# Modern analysis of automorphic forms by examples 

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## Introduction and historical notes

The aim here is persuasive proof of several important analytical results about automorphic forms, among them spectral decompositions of spaces of automorphic forms, discrete decompositions of spaces of cuspforms, meromorphic continuation of Eisenstein series, spectral synthesis of automorphic forms, a Plancherel theorem, and various notions of convergence of spectral expansions. Rather than assuming prior knowledge of the necessary analysis, or giving extensive external references, we provide customized discussions of that background, especially of ideas from 20th century analysis often neglected in the contemporary standard curriculum. Similarly, we avoid assumptions of background that would certainly be useful in studying automorphic forms, but which beginners cannot be expected to have. Therefore, we keep external references to a minimum, treating the modern analysis and other background as a significant part of the discussion.

Not only for reasons of space, the treatment of automorphic forms is deliberately neither systematic nor complete, but by three families of examples, in all cases aiming to illustrate aspects beyond the introductory case of $S L_{2}(\mathbb{Z})$ and its congruence subgroups.

The first three chapters set up three families of examples, proving essential preparatory results, proving many of the basic facts about automorphic forms, while merely stating results whose proofs are more sophisticated or difficult. The proofs of the more difficult results occupy the remainder of the book, as in many cases the arguments require various ideas not visible in the statements.

The first family of examples is introduced in chapter one, consisting of waveforms on quotients having dimensions $2,3,4,5$ with a single cusp, which is just a point. In the two-dimensional case, the space on which the functions live is the usual quotient $S L_{2}(\mathbb{Z}) \backslash \mathfrak{H}$ of the complex upper half-plane $\mathfrak{H}$. The three-dimensional case is related to $S L_{2}(\mathbb{Z}[i])$, and the four-dimensional and five-dimensional cases are similarly explicitly described. Basic discussion of the physical spaces themselves involves explication of the groups acting on them, and decompositions of these groups in terms of subgroups, and the expression of the physical spaces as $G / K$ for $K$ a maximal compact subgroup of $G$. There are natural invariant measures and integrals on $G / K$ and on $\Gamma \backslash G / K$, whose salient properties can be described quickly, with proofs deferred to a later point. Similarly, a natural Laplace-Beltrami operator $\Delta$ on $G / K$ and $\Gamma \backslash G / K$ can be described easily, but with proofs deferred. The first serious result specific to automorphic forms is about reduction theory, that is, determination of a nice set in $G / K$ that surjects to the quotient $\Gamma \backslash G / K$, for specific discrete subgroups $\Gamma$ of $G$. The four examples in this simplest scenario all admit very simple sets of representatives, called Siegel sets in every case a product of a ray and a box, with Fourier expansions possible along the box-coordinate, consonant with a decomposition of part of the group $G$ (Iwasawa decomposition). This greatly simplifies both statements and proofs of fundamental theorems.

In the simplest family of examples, the space of cuspforms consists of those functions on the quotient $\Gamma \backslash G / K$ with $0^{t h}$ Fourier coefficient identically 0 . The basic theorem, quite non-trivial to prove, is that the space of cuspforms in $L^{2}(\Gamma \backslash G / K)$ has a basis consisting of eigenfunctions for the invariant Laplacian $\Delta$. This result is one form of the discrete decomposition of cuspforms. We delay its proof, which uses many ideas not apparent in the statement of the theorem. The orthogonal complement to cuspforms in $L^{2}(\Gamma \backslash G / K)$ is readily characterized as the space of pseudo-Eisenstein series, parametrized here by test functions on $(0,+\infty)$. However, these simple, explicit automorphic forms are never eigenfunctions for $\Delta$. Rather, via Euclidean Fourier-Mellin inversion, they are expressible as integrals of (genuine) Eisenstein series, the latter eigenfunctions for $\Delta$, but unfortunately not in $L^{2}(\Gamma \backslash G / K)$. Further, it turns out that the best expression of pseudo-Eisenstein series in terms of genuine Eisenstein series $E_{s}$ involves the latter with complex parameter outside the region of convergence of the defining series. Thus arises the need to meromorphically continue the Eisenstein series in that complex parameter. Genuine proof of meromorphic continuation, with control over the behavior of the meromorphically continued function, is another basic but non-trivial result, whose proof is delayed. Granting those postponed proofs, a Plancherel theorem for the space of pseudo-Eisenstein series follows from their expansion in terms of genuine Eisenstein series, together with attention to integrals as vector-valued (rather than merely numerical), with the important corollary that such integrals commute with continuous operators on the vector space. This and other aspects of vector-valued integrals are treated at length in an appendix. Then we obtain the Plancherel theorem for the whole space of $L^{2}$ waveforms. Even for the simplest examples, these few issues illustrate the goals of this book: discrete decomposition of spaces of cuspforms, meromorphic continuation of Eisenstein series, and a Plancherel theorem.

In chapter two is the second family of examples, adele groups $G L_{2}$ over number fields. These examples subsume classical examples of quotient $\Gamma_{0}(N) \backslash \mathfrak{H}$ with several cusps, reconstituting things so that operationally there is a single cusp. Also, examples of Hilbert modular groups and Hilbert modular forms are subsumed, by rewriting things so that the vagaries of class numbers and unit groups become irrelevant. Assuming some basic algebraic number theory, we prove $p$-adic analogues of the group decomposition results proven earlier in chapter one for the purely archimedean examples. Integral operators made from $C_{c}^{o}$ functions on the $p$-adic factor groups, known as Hecke operators, are reasonable p-adic analogues of the archimedean factors' $\Delta$, although the same integral operators do make the same sense on archimedean factors. Again, the first serious result for these examples is that of reduction theory, namely, that there is a single nice set, an adelic form of a Siegel set, again nearly the product of a ray and a box, that surjects to the quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$, where $Z^{+}$is itself a ray in the center of the group. The first serious analytical result is again about discrete decomposition of spaces of cuspforms, where now relevant operators are both the invariant Laplacians and the Hecke operators. And, again, the deferred proof is much more substantial than the statement, and needs ideas not visible in the assertion itself. The orthogonal complement to cuspforms is again describable as the $L^{2}$ span of pseudo-Eisenstein series, now with a discrete parameter, a Hecke character (grossencharacter) of the ground field, in addition to the test function on $(0,+\infty)$. The pseudo-Eisenstein series are never eigenfunctions for invariant Laplacians nor for Hecke operators. Within each family, indexed by Hecke characters, every pseudo-Eisenstein series again decomposes via Euclidean Fourier-Mellin inversion as an integral of (genuine) Eisenstein series with the same discrete parameter. The genuine Eisenstein series are eigenfunctions for invariant Laplacians, and are eigenfunctions for Hecke operators at almost all finite places, but are not square-integrable. And, again, the best assertion of spectral decomposition requires a meromorphic continuation of the genuine Eisenstein series in the continuous parameter. Then a Plancherel theorem for pseudo-Eisenstein series for each discrete parameter value follows from the integral representation in terms of genuine Eisenstein series and general properties of vector-valued integrals. These are assembled into a Plancherel theorem for all $L^{2}$ automorphic forms. An appendix computes periods of Eisenstein series along copies of $G L_{1}(\widetilde{k})$ of quadratic field extensions $\widetilde{k}$ of the ground field.

Chapter three treats the most complicated of the three families of examples, including automorphic forms for $S L_{n}(\mathbb{Z})$, both purely archimedean and adelic. Again, some relatively elementary set-up regarding group decompositions is necessary, and carried out immediately. Identification of invariant differential operators and Hecke operators at finite places is generally similar to that for the previous example $G L_{2}$. A significant change is the proliferation of types of parabolic subgroups (essentially, subgroups conjugate to subgroups containing upper-triangular matrices). This somewhat complicates the notion of cuspform, although the general idea, that zeroth Fourier coefficients vanish, is still correct, if suitably interpreted. Again, the space of square-integrable cuspforms decomposes discretely, although the complexity of the proof for these examples increases significantly, and is again delayed. The increased complication of parabolic subgroups also complicates the description of the orthogonal complement to cuspforms, in terms of pseudo-Eisenstein series. For purposes of spectral decomposition, the discrete parameters now become more complicated than the $G L_{2}$ situation: cuspforms on the Levi components (diagonal blocks) in the parabolics generalize the role of Hecke characters. Further, the continuous complex parametrizations need to be over largerdimensional Euclidean spaces. Thus, we restrict attention to the two extreme cases: minimal parabolics (also called Borel subgroups) consisting exactly of upper-triangular matrices, and maximal proper parabolics, which have exactly two diagonal blocks. The minimal parabolics use no cuspidal data, but for $S L_{n}(\mathbb{Z})$ have an $(n-1)$-dimensional complex parameter. The maximal proper parabolics have just a one-dimensional complex parameter, but typically need two cuspforms on smaller groups, one on each of the two diagonal blocks. The general qualitative result that the $L^{2}$ orthogonal complement to cuspforms is spanned by pseudo-Eisenstein series of various types does still hold, and the various types of pseudo-Eisenstein series are integrals of genuine Eisenstein series with the same discrete parameters. And, again, the best description of these integrals requires the meromorphic continuation of the Eisenstein series. For non-maximal parabolics, Bochner's lemma (recalled and proven in an appendix) reduces the problem of meromorphic continuation to the maximal proper parabolic case, with cuspidal data on the Levi components. Elementary devices such as Poisson summation, that suffice for meromorphic continuation for $G L_{2}$, as we have seen in the appendix to chapter two, are inadequate to prove meromorphic continuation involving the non-elementary cuspidal data. We defer the proof. Plancherel theorems for the spectral fragments follow from the integral representations
in terms of genuine Eisenstein series, together with properties of vector-valued integrals.
The rest of the book gives proofs of those foundational analytical results, discreteness of cuspforms and meromorphic continuation of Eisenstein series, at various levels of complication, and by various devices. Perhaps surprisingly, the required analytical underpinnings are considerably more substantial than an unsuspecting or innocent bystander might imagine. Further, not everyone interested in the truth of foundational analytical facts about automorphic forms will necessarily care about their proofs, especially upon discovery that that burden is greater than anticipated. These obvious points reasonably explain the compromises made in many sources. Nevertheless, rather than either gloss over the analytical issues, or refer to encyclopedic treatments of modern analysis on a scope quite unnecessary for our immediate interests, or give suggestive but misleading neo-classical heuristics masquerading as adequate arguments for what is truly needed, the remaining bulk of the book aims to discuss analytical issues at a technical level truly sufficient to convert appealing heuristics to persuasive, genuine proofs. For that matter, one's own lack of interest in the proofs might provide all the more interest in knowing that things widely believed are in fact provable by standard methods.

Chapter four explains enough Lie theory to understand the invariant differential operators on the ambient archimedean groups $G$, both in the simplest small examples and more generally, determining the invariant Laplace-Beltrami operators explicitly in coordinates on the four simplest examples.

Chapter five explains how to integrate on quotients, without concern for explicit sets of representatives. Although in very simple situations, such as quotients $\mathbb{R} / \mathbb{Z}$ (the circle), it is easy to manipulate sets of representatives (the interval $[0,1]$ for the circle), this eventually becomes infeasible, despite the traditional example of the explicit fundamental domain for $S L_{2}(\mathbb{Z})$ acting on the upper half-plane $\mathfrak{H}$. That is, much of the picturesque detail is actually inessential, which is fortunate since that level of details is also unsustainable in all but the very simplest little examples.

Chapter six introduces natural actions of groups on spaces of functions on physical spaces on which the groups act. In some contexts, one might make a more elaborate representation theory formalism here, but it is possible to reap many of the benefits of the ideas of representation theory without the usual superstructure. That is, the idea of a linear action of a topological group on a topological vector space of functions on a physical space is the beneficial notion, with or without classification. It is true that at certain technical moments classification results are crucial, so, although we do not prove either the Borel-Casselman-Matsumoto classification in the p-adic case [Borel 1976], [Matsumoto 1977], [Casselman 1980], nor the subrepresentation theorem [Casselman 1978/80], [Casselman Miličić 1982] in the archimedean case, hopefully the roles of these results are made clear. Classification results per se, while difficult and interesting problems, do not necessarily affect the foundational analytic aspects of automorphic forms.

Chapter seven proves the discreteness of spaces of cuspforms, in various senses, in examples of varying complexity. Here, it becomes apparent that genuine proofs, as opposed to heuristics, require some sophistication concerning topologies on natural function spaces, beyond the typical Hilbert, Banach, and Fréchet spaces. Here again, there is a forward reference to the extended appendix on function spaces and classes of topological vector spaces necessary for practical analysis. Further, even less immediately apparent, but in fact already needed in the discussion of decomposition of pseudo-Eisenstein series in terms of genuine Eisenstein series, we need a coherent and effective theory of vector-valued integrals, a complete, succinct form given in the corresponding appendix, following Gelfand and Pettis, making explicit the most important corollaries on uniqueness of invariant functions, differentiation under integral signs with respect to parameters, and related.

Chapter eight fills an unobvious need, proving that automorphic forms that are of moderate growth and are eigenfunctions for Laplacians have asymptotics given by their constant terms. In the smaller examples, it is easy to make this precise. For $S L_{n}$ with $n \geq 3$, some effort is required for an accurate statement. As corollaries, $L^{2}$ cuspforms that are eigenfunctions are of rapid decay, and Eisenstein series have relatively simple asymptotics given by their constant terms. Thus, we discover again the need to prove that Eisenstein series have vector-valued meromorphic continuations, specifically, as moderate-growth functions.

Chapter nine carefully develops ideas concerning unbounded symmetric operators on Hilbert spaces, thinking especially of operators related to Laplacians $\Delta$, and especially those such that $(\Delta-\lambda)^{-1}$ is a compact-operator-valued meromorphic function of $\lambda \in \mathbb{C}$. On one hand, even a naive conception of the general behavior of Laplacians is fairly accurate, but this is due to a subtle fact that needs proof, namely,
the essential self-adjointness of Laplacians on natural spaces such as $\mathbb{R}^{n}$, multi-toruses $\mathbb{T}^{n}$, spaces $G / K$, and even spaces $\Gamma \backslash G / K$. This has a precise sense: the (invariant) Laplacian restricted to test functions has a unique self-adjoint extension, which then is necessarily its graph-closure. Thus, the naive presumption, implicit or explicit, that the graph closure is a (maximal) self-adjoint extension is correct. On the other hand, the proof of meromorphic continuation of Eisenstein series in [Colin de Verdière 1981/82/83] makes essential use of some quite counter-intuitive features of (Friedrichs') self-adjoint extensions of restrictions of self-adjoint operators, which therefore merit careful attention. In this context, the basic examples are the usual Sobolev spaces on $\mathbb{T}$ or $\mathbb{R}$, and the quantum harmonic oscillator $-\Delta+x^{2}$ on $\mathbb{R}$. An appendix recalls the proof of the spectral theorem for compact, self-adjoint operators.

Chapter ten extends the idea from [Lax-Phillips 1976] to prove that larger spaces than spaces of cuspforms decompose discretely under the action of self-adjoint extensions $\widetilde{\Delta}_{a}$ of suitable restrictions $\Delta_{a}$ of Laplacians. Namely, the space of pseudo-cuspforms $L_{a}^{2}$ at cut-off height $a$ is specified, not by requiring constant terms to vanish entirely, but by requiring that all constant terms vanish above height $a$. The discrete decomposition is proven, as expected, by showing that the resolvent $\left(\widetilde{\Delta}_{a}-\lambda\right)^{-1}$ is a meromorphic compact-operator-valued function of $\lambda$, and invoking the spectral theorem for self-adjoint compact operators. The compactness of the resolvent is a Rellich-type compactness result, proven by observing that $\left(\Delta_{a}-\lambda\right)^{-1}$ maps $L_{a}^{2}$ to a Sobolevtype space $\mathfrak{B}_{a}^{1}$ with a finer topology on $\mathfrak{B}_{a}^{1}$ than the subspace topology, and that the inclusion $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}$ is compact.

Chapter eleven uses the discretization results of chapter ten to prove meromorphic continuations and functional equations of a variety of Eisenstein series, following [Colin de Verdière 1981/82/83]'s application of the discreteness result in [Lax-Phillips 1976]. This is carried out first for the four simple examples, then for maximal proper parabolic Eisenstein series for $S L_{n}(\mathbb{Z})$, with cuspidal data. In both the simplest cases and the higher-rank examples, we identify the exotic eigenfunctions as being certain truncated Eisenstein series.

Chapter twelve uses several of the analytical ideas and methods of the previous chapters to reconsider automorphic Green's functions, and solutions to other differential equations in automorphic forms, by spectral methods. We prove a pre-trace formula in the simplest example, as an application of a comparably simple instance of a subquotient theorem, which follows from asymptotics of solutions of second-order ordinary differential equations, recalled in a later appendix. We recast the pre-trace formula as a demonstration that an automorphic Dirac $\delta$-function lies in the expected global automorphic Sobolev space. The same argument gives a corresponding result for any compact automorphic period. Subquotient/subrepresentation theorems for groups such as $G=S O(n, 1)$ (rank-one groups with abelian unipotent radicals) appeared in [Casselman-Osborne 1975], [Casselman-Osborne 1978]. For higher-rank groups $S L_{n}(\mathbb{Z})$, the corresponding subrepresentation theorem is [Casselman 1978/80], [Casselman Miliči'c 1982]. Granting that, we obtain a corresponding pre-trace formula for a class of compactly-supported automorphic distributions, showing that these distributions lie in the expected global automorphic Sobolev spaces.

Chapter thirteen is an extensive appendix with many examples of natural spaces of functions and appropriate topologies on them. One point is that too-limited types of topological vector spaces are inadequate to discuss natural function spaces arising in practice. We include essential standard arguments characterizing locally convex topologies in terms of families of seminorms. We prove the quasi-completeness of all natural function spaces, and weak duals, and spaces of maps between them. Notably, this includes spaces of distributions.

Chapter fourteen proves existence of Gelfand-Pettis vector-valued integrals of compactly-supported continuous functions taking values in locally convex, quasi-complete topological vector space. Conveniently, the previous chapter showed that all function spaces of practical interest meet these requirements. The fundamental property of Gelfand-Pettis integrals is that

$$
T\left(\int f\right)=\int T \circ f \quad(\text { for } V \text {-valued } f, T: V \rightarrow W \text { continuous linear) }
$$

at least for $f$ continuous, compactly supported, $V$-valued, where $V$ is quasi-complete and locally convex. That is, continuous linear operators pass inside the integral. In suitably-topologized natural function spaces, this situation includes differentiation with respect to a parameter. In this situation, as corollaries we can easily prove uniqueness of invariant distributions, density of smooth vectors, and similar.

Chapter fifteen carefully discusses holomorphic $V$-valued functions, using the Gelfand-Pettis integrals as well as a variant of the Banach-Steinhaus theorem. That is, weak holomorphy implies (strong) holomorphy, and the expected Cauchy integral formulas and Cauchy-Goursat theory apply almost verbatim in the vectorvalued situation. Similarly, we prove that for $f$ a $V$-valued function on an interval $[a, b], \lambda \circ f$ being $C^{k}$ for all $\lambda \in V^{*}$ implies that $f$ itself is $C^{k-1}$ as a $V$-valued function.

Chapter sixteen reviews basic results on asymptotic expansions of integrals, and of solutions to second-order ordinary differential equations. The methods are deliberately general, rather than invoking specific features of special functions, to illustrate methods that are applicable more broadly. The simple subrepresentation theorem in chapter twelve makes essential use of asymptotic expansions.

Our coverage of modern analysis does not aim to be either systematic or complete, but well-grounded and adequate for the above-mentioned issues concerning automorphic forms. In particular, several otherwiseapocryphal results are treated carefully. We want a sufficient viewpoint so that attractive heuristics, for example, from physics, can become succinct, genuine proofs. Similarly, we do not presume familiarity with Lie theory, nor algebraic groups, nor representation theory, nor algebraic geometry, and certainly not with classification of representations of Lie groups or $p$-adic groups. All these are indeed very useful, in the long run, but it is unreasonable to demand mastery of these prior to thinking about analytical issues concerning automorphic forms. Thus, we directly develop some essential ideas in these supporting topics, sufficient for immediate purposes here. [Lang 1975] and [Iwaniec 2002] are examples of the self-supporting exposition intended here.

Naturally, any novelty here is mostly in the presentation, rather than in the facts themselves, most of which have been known for several decades. Sources and origins can be most clearly described in a historical context, as follows.

The reduction theory in [1.5] is merely an imitation of the very classical treatment for $S L_{2}(\mathbb{Z})$, including some modern ideas, as in [Borel 1997]. The subtler versions in [2.2] and [3.3] are expanded versions of the first part of [Godement 1963], a more adele-oriented reduction theory than [Borel 1965/6b], [Borel 1969], and [Borel-HarishChandra 1962]. Proofs [1.9.1], [2.8.6], [3.10.1-2], [3.11.1] of convergence of Eisenstein series are due to Godement use similar ideas, reproduced for real Lie groups in [Borel 1965/6]. Convergence arguments on larger groups go back at least to [Braun 1939]'s treatment of convergence of Siegel Eisenstein series. Holomorphic Hilbert-Blumenthal modular forms were studied by [Blumenthal 1904]. What would now be called degenerate Eisenstein series for $G L_{n}$ appeared in [Epstein 1903/07]. [Picard 1882/83/84] was one of the earliest investigations beyond the elliptic modular case. Our notion of truncation is from [Arthur 1978] and [Arthur 1980].

Eigenfunction expansions and various notions of convergence are a pervasive theme here, and have a long history. The idea that periodic functions should be expressible in terms of sines and cosines is at latest from [Fourier 1822], including what we now call the Dirichlet kernel, although [Dirichlet 1829] came later. Somewhat more generally, eigenfunction expansions for Sturm-Liouville problems appeared in [Sturm 1836] and [Sturm 1833a,b/36a,b] but were not made rigorous until [Bôcher 1898/99] and [Steklov 1898] (see [Lützen 1984]). Refinements of the spectral theory of ordinary differential equations continued in [Weyl 1910], [Kodaira 1949], and others, addressing issues of non-compactness and unboundedness echoing complications in the behavior of Fourer transform and Fourier inversion on the line [Bochner 1932], [Wiener 1933]. Spectral theory and eigenfunction expansions for integral equations, which we would now call compact operators [9.A], were recognized as more tractable than direct treatment of diffferential operators soon after 1900: [Schmidt 1907], [Myller-Lebedev 1907], [Riesz 1907], [Hilbert 1909], [Riesz 1910], [Hilbert 1912]. Expansions in spherical harmonics were used in the 18 th century by S. P. Laplace and J.-L. Lagrange, and eventually subsumed in the representation theory of compact Lie groups [Weyl 1925/6] and in eigenfunction expansions on Riemannian manifolds and Lie groups, as in [Minakshisundaram-Pleijel 1949], [Povzner 1953], [Avakumović 1956], [Berezin 1956], and many others.

Spectral decomposition and synthesis of various types of automorphic forms is more recent, beginning with [Maaß 1949], [Selberg 1956], and [Roelcke 1956a,b]. The spectral decomposition for automorphic forms on general reductive groups is more complicated than might have been anticipated by the earliest pioneers. Subtleties are already manifest in [Gelfand-Fomin 1952], and then in [Gelfand-Graev 1959], [HarishChandra 1959], [Gelfand-PS 1963], [Godement 1966b], [HarishChandra 1968], [Langlands 1966], [Langlands 1967/76], [Arthur 1978], [Arthur 1980], [Jacquet 1982/83], [Moeglin-Waldspurger 1989], [Moeglin-

Waldspurger 1995], [Casselman 2005], [Shahidi 2010]. Despite various formalizations, spectral synthesis of automorphic forms seems most clearly understood in fairly limited scenarios: [Godement 1966a], [Faddeev 1967], [Venkov 1971], [Faddeev-Pavlov 1972], [Arthur 1978], [Venkov 1979], [Arthur 1980], [CogdellPS 1990], largely due to issues of convergence, often leaving discussions in an ambiguous realm of (nevertheless interesting) heuristics.

Regarding meromorphic continuation of Eisenstein series: our proof [2.B] for the case [2.9] of $G L_{2}$ is an adaptation of the Poisson summation argument from [Godement 1966a]. The essential idea already occurred in [Rankin 1939] and [Selberg 1940]. [Elstrodt-Grunewald-Mennicke 1985] treated examples including our example $S L_{2}(\mathbb{Z}[i])$, and in that context [Elstrodt-Grunewald-Mennicke 1987] treats special cases of the period computation of [2.C]. For Eisenstein series in rank one groups, compare also [CohenSarnak 1980], which treats a somewhat larger family including our simplest four examples, and then [Müller 1996]. The minimal-parabolic example in [3.12] using Bochner's lemma [3.A] essentially comes from an appendix in [Langlands 1967/76]. The arguments for the broader class of examples in chapter eleven are adaptations of [Colin de Verdière 1981/82/83], using discretization effects of pseudo-Laplacians from chapter ten, which adapts the idea of [Lax-Phillips 1976]. Certainly one should compare the arguments in [HarishChandra 1968], [Langlands 1967/76], [Wong 1990], and [Moeglin-Waldspurger 1995]. The latter gives a version of Colin de Verdière's idea due to H. Jacquet.

The discussion of group actions on function spaces in chapter six is mostly very standard. Apparently the first occurrence of the Gelfand-Kazhdan criterion idea is in [Gelfand 1950]. An extension of that idea appeared in [Gelfand-Kazhdan 1975].

The arguments for discrete decomposition of cuspforms in chapter seven are adaptations of [Godement 1966b]. The discrete decomposition examples for larger spaces of pseudo-cuspforms in chapter ten use the idea of [Lax-Phillips 1976]. The idea of this decomposition perhaps goes back to [Gelfand-Fomin 1952], and, as with many of these ideas, was elaborated-upon in the iconic sources [Gelfand-Graev 1959], [HarishChandra 1959], [Gelfand-PS 1963], [Godement 1966b], [HarishChandra 1968], [Langlands 1967/76], and [Moeglin-Waldspurger 1989].

Difficulties with pointwise convergence of Fourier series of continuous functions, and problems in other otherwise-natural Banach spaces of functions, were well appreciated in the late 19th century. There was a precedent for constructs avoiding strictly pointwise conceptions of functions in the very early 20th century, when B. Levi, G. Fubini, and D. Hilbert used Hilbert space constructs to legitimize Dirichlet's minimization principle, in essence that a non-empty closed convex set should have a (unique) point nearest a given point not in that set. The too-general form of this principle is false, in that both existence and uniqueness easily fail in Banach spaces, in natural examples, but the principle is correct in Hilbert spaces. Thus, natural Banach spaces of pointwise-valued functions, such as continuous functions on a compact set with sup norm, do not support this minimization principle. Instead, Hilbert-space versions of continuity and differentiability are needed, as in [Levi 1906]. This idea was systematically developed by [Sobolev 1937, 1938, 1950]. We recall the $L^{2}$ Sobolev spaces for circles in [9.5], for lines in [9.7], and develop various (global) automorphic versions of Sobolev spaces in chapters ten, eleven, and twelve.

For applications to analytic number theory, automorphic forms are often constructed by winding up various simpler functions containing parameters, forming Poincaré series [Cogdell-PS 1990], [Cogdell-PSSarnak 1991. Spectral expansions are the standard device for demonstration of meromorphic continuation in the parameters, if it exists at all, which is a non-trivial issue [Estermann 1928], [Kurokawa 1985a,b]. For the example of automorphic Green's functions, namely, solutions to equations $(\Delta-s(s-1)) u=\delta_{w}^{\text {afc }}$ with invariant Laplacian $\Delta$ on $\mathfrak{H}$ and automorphic Dirac $\delta$ on the right, [Huber 1955] had considered such matters in the context of lattice-point problems in hyperbolic spaces, and, independently, [Selberg 1954] had addressed this issue in lectures in Göttingen. [Neunhöffer 1973] carefully considers the convergence and meromorphic continuation of a solution of that equation formed by winding up. See also [Elstrodt 1973]. The complications or failures of pointwise convergence of the spectral synthesis expressions can often be avoided entirely by considering convergence in suitable global automorphic Sobolev spaces described in chapter twelve. See [DeCelles 2012] and [DeCelles 2016] for developments in this spirit.

Because of the naturality of the issue, and to exploit interesting idiosyncrasies, we pay considerable attention to invariant Laplace-Beltrami operators and their eigenfunctions. To have genuine proofs, rather than heuristics, chapter nine attends to rigorous notions of unbounded operators on Hilbert spaces
[vonNeumann 1929], with motivation toward [vonNeumann 1931], [Stone 1929/32], [Friedrichs 1934], [Krein 1945], [Krein 1947]. In fact, [Friedrichs 1934/5]' special construction [9.2] has several useful idiosyncracies, exploited in chapters ten and eleven. Incidentally, the apparent fact that the typically naive treatment of many natural Laplace-Beltrami operators without boundary conditions does not lead to serious mistakes is a corollary of their essential self-adjointness [9.9], [9.10]. That is, in many situations, the naive form of the operator admits a unique self-adjoint extension, and this extension is the graph closure of the original. Thus, in such situations, a naive treatment is provably reasonable. However, the Lax-Phillips discretization device, and Colin de Verdière's use of it to prove meromorphic continuation of Eisenstein series, and also to convert certain inhomogeneous differential equations to homogeneous ones, illustrate the point that restrictions of essentially self-adjoint operators need not remain essentially self-adjoint. With hindsight, this possibility is already apparent in the context of Sturm-Liouville problems [9.3].

The global automorphic Sobolev spaces of chapter twelve already enter in important auxiliary roles as the spaces $\mathfrak{B}^{1}, \mathfrak{B}_{a}^{1}$ in chapter ten's proofs of discrete decomposition of spaces of pseudo-cuspforms, and $\mathfrak{E}^{1}$ and $\mathfrak{E}_{a}^{1}$ in [11.7-11.11] proving meromorphic continuation of Eisenstein series. The basic estimate called a pre-trace formula occurred as a precursor to trace formulas, as in [Selberg 1954], [Selberg 1956], [Hejhal 1976/83], and [Iwaniec 2002]. The notion of global automorphic Sobolev spaces provides a reasonable context for discussion of automorphic Green's functions, other automorphic distributions, and solutions of partial differential equations in automorphic forms. The heuristics for Green's functions [Green 1828], [Green 1837] had repeatedly shown their utility in the 19 th century. Differential equations $(-\Delta-\lambda) u=\delta$ related to Green's functions had been used by physicists [Dirac 1928a/b, 1930], [Thomas 1935], [Bethe-Peierls 1935], with excellent corroboration by physical experiments, and are nowadays known as solvable models. At the time, and currently, in physics contexts they are rewritten as $((-\Delta+\delta)-\lambda) u=0$, viewing $-\Delta+\delta$ as a perturbation of $-\Delta$ by a singular potential $\delta$, a mathematical idealization of a very-short-range force. This was treated rigorously in [Berezin-Faddeev 1961]. The necessary systematic estimates on eigenvalues of integral operators use a subquotient theorem, which we prove for the four simple examples, as in that case the issue is about asymptotics of solutions of second-order differential equations, classically understood as recalled in an appendix (chapter sixteen). The general result is the subrepresentation theorem from [Casselman 1978/80], [Casselman Miličić 1982], improving the subquotient theorem of [Harish-Chandra 1954]. In [Varadarajan 1989] there are related computations for $S L_{2}(\mathbb{R})$.

In the discussion of natural function spaces in chapter thirteen, in preparation for the vector-valued integrals of the following chapter, the notion of quasi-completeness proves to be the correct general version of completeness. The incompleteness of weak duals has been known at least since [Grothendieck 1950], which gives a systematic analysis of completeness of various types of duals. This larger issue is systematically discussed in [Schaefer 1966/99], p. 147-8 and following. The significance of the compactness of the closure of the convex hull of a compact set appears, for example, in the discussion of vector-valued integrals in [Rudin 1991], although the latter does not make clear that this condition is fulfilled in more than Fréchet spaces, and does not mention quasi-completeness. To apply these ideas must be applicable to distributions, one might cast about for means to prove the compactness condition, eventually hitting upon the hypothesis of quasi-completeness in conjunction with ideas from the proof of the Banach-Alaoglu theorem. Indeed, in [Bourbaki 1987] it is shown (by apparently different methods) that quasi-completeness implies this compactness condition. The fact that a bounded subset of a countable strict inductive limit of closed subspaces must actually be a bounded subset of one of the subspaces, easy to prove once conceived, is attributed to Dieudonne and Schwartz in [Horvath 1966]. See also [Bourbaki 1987], III. 5 for this result. Pathological behavior of uncountable colimits was evidently first exposed in [Douady 1963].

In chapter fourteen, rather than constructing vector-valued integrals as limits following [Bochner 1935], [Birkhoff 1935], et alia, we use the [Gelfand 1936]-[Pettis 1938] characterization of integrals, which has good functorial properties and gives a forceful reason for uniqueness. The issue is existence. Density of smooth vectors follows [Gårding 1947]. Another of application of holomorphic and meromorphic vector-valued functions is to generalized functions, as in [Gelfand-Shilov 1964], studying holomorphically parametrized families of distributions. A hint appears in the discussion of holomorphic vector-valued functions in [Rudin 1991]. A variety of developmental episodes and results in the Banach-space-valued case is surveyed in [Hildebrandt 1953]. Proofs and application of many of these results are given in [Hille-Phillips 1957]. (The first edition, authored by Hille alone, is sparser in this regard.) See also [Brooks 1969] to understand the
viewpoint of those times.
Ideas about vector-valued holomorphic and differentiable functions, in chapter fifteen, appeared in [Schwartz 1950/51], [Schwartz 1952], [Schwartz 1953/4], and in [Grothendieck 1953a,b].

The asymptotic expansion results of chapter sixteen are standard. [Blaustein-Handelsman 1975] is a standard source for asymptotics of integrals. Watson's lemma and Laplace's method for integrals have been used and rediscovered repeatedly. Watson's lemma dates from at latest [Watson 1918], and Laplace's method at latest from [Laplace 1774]. [Olver 1954] notes that Carlini, [Green 1837], and [Liouville 1837] investigated relatively simple cases of asymptotics at irregular singular points of ordinary differential equations, without complete rigor. According to [Erdélyi 1956] p. 64, there are roughly two proofs that the standard argument produces genuine asymptotic expansions for solutions of the differential equation. Poincaré's approach, elaborated by J. Horn, expresses solutions as Laplace transforms and invokes Watson's lemma to obtain asymptotics. G.D. Birkhoff and his students constructed auxiliary differential equations from partial sums of the asymptotic expansion, and compared these auxiliary equations to the original [Birkhoff 1908], [Birkhoff 1909], [Birkhoff 1913]. Volterra integral operators are important in both approaches, insofar as asymptotic expansions behave better under integration than under differentiation. Our version of the Birkhoff argument is largely adapted from [Erdélyi 1956].

Many parts of this exposition are adapted and expanded from [Garrett vignettes], [Garrett mfms-notes], [Garrett fun-notes], and [Garrett alg-noth-notes]. As is surely usual in book writing, many of the issues here had plagued me for decades.
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## 1. Four small examples

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We recall basic notions related to automorphic forms on some simple arithmetic quotients, including the archetypical quotient $S L_{2}(\mathbb{Z}) \backslash \mathfrak{H}$ of the complex upper half-plane $\mathfrak{H}$ and the related quotient $S L_{2}(\mathbb{Z}) \backslash S L_{2}(\mathbb{R})$. To put this in a somewhat larger context, ${ }^{[1]}$ we consider parallel examples $\Gamma \backslash X$ and $\Gamma \backslash G$ for a few other groups $G$, discrete subgroups $\Gamma$, and spaces $X \approx G / K$ for compact subgroups $K$ of $G$. The other three examples share several of the features of $G=S L_{2}(\mathbb{R}), \Gamma=S L_{2}(\mathbb{Z}), X=\mathfrak{H} \approx G / K$ with $K=S O_{2}(\mathbb{R})$, allowing simultaneous treatment.

For many reasons, even if we are only interested in harmonic analysis on quotients $\Gamma \backslash X$, it is necessary to consider spaces of functions on the overlying spaces $\Gamma \backslash G$, on which $G$ acts by right translations, with a corresponding translation action on functions.

Some basic discussions not specific to the four examples are postponed, such as determination of invariant Laplacians in coordinates, self-adjointness properties of invariant Laplacians, proof of the formula for the left $G$-invariant measure on $X=G / K$, unwinding properties of integrals and sums, continuity of the action of $G$ on test functions on $\Gamma \backslash G$, density of test functions in $L^{2}(\Gamma \backslash X)$, vector-valued integrals, holomorphic vector-valued functions, and other generalities.

We also postpone the relatively specific proofs of the major theorems stated in the last sections of this chapter, concerning the spectral decomposition of automorphic forms, meromorphic continuation of Eisenstein series, and the theory of the constant term. Those proofs make pointed use of finer details from the more sophisticated analysis.

[^0]$$
\text { 1.1 Groups } G=S L_{2}(\mathbb{R}), S L_{2}(\mathbb{C}), S p_{1,1}^{*}, \text { and } S L_{2}(\mathbb{H})
$$

These four groups share some convenient simplifying features, which we will exploit. The first two examples $G$ are easy to describe:
$G=\left\{\begin{array}{l}\text { a special linear group over } \mathbb{R}=S L_{2}(\mathbb{R})=\text { two-by-two real matrices with determinant } 1 \\ \text { a special linear group over } \mathbb{C}=S L_{2}(\mathbb{C})=\text { two-by-two complex matrices with determinant } 1\end{array}\right.$
We will have occasion to use the general linear groups $G L_{2}(R)$ of 2-by-2 invertible matrices with entries in a ring $R$. Our other two example groups are conveniently described in terms of the Hamiltonian quaternions $\mathbb{H}=\mathbb{R}+\mathbb{R} i+\mathbb{R} j+\mathbb{R} k$, with the usual relations

$$
i^{2}=j^{2}=k^{2}=-1 \quad i j=-j i=k \quad j k=-k j=i \quad k i=-i k=j
$$

The quaternion conjugation is $\bar{\alpha}=\overline{a+b i+c j+d k}=a-b i-c j-d k$ for $\alpha=a+b i+c j+d k$, the norm is $N \alpha=\alpha \cdot \bar{\alpha}$, and $|\alpha|=(N \alpha)^{\frac{1}{2}} . \mathbb{H}$ can be modeled in two-by-two complex matrices by

$$
\rho(a+b i+c j+d k)=\left(\begin{array}{rr}
a+b i & c+d i \\
-c+d i & a-b i
\end{array}\right)
$$

with $\operatorname{det} \rho(\alpha)=N \alpha$. For a quaternion matrix $g$, let $g^{*}$ be the transpose of the entry-wise conjugate:

$$
\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right)^{*}=\left(\begin{array}{cc}
\bar{\alpha} & \bar{\gamma} \\
\bar{\beta} & \bar{\delta}
\end{array}\right) \quad(\text { for } \alpha, \beta, \gamma, \delta \in \mathbb{H})
$$

The third example group is a kind of symplectic group:

$$
\left.G=S p_{1,1}^{*}=\left\{g \in G L_{2}(\mathbb{H}): g^{*} S g=S\right\} \quad \text { (with } S=\left(\begin{array}{cc}
0 & 1 \\
1 & 0
\end{array}\right)\right)
$$

The fourth example is a special linear group $G=S L_{2}(\mathbb{H})$. In the latter, $S L_{2}$ is more convenient than $G L_{2}$, having a smaller center. However, since $\mathbb{H}$ is not commutative, the notion of determinant is problemmatical. One way to skirt the issue is to imbed $r: G L_{2}(\mathbb{H}) \rightarrow G L_{4}(\mathbb{C})$ : with quaternions $\alpha, \beta, \gamma, \delta$,

$$
r\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right)=\left(\begin{array}{cc}
\rho(\alpha) & \rho(\beta) \\
\rho(\gamma) & \rho(\delta)
\end{array}\right) \quad \text { (identified with a 4-by-4 complex matrix) }
$$

using the map $\rho$ of $\mathbb{H}$ to 2 -by- 2 complex matrices, and require that the image in $G L_{4}(\mathbb{C})$ be in the subgroup $S L_{4}(\mathbb{C})$ where determinant is 1 :

$$
S L_{2}(\mathbb{H})=\left\{g \in G L_{2}(\mathbb{H}): r(g) \in S L_{4}(\mathbb{C})\right\}
$$

Standard subgroups of any of these groups $G$ are
$P=\left\{\left(\begin{array}{cc}* & * \\ 0 & *\end{array}\right)\right\} \quad N=\left\{\left(\begin{array}{cc}1 & * \\ 0 & 1\end{array}\right)\right\} \quad M=\left\{\left(\begin{array}{cc}* & 0 \\ 0 & *\end{array}\right)\right\} \quad A^{+}=\left\{\left(\begin{array}{cc}t & 0 \\ 0 & t^{-1}\end{array}\right): t>0\right\}$
The Levi-Malcev decomposition $P=N M$ is elementary to check. By direct computation from the defining relations of the groups, one finds

$$
M=\left\{\begin{array}{cl}
\left\{\left(\begin{array}{cc}
m & 0 \\
0 & m^{-1}
\end{array}\right): m \in \mathbb{R}^{\times}\right\} & \left(\text {for } G=S L_{2}(\mathbb{R})\right) \\
\left\{\left(\begin{array}{cc}
m & 0 \\
0 & m^{-1}
\end{array}\right): m \in \mathbb{C}^{\times}\right\} & \left(\text {for } G=S L_{2}(\mathbb{C})\right) \\
\left\{\left(\begin{array}{cc}
m & 0 \\
0 & \bar{m}^{-1}
\end{array}\right): m \in \mathbb{H}^{\times}\right\} & \left(\text {for } G=S p_{1,1}^{*}\right) \\
\left\{\left(\begin{array}{cc}
a & 0 \\
0 & d
\end{array}\right): N(a d)=1, a, d \in \mathbb{H}^{\times}\right\} & \left(\text {for } G=S L_{2}(\mathbb{H})\right)
\end{array}\right.
$$

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and

$$
N=\left\{\begin{array}{cl}
\left\{\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right): x \in \mathbb{R}\right\} & \left(\text { for } G=S L_{2}(\mathbb{R})\right) \\
\left\{\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right): x \in \mathbb{C}\right\} & \left(\text { for } G=S L_{2}(\mathbb{C})\right) \\
\left\{\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right): x \in \mathbb{H}, x+\bar{x}=0\right\} & \left(\text { for } G=S p_{1,1}^{*}\right) \\
\left\{\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right): x \in \mathbb{H}\right\} & \left(\text { for } G=S L_{2}(\mathbb{H})\right)
\end{array}\right.
$$

The subgroup $P$ is the standard (proper) parabolic, $N$ is its unipotent radical, $M$ is the standard LeviMalcev component, and $A^{+}$is the standard split component. We will use these (standard) names without elaborating on their history or their connotations.

In these examples, the (spherical) Bruhat decomposition is

$$
G=\bigsqcup_{w=1, w_{o}} P w P=P \sqcup P w_{o} P=P \sqcup P w_{o} N \quad \quad\left(\text { where } w_{o}=\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)\right)
$$

with the last equality following because $w_{o}$ normalizes $M$ :

$$
P w_{o} P=P w_{o} M N=P\left(w_{o} M w_{o}^{-1}\right) w_{o} N=P w_{o} N
$$

The element $w_{o}$ is the long Weyl element. The small (Bruhat) cell is $P$ itself, and the big (Bruhat) cell is $P w_{o} P$. The (spherical, geometric) Weyl group is $\left\{1, w_{o}\right\}$. It is a group modulo the center of $G$. The proof of the Bruhat decomposition is straightforward: $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in P$ if and only if $c=0$. Otherwise, $c \neq 0$, and we try to find $p \in P$ and $n \in N$ such that $g=p w_{o} n$. To simplify, since $c \neq 0$, it is invertible, so, in a form applicable to all four cases, we can left multiply by $\left(\begin{array}{cc}c & 0 \\ 0 & c^{-1}\end{array}\right) \in M$ to make $c=1$ without loss of generality. Then try to solve

$$
\left(\begin{array}{ll}
a & b \\
1 & d
\end{array}\right)=g=p w_{o} n=\left(\begin{array}{cc}
p_{11} & p_{12} \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)\left(\begin{array}{cc}
1 & n_{12} \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
p_{12} & p_{12} n_{12}-p_{11} \\
1 & n_{12}
\end{array}\right)
$$

From the lower right entry, apparently $n_{12}=d$. For the case $G=S p_{1,1}^{*}$ the additional condition must be checked, as follows. Observe that inverting $g^{*} S g=S$ gives $g^{-1} S^{-1}\left(g^{*}\right)^{-1}=S^{-1}$, and then $S=g S g^{*}$. In particular, this gives a relation between the $c, d$ entries of $g$ :

$$
\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)=S=g S g^{*}=\left(\begin{array}{ll}
* & * \\
c & d
\end{array}\right)\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)\left(\begin{array}{cc}
* & \bar{c} \\
* & \bar{d}
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
* & c \bar{d}+d \bar{c}
\end{array}\right)
$$

For $c=1$, this gives $d+\bar{d}=0$, which is the condition for $\left(\begin{array}{ll}1 & d \\ 0 & 1\end{array}\right) \in N$ in that case. Thus, in all cases, right multiplying $g$ by $\left(\begin{array}{cc}1 & -d \\ 0 & 1\end{array}\right) \in N$ makes $d=0$, without loss of generality. Thus, it suffices to solve

$$
\left(\begin{array}{ll}
a & b \\
1 & 0
\end{array}\right)=g=p w_{o}=\left(\begin{array}{cc}
p_{11} & p_{12} \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)=\left(\begin{array}{cc}
p_{12} & -p_{11} \\
1 & 0
\end{array}\right)
$$

That is,

$$
g w_{o}^{-1}=\left(\begin{array}{rr}
-b & a \\
0 & 1
\end{array}\right)=p
$$

Since $g \in G$, the entries $a, b$ satisfy whatever relations $G$ requires, and $p \in G$. This proves the Bruhat decomposition.

### 1.2 Compact subgroups $K \subset G$, Cartan decompositions

We describe the standard maximal ${ }^{[2]}$ compact subgroups $K \subset G$ for the four examples $G$. With $\mathbb{H}^{1}$ the quaternions of norm 1, in a notation consistent with that for $S p_{1,1}^{*}$, write

$$
S p_{1}^{*}=\left\{g \in G L_{1}(\mathbb{H}): g^{*} g=1\right\}=\left\{g \in \mathbb{H}^{\times}: \bar{g} g=1\right\}=\mathbb{H}^{1}
$$

Letting $1_{2}$ be the two-by-two identity matrix, the four maximal compact subgroups are

$$
K=\left\{\begin{array}{cll}
S O_{2}(\mathbb{R}) & =\left\{g \in S L_{2}(\mathbb{R}): g^{\top} g=1_{2}\right\} & \\
& \left(\text { for } G=S L_{2}(\mathbb{R})\right) \\
S U_{2} & =\left\{g \in S L_{2}(\mathbb{C}): g^{*} g=1_{2}\right\} & \\
\left(\text { for } G=S L_{2}(\mathbb{C})\right) \\
S p_{1}^{*} \times S p_{1}^{*} & =\mathbb{H}^{1} \times \mathbb{H}^{1} & \\
S p_{2}^{*} & =\left\{g \in G L_{2}(\mathbb{H}): g^{*} g=1_{2}\right\} & \\
\left(\text { for } G=S p_{1,1}^{*}\right) \\
\left.S L_{2}(\mathbb{H})\right)
\end{array}\right.
$$

In all four cases, the indicated groups are compact. Verification of the compactness of the first three is straightforward, since their defining equations present them as spheres or products of spheres. Verification that $S p_{2}^{*}$ is compact and is a subgroup of $S L_{2}(\mathbb{H})$ merits discussion. For the fourth, observe that the defining condition

$$
\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)^{*}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=\left(\begin{array}{cc}
|a|^{2}+|c|^{2} & \bar{a} b+\bar{c} d \\
\bar{b} a+\bar{d} c & |b|^{2}+|d|^{2}
\end{array}\right)
$$

makes $S p_{2}^{*}$ a closed subset of a product of two seven-spheres, $|a|^{2}+|c|^{2}=1$ and $|b|^{2}+|d|^{2}=1$, thus, compact. Further, $S p_{2}^{*}$ lies inside $S L_{2}(\mathbb{H})$ rather than merely $G L_{2}(\mathbb{H})$. For the moment, we will prove a slightly weaker property, that the relevant determinant is $\pm 1$. Use the feature

$$
\rho(\bar{\alpha})=\varepsilon \rho(\alpha)^{\top} \varepsilon^{-1} \quad\left(\text { where } \varepsilon=\left(\begin{array}{rr}
0 & -1 \\
1 & 0
\end{array}\right), \text { for } \alpha \in \mathbb{H}\right)
$$

of the imbedding $\rho$ of $\mathbb{H}$ in 2-by-2 complex matrices, and again let

$$
r\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=\left(\begin{array}{cc}
\rho(a) & \rho(b) \\
\rho(c) & \rho(d)
\end{array}\right) \quad(\text { for } a, b, c, d \in \mathbb{H})
$$

viewed as mapping to 4 -by- 4 complex matrices. Then

$$
r\left(g^{*}\right)=J \cdot r(g) \cdot J^{-1} \quad\left(\text { where } J=\left(\begin{array}{ll}
\varepsilon & \\
& \varepsilon
\end{array}\right)=\left(\begin{array}{lll}
1 & & \\
& & \\
& & \\
& & \\
& &
\end{array}\right), \text { and } g \in G L_{2}(\mathbb{H})\right)
$$

Thus, for $g^{*} g=1_{2} \in G L_{2}(\mathbb{H})$,

$$
1_{4}=r\left(1_{2}\right)=r\left(g^{*} g\right)=r\left(g^{*}\right) \cdot r(g)=J \cdot r(g)^{\top} \cdot J^{-1} \cdot r(g)
$$

In other words, $r(g)^{\top} J r(g)=J$. ${ }^{[3]}$ Taking determinants shows $\operatorname{det} r(g)^{2}=1$, so $\operatorname{det} r(g)= \pm 1$. Thus, $g$ in the connected component of $S p_{2}^{*}$ containing 1 has $\operatorname{det} r(g)=1$.
[2] The maximality of each of these subgroups $K$ among all compact subgroups in the corresponding $G$ is not obvious, but is not used in the sequel.
[3] Thus, $r(g)$ is inside a symplectic group denoted $S p_{4}(\mathbb{C})$ or $S p_{2}(\mathbb{C})$, depending on convention.

## 1. Four small examples

The copy $K$ of $S p_{1}^{*} \times S p_{1}^{*}$ inside $S p_{1,1}^{*}$ is not immediately visible in these coordinates, which were chosen to make the parabolic $P$ visible. That is, defining $S p_{1,1}^{*}$ as the isometry group of the quaternion hermitian form $S$ obscures the nature of the (maximal) compact $K$. Changing coordinates by replacing $S$ by

$$
S^{\prime}=\frac{1}{2}\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right) S\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)^{\top}=\frac{1}{2}\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)\left(\begin{array}{rr}
1 & 1 \\
-1 & 1
\end{array}\right)=\left(\begin{array}{rr}
-1 & 0 \\
0 & 1
\end{array}\right)
$$

gives

$$
\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right) S p_{1,1}^{*}\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)^{-1}=\left\{g \in G L_{2}(\mathbb{H}): g^{*} S^{\prime} g=S^{\prime}\right\}
$$

and makes the two copies of $S p_{1}^{*}$ visible on the diagonal:

$$
\left\{k=\left(\begin{array}{cc}
* & 0 \\
0 & *
\end{array}\right): k^{*} S^{\prime} k=S^{\prime}\right\}=\left\{k=\left(\begin{array}{cc}
\mu & 0 \\
0 & \nu
\end{array}\right): \mu, \nu \in \mathbb{H}^{1}\right\}
$$

That is,

$$
K=\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)^{-1} \cdot\left\{\left(\begin{array}{rr}
\mu & 0 \\
0 & \nu
\end{array}\right): \mu, \nu \in \mathbb{H}^{1}\right\} \cdot\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)=\left\{\left(\begin{array}{ll}
\frac{\mu+\nu}{2} & \frac{-\mu+\nu}{2} \\
\frac{-\mu+\nu}{2} & \frac{\mu+\nu}{2}
\end{array}\right): \mu, \nu \in \mathbb{H}^{1}\right\}
$$

## [1.2.1] Claim:

$$
K \cap P=K \cap M=\left\{\begin{array}{cl} 
\pm 1_{2} & \left(\text { for } G=S L_{2}(\mathbb{R})\right) \\
\left\{\left(\begin{array}{cc}
\mu & 0 \\
0 & \mu^{-1}
\end{array}\right): \mu \in \mathbb{C}^{\times},|\mu|=1\right\} & \left(\text { for } G=S L_{2}(\mathbb{C})\right) \\
\left\{\left(\begin{array}{cc}
\mu & 0 \\
0 & \mu
\end{array}\right): \mu \in \mathbb{H}^{1}\right\} & \left(\text { for } G=S p_{1,1}^{*}\right) \\
\left\{\left(\begin{array}{cc}
\mu & 0 \\
0 & \nu
\end{array}\right): \mu, \nu \in \mathbb{H}^{1}\right\} & \left(\text { for } G=S L_{2}(\mathbb{H})\right)
\end{array}\right.
$$

Proof: In all but the third case, this follows from the description of $K$. For example, for $G=S L_{2}(\mathbb{R})$ and $K=S O_{2}(\mathbb{R})$, take $p=\left(\begin{array}{ll}a & b \\ 0 & a^{-1}\end{array}\right) \in P$ and examine the relation $p^{\top} p=1_{2}$ for $p$ to be in $K$ :

$$
\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)=p^{\top} p=\left(\begin{array}{ll}
a & 0 \\
b & a^{-1}
\end{array}\right)\left(\begin{array}{ll}
a & b \\
0 & a^{-1}
\end{array}\right)=\left(\begin{array}{cc}
a^{2} & \left(a+a^{-1}\right) b \\
\left(a+a^{-1}\right) b & b^{2}+a^{-2}
\end{array}\right)
$$

From the upper-left entry, $a= \pm 1$. From the off-diagonal entries, $b=0$. The arguments for $S L_{2}(\mathbb{C})$ and $S L_{2}(\mathbb{H})$ are similar. For $S p_{1,1}^{*}$, comparison to the coordinates that diagonalize $K \approx S p_{1}^{*} \times S p_{1}^{*}$ gives

$$
\begin{aligned}
\left\{\left(\begin{array}{cc}
k_{1} & 0 \\
0 & k_{2}
\end{array}\right): k_{1}, k_{2} \in \mathbb{H}^{1}\right\}= & \left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right) K\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)^{-1} \ni\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)\left(\begin{array}{cc}
a & b \\
0 & \left(a^{*}\right)^{-1}
\end{array}\right)\left(\begin{array}{rr}
1 & -1 \\
1 & 1
\end{array}\right)^{-1}= \\
& =\frac{1}{2}\left(\begin{array}{ll}
a+\left(a^{*}\right)^{-1}-b & a-\left(a^{*}\right)^{-1}+b \\
a-\left(a^{*}\right)^{-1}-b & a+\left(a^{*}\right)^{-1}+b
\end{array}\right)
\end{aligned}
$$

For example, adding the elements of the bottom row gives $a=k_{2} \in \mathbb{H}^{1}$, and also $\left(a^{*}\right)^{-1}=a$. From either off-diagonal entry, $b=0$.

In all four cases, the same discussion gives $M=A^{+} \cdot(P \cap K)=A^{+} \cdot(M \cap K)$.
The following will be essential in [7.1]:
[1.2.2] Claim: (Cartan decomposition) $G=K A^{+} K$.

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Proof: First, treat $G=S L_{2}(\mathbb{R})$. Prove that every $g \in G$ can be written as $g=s k$ with $s^{\top}=s$ and $s$ positive-definite. To find such $s$, assume for the moment that it exists, and consider

$$
g \cdot g^{\top}=(s k) \cdot(s k)^{\top}=s k \cdot k^{-1} s=s^{2}
$$

Certainly $g g^{\top}$ is symmetric and positive-definite, so having a positive-definite symmetric square root of positive-definite symmetric $t$ would produce $s$. Such $t$ gives a positive, symmetric operator on $\mathbb{R}^{2}$, which by the spectral theorem has an orthonormal basis of eigenvalues. That is, there is $h \in K$ such that $h t h^{\top}=\delta$ is diagonal, necessarily with positive diagonal entries. With $\delta^{\frac{1}{2}}$ be the positive diagonal square root of $\delta$,

$$
\left(h^{\top} \delta^{\frac{1}{2}} h\right)^{2}=h^{\top} \delta^{\frac{1}{2}} h \cdot h^{\top} \delta^{\frac{1}{2}} h=h^{\top} \delta^{\frac{1}{2}} \cdot \delta^{\frac{1}{2}} h=h^{\top} \delta h=t
$$

Thus, take $s=h^{\top} \delta^{\frac{1}{2}} h$, and every $g \in G$ can be written as $g=k s$. Indeed, we have more:

$$
g=k s=k \cdot h^{\top} \delta^{\frac{1}{2}} h=\left(k \cdot h^{\top}\right) \cdot \delta^{\frac{1}{2}} \cdot h \in K \cdot A^{+} \cdot K
$$

giving the claim in this case. The cases of $G=S L_{2}(\mathbb{C})$ is similar, using $g=s k$ with $s=s^{*}$ hermitian positive-definite and $k^{*}=k^{-1} \in K$, invoking the spectral theorem for hermitian positive-definite operators. The same argument succeeds for $G=S L_{2}(\mathbb{H})$ with quaternion conjugation replacing complex, with a suitably adapted spectral theorem for $s \in G L_{2}(\mathbb{H})$ with $s^{*}=s$ and $x^{*} s c$ real and positive for all non-zero 2-by-1 quaternion matrices $x$. ${ }^{[4]}$

The case of $G=S p_{1,1}^{*}$ essentially reduces to the case of $S L_{2}(\mathbb{H})$, as follows. Since $g^{*} S g=S$, $S g S^{-1}=\left(g^{*}\right)^{-1}$. Anticipating the Cartan decomposition $g=s k$, from $g g^{*}=s s^{*}=s^{2}$, by the quaternionic version of the spectral theorem, there is $k \in S p_{2}^{*}$ such that $k^{-1} g g^{*} k=\Lambda$ with $\Lambda$ positive real diagonal. We want to adjust $k$ to be in $S p_{1,1}^{*} \cap S p_{2}^{*}$, while preserving the property $k^{-1} g g^{*} k=\Lambda$. Unless $g g^{*}$ is scalar, the diagonal entries are distinct. By $S g S^{-1}=\left(g^{*}\right)^{-1}$ and $S g^{*} S^{-1}=g^{-1}$ for $g \in G$,

$$
\begin{gathered}
\Lambda^{-1}=\left(\Lambda^{*}\right)^{-1}=S \Lambda S^{-1}=S\left(k^{-1} g g^{*} k\right) S^{-1}=\left(S k S^{-1}\right)^{-1} \cdot S g g^{*} S^{-1} \cdot S k S^{-1} \\
=\left(S k S^{-1}\right)^{-1} \cdot\left(g g^{*}\right)^{-1} \cdot S k S^{-1}
\end{gathered}
$$

Inverting gives $\Lambda=\left(S k S^{-1}\right)^{-1} \cdot g g^{*} \cdot S k S^{-1}$. Also $\Lambda=k^{-1} g g^{*} k$, so

$$
\left(S k S^{-1}\right) \cdot \Lambda \cdot\left(S k S^{-1}\right)^{-1}=g g^{*}=k \cdot \Lambda \cdot k^{-1}
$$

That is, $k^{-1} \cdot S k S^{-1}$ commutes with $\Lambda$, and $\delta=k^{-1} \cdot S k S^{-1}$ is at worst diagonal:

$$
S k S^{-1}=k \cdot \delta=k \cdot\left(\begin{array}{ll}
a & 0 \\
0 & d
\end{array}\right)
$$

Since $\delta \in S p_{2}^{*}$, $a \cdot \bar{a}=1$ and $d \cdot \bar{d}=1$. To preserve $k^{-1} g g^{*} k=\Lambda$, to adjust $k$ to be in $K=S p_{2}^{*} \cap S p_{1,1}^{*}$, adjust $k$ by diagonal matrices $\varepsilon$ in $S p_{2}^{*}$. The condition for $k \varepsilon$ to be in $K$ is

$$
(k \cdot \varepsilon)=\left((k \varepsilon)^{*}\right)^{-1}=S(k \varepsilon) S^{-1}=S k S^{-1} \cdot S \varepsilon S^{-1}=k \cdot \delta \cdot S \varepsilon S^{-1}
$$

so take $\varepsilon=S^{-1} \delta S$. The rest of the argument runs as in the first three cases.
[4] In all three of these cases, a Rayleigh-Ritz approach gives a sufficient spectral theorem, as follows. Let $F$ be $\mathbb{R}$, $\mathbb{C}$, or $\mathbb{H}$. Let $\langle x, y\rangle=y^{*} x$ for 2-by-1 matrices $x, y$ over $F$. Let $T: F^{2} \rightarrow F^{2}$ be right $F$-linear, and positive hermitian in the sense that $\langle T x, x\rangle$ is positive, real for $x \neq 0$. Then $x$ with $\langle x, x\rangle=1$ maximizing $\langle T x, x\rangle$ is an eigenvector for $T$. For non-scalar $T$, the unit vector $y$ minimizing $\langle T y, y\rangle$ is an eigenvector for $T$ orthogonal to $x$. Letting $k$ be the matrix with columns $x, y$, the conjugated matrix $k^{-1} T k$ is diagonal.

### 1.3 Iwasawa decomposition $G=P K=N A^{+} K$

The subgroups $P$ and $K$ are not normal in $G$, so the Iwasawa decompositions $G=P K=\{p k: p \in P, k \in$ $K\}$ do not express $G$ as a product group. Nevertheless, these decompositions are essential.
[1.3.1] Claim: (Iwasawa decomposition) $G=P K=N A^{+} K$. In particular, the map $N \times A^{+} \times K \longrightarrow G$ by $n \times a \times k \longrightarrow n a k$ is an injective set map (and is a diffeomorphism).
Proof: For $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G$, in the easy case that $c=0$, then $g \in P$. In all cases, once we have $g=n m \in P$, we can adjust $g$ on the right by $M \cap K$ to put the Levi component $m$ into $A^{+}$.

One approach is to think of right multiplication by $K$ as rotating the lower row $(c d)$ of $g \in G$ to put it into the form $(0 *)$ of the lower row of an element of $P$. For $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G=S L_{2}(\mathbb{R})$ : right multiplication by the explicit element

$$
k=\left(\begin{array}{cc}
\frac{d}{\sqrt{c^{2}+d^{2}}} & \frac{c}{\sqrt{c^{2}+d^{2}}} \\
\frac{-c}{\sqrt{c^{2}+d^{2}}} & \frac{d}{\sqrt{c^{2}+d^{2}}}
\end{array}\right) \in K=\operatorname{SO}_{2}(\mathbb{R})
$$

puts $g k \in P$ :

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left(\begin{array}{cc}
\frac{d}{\sqrt{c^{2}+d^{2}}} & \frac{c}{\sqrt{c^{2}+d^{2}}} \\
\frac{-c}{\sqrt{c^{2}+d^{2}}} & \frac{d}{\sqrt{c^{2}+d^{2}}}
\end{array}\right)=\left(\begin{array}{ll}
* & * \\
0 & *
\end{array}\right)
$$

Similarly, for $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G=S L_{2}(\mathbb{C})$, right multiplication by

$$
k=\left(\begin{array}{cc}
\frac{d}{\sqrt{|c|^{2}+|d|^{2}}} & \frac{\bar{c}}{\sqrt{|c|^{2}+|d|^{2}}} \\
\frac{-c}{\sqrt{|c|^{2}+|d|^{2}}} & \frac{\bar{d}}{\sqrt{|c|^{2}+|d|^{2}}}
\end{array}\right) \in K=S U(2)
$$

gives $g k \in P$. Likewise, for $G=S L_{2}(\mathbb{H})$, nearly the same explicit expression as for $S L_{2}(\mathbb{C})$ succeeds, with complex conjugation replaced by quaternion conjugation, accommodating the non-commutativity: [5]

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left(\begin{array}{cc}
\frac{c^{-1} d}{\sqrt{1+\left|c^{-1} d\right|^{2}}} & \frac{1}{\sqrt{1+\left|c^{-1} d\right|^{2}}} \\
\frac{-1}{\sqrt{1+\left|c^{-1} d\right|^{2}}} & \frac{\frac{c^{-1} d}{\sqrt{1+\left|c^{-1} d\right|^{2}}}}{\sqrt{ }}
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
0 & *
\end{array}\right) \in P
$$

For $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G=S p_{1,1}^{*}$, we hope that a matrix $k$ of a similar form lies in $K \approx S p_{1}^{*} \times S p_{1}^{*}$, and then $g k \in P$. To be sure that the defining relation for $S p_{1,1}^{*}$ is fulfilled, use the more explicit coordinates

$$
K=\left\{\left(\begin{array}{cc}
\frac{\mu+\nu}{2} & \frac{-\mu+\nu}{2} \\
\frac{-\mu+\nu}{2} & \frac{\mu+\nu}{2}
\end{array}\right): \mu, \nu \in \mathbb{H}^{1}\right\}
$$

To reduce the issue to more manageable pieces, left multiply $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ by $\left(\begin{array}{cc}c^{*} & 0 \\ 0 & c^{-1}\end{array}\right)$ to make $c=1$. As earlier, $g^{*} S g=S$ implies $g S g^{*}=S$, so $c \bar{d}+d \bar{c}=0$, and with $c=1$ we have $d+\bar{d}=0$. Also, $|1+d|^{2}=1+|d|^{2}$.
[5] This explicit element lies in the connected component of $S p_{2}^{*}$ containing 1, so this argument for the Iwasawa decomposition is complete whether or not we have verified that $S p_{2}^{*} \subset S L_{2}(\mathbb{H})$.

Thus, with $\mu=\frac{d+1}{|d+1|}$ and $\nu=\frac{d-1}{|d-1|}, K$ contains

$$
\left(\begin{array}{cc}
\frac{\mu+\nu}{2} & \frac{-\mu+\nu}{2} \\
\frac{-\mu+\nu}{2} & \frac{\mu+\nu}{2}
\end{array}\right)=\left(\begin{array}{cc}
\frac{d}{\sqrt{1+|d|^{2}}} & \frac{-1}{\sqrt{1+|d|^{2}}} \\
\frac{-1}{\sqrt{1+|d|^{2}}} & \frac{d}{\sqrt{1+|d|^{2}}}
\end{array}\right)
$$

Then $g k \in P$, giving the Iwasawa decomposition in this case. In all cases, the fact that $N \cap A^{+}=\{1\}$ and $N A^{+} \cap K=\{1\}$ proves the injectivity of the multiplication $n \times a \times k \rightarrow n a k$.

The following assertion is a generalization of the standard fact that

$$
\operatorname{Im}(g z)=\frac{y}{|c z+d|^{2}} \quad\left(\text { for } z=x+i y \in \mathfrak{H} \text { and } g=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in S L_{2}(\mathbb{R})\right)
$$

This is the foundation for reduction theory for these examples, that is, for determination of the behavior of images $\gamma \cdot g K$ as $\gamma$ varies in $\Gamma$, as below. Let

$$
a_{y}=\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{y}
\end{array}\right) \quad(\text { with } y>0)
$$

[1.3.2] Claim: For Iwasawa decomposition $g=n a_{y} k$ with $n \in N, y>0$, and $k \in K$, say that $y$ is the height of $g$. In all four cases,

$$
\operatorname{height}\left(g \cdot n_{x} a_{y}\right)=\frac{y}{|c y|^{2}+|c x+d|^{2}} \quad\left(\text { for } g=\left(\begin{array}{cc}
a & b \\
c & d
\end{array}\right) \in G, n_{x} \in N, \text { and } y>0\right)
$$

Proof: This is a direct computation.

$$
g n_{x} a_{y}=\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{y}
\end{array}\right)=\left(\begin{array}{ll}
a \sqrt{y} & \frac{a x+b}{\sqrt{y}} \\
c \sqrt{y} & \frac{c x+d}{\sqrt{y}}
\end{array}\right)=\left(\begin{array}{cc}
* & 0 \\
0 & c \sqrt{y}
\end{array}\right)\left(\begin{array}{cc}
* & * \\
1 & \frac{x+c^{-1} d}{y}
\end{array}\right)
$$

For $G=S L_{2}(\mathbb{R}), S L_{2}(\mathbb{C})$, and $S L_{2}(\mathbb{H})$ with respect compact subgroups $K$, for $D$ in $\mathbb{R}, \mathbb{C}, \mathbb{H}$, respectively,

$$
k=\left(\begin{array}{cc}
\frac{D}{\sqrt{1+|D|^{2}}} & \frac{1}{\sqrt{1+|D|^{2}}} \\
\frac{-1}{\sqrt{1+|D|^{2}}} & \frac{\bar{D}}{\sqrt{1+|D|^{2}}}
\end{array}\right) \in K
$$

In those three cases, letting $D=\frac{x+c^{-1} d}{y}$,

$$
\begin{gathered}
g n_{x} a_{y} \cdot k=\left(\begin{array}{cc}
* & 0 \\
0 & c \sqrt{y}
\end{array}\right)\left(\begin{array}{cc}
* & * \\
1 & \frac{x+c^{-1} d}{y}
\end{array}\right)\left(\begin{array}{cc}
\frac{D}{\sqrt{1+|D|^{2}}} & \frac{1}{\sqrt{1+|D|^{2}}} \\
\frac{-1}{\sqrt{1+|D|^{2}}} & \frac{\bar{D}}{\sqrt{1+|D|^{2}}}
\end{array}\right) \\
=\left(\begin{array}{cc}
* & 0 \\
0 & c \sqrt{y}
\end{array}\right)\left(\begin{array}{cc}
* & * \\
0 & \frac{1+|D|^{2}}{\sqrt{1+|D|^{2}}}
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
0 & |c| \sqrt{y} \sqrt{1+|D|^{2}}
\end{array}\right) \cdot\left(\begin{array}{cc}
* & 0 \\
0 & \frac{c}{|c|}
\end{array}\right)
\end{gathered}
$$

noting that $\left(\begin{array}{cc}* & 0 \\ 0 & \frac{c}{|c|}\end{array}\right) \in K$. Simplifying,

$$
|c| \sqrt{y} \sqrt{1+|D|^{2}}=|c| \sqrt{y} \sqrt{1+\left|\frac{x+c^{-1} d}{y}\right|^{2}}=\sqrt{\frac{|c y|^{2}+|c x+d|^{2}}{y}}
$$

Thus, in these three cases,

$$
g n_{x} a_{y} \in N\left(\begin{array}{cc}
\sqrt{y^{\prime}} & 0 \\
0 & 1 / \sqrt{y^{\prime}}
\end{array}\right) K \quad \text { with } \quad y^{\prime}=\frac{y}{|c y|^{2}+|c x+d|^{2}}
$$

## 1. Four small examples

For $G=S p_{1,1}^{*}$, the explicit element of $K$ is slightly different

$$
k=\left(\begin{array}{cc}
\frac{D}{\sqrt{1+|D|^{2}}} & \frac{-1}{\sqrt{1+|D|^{2}}} \\
\frac{-1}{\sqrt{1+|D|^{2}}} & \frac{D}{\sqrt{1+|D|^{2}}}
\end{array}\right) \in K
$$

but the conclusion will be the same: with $D=\frac{x+c^{-1} d}{y}$

$$
g n_{x} a_{y} \cdot k=\left(\begin{array}{cc}
* & 0 \\
0 & c \sqrt{y}
\end{array}\right)\left(\begin{array}{cc}
* & * \\
1 & \frac{x+c^{-1} d}{y}
\end{array}\right)\left(\begin{array}{cc}
\frac{D}{\sqrt{1+|D|^{2}}} & \frac{-1}{\sqrt{1+|D|^{2}}} \\
\frac{-1}{\sqrt{1+|D|^{2}}} & \frac{D}{\sqrt{1+|D|^{2}}}
\end{array}\right)=\left(\begin{array}{cc}
* & 0 \\
0 & c \sqrt{y}
\end{array}\right)\left(\begin{array}{cc}
* & * \\
0 & \frac{-1+D^{2}}{\sqrt{1+|D|^{2}}}
\end{array}\right)
$$

For $S p_{1,1}^{*}$, as in earlier computations, the relation $h^{*} S h=S$ gives $h S h^{*}=S$, so for $h=\left(\begin{array}{cc}* & * \\ 1 & D\end{array}\right)$ we find $D+\bar{D}=0$. That is, $D$ is purely imaginary, so $D^{2}=-|D|^{2}$, and

$$
g n_{x} a_{y} \cdot k=\left(\begin{array}{cc}
* & * \\
0 & -c \sqrt{y} \sqrt{1+|D|^{2}}
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
0 & |c| \sqrt{y} \sqrt{1+|D|^{2}}
\end{array}\right) \cdot\left(\begin{array}{cc}
* & 0 \\
0 & \frac{-c}{|c|}
\end{array}\right)
$$

The remainder of the computation is identical to the other three cases.

### 1.4 Some convenient Euclidean rings

We recall proofs that, just as the ordinary integers are Euclidean, the Gaussian integers $\mathbb{Z}[i]$ and Hurwitz quaternion integers are Euclidean. This will greatly simplify the geometry of quotients $\Gamma \backslash X$ in [1.5.1] by assuring that there is just a single cusp.

Recall the simplest version of Euclidean-ness for a ring $R$ with 1: there is a function $\|\cdot\|: R \rightarrow \mathbb{Z}$ such that $\|r\| \geq 0$ and $\|r\|=0$ implies $r=0$, such that $\left\|r r^{\prime}\right\|=\|r\| \cdot\left\|r^{\prime}\right\|$, and, for every $a \in R$ and every $0 \neq d \in R$, there is $q \in R$ such that $\|a-q d\|<\|d\|$.

Since $\|1\|=\left\|1^{2}\right\|=\|1\| \cdot\|1\|$ and $0<\|1\|$, necessarily $\|1\|=1$. Units $r \in R^{\times}$have $\|r\|=1$, because $r s=1$ gives $\|r\| \cdot\|s\|=\|r s\|=\|1\|=1$, and $\|\cdot\|$ takes non-negative integer values.

Euclidean-ness implies that every left ideal is principal: let $d$ be an element having the smallest norm in a given non-zero left ideal $I$. For any $a \in I$, there is $q \in R$ such that $\|a-q d\|<\|d\|$. Thus, $\|a-q d\|=0$, and $a=q d$.

To show that $R=\mathbb{Z}[i]$ is Euclidean with respect to the square of the usual complex absolute value $\|\cdot\|=|\cdot|^{2}$, for $a \in \mathbb{Z}[i]$ and given $0 \neq d \in \mathbb{Z}[i]$, we need to find $q \in \mathbb{Z}[i]$ such that $\|a-d q\|<\|d\|$. The requirement $\|a-q d\|<\|d\|$ is equivalent to $\|a / d-q\|<1$. Thus, given $a / d \in \mathbb{Q}(i)$, we want $q \in \mathbb{Z}[i]$ within distance-squared 1. With $a / d=u+i v$ with $u, v \in \mathbb{Q}$, taking $u^{\prime}, v^{\prime} \in \mathbb{Z}$ such that $\left|u-u^{\prime}\right| \leq \frac{1}{2}$ and $\left|v-v^{\prime}\right| \leq \frac{1}{2}$ gives the desired $\left\|a / d-\left(u^{\prime}+i v^{\prime}\right)\right\| \leq\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{2}\right)^{2}<1$.

In the rational quaternions $\mathbb{H}_{\mathbb{Q}}=\mathbb{Q}+\mathbb{Q} i+\mathbb{Q} j+\mathbb{Q} k$, the natural choice $\mathbb{Z}+\mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k$ for integers is not optimal. Instead, we use the slightly larger ring of Hurwitz integers:

$$
\mathfrak{o}=(\mathbb{Z}+\mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k)+\mathbb{Z} \cdot \frac{1+i+j+k}{2}
$$

We prove that the Hurwitz integers are Euclidean, using the square of the quaternion norm: $\|\cdot\|=|\cdot|^{2}$. To see that the norm-squared takes integer values on $\mathfrak{o}$, the only possible difficulty might be a denominator of 4, which does not occur, since

$$
(2 a+1)^{2}+(2 b+1)^{2}+(2 c+1)^{2}+(2 d+1)^{2}=0 \bmod 4 \quad(\text { for all } a, b, c, d \in \mathbb{Z})
$$

Given $a \in \mathfrak{o}$ and $0 \neq d \in \mathfrak{o}$, to show that there is $q \in \mathfrak{o}$ such that $\|a-q d\|<\|d\|$ is equivalent to $\left\|a d^{-1}-q\right\|<1$. For $a d^{-1}=x+y i+z j+w k$ with $x, y, z, w \in \mathbb{Q}$, there are $x^{\prime}, y^{\prime}, z^{\prime}, w^{\prime} \in \mathbb{Z}$ differing by at most $\frac{1}{2}$ in absolute value from the respective $x, y, z, w$. However, the resulting estimate

$$
\left\|(x+y i+z j+w k)-\left(x^{\prime}+y^{\prime} i+z^{\prime} j+w^{\prime} k\right)\right\| \leq\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{2}\right)^{2}=1
$$

is insufficient. Nevertheless, being slightly more precise, if $\left|x-x^{\prime}\right|<\frac{1}{2}$ or $\left|y-y^{\prime}\right|<\frac{1}{2}$ or $\left|z-z^{\prime}\right|<\frac{1}{2}$ or $\left|w-w^{\prime}\right|<\frac{1}{2}$, then we do have the desired

$$
\left\|(x+y i+z j+w k)-\left(x^{\prime}+y^{\prime} i+z^{\prime} j+w^{\prime} k\right)\right\|<1
$$

That is, the only case of failure is $\left|x-x^{\prime}\right|=\left|y-y^{\prime}\right|=\left|z-z^{\prime}\right|=\left|w-w^{\prime}\right|=\frac{1}{2}$. Subtracting 1 from $x, y, z, w$ if necessary, without loss of generality $x-x^{\prime}=y-y^{\prime}=z-z^{\prime}=w-w^{\prime}=\frac{1}{2}$. In that case,

$$
(x+y i+z j+w k)-\left(\left(x^{\prime}+y^{\prime} i+z^{\prime} j+w^{\prime} k\right)+\frac{1+i+j+k}{2}\right)=0
$$

proving that the Hurwitz integers are Euclidean. A qualitative version of the Euclidean-ness of o will be useful in one of the proofs of unicuspidality: for $\alpha=x+y i+z j+w k$ with $x, y, z, w \in \mathbb{Q}$, there is $x^{\prime}+y^{\prime} i+z^{\prime} j+w^{\prime} k \in \mathfrak{o}$ such that

$$
\left|(x+y i+z j+w k)-\left(x^{\prime}+y^{\prime} i+z^{\prime} j+w^{\prime} k\right)\right| \leq \frac{\sqrt{13}}{4}
$$

Adjust $\alpha$ by an element of $\mathbb{Z}+\mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k$ so that, without loss of generality, all coefficients are of absolute value at most $\frac{1}{2}$. If any one coefficient is smaller than $1 / 4$, then $|\alpha|^{2} \leq\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{2}\right)^{2}+\left(\frac{1}{4}\right)^{2}=13 / 16$ as desired. When all coefficients are between $1 / 4$ and $1 / 2$ in absolute value, make them all of the same sign by adding or subtracting 1 to either one or two, paying the price that those one or two are of absolute value between $1 / 2$ and $3 / 4$, while the others are still of absolute value between $1 / 4$ and $1 / 2$. Adding or subtracting $(1+i+j+k) / 2$ depending on sign, all coefficients are between $-1 / 4$ and $1 / 4$, and the quaternion norm of the result is at most $\frac{1}{2} \leq \frac{\sqrt{13}}{4}$.

### 1.5 Discrete subgroups $\Gamma \subset G$, reduction theory

We specify discrete ${ }^{[6]}$ subgroups $\Gamma$ of each of the examples $G$, so that $\Gamma \backslash G / K$ has just one cusp, in a sense made precise below. Reduction theory is the exhibition of a simple approximate collection of representatives for the quotient $\Gamma \backslash G / K$ sufficient to understand the most basic geometric features of that quotient. [7] The simple outcome in the present examples, unicuspidality, simplifies meromorphic continuation of Eisenstein series and simplifies the form of the spectral decomposition of the space of square-integrable automorphic forms on $\Gamma \backslash G / K$. The four cases are ${ }^{[8]}$

$$
\Gamma=\left\{\begin{array}{lll}
S L_{2}(\mathbb{Z}) & \left(\text { for } G=S L_{2}(\mathbb{R})\right) & \text { (elliptic modular group) } \\
S L_{2}(\mathbb{Z}[i]) & \left(\text { for } G=S L_{2}(\mathbb{C})\right) & \text { (a Bianchi modular group) } \\
S p_{1,1}^{*}(\mathfrak{o}) & \left(\text { for } G=S p_{1,1}^{*}\right) \\
S L_{2}(\mathfrak{o}) & \left(\text { for } G=S L_{2}(\mathbb{H})\right)
\end{array}\right.
$$

[6] As usual, a subset $D$ of a topological space $X$ is discrete when every point $x \in D$ has a neighborhood $U$ such that $U \cap D=\{x\}$. The topologies on our groups $G$ are the subspace topologies from the ambient real vector spaces of 2-by-2 real, complex, or quaternion matrices.
[7] In some contexts, the goal of determination of an exact, explicit collection of representatives in $G / K$ for the quotient $\Gamma \backslash G / K$ is given high priority. A precise collection of representative is often called a fundamental domain. However, in general determination of an explicit fundamental domain is infeasible. Fortunately, it is also inessential. [8] The elliptic modular group has its origins in dim antiquity. [Picard 1883] and [Picard 1884] looked at similar subgroups of small non-compact unitary groups. L. Bianchi [Bianchi 1892] looked at a family of discrete subgroups of $S L_{2}(\mathbb{C})$, such as $S L_{2}(\mathbb{Z}[i])$. W. de Sitter proposed a model of space-time in which the cosmological constant dominates and matter is negligible, with symmetry group $S O(4,1)$, and $S p_{1,1}^{*}$ is a two-fold cover of $S O(4,1)$. No automorphic forms directly entered his work, but his attention to specific groups, as in the more theoretical work of [Bargman 1947] and [Wigner 1939], provided examples which eventually were appreciated for their illustration of phenomena with mathematical significance beyond physics itself. [Hurwitz 1898] studied the quaternion integers $\mathfrak{o}$ which bear his name. See also [Hurwitz 1919] and [Conway-Smith 2003].

## 1. Four small examples

where $S p_{1,1}^{*}(\mathfrak{o})$ and $S L_{2}(\mathfrak{o})$ denote the elements of $S p_{1,1}^{*}$ and $S L_{2}(\mathbb{H})$ with entries in the ring of Hurwitz integers $\mathfrak{o}$.

In all examples, $\Gamma \cap P=(\Gamma \cap M) \cdot(\Gamma \cap N)$. We have

$$
\Gamma \cap N=\left\{( \begin{array} { l l } 
{ 1 } & { x } \\
{ 0 } & { 1 }
\end{array} ) \text { where } \left\{\begin{array}{ll}
x \in \mathbb{Z} & \left(\text { for } \Gamma=S L_{2}(\mathbb{Z})\right) \\
x \in \mathbb{Z}[i] & \left(\text { for } \Gamma=S L_{2}(\mathbb{Z}[i])\right) \\
x \in \mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k & \left(\text { for } \Gamma=S p_{1,1}^{*}(\mathfrak{o})\right) \\
x \in \mathfrak{o} & \left(\text { for } \Gamma=S L_{2}(\mathbb{H})\right)
\end{array}\right.\right.
$$

As in the discussion of Euclidean-ness, the quotients $(\Gamma \cap N) \backslash N$ have (redundant) representatives

$$
(\Gamma \cap N) \backslash N=\left\{\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)\right\} \text { where } \begin{cases}x \in \mathbb{R},|x| \leq \frac{1}{2} & \left(\text { for } \Gamma=S L_{2}(\mathbb{Z})\right) \\
x \in \mathbb{C},|x| \leq \frac{1}{\sqrt{2}} & \left(\text { for } \Gamma=S L_{2}(\mathbb{Z}[i])\right) \\
x=a i+b j+c k,|x| \leq \frac{\sqrt{3}}{2} & \left(\text { for } \Gamma=S p_{1,1}^{*}(\mathfrak{o})\right) \\
x \in \mathbb{H},|x| \leq \frac{\sqrt{13}}{4} & \left(\text { for } \Gamma=S L_{2}(\mathbb{H})\right)\end{cases}
$$

In particular, $(\Gamma \cap N) \backslash N$ is compact. We have

$$
\Gamma \cap M=\left\{\left(\begin{array}{ll}
a & 0 \\
0 & d
\end{array}\right)\right\} \text { where } \begin{cases}a=d^{-1} \in \mathbb{Z}^{\times}=\{ \pm 1\} & \left(\text { for } \Gamma=S L_{2}(\mathbb{Z})\right) \\
a=d^{-1} \in \mathbb{Z}[i]^{\times}=\{ \pm 1, \pm i\} & \left(\text { for } \Gamma=S L_{2}(\mathbb{Z}[i])\right) \\
a=\left(d^{*}\right)^{-1} \in \mathfrak{o}^{\times} & \left(\text {for } \Gamma=S p_{1,1}^{*}(\mathfrak{o})\right) \\
a, d \in \mathfrak{o}^{\times} & \left(\text {for } \Gamma=S L_{2}(\mathbb{H})\right)\end{cases}
$$

The groups of units $\mathbb{Z}^{\times}$and $\mathbb{Z}[i]^{\times}$are well-known, and finite. The group $\mathfrak{o}^{\times}$is also finite, but less trivial. As noted earlier, $\alpha \in \mathfrak{o}^{\times}$implies $|\alpha|=1$. Certainly $\mathfrak{o} \subset \frac{1}{2} \cdot(\mathbb{Z}+\mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k)$ and

$$
\frac{a^{2}+b^{2}+c^{2}+d^{2}}{4}=\left|\frac{a+b i+c j+d k}{2}\right|^{2} \leq 1
$$

implies $|a| \leq 2,|b| \leq 2,|c| \leq 2$, and $|d| \leq 2$, giving a crude bound on the number of possibilities for integers $a, b, c, d$.

For compact $C \subset N$, a standard Siegel set is a subset of $G$ of the form

$$
\mathfrak{S}_{t, C}=\left\{n a_{y} k: n \in C, k \in K, y \geq t\right\}
$$

This is essentially a half-infinite rectangle right-multiplied by $K$. On other occasions, a Siegel set is construed as a subset of the quotient $(\Gamma \cap N) \backslash G$, or as a subset of $G / K$. These distinctions are inessential. Let $\Gamma_{\infty}=P \cap \Gamma$. Left multiplication by $N$ does not change heights on $G$. Since $\Gamma_{\infty} \subset N \times(M \cap \Gamma)$ and $M \cap \Gamma \subset M \cap K$, left multiplication by $\Gamma_{\infty}$ does not change heights. Siegel sets are a simple type of set among which, as it turns out, we can find approximate sets of representatives for the quotient $\Gamma \backslash G / K$. That is, reduction theory for these examples is relatively simple:
[1.5.1] Theorem: For all four examples, $\Gamma \backslash G$ is unicuspidal, in the sense that there is $t>0$ and compact $C \subset N$ such that a single Siegel set covers $G$ :

$$
\bigcup_{\gamma \in \Gamma} \gamma \cdot \mathfrak{S}_{t, C}=G
$$

Proof: For fixed $x, y$, the function $q(c, d)=q_{x, y}(c, d)=|c x+d|^{2}+|c y|^{2}$ is a homogeneous real-valued quadratic polynomial function on $\mathbb{R} \oplus \mathbb{R}, \mathbb{C} \oplus \mathbb{C}$, or $\mathbb{H} \oplus \mathbb{H}$, in the respective cases, with $y>0$ and appropriate $x$. It is positive definite: $q(c, d)=0$ implies $c=0=d$. Thus, $q(c, d)$ is comparable to $|c|^{2}+|d|^{2}$ : there are positive constants $A, B$ depending on $x, y$ such that

$$
A \cdot\left(|c|^{2}+|d|^{2}\right) \leq q(c, d) \leq B \cdot\left(|c|^{2}+|d|^{2}\right)
$$

The number of points $(c, d)$ in a lattice $\mathbb{Z} \oplus \mathbb{Z}, \mathbb{Z}[i] \oplus \mathbb{Z}[i]$, or $\mathfrak{o} \oplus \mathfrak{o}$ inside a ball of finite radius is finite. In particular, in the orbit $\Gamma \cdot n_{x} a_{y}$ there are only finitely-many values height $\left(\gamma \cdot n_{x} a_{y}\right)$ above a given bound $t>0$. In particular, the supremum of these heights is attained. Thus, every $\Gamma$-orbit contains (at least one) $n_{x} a_{y}$ of maximum height, and $|c x+d|^{2}+|c y|^{2} \geq 1$ for all lower rows $(c d)$ of $\gamma \in \Gamma$. In particular, with $c=1$ and $d=0,|x|^{2}+|y|^{2} \geq 1$.

Thus, given $n_{x} a_{y} K \in G / K$, adjust on the left by $\gamma \in \Gamma$ so that $\gamma n_{x} a_{y} K$ is (one of) the highest in its orbit on $G / K$. In particular, this makes $|x|^{2}+|y|^{2} \geq 1$. From the specific estimates on parameters $\xi$ for representatives $n_{\xi}$ of $(\Gamma \cap N) \backslash N$, in all cases there is $0<t<1$ such that $|\xi| \leq 1-t$ for all representatives. Thus, if $|x|>1-t$, further adjust on the left by $\gamma \in \Gamma \cap N$ so that $|x| \leq 1-t$, without altering the height. Thus, the new $n_{x} a_{y}$ is still among the highest in its orbit, and $|x|^{2}+|y|^{2} \geq 1$ still holds. Thus

$$
|y|^{2} \geq 1-|x|^{2} \geq 1-(1-t)^{2}=t(2-t) \geq t
$$

Thus, every $\Gamma$-orbit has a representative in the Siegel set $\mathfrak{S}_{t, C}$, where $C=\left\{n_{x} \in N:|x| \leq 1-t\right\}$.
The following part of reduction theory is more technical, but essential.
[1.5.2] Theorem: For given $t, t^{\prime}>0$ and compact subsets $C, C^{\prime}$ of $N$, there are only finitely-many $\gamma \in \Gamma$ such that $\mathfrak{S}_{t, C} \cap \gamma \cdot \mathfrak{S}_{t^{\prime}, C^{\prime}} \neq \phi$. Further, given $t>0$, for sufficiently large $t^{\prime}>0, \mathfrak{S}_{t, C} \cap \gamma \mathfrak{S}_{t^{\prime}, C^{\prime}} \neq \phi$ implies $\gamma \in \Gamma_{\infty}$.
Proof: Continue in the context of the proof of the previous theorem. Given $t>0$, take $t^{\prime}>1 / t$. For $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \notin \Gamma_{\infty}, c \neq 0$, so $|c| \geq 1$. For $y \geq t^{\prime}$ and arbitrary $x$,

$$
\operatorname{height}\left(\gamma \cdot n_{x} a_{y}\right)=\frac{y}{|c x+d|^{2}+|c y|^{2}} \leq \frac{y}{|c|^{2} \cdot y^{2}} \leq \frac{1}{1 \cdot y} \leq \frac{1}{t^{\prime}}<t
$$

Thus, $\mathfrak{S}_{t, C} \cap \gamma \cdot \mathfrak{S}_{t^{\prime}, C^{\prime}} \neq \phi$ implies $\gamma \in \Gamma_{\infty}$ for such $t, t^{\prime}$.
For arbitrary $0<t \leq t^{\prime}$, to show finiteness of the set of $\gamma$ so that $\mathfrak{S}_{t, C} \cap \gamma \mathfrak{S}_{t^{\prime}, C^{\prime}} \neq \phi$, take $t^{\prime \prime}$ strictly larger than $t, t^{\prime}, 1 / t$, and $1 / t^{\prime}$. The two sets

$$
\Omega=\left\{n_{x} a_{y} k: n \in C, k \in K, t \leq y \leq t^{\prime \prime}\right\} \quad \text { and } \quad \Omega^{\prime}=\left\{n_{x} a_{y} k: n \in C^{\prime}, k \in K, t^{\prime} \leq y \leq t^{\prime \prime}\right\}
$$

are compact, and $\mathfrak{S}_{t, C}=\mathfrak{S}_{t^{\prime \prime}, C} \cup \Omega$ and $\mathfrak{S}_{t^{\prime}, C^{\prime}}=\mathfrak{S}_{t^{\prime \prime}, C^{\prime}} \cup \Omega^{\prime}$. For the asserted finiteness, it suffices to treat the pieces separately.

By the previous paragraph, $\mathfrak{S}_{t, C}=\mathfrak{S}_{t^{\prime \prime}, C} \cup \Omega$ meets $\gamma \mathfrak{S}_{t^{\prime \prime}, C^{\prime}}$ only for $\gamma \in \Gamma_{\infty}$. Since $\Gamma_{\infty}=(\Gamma \cap N) \cdot(\Gamma \cap M)$ and $\Gamma \cap M$ is finite, it suffices to consider $\gamma=n_{x} \in \Gamma \cap N$. In that case, $\gamma \mathfrak{S}_{t^{\prime \prime}, C^{\prime}}=\mathfrak{S}_{t^{\prime \prime}, C^{\prime}+x}$. By the Iwasawa decomposition, $\mathfrak{S}_{t, C} \cap \mathfrak{S}_{t^{\prime \prime}, C^{\prime}+x} \neq \phi$ if and only if $C \cap\left(C^{\prime}+x\right) \neq \phi$. Equivalently, $x \in C-C^{\prime}$ and $x$ is in a lattice $\mathbb{Z}, \mathbb{Z}[i], \mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k$, or $\mathfrak{o}$, respectively. The set $C-C^{\prime}$ of element-wise differences is compact. Either by the lemma below, or by more elementary considerations, the set of such $x$ is finite.

Similarly $\mathfrak{S}_{t^{\prime \prime}, C}$ meets $\gamma \Omega^{\prime}$ only for $\gamma \in \Gamma_{\infty}$, and there are only finitely many possibilities.
The interaction of $\Omega$ and $\Omega^{\prime}$ is subtler. For $\gamma$ such that $\Omega \cap \gamma \Omega^{\prime} \neq \phi$, there are $\omega \in \Omega$ and $\omega^{\prime} \in \Omega^{\prime}$ such that $\omega=\gamma \omega^{\prime}$. That is,

$$
\gamma=\omega\left(\omega^{\prime}\right)^{-1} \cap \Gamma \subset \Omega \cdot \Omega^{\prime-1} \cap \Gamma
$$

where $\Omega^{\prime-1}$ is element-wise inversion. Inversion and multiplication are continuous maps ${ }^{[9]} G \rightarrow G$ and $G \times G \rightarrow G$, so they map compacts to compacts, so $\Omega \Omega^{\prime-1}$ is compact. By the following lemma, such a set is finite.
[9] The multiplication of real, complex, or quaternion matrices is polynomial in the entries, so is continuous. The continuity of inversion can be seen via the explicit formula in terms of determinants of minors over a field $k$, for example: for $g \in S L_{n}(k)$, letting $A_{i j}$ be the $(n-1)$-by- $(n-1)$ matrix obtained by deleting the $i^{t h}$ row and $j^{t h}$ column, $(-1)^{i+j} \operatorname{det} A_{i j}$ is the $i j^{t h}$ entry of $g^{-1}$.

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[1.5.3] Lemma: The intersection $\Omega \cap \Gamma$ of a compact subset $\Omega$ of a topological group ${ }^{[10]} G$ and a discrete subgroup $\Gamma \subset G$ is finite.
Proof: First, we prove that $\Gamma$ is closed, which would not necessarily hold for a discrete subset. For us, a topological group is locally compact, Hausdorff, and countably-based. Let $U$ be a neighborhood of 1 such that $U \cap \Gamma=\{1\}$. By continuity of inversion and multiplication, there is a neighborhood $U_{1}$ of 1 such that $U_{1}^{-1} \cdot U_{1} \subset U$. For $g \notin \Gamma$ but $g$ in the closure of $\Gamma$ in $G$, the neighborhood $g U_{1}$ of $g$ contains infinitely-many elements of $\Gamma$. For $\gamma \neq \delta$ two such,

$$
1 \neq \gamma^{-1} \cdot \delta \in\left(g U_{1}\right)^{-1} \cdot\left(g U_{1}\right)=U_{1}^{-1} \cdot U_{1} \subset U
$$

contradiction. Thus, $\Gamma$ is closed in $G$.
In a Hausdorff space $G$, a compact subset $C$ is closed, so $C \cap \Gamma$ is closed. A closed subset of a compact set is compact. Thus, $C \cap \Gamma$ is compact, and it is (still) discrete. Discrete compact sets are finite, proven as follows. For each $\gamma \in C \cap \Gamma$, let $N_{\gamma}$ be a neighborhood of $\gamma$ in $G$ containing no other element of $\Gamma$. The open cover $\left\{N_{\gamma}: \gamma \in C \cap \Gamma\right\}$ of $C \cap \Gamma$ has a finite subcover $N_{\gamma_{1}} \cup \ldots \cup N_{\gamma_{n}}$. Since $N_{\gamma_{i}} \cap \Gamma=\gamma_{i}$, necessarily $C \cap \Gamma$ is finite.

### 1.6 Invariant measures, invariant Laplacians

Proofs of the assertions in this section require substantial preparation, and succeed for very general reasons, so are postponed to [5.2] and [4.2]. In all four examples, the subgroup $N A^{+}$of $P$ is transitive on $X=G / K$, by the Iwasawa decomposition $G=N A^{+} K$, giving a bijection

$$
X=G / K=\left(N A^{+} K\right) / K \approx\left(N A^{+}\right) /\left(N A^{+} \cap K\right)=N A^{+}
$$

In coordinate-independent formulations, notation $x \in X$ is reasonable, despite the fact that, somewhat incompatibly, when convenient we will use coordinates

$$
(x, y) \longrightarrow n_{x} a_{y}=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{y}
\end{array}\right)
$$

on $N A^{+} \approx G / K=X$, as above, with $y>0$, and $x$ in $\mathbb{R}, \mathbb{C} \approx \mathbb{R}^{2}, \mathbb{R} i+\mathbb{R} j+\mathbb{R} k \approx \mathbb{R}^{3}$, or $\mathbb{H} \approx \mathbb{R}^{4}$, respectively. These coordinates $(x, y) \in \mathbb{R}^{\ell-1} \times(0,+\infty)$ are standard coordinates on real hyperbolic $\ell$-space, with $\ell=2,3,4,5$.

Although we will eventually need the right translation action of $G$ on $G$ and on functions on $G$, for the moment we are considering the quotient $X=G / K$. Since $K$ is not a normal subgroup, there is no sensible right translation action of $G$ on $X=G / K$, only the left translation action $g \cdot\left(g_{o} K\right)=\left(g g_{o}\right) K$.

The group $G$ acts on the collection $C_{c}^{o}(X)$ of continuous, compactly-supported functions on $X=G / K$ by left translation

$$
L_{g} f(x)=f\left(g^{-1} x\right)
$$

with the inverse inserted to have the associativity

$$
L_{g_{1}} L_{g_{2}} f=L_{g_{1} g_{2}} f
$$

A $G$-invariant measure/integral $\mu$ on the quotient $X=G / K$ is characterized by the property

$$
\int_{X} L_{g} f d \mu=\int_{X} f d \mu \quad\left(\text { for all } g \in G, f \in C_{c}^{o}(X)\right)
$$

[10] As usual, a topological group is a locally compact, Hausdorff topological space $G$ with a countable basis, and so that the inverse map $g \rightarrow g^{-1}$ and multiplication $g_{1} \times g_{2} \rightarrow g_{1} g_{2}$ are continuous.

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In [14.4] we will see that such an invariant measure/integral is unique up to scalar multiples, and is given in the $x, y \rightarrow n_{x} a_{y}$ coordinates by

$$
f \longrightarrow \int_{0}^{\infty} \int_{\mathbb{R}^{\ell-1}} f(x, y) \frac{d x d y}{|y|^{\ell}} \quad(\text { where } \ell=2,3,4,5, \text { respectively) }
$$

[1.6.1] Corollary: (of reduction theory) The invariant volume of $\Gamma \backslash X=\Gamma \backslash G / K$ is finite.
Proof: Since there is a Siegel set $\mathfrak{S}_{t, C}$ that surjects to $\Gamma \backslash X$ for some compact $C \subset N$ and some $t>0$, it suffices to show that the invariant measure of a Siegel set is finite. In the $x, y \rightarrow n_{x} a_{y}$ coordinates,

$$
\int_{\mathfrak{S}_{t, C}} 1 d \mu=\int_{t}^{\infty} \int_{\left\{x: n_{x} \in C\right\}} 1 \frac{d x d y}{y^{\ell}}=(N \text {-volume of } C) \cdot \int_{t}^{\infty} \frac{d y}{y^{\ell}}
$$

where $\ell=2,3,4,5$ in the respective examples. Each of these integrals is finite.
Test functions $C_{c}^{\infty}(G / K)$ should be compactly-supported, infinitely-differentiable functions on $G / K$. However, while these groups $G$ are smooth manifolds, it is less clear whether $G / K$ is a smooth manifold. This potential issue is rendered irrelevant by taking

$$
C_{c}^{\infty}(G / K)=\{\text { right } K \text {-invariant test functions on } G\}=C_{c}^{\infty}(G)^{K}
$$

where right $K$-invariance means $f(g k)=f(g)$ for all $g \in G$ and $k \in K$. The invariance of a $G$-invariant Laplacian $\Delta$ on the quotient $X=G / K$ is the property

$$
\Delta\left(L_{g} f\right)=L_{g}(\Delta f) \quad\left(\text { for all } g \in G, f \in C_{c}^{\infty}(X)\right)
$$

In [4.2] we will see that such a Laplacian is essentially canonical, and in the $x, y$ coordinates is, up to constants which we might want to adjust later for notational convenience,

$$
\Delta= \begin{cases}y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) & \left(\text { for } G=S L_{2}(\mathbb{R}), x \in \mathbb{R}\right) \\ y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-y \frac{\partial}{\partial y} & \left(\text { for } G=S L_{2}(\mathbb{C}), x=x_{1}+i x_{2} \in \mathbb{C}\right) \\ y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial x_{3}^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-2 y \frac{\partial}{\partial y} & \left(\text { for } G=S p_{1,1}^{*}, x=x_{1} i+x_{2} j+x_{3} k\right) \\ y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial x_{3}^{2}}+\frac{\partial^{2}}{\partial x_{4}^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-3 y \frac{\partial}{\partial y} & \left(\text { for } G=S L_{2}(\mathbb{H}), x=x_{1}+x_{2} i+x_{3} j+x_{4} k\right)\end{cases}
$$

In [6.6] we will see the symmetry property of $\Delta$ :

$$
\int_{X} \Delta f \cdot \bar{F} d \mu=\int_{X} f \cdot \overline{\Delta F} d \mu \quad\left(\text { for } f, F \in C_{c}^{\infty}(X)\right)
$$

Also, the negative-semi-definite property

$$
\int_{X} \Delta f \cdot \bar{f} d \mu \leq 0 \quad\left(\text { for } f \in C_{c}^{\infty}(X)\right)
$$

For $G=S L_{2}(\mathbb{R})$, where $\ell=2$, by chance the first-order term in $y$ in $\Delta$ disappears, the powers of $y$ in $\Delta$ and $\mu$ cancel, and the symmetry property reduces to symmetry of the Euclidean Laplacian, just integration by parts. Although attractive, this coincidence is inessential.

The left $G$-invariance of $\mu$ and $\Delta$ assure that they descend to the quotient $\Gamma \backslash X \approx \Gamma \backslash G / K$. We will use the same symbols for the versions on $\Gamma \backslash X$. As we see in [5.2], uniqueness of invariant measure/integral entails unwinding identities such as

$$
\int_{\Gamma \backslash X}\left(\sum_{\gamma \in \Gamma} L_{\gamma} f\right) d \mu=\int_{\Gamma \backslash X}\left(\sum_{\gamma \in \Gamma} f\left(\gamma^{-1} x\right)\right) d \mu(x)=\int_{X} f d \mu \quad\left(\text { for } f \in C_{c}^{o}(X)\right)
$$

## 1. Four small examples

For the Laplacian,

$$
\Delta f \circ q=\Delta(f \circ q) \quad\left(\text { for } f \in C_{c}^{\infty}(\Gamma \backslash X)\right)
$$

The symmetry and negative semi-definiteness of $\Delta$ descend to $C^{\infty}(\Gamma \backslash X)=C^{\infty}(\Gamma \backslash G)^{K}$, where we take advantage of the fact that $\Gamma \backslash G$ is a smooth manifold. Again, see [6.6] for proofs.

As expected, with invariant measure $\mu$ descended to $\Gamma \backslash X$, the usual hermitian inner product is ${ }^{[11]}$

$$
\langle f, F\rangle_{\Gamma \backslash X}=\int_{\Gamma \backslash X} f \cdot \bar{F} d \mu
$$

with associated norm

$$
|f|_{L^{2}(\Gamma \backslash X)}=\langle f, f\rangle_{\Gamma \backslash X}^{\frac{1}{2}}
$$

As usual, in the characterization

$$
L^{2}(\Gamma \backslash X)=\left\{\text { measurable } f:|f|_{L^{2}(\Gamma \backslash X)}<\infty\right\}
$$

elements of $L^{2}(\Gamma \backslash X)$ are equivalence classes of measurable functions, with equivalence being equality almosteverywhere. In $[6.5]$ we will see that this characterization is equivalent to a characterization as $L^{2}$-completion of test functions $C_{c}^{\infty}(\Gamma \backslash X)$.

### 1.7 Discrete decomposition of $L^{2}(\Gamma \backslash G / K)$ cuspforms

The theorems stated below will be proven later, in [7.1-7.7], but we can set up precise statements.
In this section and much of the sequel, waveform, automorphic form, and automorphic function will be used roughly as synonyms, referring to $\mathbb{C}$-valued functions on $\Gamma \backslash X$, meeting further conditions depending on the situation. Such functions are identifiable with $\Gamma$-invariant functions on $X$, by composing with the quotient $\operatorname{map} X \rightarrow \Gamma \backslash X$.

The constant term $c_{P} f$ of a waveform $f$ on $\Gamma \backslash X$ is a function on $X=G / K$ defined by

$$
(\text { constant term }) f(x)=c_{P} f(x)=\int_{(N \cap \Gamma) \backslash N} f(n \cdot x) d n
$$

Here the group $N$ is abelian, isomorphic to $\mathbb{R}^{\ell-1}$ for $\ell=2,3,4,5$, and $N \cap \Gamma$ is a discrete subgroup with compact quotient $(N \cap \Gamma) \backslash N$. We give $N$ the measure from the coordinate $x \rightarrow n_{x}$ with $x \in \mathbb{R}^{\ell-1}$, and as above and in [5.2] give the quotient the unique compatible measure for unwindings

$$
\left.\int_{(N \cap \Gamma) \backslash N}\left(\sum_{\gamma \in N \cap \Gamma} \varphi(\gamma n)\right) d n=\int_{N} \varphi(n) d n \quad \text { (for all } \varphi \in C_{c}^{o}(N)\right)
$$

By changing variables, we see that, although the constant term has probably lost left $\Gamma$-invariance, $c_{P} f$ is a left $N$-invariant function on $X=G / K$ :
$c_{P} f\left(n^{\prime} x\right)=\int_{(N \cap \Gamma) \backslash N} f\left(n \cdot n^{\prime} x\right) d n=\int_{(N \cap \Gamma) \backslash N} f\left(\left(n n^{\prime}\right) \cdot x\right) d n=\int_{(N \cap \Gamma) \backslash N} f(n \cdot x) d n \quad\left(\right.$ for $\left.n^{\prime} \in N\right)$
Thus, constant terms of functions $f$ on $\Gamma \backslash G / K$ can be viewed as functions on the ray

$$
\left.N \backslash X=N \backslash G / K=N \backslash\left(N A^{+} K\right) / K\right) \approx A^{+} \approx(0, \infty)
$$

[11] While integrals of $\Gamma$-invariant functions on $\mathfrak{H}$ on the quotient $\Gamma \backslash X$ can be understood in an elementary way as integrals over explicit fundamental domains, such a viewpoint impedes understanding of integration by parts on $C_{c}^{\infty}(\Gamma \backslash X)$. It is better to use an intrinsic integral on the quotient, characterized by the unwinding relation above, as in [5.2].

Similarly, since $\Gamma_{\infty}=P \cap \Gamma$ normalizes $N \cap \Gamma$, the constant term is left $\Gamma_{\infty}$-invariant. ${ }^{[12]}$ Altogether, $c_{P} f$ is left invariant by the group $N \Gamma_{\infty}$.

All this presumes that $c_{P} f$ has at least as much sense as a function with point-wise values as did $f$, but we need more than that. For example, unfortunately, it turns out that $f \in L^{2}(\Gamma \backslash G / K)$ does not imply that $c_{P} f \in L^{2}(N \backslash G / K)$. More cautiously, suppose $f$ is locally $L^{1}$, meaning that $|f|$ has finite integrals over compact subsets of $\Gamma \backslash G$. Fubini's theorem implies that a compactly-supported integral of $f$ in one of several variables is again locally $L^{1}$. This applies to $n \times y \rightarrow f\left(n a_{y}\right)$. The nature of the constant term map is clarified in [1.8].

Cuspforms are waveforms $f$ meeting the Gelfand condition $c_{P} f=0$. In some contexts, the term cuspform further connotes $\Delta$-eigenfunctions in $L^{2}(\Gamma \backslash G / K)$, but for present purposes the latter usage is too-restrictive. A genuine minor complication is that $L^{2}$ functions do not have good pointwise values, so vanishing of the constant term must mean almost everywhere for $L^{2}$ functions. Thus, it is often better to consider the constantterm map as a map on distributions, and the Gelfand condition as a distributional vanishing condition on distributions, as below in [1.8]. As usual, put

$$
L_{o}^{2}(\Gamma \backslash X)=\left\{L^{2} \text {-cuspforms }\right\}=\left\{f \in L^{2}(\Gamma \backslash G / K): c_{P} f=0\right\}
$$

The first main theorem, proven in [7.1-7.7], is the discrete decomposition of the space of cuspforms: one version is
[1.7.1] Theorem: The space $L_{o}^{2}(\Gamma \backslash G / K)$ of square-integrable cuspforms is a closed subspace of $L^{2}(\Gamma \backslash G / K)$, and has an orthonormal basis of $\Delta$-eigenfunctions. Each eigenspace is finite-dimensional, and the number of eigenvalues below a given bound is finite. (Proof in [7.1-7.7].)

The closed-ness of the space of $L^{2}$ cuspforms comes from recharacterization of it in terms of pseudoEisenstein series, in [1.8].

In contrast, the full space $L^{2}(\Gamma \backslash X)$ does not have a basis of $\Delta$-eigenfunctions: as proven in [1.12], the orthogonal complement of cuspforms in $L^{2}(\Gamma \backslash X)$ mostly consists of integrals of non- $L^{2}$ eigenfunctions for $\Delta$, the Eisenstein series $E_{s}$, introduced just below in [1.9].

The operator $\Delta$ presents some technical issues. For example, while $L^{2}(\Gamma \backslash X)$ lies inside the collection of distributions on $\Gamma \backslash X$, and interpreting $\Delta$ distributionally would make it well-defined on all of $L^{2}(\Gamma \backslash X)$, it would not stabilize $L^{2}(\Gamma \backslash X)$. This would seem to obstruct use of its symmetry or self-adjointness as an (unbounded) operator on a Hilbert space. On the other hand, indeed, no version of $\Delta$ can be defined on all of $L^{2}(\Gamma \backslash X)$ while retaining the symmetry $\langle\Delta f, F\rangle=\langle f, \Delta F\rangle$ for test functions $f, F$ in $L^{2}(\Gamma \backslash X)$. This situation requires careful treatment of unbounded, densely-defined operators on Hilbert spaces, as in [9.1-9.2].

### 1.8 Pseudo-Eisenstein series

Returning to $L^{2}(\Gamma \backslash X)$, we want to express the orthogonal complement of cuspforms $L_{o}^{2}(\Gamma \backslash X)$ in terms of $\Delta$-eigenfunctions, as discussed below in [1.12] and [1.13]. To exhibit explicit $L^{2}$ functions demonstrably spanning the orthogonal complement to cuspforms, we intend to recast the Gelfand vanishing condition. First, for $f \in L^{2}(\Gamma \backslash X)$, the constant term $c_{P} f$ is a left $N \Gamma_{\infty}$-invariant function on $G$. It vanishes as a distribution if and only if

$$
\int_{N \Gamma_{\infty} \backslash G} \varphi \cdot c_{P} f=0 \quad\left(\text { for all } \varphi \in C_{c}^{\infty}\left(N \Gamma_{\infty} \backslash G\right)\right)
$$

with right $G$-invariant measure on $N \Gamma_{\infty} \backslash G$ as in [5.2]. In fact, since $f$ is right $K$-invariant, $c_{P} f$ is right $K$-invariant, so we need only test against $\varphi \in C_{c}^{\infty}\left(N \Gamma_{\infty} \backslash G\right)^{K}$. The isomorphisms

$$
N \backslash X \approx N \backslash G / K \approx N \backslash\left(N A^{+} K\right) / K \approx A^{+}
$$

identify $N \backslash X$ with the ray $A^{+} \approx(0,+\infty)$, and identify right $K$-invariant functions $\varphi$ on $N \backslash G$ with functions of $y=$ height $\left(n a_{y} k\right)$. As in the previous section, for $f$ in $L^{2}$, since $f$ is locally integrable its constant term is locally integrable, by Fubini's theorem. Thus, $c_{P} f$ can be integrated against test functions on $N \backslash G / K$.
[12] In the present examples, $\Gamma_{\infty}=P \cap \Gamma$ is only finite index larger than $N \cap \Gamma$, but in other examples this index can be infinite.

## 1. Four small examples

Given $\varphi$ in $C_{c}^{\infty}\left(N \Gamma_{\infty} \backslash G\right)^{K}$, the corresponding pseudo-Eisenstein series $\Psi_{\varphi}$ should be a function in $C_{c}^{\infty}(\Gamma \backslash X)$ fitting into an adjunction:

$$
\int_{N \Gamma_{\infty} \backslash G} \varphi \cdot c_{P} f=\int_{\Gamma \backslash G} \Psi_{\varphi} \cdot f \quad\left(\text { for } f \in L^{2}(\Gamma \backslash X)\right)
$$

This adjunction will involve an unwinding/winding-up, so we might prefer that $c_{P} f$ be continuous, to easily invoke properties of vector-valued integrals [14.1]. For general reasons [6.1], $C_{c}^{o}(\Gamma \backslash G)$ is dense in $L^{2}(\Gamma \backslash G)$ in the $L^{2}$ topology, and for general reasons [5.1] the left action of $(N \cap \Gamma) \backslash N$ on the Fréchet space $C^{o}((N \cap \Gamma) \backslash G)$ is a continuous map $(N \cap \Gamma) \backslash N \times C^{o}((N \cap \Gamma) \backslash G) \rightarrow C^{o}(N \backslash G)$, so $c_{P} f$ exists as a $C^{o}(N \backslash G)$-valued GelfandPettis integral [14.1]. For $f \in C^{o}(\Gamma \backslash G)$, the integral of $c_{P} f$ against $\varphi \in C_{c}^{\infty}(N \backslash G / K)$ is the integral of a compactly-supported, continuous function.

Direct computation yields a canonical expression for the desired $\Psi_{\varphi}$, using the left $N \Gamma_{\infty}$-invariance of $\varphi$ and the left $\Gamma$-invariance of $f$, as follows. First, unwinding as in [5.2],

$$
\int_{N \Gamma_{\infty} \backslash G} \varphi \cdot c_{P} f=\int_{N \Gamma_{\infty} \backslash G} \varphi(g)\left(\int_{N \cap \Gamma \backslash N} f(n g) d n\right) d \mu(g)=\int_{\Gamma_{\infty} \backslash G} \varphi(g) f(g) d \mu(g)
$$

Winding up, using the left $\Gamma$-invariance of $f$,

$$
\int_{\Gamma_{\infty} \backslash G} f(g) \varphi(g) d \mu(g)=\int_{\Gamma \backslash G} \sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} f(\gamma \cdot g) \varphi(\gamma \cdot g) d \mu(g)=\int_{\Gamma \backslash G} f(g)\left(\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi(\gamma g)\right) d \mu(g)
$$

The inner sum in the last integral is the pseudo-Eisenstein series ${ }^{[13]}$ attached to $\varphi$ :

$$
\Psi_{\varphi}(g)=\sum_{\gamma \in \Gamma_{\infty \backslash \Gamma}} \varphi(\gamma g)
$$

The convergence of the sum needs attention:
[1.8.1] Claim: The series for a pseudo-Eisenstein series $\Psi_{\varphi}$ is locally finite, meaning that for $g$ in a fixed compact in $G$, there are only finitely-many non-zero summands in $\Psi_{\varphi}(g)=\sum_{\gamma} \varphi(\gamma g)$. Thus, $\Psi_{\varphi} \in C_{c}^{\infty}(\Gamma \backslash X)$.
Proof: Given $\varphi \in C_{c}^{\infty}(N \backslash G / K)$, let $C \subset G$ be compact so that $N \cdot C$ contains the support of $\varphi$. Fix compact $C_{o} \subset G$ in which $g \in G$ is constrained to lie. Then a summand $\varphi(\gamma g)$ is non-zero only if $\gamma g \in N \cdot C$, which holds only if

$$
\gamma \in \Gamma_{\infty} \cdot C \cdot g^{-1}
$$

so

$$
\gamma \in \Gamma \cap \Gamma_{\infty} \cdot C \cdot C_{o}^{-1}
$$

In the quotient $G \rightarrow \Gamma_{\infty} \backslash G$, the image of $\Gamma$ is closed and discrete. The image of the compact set $N \cdot C \cdot C_{o}^{-1}$ under the continuous quotient map is compact, since $(\Gamma \cap N) \backslash N$ is compact, and continuous images of compacts are compact. Thus, left modulo $\Gamma_{\infty}$, that intersection is the intersection of a closed discrete set and a compact set, so finite. (Compare the [1.5.2] from reduction theory.) Therefore, the series is locally finite, and defines a smooth function on $\Gamma \backslash G$. Summing over left translates certainly retains right $K$-invariance.

To show that $\Psi_{\varphi}$ has compact support in $\Gamma \backslash G$, proceed similarly. That is, for a summand $\varphi(\gamma g)$ to be non-zero, it must be that $g \in \Gamma \cdot C$. The image $\Gamma \backslash(\Gamma \cdot C)$ is compact, being the continuous image of the compact set $C$ under the continuous map $G \rightarrow \Gamma \backslash G$, proving the compact support.
[1.8.2] Corollary: Square-integrable cuspforms are the orthogonal complement in $L^{2}(\Gamma \backslash X)$ to the subspace of $L^{2}(\Gamma \backslash X)$ spanned by the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}(N \backslash X)$. The map $f \rightarrow c_{P} f$ is continuous from $L^{2}(\Gamma \backslash G / K)$ to distributions on $N \backslash G / K$.
[13] In 1966 Godement called these incomplete theta series. More recently Moeglin-Waldspurger reinforced the precedent of calling them pseudo-Eisenstein series

Proof: Again, as above, for general reasons [6.1] $C_{c}^{o}(\Gamma \backslash G / K)$ is dense in $L^{2}(\Gamma \backslash G / K)$, and the constant terms $c_{P} f$ are continuous for such $f$, so integrals against $\varphi \in C_{c}^{\infty}(N \backslash G / K)$ exist. Then the adjunction gives

$$
\left|\int_{N \backslash G / K} c_{P} f \cdot \bar{\varphi}\right|=\left|\left\langle f, \Psi_{\varphi}\right\rangle\right| \leq|f|_{L^{2}} \cdot\left|\Psi_{\varphi}\right|_{L^{2}}
$$

Thus, $f \rightarrow \int_{N \backslash G / K} c_{P} f \cdot \bar{\varphi}$ is a continuous linear functional on $L^{2}$. In particular, the kernels are closed, and the intersection of all these is the space of $L^{2}$ cuspforms. The inequality is exactly the continuity of $f \rightarrow c_{P} f$ with the weak dual topology [13.14] on distributions on $(0, \infty) \approx N \backslash G / K$.

Since $\Delta$ commutes with the group action, the effect of $\Delta$ on a pseudo-Eisenstein series is reflected entirely in its effect on the data: the sum is locally finite, so interchange of the operator and the sum is easy, giving

$$
\Delta \Psi_{\varphi}=\Delta \sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi \circ \gamma=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \Delta(\varphi \circ \gamma)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}(\Delta \varphi) \circ \gamma=\Psi_{\Delta \varphi}
$$

This correctly suggests that a suitable dense subspace of $L_{o}^{2}(\Gamma \backslash X)$ is indeed stable under $\Delta$. However, at this point we do not have a good device to prove density of smooth cuspforms with sufficient decay to prove symmetry $\langle\Delta f, F\rangle=\langle f, \Delta F\rangle$. For that matter, there is no reason to expect test functions in $L_{o}^{2}(\Gamma \backslash X)$ to be dense, since smooth-truncation to arrange compact support can succeed directly in $y \gg 1$, but disturbs the constant term as $y \rightarrow 0^{+}$. A convincing argument for smoothness of cuspforms and behavior of $\Delta$ on them can be given after the decomposition result of [7.1-7.7].

### 1.9 Eisenstein series

We can attempt to make a pseudo-Eisenstein series $\Psi_{\varphi}$ which is a $\Delta$-eigenfunction, by using a function $\varphi$ on $N \backslash X=N \backslash G / K$ which is a $\Delta$-eigenfunction. Using the $y$-coordinate on $N \backslash G / K \approx A^{+}$, the differential equation is

$$
\lambda \cdot \varphi=\Delta \varphi=\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-(\ell-2) y \frac{\partial}{\partial y}\right) \varphi=y^{2} \varphi^{\prime \prime}-(\ell-2) y \varphi^{\prime}
$$

The differential equation $y^{2} \varphi^{\prime \prime}-(\ell-2) y \varphi^{\prime}-\lambda \varphi=0$ is of Euler type, that is, will have solutions of the form $y^{\alpha}$, with $\alpha$ determined by

$$
0=y^{2} \cdot \alpha(\alpha-1) y^{\alpha-2}-(\ell-2) y \cdot \alpha y^{\alpha-1}-\lambda \cdot y^{\alpha}=y^{\alpha} \cdot(\alpha(\alpha-1)-(\ell-2) \alpha-\lambda)
$$

That is, for given $\lambda$, the corresponding exponents $\alpha$ are found by solving the indicial equation $\alpha(\alpha-1)-(\ell-2) \alpha-\lambda=0$, so $\lambda=\alpha(\alpha-(\ell-1))$. This computation suggests incorporating the factor $\ell-1$ into the exponent. Thus, with a function $\eta$ on $N \backslash G / K$ defined by

$$
\eta\left(n a_{y} k\right)=y^{\ell-1} \quad(\text { with } n \in N, k \in K, \text { and } \ell=2,3,4,5, \text { respectively })
$$

we have

$$
\Delta \eta^{s}=(\ell-1)^{2} \cdot s(s-1) \cdot \eta^{s}
$$

Unfortunately, $\eta^{s}$ is not in $C_{c}^{\infty}(N \backslash G / K)$, although it is smooth. The genuine Eisenstein series $E_{s}$ on $\Gamma \backslash X$ is ${ }^{[14]}$

$$
E_{s}(g)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \eta(\gamma \cdot g)^{s}
$$

The following claim has a much more elementary proof in the two simplest cases $\Gamma=S L_{2}(\mathbb{Z}), S L_{2}(\mathbb{Z}[i])$, but something more is required for $\Gamma=S p_{1,1}^{*}$ and $S L_{2}(\mathbb{H})$. We give an argument that applies uniformly to all four:

[^1]
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[1.9.1] Claim: For $\operatorname{Re}(s)>1$, the series expression for $E_{s}$ converges absolutely and uniformly on compacts, to a continuous function on $\Gamma \backslash X$, of moderate growth

$$
\left|E_{s}(g)\right|<_{t, C} \eta^{\operatorname{Re}(s)} \quad\left(\text { on } \mathfrak{S}_{t, C}, \text { implied constant depending on } t, C\right)
$$

Proof: It suffices to consider $s=\sigma$ real. From [1.3.2], in $x, y \rightarrow n_{x} a_{y}$ coordinates,

$$
\eta\left(\left(\begin{array}{ll}
* & * \\
c & d
\end{array}\right) n_{x} a_{y}\right)=\left(\frac{y}{|c x+d|^{2}+|c y|^{2}}\right)^{\ell-1}
$$

and $(c d) \rightarrow|c x+d|^{2}+|c y|^{2}$ is a positive-definite quadratic function on the real vector space in which ( $c d$ ) lies. The coefficients of the quadratic function depend continously on $x, y$, so for $x, y$ in a fixed compact there are uniform constants $A, B$ such that

$$
A \cdot\left(|c|^{2}+|d|^{2}\right) \leq|c x+d|^{2}+|c y|^{2} \leq B \cdot\left(|c|^{2}+|d|^{2}\right)
$$

In particular, for another pair $x^{\prime}, y^{\prime}$ in the same compact,

$$
\left|c x^{\prime}+d\right|^{2}+\left|c y^{\prime}\right|^{2} \leq B \cdot\left(|c|^{2}+|d|^{2}\right)=\frac{B}{A} \cdot\left(A \cdot\left(|c|^{2}+|d|^{2}\right)\right) \leq \frac{B}{A} \cdot\left(|c x+d|^{2}+|c y|^{2}\right)
$$

Thus, convergence of the series is equivalent to convergence of an averaged form, namely,

$$
\int_{C} \sum_{\Gamma_{\infty} \backslash \Gamma} \eta\left(\gamma n_{x} a_{y}\right)^{\sigma} \frac{d x d y}{y^{\ell}}
$$

Similarly, since the inf of lengths of non-zero vectors in a lattice in a real vector space is positive, there is a uniform non-zero lower bound for $|c x+d|^{2}+|c y|^{2}$ for $n_{x} a_{y} \in C$ and $(c d)$ a lower row in $\Gamma$. That is, the sup of $\eta(\gamma g)$ over $\gamma \in \Gamma$ and $g \in C$ is finite, and is attained. Let the sup be $\mu^{\ell-1}$ for $\mu>0$. Then $\Gamma \cdot C$ is contained in

$$
Y=\left\{n_{x} a_{y} \in X: y^{\ell-1}=\eta\left(n_{x} a_{y}\right) \leq \mu^{\ell-1}\right\}=\left\{n_{x} a_{y} \in X: y \leq \mu\right\}
$$

By discreteness of $\Gamma$ in $G$, we can shrink $C$ so that, for $\gamma$ in $\Gamma$, if $\gamma C \cap C \neq \phi$ then $\gamma=1$, so that

$$
\int_{C} E_{s}=\int_{\Gamma \backslash \Gamma \cdot C} E_{s}
$$

Unwind:

$$
\begin{gathered}
\int_{\Gamma \backslash \Gamma \cdot C} E_{s}=\int_{\Gamma_{\infty} \backslash \Gamma \cdot C} \eta\left(n_{x} a_{y}\right)^{\sigma} \frac{d x d y}{y^{\ell}} l e \int_{\Gamma_{\infty} \backslash Y} \eta\left(n_{x} a_{y}\right)^{\sigma} \frac{d x d y}{y^{\ell}}=\int_{N \cap \Gamma \backslash N} 1 \cdot \int_{0}^{\mu} \eta^{\sigma} \frac{d y}{y^{\ell}} \\
\ll \int_{0}^{\mu} y^{(\ell-1) \sigma} \frac{d y}{y^{\ell}}=\int_{0}^{\mu} y^{(\ell-1)(\sigma-1)} \frac{d y}{y}
\end{gathered}
$$

This is convergent for $\sigma-1>0$. This argument also proves the uniform convergence on compacts.
To see the moderate growth property, without yet attempting to prove that $E_{s}$ is smooth, differentiating the summands with $c \neq 0$

$$
\frac{\partial}{\partial y}\left|\frac{1}{\left(|c x+d|^{2}+|c y|^{2}\right)^{s(\ell-1)}}\right|=\frac{\partial}{\partial y} \frac{1}{\left(|c x+d|^{2}+|c y|^{2}\right)^{\operatorname{Re}(s) \cdot(\ell-1)}}<0
$$

shows that they all strictly decrease as $\eta(g)$ increases. Precisely, $\left|E_{s}\left(n a_{y}\right)\right|<\left|E_{s}\left(n a_{y^{\prime}}\right)\right|$ for $0<y<y^{\prime}$, for every $n \in N$. Since $E_{s}$ is continuous, it has a bound $B$ on a compact set $\left\{g \in \mathfrak{S}_{t, C}: \eta(g) \leq T\right\}$. Thus, $\left|y^{-s} E_{s}\right| \leq B$ on $\mathfrak{S}_{t, C}$.

Of course, we want convergence to a smooth function:
[1.9.2] Claim: The series for $E_{s}$ converges in the $C^{\infty}$ topology for $\operatorname{Re}(s)>1$, and produces a $C^{\infty}$ moderategrowth function on $\Gamma \backslash X=\Gamma \backslash G / K$. (Proof in [11.5].)

As in [13.5], the idea of the $C^{\infty}$ topology is that it is given by the collection of seminorms given by sups on compacts of all derivatives. One issue is which derivatives to use, and how to estimate them. In the present setting, one might be tempted to use derivatives with respect to coordinates $x, y$, but there is a significant disadvantage: $\partial / \partial x$ and $\partial / \partial y$ do not commute with the action of $\Gamma$ on $X$, so that $\partial E_{s} / \partial x$ is unlikely to be left $\Gamma$-invariant. For that matter, the effect of differentiating with respect to $y$ (after removing the common factor $\left.y^{s(\ell-1)}\right)$

$$
\frac{\partial}{\partial y} \frac{1}{\left(|c x+d|^{2}+|c y|^{2}\right)^{s(\ell-1)}}=-s(\ell-1) \cdot \frac{2 y \cdot|c|^{2}}{\left(|c x+d|^{2}+|c y|^{2}\right)^{s(\ell-1)+1}}
$$

on convergence of the series is difficult to appraise. The less-elementary approach in [11.5] uses left- $G$ invariant derivatives on $G$, which preserve left $\Gamma$-invariance, and which stabilize a somewhat-larger class of Eisenstein series.
[1.9.3] Corollary: $\operatorname{In} \operatorname{Re}(s)>1, E_{s}$ inherits the eigenvalue property from $\eta^{s}$ :

$$
\Delta E_{s}=(\ell-1)^{2} \cdot s(s-1) \cdot E_{s}
$$

Proof: Granting the convergence in the $C^{\infty}$ topology, in $\operatorname{Re}(s)>1$, and using the fact that $\Delta$ commutes with translations by $\Gamma$, letting $\Delta \eta^{s}=\lambda_{s} \cdot \eta^{s}$,

$$
\Delta \sum_{\gamma} \eta(\gamma g)^{s}=\sum_{\gamma} \Delta\left(\eta(\gamma g)^{s}\right)=\sum_{\gamma}\left(\Delta \eta^{s}\right)(\gamma g)=\lambda_{s} \cdot \sum_{\gamma} \eta^{s}(\gamma g)
$$

as claimed.
However, as we see below, $E_{s}$ is never in $L^{2}(\Gamma \backslash G)$.
Granting adequate convergence of $E_{s}$, and granting that the differential operator $\Delta$ can move inside the integral (see [14.1]) expressing the constant term, the constant term $c_{P} E_{s}$ is a $\Delta$-eigenfunction with the same eigenvalue:

$$
\begin{aligned}
\Delta\left(c_{P} E_{s}\right)(g)= & \Delta \int_{N \cap \Gamma \backslash N} E_{s}(n g) d n=\int_{N \cap \Gamma \backslash N} \Delta E_{s}(n g) d n=\int_{N \cap \Gamma \backslash N}(\ell-1)^{2} \cdot s(s-1) \cdot E_{s}(n g) d n \\
& =(\ell-1)^{2} \cdot s(s-1) \cdot \int_{N \cap \Gamma \backslash N} E_{s}(n g) d n=(\ell-1)^{2} \cdot s(s-1) \cdot c_{P} E_{s}(g)
\end{aligned}
$$

[1.9.4] Claim: The constant term of the Eisenstein series $E_{s}$ is of the form

$$
c_{P} E_{s}=\eta^{s}+c_{s} \eta^{1-s}
$$

Proof: From the very beginning of this section, at least for $s \neq \frac{1}{2}, \eta^{s}$ and $\eta^{1-s}$ are a basis for the space of $\Delta$-eigenfunctions with eigenvalue $(\ell-1)^{2} \cdot s(s-1)$ on $N \backslash G / K$. Thus, the constant term $c_{P} E_{s}$ is a linear combination of $\eta^{s}$ and $\eta^{1-s}$. The term $\eta^{s}$ comes from the representative $1 \in \Gamma_{\infty} \backslash \Gamma$. Every other representative $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ has $c \neq 0$. Implicitly recapitulating the computation in the Bruhat decomposition from [1.1],

$$
\int_{N \cap \Gamma \backslash N} \sum_{c, d}\left(\frac{y}{|c x+d|^{2}+|c y|^{2}}\right)^{s(\ell-1)} d n_{x}=\sum_{0 \neq c}^{\prime} \frac{1}{|c|^{s(\ell-1)}} \int_{N \cap \Gamma \backslash N} \sum_{d}\left(\frac{y}{\left|x+\frac{d}{c}\right|^{2}+|y|^{2}}\right)^{s(\ell-1)} d n_{x}
$$

where the sum over $c$ is over all possible lower-left entries of $\gamma \in \Gamma$, modulo $M \cap \Gamma$, and the inner sum over $d$ is over possible lower-right entries given lower-left entry $c$. With $\Xi=\left\{\xi: n_{\xi} \in \Gamma \cap N\right\}$, this is

$$
\sum_{0 \neq c} \frac{1}{|c|^{s(\ell-1)}} \sum_{d \bmod c} \int_{N \cap \Gamma \backslash N} \sum_{\xi \in \Xi}\left(\frac{y}{\left|x+\xi+\frac{d}{c}\right|^{2}+|y|^{2}}\right)^{s(\ell-1)} d n_{x}
$$

## 1. Four small examples

Note that if $\gamma=\left(\begin{array}{ll}* & * \\ c & d\end{array}\right) \in \Gamma$ and $\xi \in \Xi$, then

$$
\Gamma \ni \gamma \cdot n_{\xi}=\left(\begin{array}{ll}
* & * \\
c & d
\end{array}\right) \cdot\left(\begin{array}{ll}
1 & \xi \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
c & d+c \xi
\end{array}\right)
$$

The integral unwinds, giving

$$
\sum_{0 \neq c} \frac{1}{|c|^{s(\ell-1)}} \sum_{d \bmod c} \int_{N}\left(\frac{y}{\left|x+\frac{d}{c}\right|^{2}+|y|^{2}}\right)^{s(\ell-1)} d n_{x}
$$

For each suitable $d \bmod c$, replace $x$ by $x-\frac{d}{c}$, and let $\nu(c)$ be the number of such $d$, so the whole becomes

$$
\sum_{0 \neq c} \frac{\nu(c)}{|c|^{s(\ell-1)}} \int_{\mathbb{R}^{\ell-1}}\left(\frac{y}{|x|^{2}+|y|^{2}}\right)^{s(\ell-1)} d x=y^{(1-s)(\ell-1)} \sum_{0 \neq c} \frac{\nu(c)}{|c|^{s(\ell-1)}} \int_{\mathbb{R}^{\ell-1}}\left(\frac{1}{|x|^{2}+1}\right)^{s(\ell-1)} d x
$$

upon replacing $x$ by $y x$. This demonstrates the asserted shape of the constant term.
In fact, as usual, $c_{s}$ has an Euler product, and is identifiable as a ratio of $L$-functions, as we consider further in the sequel.
[1.9.5] Corollary: The Eisenstein series is not in $L^{2}(\Gamma \backslash X)$.
Proof: By reduction theory [1.5], it suffices to show that it is not square-integrable on a quotient $(N \cap \Gamma) \backslash\left\{n_{x} a_{y} \in X: y \geq t_{o}\right\}$ for $t_{o}$ large enough. Functions on such a set are left $N \cap \Gamma$-invariant functions on $N \times A^{+}$, so have Fourier expansions on the product of circles $(N \cap \Gamma) \backslash N$, with Fourier coefficients depending on $a \in A^{+}$. Specifically, let $\Psi$ be the collection of $N \cap \Gamma$-invariant continuous group homomorphisms $\psi: N \rightarrow \mathbb{C}^{\times}$. A function $f$ in $L^{2}((N \cap \Gamma) \backslash N)$ has a Fourier expansion converging (at least) in $L^{2}$ : the $\psi^{t h}$ Fourier coefficient is

$$
\widehat{f}(\psi)=\int_{(N \cap \Gamma) \backslash N} \bar{\psi}(n) \cdot f(n) d n
$$

giving $(N \cap \Gamma) \backslash N$ total measure 1, and

$$
f(n)=\sum_{\psi} \widehat{f}(\psi) \cdot \psi(n) \quad\left(n \in N, \text { convergent in an } L^{2} \text { sense }\right)
$$

The Plancherel theorem for $L^{2}((N \cap \Gamma) \backslash N)$ is

$$
\int_{(N \cap \Gamma) \backslash N}|f(n)|^{2} d n=\sum_{\psi}|\widehat{f}(\psi)|^{2}
$$

For a function $f$ on $(N \cap \Gamma) \backslash N \times A^{+}$, the Fourier coefficients are functions of $a \in A^{+}$, and

$$
f(n a)=\sum_{\psi} \widehat{f}(\psi)(a) \cdot \psi(n)
$$

Plancherel for $L^{2}((N \cap \Gamma) \backslash N)$ now gives

$$
\int_{(N \cap \Gamma) \backslash N}|f(n a)|^{2} d n=\sum_{\psi}|\widehat{f}(\psi)(a)|^{2}
$$

In particular, Bessel's inequality applied to the trivial character $\psi=1$ gives

$$
\int_{(N \cap \Gamma) \backslash N}|f(n a)|^{2} d n \geq|\widehat{f}(1)(a)|^{2}
$$

For $f(n a)=E_{s}(n a)$, that Fourier component for the trivial character in is exactly the constant term $c_{P} E_{s}=\eta^{s}+c_{s} \eta^{1-s}$, so

$$
\int_{(N \cap \Gamma) \backslash N} \int_{t_{o}}^{\infty}\left|E_{s}(n a)\right|^{2} d n d a \geq \int_{t_{o}}^{\infty}\left|\eta^{s}+c_{s} \eta^{1-s}\right|^{2} \frac{d y}{y^{\ell}}
$$

With $\sigma=\operatorname{Re}(s)>1$,

$$
\begin{aligned}
& \int_{t_{o}}^{T}\left|\eta^{s}+c_{s} \eta^{1-s}\right|^{2} \frac{d y}{y^{\ell}} \gg_{s} \int_{t_{o}}^{T} \eta^{2 \sigma} \frac{d y}{y^{\ell}}=\int_{t_{o}}^{T} y^{2 \sigma(\ell-1)} \frac{d y}{y^{\ell}} \\
& \quad=\int_{t_{o}}^{T} y^{(2 \sigma-1)(\ell-1)} \frac{d y}{y}=T^{(2 \sigma-1)(\ell-1)}-t_{o}^{(2 \sigma-1)(\ell-1)}
\end{aligned}
$$

This blows up as $T \rightarrow \infty$ for $2 \sigma-1>0$.

### 1.10 Meromorphic continuation of Eisenstein series

Although special tricks [2.B] applicable to $\Gamma=S L_{2}(\mathbb{Z})$ and $\Gamma=S L_{2}(\mathbb{Z}[i])$ have been known for almost 100 years, those tricks almost immediately fail in any larger context. For example, they do not apply to $\Gamma=S p_{1,1}^{*}(\mathfrak{o})$ or $\Gamma=S L_{2}(\mathbb{Z})$. [Selberg 1956] and [Roelcke 1956] first approached more general cases.

In [11.4] we will give a proof applying uniformly to our four example cases:
[1.10.1] Theorem: $E_{s}$ has a meromorphic continuation in $s \in \mathbb{C}$, as a smooth function of moderate growth on $\Gamma \backslash G$. As a function of $s, E_{s}(g)$ it is of at most polynomial growth vertically, uniformly in bounded strips, uniformly for $g$ in compacts. (Proof in [11.4].)

Although we give further details in a somewhat different logical order in [11.4], some consequences of the meromorphic continuation can be discussed directly:
[1.10.2] Corollary: The eigenfunction property $\Delta E_{s}=\lambda_{s} \cdot E_{s}$ with $\lambda_{s}=(\ell-1)^{2} \cdot s(s-1)$ persists under meromorphic continuation.
Proof: Both $\Delta E_{s}$ and $\lambda_{s} \cdot E_{s}$ are holomorphic function-valued functions of $s$, taking values in the topological vector space of smooth functions. They agree in the region of convergence $\operatorname{Re}(s)>1$, so the vector-valued form [15.2] of the Identity Principle from complex analysis gives the result.
[1.10.3] Corollary: The meromorphic continuation of $E_{s}$ implies the meromorphic continuation of the constant term $c_{P} E_{s}=\eta^{s}+c_{s} \eta^{1-s}$, in particular, of the function $c_{s}$.
Proof: Since $E_{s}$ meromorphically continues at least as a smooth function, the integral over the compact set $(N \cap \Gamma) \backslash N$ expressing a pointwise value $c_{P} E_{s}(g)$ of the constant term certainly converges absolutely. In fact, the function-valued function $n \rightarrow\left(g \rightarrow E_{s}(n g)\right)$ is a continuous smooth-function-valued function, and has a smooth-function-valued Gelfand-Pettis integral $g \rightarrow c_{P} E_{s}(g)$ [14.1].

In particular, the constant term $c_{P} E_{s}$ of the continuation of $E_{s}$ must still be of the form $A_{s} \eta^{s}+B_{s} \eta^{1-s}$ for some functions $A_{s}, B_{s}$, since (at least for $s \neq \frac{1}{2}$ ) $\eta^{s}$ and $\eta^{1-s}$ are the two linearly independent solutions of $\Delta f=\lambda_{s} \cdot f$ for functions $f$ on $N \backslash G / K \approx A^{+}$. In the region of convergence $\operatorname{Re}(s)>1$, the linear independence of $\eta^{s}$ and $\eta^{1-s}$ gives $A_{s}=1$ and $B_{s}=c_{s}$. The vector-valued form of the Identity Principle from complex analysis implies that $A_{s}=1$ throughout, and that $B_{s}=c_{s}$ throughout. In particular, this gives the meromorphic continuation of $c_{s}$.

The theory of the constant term in [8.1] asserts that a $\Delta$-eigenfunction of moderate growth is asymptotic to its constant term. For example,
[1.10.4] Claim: For every $s$ away from poles of $s \rightarrow E_{s}$, in a fixed Siegel set $\mathfrak{S}_{t, C}$,

$$
E_{s}\left(n a_{y} k\right)-\left(\eta^{s}+c_{s} \eta^{1-s}\right) \ll y^{-B}
$$

for every $B>0$, with the implied constant depending on $t, s$ and $B$. That is, $E_{s}-c_{P} E_{s}$ is rapidly decreasing in a Siegel set. More generally, for $s_{o}$ a pole of $s \rightarrow E_{s}$ of order $\nu$,

$$
\left.\left(s-s_{o}\right)^{\nu} E_{s}\left(n a_{y} k\right)\right|_{s=s_{o}}-\left.\left(s-s_{o}\right)^{\nu}\left(\eta^{s}+c_{s} \eta^{1-s}\right)\right|_{s=s_{o}} \ll y^{-B}
$$

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Proof: Since $E_{s}$ is a $\Delta$-eigenfunction of moderate growth, the theory of the constant term [8.1] exactly assures that $E_{s}$ is asymptotic to its constant term, in the sense of the first assertion. Near a pole $s_{o}$ of order $\nu$, writing a vector-valued Laurent expansion

$$
E_{s}=\frac{C_{\nu}}{\left(s-s_{o}\right)^{\nu}}+\frac{C_{\nu-1}}{\left(s-s_{o}\right)^{\nu-1}}+\ldots
$$

as in [15.2], where the coefficients $C_{j}$ are moderate-growth automorphic forms. Application of $\Delta$ termwise is justified, for example by invocation of the vector-valued form of Cauchy's formulas [15.2]: with $\lambda_{s}=$ $(\ell-1)^{2} \cdot s(s-1)$,

$$
\lambda_{s} \cdot\left(\frac{C_{\nu}}{\left(s-s_{o}\right)^{\nu}}+\frac{C_{\nu-1}}{\left(s-s_{o}\right)^{\nu-1}}+\ldots\right)=\lambda_{s} \cdot E_{s}=\Delta E_{s}=\frac{\Delta C_{\nu}}{\left(s-s_{o}\right)^{\nu}}+\frac{\Delta C_{\nu-1}}{\left(s-s_{o}\right)^{\nu-1}}+\ldots
$$

Multiplying through by $\left(s-s_{o}\right)^{\nu}$ and evaluating at $s=s_{o}, \lambda_{s_{o}} \cdot C_{\nu}=\Delta C_{\nu}$ as claimed. Then apply the theory of the constant term [8.1].

Granting the meromorphic continuation and the asymptotic estimation of the Eisenstein series by its constant term, the functional equation of $E_{s}$ is determined by its constant term:
[1.10.5] Corollary: $E_{s}$ has the functional equation $E_{1-s}=c_{1-s} E_{s}$, and $c_{s} \cdot c_{1-s}=1$. In particular, $\left|c_{s}\right|=1$ on $\operatorname{Re}(s)=\frac{1}{2}$.
Proof: Take $\operatorname{Re}(s)=\sigma>\frac{1}{2}$ and $s \notin \mathbb{R}$. Then $f=E_{1-s}-c_{1-s} E_{s}$ has constant term

$$
c_{P} f=\left(\eta^{1-s}+c_{1-s} \eta^{s}\right)-c_{1-s}\left(\eta^{s}+c_{s} \eta^{1-s}\right)=\left(1-c_{1-s} c_{s}\right) \cdot \eta^{1-s}
$$

For $\operatorname{Re}(s)>\frac{1}{2}, \eta^{1-s}$ is square-integrable on $\mathfrak{S}_{t, C}$ :

$$
\int_{\mathfrak{S}_{t, C}}\left|\eta^{1-s}\right|^{2} \frac{d x d y}{y^{\ell}}=(N \text {-measure } C) \cdot \int_{t}^{\infty}\left|y^{(1-s)(\ell-1)}\right|^{2} \frac{d y}{y^{\ell}}=(N \text {-measure } C) \cdot \int_{t}^{\infty} y^{(1-2 \sigma)(\ell-1)} \frac{d y}{y}
$$

Since $1-2 \sigma<0$, the integral is absolutely convergent. By the theory of the constant term [8.1], on a standard Siegel set

$$
\left.f=c_{P} f+(\text { rapidly decreasing })<_{s} \eta^{1-\sigma}+\text { (rapidly decreasing }\right)
$$

Thus, on $\mathfrak{S}_{t, C}$,

$$
|f|^{2} \ll \mid \eta^{1-\sigma}+\left.(\text { rapidly decreasing })\right|^{2}
$$

$=\eta^{2(1-\sigma)}+2 \cdot \eta^{1-\sigma} \cdot($ rapidly decreasing $)+(\text { rapidly decreasing })^{2}=\eta^{2(1-\sigma)}+($ rapidly decreasing $)$
Thus, $f=E_{1-s}-c_{1-s} E_{s} \in L^{2}(\Gamma \backslash X)$. It is a $\Delta$-eigenfunction with eigenvalue $\lambda_{s}=(\ell-1)^{2} \cdot s(s-1)$, which is not real for $\operatorname{Re}(s)>\frac{1}{2}$ and $s \notin \mathbb{R}$. But

$$
\lambda_{s} \cdot\langle f, f\rangle=\left\langle\lambda_{s} f, f\right\rangle=\langle\Delta f, f\rangle=\overline{\langle f, \Delta f\rangle}=\overline{\left\langle f, \lambda_{s} f\right\rangle}=\overline{\lambda_{s}} \cdot \overline{\langle f, f\rangle}=\overline{\lambda_{s}} \cdot\langle f, f\rangle
$$

Note that we did not use symmetry properties of $\Delta$, but only that $\langle f, F\rangle=\overline{\langle F, f\rangle}$. Thus, necessarily $E_{1-s}-c_{1-s} E_{s}=0$ for such $s$. For all $g \in G$, by the Identity Principle applied to the $\mathbb{C}$-valued meromorphic functions $s \longrightarrow\left(E_{1-s}(g)-c_{1-s} E_{s}(g)\right)$, the same identity applies for all $s$ away from poles.

Since the constant term $\left(1-c_{s} c_{1-s}\right) \cdot \eta^{1-s}$ of $E_{1-s}-c_{1-s} E_{s}=0$ is identically 0 , necessarily $c_{s} c_{1-s}=1$. Further, $s \rightarrow \overline{c_{\bar{s}}}$ is holomorphic and equal to $c_{s}$ for $\operatorname{Re}(s) \gg 1$, so the Identity Principle gives equality everywhere. Then $\overline{c_{\frac{1}{2}+i t}}=c_{1-\left(\frac{1}{2}+i t\right)}=c_{\frac{1}{2}-i t}$, and $\left|c_{\frac{1}{2}+i t}\right|^{2}=c_{\frac{1}{2}+i t} c_{\frac{1}{2}-i t}=1$.
[1.10.6] Claim: For $\operatorname{Re}(s) \neq \frac{1}{2}$ and $s \notin \mathbb{R}$, so that $\lambda_{s_{o}} \notin \mathbb{R}$, the poles of $s \rightarrow E_{s}$ are exactly the poles of $c_{s}$, and of the same order.

Proof: For $s_{o}$ such a pole, of order $\nu \geq 1$, corollary [1.10.2] showed that the leading Laurent coefficient is a $\Delta$-eigenfunction with eigenvalue $\lambda_{s_{o}}$ and is of moderate growth, so is asymptotic to its constant term. The Laurent expansion of the constant term is the constant term of the Laurent expansion, from the vectorvalued version of Cauchy's formula [15.2], and using the good behavior of continuous, compactly-supported vector-valued integrals [14.1]. Thus, if $c_{s}$ failed to have a pole at $s_{o}$, then the leading Laurent coefficient of $E_{s}$ at $s_{o}$ would have vanishing constant term, so (by the theory of the constant term) would be in $L^{2}(\Gamma \backslash X)$. Then $\lambda_{s_{o}} \in \mathbb{R}$, which is impossible for $s_{o}$ as in the hypotheses.

### 1.11 Truncation and Maaß-Selberg relations

The genuine Eisenstein series are not in $L^{2}(\Gamma \backslash X)$, but from the theory of the constant term [8.1] the only obstruction is the constant term, which is subtly altered by truncation, sufficiently removing this obstacle. The Maaß-Selberg relations are computation of the $L^{2}$ inner products of the resulting truncated Eisenstein series. As corollaries, we show that $E_{s}$ has only finitely-many poles in $\operatorname{Re}(s) \geq \frac{1}{2}$, that these are simple, lie on $\left(\frac{1}{2}, 1\right]$, and the residues are in $L^{2}(\Gamma \backslash X)$. Granting the spectral decomposition of cuspforms [7.1], and from the theory of the constant term [8.1] that the $\Delta$-eigenfunction cuspforms are of rapid decay, we prove that these residues of Eisenstein series are orthogonal to cuspforms.

The truncation operators $\wedge^{T}$ for large positive real $T$ act on an automorphic form $f$ by killing off $f$ 's constant term for large $y$. Thus, for a right $K$-invariant function, with a normalized version of height given by $\eta\left(n a_{y} k\right)=y^{\ell-1}$, one might imagine

$$
\text { (naive } T \text {-truncation of } f)(g)= \begin{cases}f(g) & \text { for } \eta(g) \leq T \\ f(g)-c_{P} f(g) & \text { for } \eta(g)>T\end{cases}
$$

This is not quite right. On a standard Siegel set $\mathfrak{S}_{t, C}$ this description is accurate, but it fails to correctly describe the truncated function on the whole domain $X$ or whole group $G$, in the sense that the truncation is not properly described as an automorphic form, that is, as a left $\Gamma$-invariant function. We want truncation to produce automorphic forms. For sufficiently large (depending on the reduction theory) $T$ we achieve the same effect by first defining the tail $c_{P}^{T} f$ of the constant term $c_{P} f$ of $f$ to be

$$
c_{P}^{T} f(g)=\left\{\begin{array}{cc}
0 & (\text { for } \eta(g) \leq T) \\
c_{P} f(y) & (\text { for } \eta(g)>T)
\end{array}\right.
$$

For legibility, we may replace a subscript by an argument in parentheses in the notation for pseudo-Eisenstein series: write

$$
\Psi(\varphi)=\Psi_{\varphi}
$$

Although $c_{P}^{T} f$ need not be smooth, nor compactly supported, by design (that is, for $T$ sufficiently large) its support is sufficiently high so that we have control over the analytical issues:
[1.11.1] Claim: For $T$ sufficiently large, the pseudo-Eisenstein series $\Psi\left(c_{P}^{T} f\right)$ is a locally finite sum, hence, uniformly convergent on compacts.
Proof: The tail $c_{P}^{T} f$ is left $N$-invariant. The reduction theory of [1.5] shows that, given $t_{o}$, for large-enough $t$, a set $\left\{n a_{y} k: y>t_{o}\right\}$ does not meet $\gamma \cdot\left\{n a_{y} k: y>t\right\}$ unless $\gamma \in \Gamma_{\infty}$. Thus, for large-enough $T$, $\left\{n a_{y} k: y>T\right\}$ does not meet $\gamma \cdot\left\{n a_{y} k: y>T\right\}$ unless $\gamma \in \Gamma_{\infty}$. Thus, $\gamma_{1} \cdot\left\{n a_{y} k: y>T\right\}$ does not meet $\gamma_{2} \cdot\left\{n a_{y} k: y>T\right\}$ unless $\gamma_{1} \Gamma_{\infty}=\gamma_{2} \Gamma_{\infty}$.

Similarly,
[1.11.2] Claim: On a standard Siegel set $\mathfrak{S}_{t, C}, \Psi\left(c_{P}^{T} f\right)=c_{P}^{T} f$ for all $T$ sufficiently large depending on $t$.
Proof: By reduction theory, a set $\left\{n a_{y} k: y>t_{o}\right\}$ does not meet $\gamma \cdot\left\{n a_{y} k: y>T\right\}$ unless $\gamma \in \Gamma_{\infty}$, for large-enough $T$ depending on $t_{o}$. Thus, for large-enough $T,\left\{n a_{y} k: y>T\right\}$ does not meet $\mathfrak{S}_{t_{o}, C}$ unless $\gamma \in \Gamma_{\infty}$. That is, the only non-zero summand in $\Psi\left(c_{P}^{T} f\right)$ is the term $c_{P}^{T} f$ itself.

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Thus, we find that the proper definition of the truncation operator $\wedge^{T}$ is

$$
\wedge^{T} f=f-\Psi\left(c_{P}^{T} f\right)
$$

As desired, a critical effect of the truncation procedure is:
[1.11.3] Claim: For $s$ away from poles, the truncated Eisenstein series $\wedge^{T} E_{s}$ is of rapid decay in all Siegel sets.
Proof: From the theory of the constant term [8.1], $E_{s}-c_{P} E_{s}$ is of rapid decay in a standard Siegel set. By the previous claim, $\left(E_{s}-c_{P}^{T} E_{s}\right)(g)=\left(E_{s}-c_{P} E_{s}\right)(g)$ for $\eta(g) \geq T$, so it is is also of rapid decay.
[1.11.4] Theorem: (Maaß-Selberg relation) Up to a uniform constant depending on normalization of measure,

$$
\frac{1}{\ell-1} \int_{\Gamma \backslash X} \wedge^{T} E_{s} \cdot \wedge^{T} E_{r}=\frac{T^{s+r-1}}{s+r-1}+c_{s} \frac{T^{(1-s)+r-1}}{(1-s)+r-1}+c_{r} \frac{T^{s+(1-r)-1}}{s+(1-r)-1}+c_{s} c_{r} \frac{T^{(1-s)+(1-r)-1}}{(1-s)+(1-r)-1}
$$

Proof: First,

$$
\int_{\Gamma \backslash X} \wedge^{T} E_{s} \cdot \wedge^{T} E_{r}=\int_{\Gamma \backslash X} \wedge^{T} E_{s} \cdot E_{r}
$$

because the tail of the constant term of $E_{r}$ is orthogonal to the truncated version $\wedge^{T} E_{s}$ of $E_{s}$. Then

$$
\begin{aligned}
\int_{\Gamma \backslash X} \wedge^{T} E_{s} \cdot \wedge^{T} E_{r} & =\int_{\Gamma \backslash X}\left(\Psi\left(\eta^{s}\right)-\Psi\left(\begin{array}{ll}
0 & (\text { for } \eta<T) \\
\eta^{s}+c_{s} \eta^{1-s} & (\text { for } \eta \geq T)
\end{array}\right)\right) \cdot E_{r} \\
& =\int_{\Gamma \backslash X} \Psi\left(\begin{array}{ll}
\eta^{s} & (\text { for } \eta<T) \\
-c_{s} \eta^{1-s} & (\text { for } \eta \geq T)
\end{array}\right) \cdot E_{r}
\end{aligned}
$$

Unwinding the awkward pseudo-Eisenstein series, noting that $\Gamma_{\infty}$ differs from $N \cap \Gamma$ only by the finite group $M \cap K$ which commutes with $A^{+}$, and the integrand is right $K$-invariant,

$$
\begin{gathered}
\int_{\Gamma_{\infty} \backslash X}\left\{\begin{array}{ll}
\eta^{s} & (\text { for } \eta<T) \\
-c_{s} \eta^{1-s} & (\text { for } \eta \geq T)
\end{array}\right\} \cdot E_{r}=\int_{N \backslash X} \int_{N \cap \Gamma \backslash N}\left\{\begin{array}{ll}
\eta^{s} & (\text { for } \eta<T) \\
-c_{s} \eta^{1-s} & (\text { for } \eta \geq T)
\end{array}\right\} \cdot E_{r} \\
=\int_{N \backslash X}\left\{\begin{array}{ll}
\eta^{s} & (\text { for } \eta<T) \\
-c_{s} \eta^{1-s} & (\text { for } \eta \geq T)
\end{array}\right\} \cdot\left(\int_{N \cap \Gamma \backslash N} E_{r}(n g) d n\right) d g \\
=\int_{N \backslash X}\left\{\begin{array}{ll}
\eta^{s} & (\text { for } \eta<T) \\
-c_{s} \eta^{1-s} & (\text { for } \eta \geq T)
\end{array}\right\} \cdot\left(\eta^{r}+c_{r} \eta^{1-r}\right)=\int_{0}^{\infty} \begin{cases}\eta^{s}\left(\eta^{r}+c_{r} \eta^{1-r}\right) & (\text { for } \eta<T) \\
-c_{s} \eta^{1-s}\left(\eta^{r}+c_{r} \eta^{1-r}\right) & (\text { for } \eta \geq T) \\
= & \int_{0}^{T} \eta^{s} \cdot\left(\eta^{r}+c_{r} \eta^{1-r}\right) \frac{d y}{y^{\ell}}-\int_{T}^{\infty} c_{s} \eta^{1-s}\left(\eta^{r}+c_{r} \eta^{1-r}\right) \frac{d y}{y^{\ell}}\end{cases}
\end{gathered}
$$

Note that the measure $d y / y^{\ell}$ is descended from the right $G$-invariant measure on $N \backslash G$. Assume that $\operatorname{Re}(r)$ is bounded above and below, so that $\operatorname{Re}(1-r)$ is also bounded, and take $\operatorname{Re}(s)$ sufficiently large so that all the integrals converge. The above becomes

$$
\int_{0}^{T} \eta^{s+r-1} \frac{d y}{y}+c_{r} \int_{0}^{T} \eta^{s+(1-r)-1} \frac{d y}{y}-c_{s} \int_{T}^{\infty} \eta^{(1-s)+r-1} \frac{d y}{y}-c_{s} c_{r} \int_{T}^{\infty} \eta^{(1-s)+(1-r)-1} \frac{d y}{y}
$$

Since $d \eta / \eta=(\ell-1) \cdot d y / y$, this gives the expression of the theorem. Note that $\ell-1=1$ in the most familiar case of $\Gamma=S L_{2}(\mathbb{Z})$. By analytic continuation (in $s$ and in $r$ ) it is valid everywhere it makes sense.

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The following corollaries can be proven directly in special cases by use of explicit details about Fourier expansion of the Eisenstein series. However, the arguments here generalize.
[1.11.5] Corollary: There are only finitely-many poles of $E_{s}$ in the region $\operatorname{Re}(s) \geq \frac{1}{2}$, all on the segment $\left(\frac{1}{2}, 1\right]$, and these are poles of $c_{s}$. Any such pole is simple, with residue in $L^{2}(\Gamma \backslash X)$. Specifically, with measure normalized as in the previous proof,

$$
\int_{\Gamma \backslash X}\left|\operatorname{Res}_{\sigma_{o}} E_{s}\right|^{2}=(\ell-1) \cdot \operatorname{Res}_{\sigma_{o}} c_{s}<+\infty
$$

Such a residue is also smooth, and moderate growth, and has eigenvalue $\lambda_{s_{o}}=(\ell-1)^{2} \cdot s_{o}\left(s_{o}-1\right)$.
Proof: The Eisenstein series is indeed treated as a meromorphic function-valued function, as in [15.2], so its Laurent coefficients or power series coefficients are functions in the same topological vector space, by the vector-valued form of Cauchy's formulas [15.2]. From the identity principle, since $\overline{E_{\bar{s}}}=E_{s}$ for $\operatorname{Re}(s)>1$, we have $\overline{E_{s}}=E_{\bar{s}}$ for all $s$ away from poles, and similarly for truncated Eisenstein series. Thus, taking $r=\bar{s}=\sigma-i t$ in the theorem,

$$
\frac{1}{\ell-1} \int_{\Gamma \backslash X}\left|\wedge^{T} E_{s}\right|^{2}=\frac{T^{2 \sigma-1}}{2 \sigma-1}+c_{s} \frac{T^{-2 i t}}{-2 i t}+c_{\bar{s}} \frac{T^{2 i t}}{2 i t}+c_{s} c_{\bar{s}} \frac{T^{1-2 \sigma}}{1-2 \sigma}
$$

Suppose $E_{s}$ has a pole $s_{o}=\sigma_{o}+i t_{o}$ of order $\nu$ with $t_{o} \neq 0$ and $\sigma_{o}>\frac{1}{2}$.
From corollary [1.10.4] to the theory of the constant term [8.1], with non-real eigenvalue, this is equivalent to the assertion that $c_{s}$ has a pole at $s_{o}$ of order $\nu$. Also, $c_{\bar{s}}=\overline{c_{s}}$, so $c_{s}$ has a pole at $\overline{s_{o}}$ as well, of the same order, with leading Laurent term the complex conjugate of that at $s_{o}$. Thus, the function $\wedge^{T} E_{s}$ also has a pole exactly at poles of $c_{s}$, of the same order, for non-real $\lambda_{s}$.

Take $s=\sigma_{o}+i t$ in the above. In the real variable $t$, the left-hand side of the Maaß-Selberg relation is asymptotic to a positive constant multiple of $\left(t-t_{o}\right)^{-2 \nu}$ as $t \rightarrow t_{o}$, since the pole is of order $\nu$ and inner products are positive. The first term on the right-hand side is bounded as $t \rightarrow t_{o}$, and the second and third terms are asymptotic to non-zero constant multiples of $\left(t-t_{o}\right)^{-\nu}$. Thus, the first three terms on the right can be ignored as $t \rightarrow t_{o}$. The fourth term on the right-hand side is asymptotic to a positive constant multiple of $\left(t-t_{o}\right)^{-2 \nu}$ from $c_{s} c_{\bar{s}}$, multiplied by $T^{1-2 \sigma_{o}} /\left(1-2 \sigma_{o}\right)$. The denominator is negative, so that, altogether, the fourth term on the right-hand side is asymptotic to a negative constant multiple of $\left(t-t_{o}\right)^{-2 \nu}$. The positivity of the left-hand side of the Maaß-Selberg, and negativity of the right-hand side ( as $t \rightarrow t_{o}$ ), contradict the hypothesized pole. Thus, $E_{s}$ and $c_{s}$ have no poles off the real axis in the region $\operatorname{Re}(s)>1 / 2$.

Next, let $s_{o}=\sigma_{o}$ be a pole of $E_{s}$ of order $\nu \geq 1$ on $\left(\frac{1}{2}, 1\right]$. Take $r=s=\sigma_{o}+i t$, obtaining

$$
\frac{1}{\ell-1} \int_{\Gamma \backslash X}\left|t^{\nu} \cdot \wedge^{T} E_{s}\right|^{2}=t^{2 \nu} \cdot\left(\frac{T^{2 \sigma_{o}-1}}{2 \sigma_{o}-1}+c_{s} \frac{T^{-2 i t}}{-2 i t}+c_{\bar{s}} \frac{T^{2 i t}}{2 i t}+c_{s} c_{\bar{s}} \frac{T^{1-2 \sigma_{o}}}{1-2 \sigma_{o}}\right)
$$

As $t \rightarrow t_{o}=0$, the right-hand side goes to 0 unless $c_{s}$ also has a pole of order $\nu$ at $s_{o}$. The fourth term is negative, and if $\nu>1$ is the only term that survives on the right-hand side as $t \rightarrow 0$, contradicting the non-negativity of the left-hand side. Thus, $\nu=1$, in which case the second and third terms' blow-up is of the same order as the left-hand side and the fourth term on the right-hand side. This proves that any pole on $\left(\frac{1}{2}, 1\right]$ is simple.

Letting $t \rightarrow 0$,

$$
\frac{1}{\ell-1} \int_{\Gamma \backslash X}\left|\operatorname{Res}_{\sigma_{o}} E_{s}^{T}\right|^{2}=\frac{\operatorname{Res}_{\sigma_{o}} c_{s}}{2}+\frac{\operatorname{Res}_{\sigma_{o}} c_{\bar{s}}}{2}+\operatorname{Res}_{\sigma_{o}} c_{s} \cdot \overline{\operatorname{Res}_{\sigma_{o}} c_{s}} \frac{T^{1-2 \sigma_{o}}}{1-2 \sigma_{o}}
$$

Since $1-2 \sigma_{o}<0$, the limit of the last term is 0 as $T \rightarrow+\infty$, given the square-integrability of the residue. General considerations about meromorphic vector-valued functions [15.2] and Gelfand-Pettis integrals [14.1] assure that taking residues commutes with taking the limit as $T \rightarrow \infty$. The two remaining terms are equal, since the pole is on the real line.

Now suppose $s_{o}$ is a pole of $E_{s}$ of order $\nu \geq 1$ on the line $\operatorname{Re}(s)=\frac{1}{2}$ and off $\mathbb{R}$. The leading Laurent coefficient $C_{\nu}$ of $E_{s}$ at $s_{o}$ is a $\Delta$-eigenfunction with eigenvalue $\lambda_{s_{o}}$, and is of moderate growth, again by the

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vector-valued form of Cauchy's formulas. Thus, by the theory of the constant term [8.1], $C_{\nu}$ is asymptotic to its constant term $c_{P} C_{\nu}$, which, again, is the the leading Laurent coefficient of the constant term $\eta^{s}+c_{s} \eta^{1-s}$ of $E_{s}$. The property $\left|c_{s}\right|=1$ on $\operatorname{Re}(s)=\frac{1}{2}$ shows that $c_{s}$ has no pole there. Then $c_{P} C_{\nu}=0$, so $\wedge^{T} C_{\nu}=C_{\nu}$, and $C_{\nu}$ is in $L^{2}(\Gamma \backslash X)$. By Maaß-Selberg with $r=\bar{s}=\sigma-i t$,

$$
\begin{aligned}
& \frac{1}{\ell-1} \int_{\Gamma \backslash X}\left|C_{\nu}\right|^{2}=\frac{1}{\ell-1} \int\left|\lim _{s \rightarrow s_{o}}\left(s-s_{o}\right)^{2 \nu} \wedge^{T} E_{s}\right|^{2} \\
= & \lim _{s \rightarrow s_{o}}\left(s-s_{o}\right)^{2 \nu} \cdot\left(\frac{T^{2 \sigma-1}}{2 \sigma-1}+c_{s} \frac{T^{-2 i t}}{-2 i t}+c_{\bar{s}} \frac{T^{2 i t}}{2 i t}+c_{s} c_{\bar{s}} \frac{T^{1-2 \sigma}}{1-2 \sigma}\right)
\end{aligned}
$$

Since $\nu \geq 1$, approaching $s_{o}$ from off the line, the limit of $\left(s-s_{o}\right)^{2 \nu} /(2 \sigma-1)$ is 0 . Since $\left|c_{s}\right| \rightarrow 1$ as $s \rightarrow s_{o} \in \frac{1}{2}+i \mathbb{R}$, the whole limit is 0 . Thus, $E_{s}$ has no pole on $\frac{1}{2}+i \mathbb{R}$.

Finally, we see that the residues are not only in $L^{2}(\Gamma \backslash X)$, but are also smooth (and moderate growth pointwise) $\Delta$-eigenfunctions with the indicated eigenvalues. By the vector-valued form of Cauchy's formula for residues,

$$
\operatorname{Res}_{s=s_{o}} E_{s}=\frac{1}{2 \pi i} \int_{\gamma} \frac{E_{w}}{w-s} d w
$$

where $\gamma$ is a small circle around $s_{o}$, traversed counter-clockwise. Since $w \rightarrow E_{w} /(s-w)$ is a continuous moderate-growth-function-valued function, Gelfand-Pettis [14.1] assures that the integral is in the same space. In particular, the residue is smooth. Because the pole is simple, the function $f_{s}=\left(s-s_{o}\right) E_{s}$ has a removable singularity at $s_{o}$, and its value there is the residue. In the topology on moderate-growth functions (as in [13.10]), $\Delta$ is a continuous map. From the theory of vector-valued holomorphic functions [15.2] and Gelfand-Pettis integrals [14.1], evaluation commutes with continuous linear maps, so

$$
\Delta\left(\operatorname{Res}_{s_{o}} E_{s}\right)=\Delta\left(\left.f_{s}\right|_{s=s_{o}}\right)=\left.\left(\Delta f_{s}\right)\right|_{s=s_{o}}=\lambda_{s_{o}} \cdot f_{s_{o}}
$$

demonstrating that residues are (smooth and) $\Delta$-eigenfunctions.
For $\Gamma=S L_{2}(\mathbb{Z})$ and $S L_{2}(\mathbb{Z}[i])$, there are special arguments that show that the only relevant residues of Eisenstein series are at $s=1$. The eigenvalue $\lambda_{s}=s(s-1)$ of $\Delta$ at $s=1$ is 0 .
[1.11.6] Claim: Any smooth $f \in L^{2}(\Gamma \backslash G / K)$ such that $\Delta f=0$ is constant.
Proof: Let $\nabla$ be the tangent-space-valued gradient on $\Gamma \backslash G / K$, as developed in more detail in [10.7],

$$
\int_{\Gamma \backslash G} \Delta f \cdot F=-\int_{\Gamma \backslash G}\langle\nabla f, \nabla F\rangle
$$

where, for the moment, $\langle-,-\rangle$ is a inner product on the tangent space. For $\Delta f=0$, this gives

$$
0=\int_{\Gamma \backslash G} 0 \cdot f=\int_{\Gamma \backslash G} \Delta f \cdot f=-\int_{\Gamma \backslash G}\langle\nabla f, \nabla f\rangle
$$

Thus, $\nabla f$ is identically 0 , so $f$ is constant.
The other two current examples, $\Gamma=S p_{1,1}^{*}(\mathfrak{o})$ and $S L_{2}(\mathfrak{o})$, do not admit those special arguments to decisively locate poles, although they still do have poles at $s=1$, with constant residues, by the same argument. To treat residues in $\operatorname{Re}(s)>\frac{1}{2}$ generally:
[1.11.7] Corollary: Residues of Eisenstein series at distinct poles $s_{1}, s_{2}$ in $\left(\frac{1}{2}, 1\right]$ are mutually orthogonal.
Proof: Let $f_{j}$ be the residue at $s_{j}$, with eigenvalue $\lambda_{j}$. The eigenvalues are real, since $s_{1}, s_{2} \in \mathbb{R}$. It is reasonable to think that $\Delta$ has the symmetry $\left\langle\Delta f_{1}, f_{2}\right\rangle=\left\langle f_{1}, \Delta f_{2}\right\rangle$ so that the usual argument

$$
\lambda_{1} \cdot\left\langle f_{1}, f_{2}\right\rangle=\left\langle\Delta f_{1}, f_{2}\right\rangle=\left\langle f_{1}, \Delta f_{2}\right\rangle=\lambda_{2} \cdot\left\langle f_{1}, f_{2}\right\rangle
$$

would give $\left(\lambda_{1}-\lambda_{2}\right) \cdot\left\langle f_{1}, f_{2}\right\rangle=0$, and then $\left\langle f_{1}, f_{2}\right\rangle=0$. However, the defensible starting-point [6.6] for this symmetry property of $\Delta$ is that it holds for functions in $C_{c}^{\infty}(\Gamma \backslash G / K)$, in effect avoiding any boundary

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terms in integration by parts. To preserve symmetry in an extension requires care. In fact, the method of [9.10] shows that $\Delta$ is essentially self-adjoint in the sense of having a unique self-adjoint extension to an unbounded operator densely-defined in $L^{2}(\Gamma \backslash G / K)$. The domain of that extension does include the residues $f_{j}$, but demonstration of the latter fact is more a consequence than starting-point.

Instead, the Maaß-Selberg relation (with $s_{1} \neq s_{2}$ both real, eliminating some complex conjugations) gives

$$
\begin{gathered}
\frac{1}{\ell-1} \int_{\Gamma \backslash X} \wedge^{T} E_{s_{1}} \cdot \wedge^{T} E_{s_{2}} \\
=\frac{T^{s_{1}+s_{2}-1}}{s_{1}+s_{2}-1}+c_{s_{1}} \frac{T^{\left(1-s_{2}\right)+s_{2}-1}}{\left(1-s_{1}\right)+s_{2}-1}+c_{s_{2}} \frac{T^{s_{1}+\left(1-w_{2}\right)-1}}{s_{1}+\left(1-s_{2}\right)-1}+c_{s_{1}} c_{s_{2}} \frac{T^{\left(1-s_{1}\right)+\left(1-s_{2}\right)-1}}{\left(1-s_{1}\right)+\left(1-s_{2}\right)-1}
\end{gathered}
$$

With simple poles of $E_{s}$ at $s_{1}$ and $s_{2}$, multiplying through by $\left(s-s_{1}\right)\left(s^{\prime}-s_{2}\right)$ and taking the limit as $s \rightarrow s_{1}$ and $s^{\prime} \rightarrow s_{2}$ gives

$$
\frac{1}{\ell-1} \int_{\Gamma \backslash X} \wedge^{T} \operatorname{Res}_{s_{1}} E_{s_{1}} \cdot \wedge^{T} \operatorname{Res}_{s_{2}} E_{s_{2}}=0+0+0+\operatorname{Res}_{s=s_{1}} c_{s} \cdot \operatorname{Res}_{s=s_{2}} c_{s} \cdot \frac{T^{\left(1-s_{1}\right)+\left(1-s_{2}\right)-1}}{\left(1-s_{1}\right)+\left(1-s_{2}\right)-1}
$$

Since $\left(1-s_{1}\right)+\left(1-s_{2}\right)-1<\left(1-\frac{1}{2}\right)+\left(1-\frac{1}{2}\right)-1<0$, letting $T \rightarrow+\infty$ gives the orthogonality. ///
[1.11.8] Corollary: The residues of $E_{s}$ for $s \in\left(\frac{1}{2}, 1\right]$ are orthogonal to cuspforms.
Proof: This uses the spectral decomposition of cuspforms [1.7]: there is an orthonormal basis for $L^{2}$ cuspforms consisting of $\Delta$-eigenfunctions, and each eigenspace is finite-dimensional. The theory of the constant term [8.1] shows that any such eigenfunction is asymptotic to its constant term. Since constant terms of cuspforms are 0, cuspform-eigenfunctions are of rapid decay in Siegel sets.

Thus, for a cuspform-eigenfunction $f$, granting [1.9.1] that Eisenstein series $E_{s}$ are of moderate growth on Siegel sets, the literal integrals $\left\langle f, E_{s}\right\rangle=\int_{\Gamma \backslash X} f \cdot \overline{E_{s}}$ are absolutely convergent for all $s$ away from poles. These are not $L^{2}$ inner products, since $E_{s}$ is never in $L^{2}$, but we use the same notation for brevity. In the region of convergence $\operatorname{Re}(s)>1$, any integral $\int_{\Gamma \backslash X} f \cdot E_{s}$ unwinds to compute an integral against the constant term of $f$, and the latter is 0 :

$$
\begin{gathered}
\int_{\Gamma \backslash X} f \cdot E_{s}=\int_{\Gamma_{\infty} \backslash X} f \cdot \eta^{s}=\int_{N \Gamma_{\infty} \backslash X}\left(\int_{(N \cap \Gamma) \backslash N} f(n g) \eta(n g)^{s} d n\right) d g \\
=\int_{N \Gamma_{\infty} \backslash X} \eta^{s}(g)\left(\int_{(N \cap \Gamma) \backslash N} f(n g) d n\right) d g=\int_{N \Gamma_{\infty} \backslash X} \eta^{s}(g) \cdot c_{P} f(g) d g=\int_{N \Gamma_{\infty} \backslash X} \eta^{s}(g) \cdot 0 d g=0
\end{gathered}
$$

Because $s \rightarrow E_{s}$ is a meromorphic function-valued function taking values in (at least continuous) functions of moderate growth, the function $s \rightarrow\left\langle f, E_{s}\right\rangle$ is meromorphic on $\mathbb{C}$. By the Identity Principle, since this function is 0 on $\operatorname{Re}(s)>1$, it is identically 0 . The vector-valued form of Cauchy's formula expresses the residue at $s=s_{o}$ as an integral:

$$
\operatorname{Res}_{s=s_{o}} E_{s}=\frac{1}{2 \pi i} \int_{\gamma} \frac{E_{w}}{w-s} d w
$$

where $\gamma$ is a small circle around $s_{o}$, traversed counter-clockwise. Then

$$
\left\langle f, \operatorname{Res}_{s_{o}} E_{s}\right\rangle=\left\langle f, \frac{1}{2 \pi i} \int_{\gamma} \frac{E_{w}}{w-s} d w\right\rangle
$$

The functional $u \rightarrow\langle f, u\rangle$ is a continuous linear functional on functions of moderate growth, and $w \rightarrow E_{w} /(s-w)$ is a continuous, compactly-supported moderate-growth-function-valued function, so by Gelfand-Pettis [14.1] the inner product passes inside the integral:

$$
\left\langle f, \operatorname{Res}_{s_{o}} E_{s}\right\rangle=\frac{1}{2 \pi i} \int_{\gamma}\left\langle f, E_{w}\right\rangle \frac{1}{w-s} d w=\frac{1}{2 \pi i} \int_{\gamma} 0 \cdot \frac{1}{w-s} d w=0
$$

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Again, $\left\langle f, E_{s}\right\rangle$ is not an $L^{2}$ pairing, because $E_{s}$ is not in $L^{2}$. Nevertheless, because of the rapid decay of $f$, the implied integral is absolutely convergent. This proves that the residues of $E_{s}$ for $s \in\left(\frac{1}{2}, 1\right]$, all of which are in $L^{2}$, are $L^{2}$-orthogonal to cuspforms.

### 1.12 Decomposition of pseudo-Eisenstein series

We saw in [1.8] that the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}(N \backslash G / K)$ generate the orthogonal complement to cuspforms in $L^{2}(\Gamma \backslash G / K)$ : since the orthogonal complement of these pseudo-Eisenstein series is the space of cuspforms, the orthogonal complement to cuspforms is the $L^{2}$-closure of the set of these pseudo-Eisenstein series.

To express such pseudo-Eisenstein series as superpositions of $\Delta$-eigenfunctions in the four examples at hand, once we know the meromorphic continuation and functional equation of the genuine Eisenstein series $E_{s}$, the essential harmonic analysis is Fourier transform on the real line, in coordinates in which it is known as a Mellin transform. That is, the non-cuspidal part of harmonic analysis on $\Gamma \backslash X$ in each of these four examples reduces to harmonic analysis on $\mathbb{R}$.

For $\varphi \in C_{c}^{\infty}(\Gamma \backslash G / K)=C_{c}^{\infty}(N \backslash G)^{K}$, the pseudo-Eisenstein series $\Psi_{\varphi}$ is in $C_{c}^{\infty}(\Gamma \backslash G)^{K}$, so its integral against $E_{s}$ converges absolutely, since $E_{s}$ is continuous, even after meromorphic continuation. Thus, by abuse of notation, we may write

$$
\left\langle\Psi_{\varphi}, E_{s}\right\rangle=\int_{\Gamma \backslash X} \Psi_{\varphi} \cdot \overline{E_{s}}
$$

even though this $\langle$,$\rangle cannot be the L^{2}$ pairing, since $E_{s} \notin L^{2}(\Gamma \backslash X)$. The following is a preliminary version of a spectral decomposition of the $L^{2}$ closure of the space containing pseudo-Eisenstein series, insofar as it only treats $\Psi_{\varphi}$ with test-function $\varphi$, only computes point-wise values, so does not consider the integral of genuine Eisenstein series as a function-valued integral, and omits a Plancherel assertion.
[1.12.1] Theorem: (Numerical form) Let $s_{o}$ run over poles of $E_{s}$ in $\operatorname{Re}(s)>\frac{1}{2}$. For $\varphi \in C_{c}^{\infty}(N \backslash G / K)$, the pseudo-Eisenstein series $\Psi_{\varphi}$ is expressible in terms of genuine Eisenstein series, by an integral converging absolutely and uniformly on compacts in $\Gamma \backslash G / K$ :

$$
\Psi_{\varphi}(g)=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \cdot E_{s}(g) d s+(\ell-1) \sum_{s_{o}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s}(g)
$$

where we abuse notation by writing $\left\langle\Psi_{\varphi}, E_{s}\right\rangle=\int_{\Gamma \backslash G} \Psi_{\varphi} \cdot \bar{E}_{s}$ even though $E_{s}$ is not in $L^{2}$.
Proof: One form of Fourier inversion for Schwartz functions ${ }^{[15]} f$ on the real line is

$$
f(x)=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\left(\int_{-\infty}^{\infty} f(u) e^{-i \xi u} d u\right) e^{i \xi x} d \xi
$$

Both outer and inner integrals converge very well, uniformly pointwise. The inner integral is a Schwartz function in $\xi$. Fourier transforms on $\mathbb{R}$ put into multiplicative coordinates are Mellin transforms: for $\varphi \in C_{c}^{\infty}(0,+\infty)$, take $f(x)=\varphi\left(e^{x}\right)$. Let $y=e^{x}$ and $r=e^{u}$, and rewrite Fourier inversion as

$$
\varphi(y)=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\left(\int_{0}^{\infty} \varphi(r) r^{-i \xi} \frac{d r}{r}\right) y^{i \xi} d \xi
$$

The Fourier transform in these coordinates is a Mellin or Laplace transform. For compactly-supported $\varphi$, as we use throughout this discussion, the integral definition extends to all complex $s$ in place of $i \xi$, and $d \xi$ replaced by $-i d s$. The variant Fourier inversion identity gives Mellin inversion

$$
\varphi(y)=\frac{1}{2 \pi i} \int_{-i \infty}^{i \infty}\left(\int_{0}^{\infty} \varphi(r) r^{-s} \frac{d r}{r}\right) y^{s} d s
$$

[15] As usual, Schwartz functions $\mathscr{S}(\mathbb{R})$ on $\mathbb{R}$ or any copy of it are smooth functions $f$ such that $f$ and all its derivatives are rapidly decreasing, in the sense that $\left(1+x^{2}\right)^{N} \cdot\left|f^{(k)}(x)\right|$ is bounded on $x \in \mathbb{R}$ for every $k$ and $N$. These sups are a family of seminorms that give $\mathscr{S}(\mathbb{R})$ a Fréchet space structure. See chapter 12.

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By an easy part of the Paley-Wiener theorem [13.16], for $f \in C_{c}^{\infty}(\mathbb{R})$ the Fourier transform is an entire function in $s$, of rapid decay on horizontal lines, uniformly so on strips of finite width, so the Mellin transform of $\varphi$ has rapid decay vertically. This allows movement of the contour: for compactly-supported $\varphi$, Mellin inversion is

$$
\varphi(y)=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left(\int_{0}^{\infty} \varphi(r) r^{-s} \frac{d r}{r}\right) y^{s} d s \quad \text { (for any real } \sigma \text { ) }
$$

In the present context, adjust the coordinates so that the Mellin transform is an integral against $\eta\left(a_{y}\right)^{s}=$ $y^{(\ell-1) s}$, and inversion likewise: replace $s$ by $s(\ell-1)$ (and re-adjust the contour):

$$
\begin{aligned}
\varphi(y) & =\frac{(\ell-1)}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left(\int_{0}^{\infty} r^{-(\ell-1) s} \varphi(r) \frac{d r}{r}\right) y^{s(\ell-1)} d s \\
& =\frac{(\ell-1)}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left(\int_{0}^{\infty} \eta\left(a_{r}\right)^{-s} \varphi(r) \frac{d r}{r}\right) \eta\left(a_{y}\right)^{s} d s
\end{aligned}
$$

Identifying $N \backslash G / K \approx A^{+} \approx(0,+\infty)$, this is

$$
\varphi\left(a_{y}\right)=\frac{(\ell-1)}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left(\int_{0}^{\infty} \eta\left(a_{r}\right)^{-s} \varphi\left(a_{r}\right) \frac{d r}{r}\right) \eta\left(a_{y}\right)^{s} d s
$$

Thus,

$$
\varphi(g)=\frac{(\ell-1)}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left(\int_{0}^{\infty} \eta\left(a_{r}\right)^{-s} \varphi\left(a_{r}\right) \frac{d r}{r}\right) \eta(g)^{s} d s \quad(\text { for all } g \in G)
$$

The Mellin transform useful here is

$$
\mathcal{M} \varphi(s)=\int_{0}^{\infty} \eta\left(a_{r}\right)^{-s} \varphi\left(a_{r}\right) \frac{d r}{r}
$$

and the pseudo-Eisenstein series is

$$
\Psi_{\varphi}(g)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi(\gamma g)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi\left(a_{\gamma g}\right)=\frac{(\ell-1)}{2 \pi i} \sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot \eta(\gamma g)^{s} d s
$$

Taking $\sigma=0$ would be natural, but with $\sigma=0$ the double integral (sum and integral) is not absolutely convergent, and the two integrals cannot be interchanged. For $\sigma>1$, the Eisenstein series is absolutely convergent, so the rapid vertical decrease of $\mathcal{M} \varphi$ makes the double integral absolutely convergent, and by Fubini the two integrals can be interchanged:

$$
\Psi_{\varphi}(g)=\frac{(\ell-1)}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s)\left(\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \eta(\gamma g)^{s}\right) d s \quad(\text { with } \sigma>1)
$$

The inner sum is the Eisenstein series $E_{s}(g)$, so

$$
\Psi_{\varphi}(g)=\frac{(\ell-1)}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot E_{s}(g) d s \quad(\text { for } \sigma>1)
$$

Although this does express $\Psi_{\varphi}$ as a superposition of $\Delta$-eigenfunctions, it is unsatisfactory, because it should refer to $\mathcal{M} c_{P} \Psi_{\varphi}$, not to $\mathcal{M} \varphi$, in order to give a direct decomposition formula for functions in the span of the pseudo-Eisenstein series.

We want to move the line of integration to the left, to $\sigma=1 / 2$, stabilized by the functional equation of $E_{s}$. From the corollary [1.11.5] to the Maaß-Selberg relations, there are only finitely-many poles of $E_{s}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$, removing one possible obstacle to the contour move. From the theorem [1.10.1] on meromorphic continuation, even the meromorphically continued $E_{s}\left(g_{o}\right)$ is of polynomial growth vertically in $s$, uniformly in

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bounded strips in $s$, uniformly for $g_{o}$ in compacts. Thus, we may move the contour, picking up finitely-many residues:

$$
\Psi_{\varphi}=\frac{(\ell-1)}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s} d s+(\ell-1) \sum_{s_{o}} \mathcal{M} \varphi\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}} E_{s}
$$

since the poles of $E_{s}$ are simple and $\mathcal{M} \varphi$ is entire. The $1 / 2 \pi i$ from inversion cancels the $2 \pi i$ in the residue formula. By the adjunction/unwinding property of $\Psi_{\varphi}$, on $\operatorname{Re}(s)=\frac{1}{2}$,

$$
\begin{gathered}
\left\langle\Psi_{\varphi}, E_{s}\right\rangle=\int_{\Gamma \backslash X} \Psi_{\varphi} \cdot E_{1-s}=\int_{\Gamma_{\infty} \backslash X} \varphi \cdot c_{P} E_{1-s}=\int_{\Gamma_{\infty} \backslash X}\left(\eta^{1-s}+c_{1-s} \eta^{s}\right) \cdot \varphi \\
=\int_{0}^{\infty}\left(\eta^{1-s}+c_{1-s} \eta^{s}\right) \cdot \varphi(y) \frac{d y}{y^{\ell}}=\int_{0}^{\infty}\left(\eta^{-s}+c_{1-s} \eta^{-(1-s)}\right) \cdot \varphi(y) \frac{d y}{y}=\mathcal{M} \varphi(s)+c_{1-s} \mathcal{M} \varphi(1-s)
\end{gathered}
$$

The integral part of the expression of $\Psi_{\varphi}$ in terms of Eisenstein series can be folded in half, integrating from $\frac{1}{2}+i 0$ to $\frac{1}{2}+i \infty$ rather than from $\frac{1}{2}-i \infty$ to $\frac{1}{2}+i \infty$ :

$$
\begin{gathered}
\Psi_{\varphi}-(\text { residual part })=\frac{(\ell-1)}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot E_{s}(g) d s=\frac{(\ell-1)}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s}+\mathcal{M} \varphi(1-s) E_{1-s} d s \\
=\frac{(\ell-1)}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s}+\mathcal{M} \varphi(1-s) c_{1-s} E_{s} d s=\frac{(\ell-1)}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle E_{s} d s
\end{gathered}
$$

by the functional equation and the computation of $\left\langle\Psi_{\varphi}, E_{s}\right\rangle$ just above. The integral can be written as an integral over the whole line $\operatorname{Re}(s)=\frac{1}{2}$, by the functional equation of $E_{s}$ and dividing by 2 :

$$
\Psi_{\varphi}-(\text { residual part })=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \cdot E_{s} d s
$$

It remains to explicate the finitely-many residues which appear. The notation is normalized so that in all these examples there is a pole at $s=1$. The coefficient $\mathcal{M} \varphi(1)$ is

$$
\mathcal{M} \varphi(1)=\int_{0}^{+\infty} \varphi\left(a_{y}\right) \eta^{-1} \frac{d y}{y}=\int_{0}^{+\infty} \varphi\left(a_{y}\right) \frac{d y}{y^{\ell}}=\int_{\Gamma_{\infty} \backslash X} \varphi\left(n_{x} a_{y}\right) \frac{d x d y}{y^{\ell}}
$$

giving $(N \cap \Gamma) \backslash N$ total measure 1. Winding up,
$\mathcal{M} \varphi(1)=\int_{\Gamma \backslash \mathfrak{H}} \sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi\left(\gamma n_{x} a_{y}\right) \frac{d x d y}{y^{\ell}}=\int_{\Gamma \backslash \mathfrak{H}} \Psi_{\varphi}\left(n_{x} a_{y}\right) \frac{d x d y}{y^{\ell}}=\int_{\Gamma \backslash \mathfrak{H}} \Psi_{\varphi}\left(n_{x} a_{y}\right) \cdot 1 \frac{d x d y}{y^{\ell}}=\left\langle\Psi_{\varphi}, 1\right\rangle$
That is, $\mathcal{M} \varphi(1)$ is the inner product of $\Psi_{\varphi}$ with the constant function 1 . For $\Gamma=S L_{2}(\mathbb{Z})$ and $\Gamma=S L_{2}(\mathbb{Z}[i])$, special arguments [2.B] easily show that the only pole of $E_{s}$ in the half-plane $\operatorname{Re}(s) \geq 1 / 2$ is at $s_{o}=1$, is simple, and the residue is a constant function. However, these special arguments do not easily extend to $S p_{1,1}^{*}(\mathfrak{o})$ or $S L_{2}(\mathfrak{o})$, and, in any case, these are meant to be examples toward a larger context. As the pseudo-Eisenstein series do, $E_{s}$ fits into an adjunction

$$
\int_{\Gamma \backslash X} E_{s} \cdot f=\int_{\Gamma_{\infty} \backslash X} \eta^{s} \cdot c_{P} f \quad(\text { for } f \text { on } \Gamma \backslash X)
$$

whenever the implied integrals converge absolutely. Via the analytic continuation of $E_{s}$, the adjunction asserts that integrals against Eisenstein series are Mellin transforms of constant terms:

$$
\int_{\Gamma \backslash X} E_{s} \cdot f=\int_{0}^{\infty} c_{P} f\left(a_{y}\right) \eta^{s} \frac{d y}{y^{\ell}}=\int_{0}^{\infty} c_{P} f\left(a_{y}\right) \eta^{-(1-s)} \frac{d y}{y}=\mathcal{M} c_{P} f(1-s)
$$

Again, at a pole $s_{o}$ of $E_{s}$ in $\operatorname{Re}(s)>\frac{1}{2}, c_{s}$ also has a pole of the same order. Since $c_{s} \cdot c_{1-s}=1$, necessarily $c_{1-s}$ has a zero at $s_{o}$. Thus, from

$$
\mathcal{M} c_{P} \Psi_{\varphi}(s)=\mathcal{M} \varphi(s)+c_{1-s} \mathcal{M} \varphi(1-s)
$$

at a pole $s_{o}$ of $E_{s}$ we have

$$
\mathcal{M} c_{P} \Psi_{\varphi}\left(s_{o}\right)=\mathcal{M} \varphi\left(s_{o}\right)+c_{1-s_{o}} \mathcal{M} \varphi\left(1-s_{o}\right)=\mathcal{M} \varphi\left(s_{o}\right)+0 \cdot \mathcal{M} \varphi\left(1-s_{o}\right)=\mathcal{M} \varphi\left(s_{o}\right)
$$

That is, the value $\mathcal{M} c_{P} \Psi_{\varphi}$ at $s_{o}$ is just the value of $\mathcal{M} \varphi$, so the coefficients appearing in the decomposition of $\Psi_{\varphi}$ are intrinsic. Thus, the decomposition above has an intrinsic form as in the statement of the theorem.

To have an $L^{2}$ assertion and Plancherel require somewhat more care in the argument, as in the following section.

### 1.13 Plancherel for pseudo-Eisenstein series

A refined form of the previous theorem, proving convergence of the integral as a $C^{\infty}(\Gamma \backslash G / K)$-valued integral, from a corresponding result for behavior of Fourier inversion integrals, gives an immediate proof of a Plancherel theorem for pseudo-Eisenstein series.
[1.13.1] Theorem: (Function-valued form) Let $s_{o}$ run over poles of $E_{s}$ in $\operatorname{Re}(s)>\frac{1}{2}$. For $\varphi \in C_{c}^{\infty}(N \backslash X)=$ $C_{c}^{\infty}(N \backslash G)^{K}$, the pseudo-Eisenstein series $\Psi_{\varphi}$ is expressible in terms of genuine Eisenstein series, by an integral converging as a Gelfand-Pettis $C^{\infty}(\Gamma \backslash G / K)$-valued integral:

$$
\Psi_{\varphi}=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \cdot E_{s} d s+(\ell-1) \sum_{s_{o}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s}
$$

writing $\left\langle\Psi_{\varphi}, E_{s}\right\rangle=\int_{\Gamma \backslash G} \Psi_{\varphi} \cdot \bar{E}_{s}$ even though $E_{s}$ is not in $L^{2}$.
Proof: Let $\psi_{\xi}(x)=e^{i \xi x}$. The integral expressing Fourier inversion for Schwartz functions $f$ on the real line

$$
f(x)=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\left(\int_{-\infty}^{\infty} f(u) \bar{\psi}_{\xi}(u) d u\right) \psi_{\xi}(x) d \xi=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \psi_{\xi}(x) \cdot \widehat{f}(\xi) d \xi
$$

does not express $f$ as a superposition of Schwartz functions, but as a superposition of exponentials $x \rightarrow e^{2 \pi i \xi x}$. These exponentials are not Schwartz, and are not $L^{2}$. But the Fourier inversion integral does converge as a Gelfand-Pettis integral with values in the Fréchet space $C^{\infty}(\mathbb{R})$, by [14.3]. Changing coordinates to give Mellin inversion for functions on $(0,+\infty) \approx N \backslash G / K$ gives convergence as Gelfand-Pettis integral with values in $C^{\infty}(0,+\infty) \approx C^{\infty}(N \backslash G / K) \subset C^{\infty}(G)$, with its Fréchet-space structure as in [13.5].

By the same steps as in the proof of the numerical form of the theorem,

$$
\Psi_{\varphi}=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s} d s+(\ell-1) \sum_{s_{o}} \mathcal{M} \varphi\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}} E_{s}
$$

as a $C^{\infty}(G)$-valued Gelfand-Pettis integral. As in [13.6] and the analogue for $G$ as in [6.2], [6.4], left and right translation by $G$ are continuous maps on $C^{\infty}(G)$, so the linear operators of left translation by $\Gamma$ and right translation by $K$ commute with the integral, so the integral converges as a Gelfand-Pettis integral with values in $C^{\infty}(\Gamma \backslash G / K)$. Similarly, the rearrangement to the statement of the theorem preserves this convergence.
[1.13.2] Corollary: For $\varphi, \psi \in C_{c}^{\infty}(N \backslash G / K) \approx C_{c}^{\infty}\left(A^{+}\right) \approx C_{c}^{\infty}(0,+\infty)$,

$$
\left\langle\Psi_{\varphi}, \Psi_{\psi}\right\rangle=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \overline{\left\langle\Psi_{\psi}, E_{s}\right\rangle} d s+(\ell-1) \sum_{s_{o}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \overline{\left\langle\Psi_{\psi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle}
$$

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Proof: For $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$, the map $F \rightarrow \int_{\Gamma \backslash G} F \cdot \bar{f}$ is a continuous linear functional on $F \in C^{o}(\Gamma \backslash G / K)$, so the Gelfand-Pettis property legitimizes the obvious interchange:

$$
\begin{aligned}
\left\langle\Psi_{\varphi}, f\right\rangle & =\left\langle\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle E_{s} d s+(\ell-1) \sum_{s_{o}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s}, \quad f\right\rangle \\
& =\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \overline{\left\langle f, E_{s}\right\rangle} d s+\sum_{s_{o}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \overline{\left\langle f, \operatorname{Res}_{s_{o}} E_{s}\right\rangle}
\end{aligned}
$$

where $\left\langle E_{s}, f\right\rangle$ converges because $f \in C_{c}^{o}(\Gamma \backslash G / K)$. Taking $f=\Psi_{\psi}$ for $\psi \in C_{c}^{\infty}(N \backslash G / K)$ gives the asserted equality.

This corollary looks like an assertion of a Plancherel theorem, that is, inducing (extending by continuity) an isometry from the $L^{2}$ closure of the span of pseudo-Eisenstein series $\Psi_{\varphi}$ with test function data $\varphi$ to spaces of functions on $\frac{1}{2}+i \mathbb{R}$ and a finite-dimensional space generated by residues. What remains to show is surjectivity to a clearly specified space, and orthogonality of the residues to the integrals on $\frac{1}{2}+i \mathbb{R}$, neither of which is surprising.
[1.13.3] Claim: The residues of $E_{s}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$ are in the closure of the space of pseudo-Eisenstein series. Proof: The residues $\operatorname{Res}_{s_{j}} E_{s}$ are in $L^{2}$ by [1.11.5], mutually orthogonal by [1.11.7], and orthogonal to cuspforms by [1.11.8]. By the adjunction property [1.8.2] they are in the closure of the span of the pseudoEisenstein series.

Thus, for test function $\varphi$ and expansion

$$
\Psi_{\varphi}=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \cdot E_{s} d s+(\ell-1) \sum_{s_{o}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s}
$$

the integral is itself in the closure of the span of the pseudo-Eisenstein series. The functions $F$ on $\frac{1}{2}+i \mathbb{R}$ possibly arising as $F(s)=\left\langle\Psi_{\varphi}, E_{s}\right\rangle$ are constrained by the functional equation $E_{1-s}=c_{1-s} E_{s}$ :

$$
\left\langle\Psi_{\varphi}, E_{1-s}\right\rangle=\left\langle\Psi_{\varphi}, c_{1-s} E_{s}\right\rangle=\overline{c_{1-s}} \cdot\left\langle\Psi_{\varphi}, E_{s}\right\rangle=c_{s} \cdot\left\langle\Psi_{\varphi}, E_{s}\right\rangle \quad\left(\text { on } \operatorname{Re}(s)=\frac{1}{2}\right)
$$

Let

$$
V=\left\{F \in L^{2}\left(\frac{1}{2}+i \mathbb{R}\right): F(1-s)=c_{s} F(s)\right\}
$$

[1.13.4] Claim: The images $\left\langle\Psi_{\varphi}, E_{s}\right\rangle \oplus\left(\ldots,\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{j}} E_{s}\right\rangle, \ldots\right)$ are dense in $V \oplus \mathbb{C}^{n}$.
Proof: The residues are in the closure of pseudo-Eisenstein series, so the integral parts of the spectral decompositions are in the closure, as well, by subtraction. The remaining question is identification of the $L^{2}$ closure of the functions $s \rightarrow\left\langle\Psi_{\varphi}, E_{s}\right\rangle$. Test functions $\varphi$ are dense in the Schwartz space, and the map $\varphi \rightarrow \mathcal{M} \varphi$, essentially Fourier transform, is an isomorphism to the Schwartz space on $\frac{1}{2}+i \mathbb{R}$, so the images $\mathcal{M} \varphi$ of test functions are dense in that Schwartz space, which is dense in $L^{2}$. Noting that $\left|c_{s}\right|=1$ on $\operatorname{Re}(s)=\frac{1}{2}$, the averaging map $F(s) \longrightarrow F(s)+c_{1-s} F(1-s)$ is a surjection of $L^{2}\left(\frac{1}{2}+i \mathbb{R}\right)$ to $V$, so the images $\left\langle\Psi_{\varphi}, E_{s}\right\rangle=\mathcal{M} c_{P} \Psi_{\varphi}$ are dense there, so are dense in $V$.

A typical polarization argument finishes the proof of Plancherel. Recall
[1.13.5] Lemma: Let $V$ be a Hilbert space with subspaces $V_{1}$ and $V_{2}$. If $\left|v_{1}+v_{2}\right|^{2}=\left|v_{1}\right|^{2}+\left.v_{2}\right|^{2}$ for every $v_{1} \in V_{1}$ and $v_{2} \in V_{2}$, then $V_{1}$ and $V_{2}$ are orthogonal.
Proof: We aim to show that $\left\langle v_{1}, v_{2}\right\rangle=0$. Adjusting $v_{2}$ by a complex number of absolute value 1 , we may suppose that this inner product is real. Then

$$
4\left\langle v_{1}, v_{2}\right\rangle=\left|v_{1}+v_{2}\right|^{2}-\left|v_{1}-v_{2}\right|^{2}=\left|v_{1}\right|^{2}+\left|v_{2}\right|^{2}-\left(\left|v_{1}\right|^{2}+\left|v_{2}\right|^{2}\right)=0
$$

as claimed.

Thus, we can distinguish the residual part of $\Psi_{\varphi}$ by

$$
\Psi_{\varphi}^{R}=(\ell-1) \sum_{s_{j}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{j}} E_{s}\right\rangle \cdot \operatorname{Res}_{s_{j}} E_{s}
$$

and the continuous part

$$
\Psi_{\varphi}^{C}=\Psi_{\varphi}-\Psi_{\varphi}^{R}=\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s}\right\rangle \cdot E_{s} d s
$$

Both parts are in the closure of the images of pseudo-Eisenstein series, from above. Extending by continuity the relation [1.13.2],

$$
\left|\Psi_{\varphi}\right|_{L^{2}}^{2}=\left|\Psi_{\varphi}^{C}\right|_{L^{2}}^{2}+\left|\Psi_{\varphi}^{R}\right|_{L^{2}}^{2}
$$

and these two parts are orthogonal. We have the corresponding Plancherel theorem:
[1.13.6] Corollary: The map $\Psi_{\varphi} \longrightarrow\left\langle\Psi_{\varphi}, E_{s}\right\rangle \oplus\left(\ldots,\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{j}} E_{s}\right\rangle, \ldots\right)$ with test functions $\varphi$ is an $L^{2}$ isometry to its image in $V \oplus \mathbb{C}^{n}$, and that image is dense in $V \oplus \mathbb{C}^{n}$. Extending by continuity gives an isometry of the $L^{2}$ closure of the space of $\Psi_{\varphi}$ 's to $V \oplus \mathbb{C}^{n}$.
[1.13.7] Remark: Except on smaller subspaces, such as the span of the pseudo-Eisenstein series with testfunction data, the integrals above are not literal, but are the extension-by-continuity of those integrals, as with Fourier transform and Fourier inversion on $L^{2}(\mathbb{R})$.

### 1.14 Automorphic spectral expansion and Plancherel theorem

Combining the decomposition of cuspforms and the decomposition of their orthogonal complement: letting $s_{o}$ run over the poles of $E_{s}$ in $\operatorname{Re}(s)>\frac{1}{2}$, and letting $F$ run over an orthonormal basis for the space of cuspforms on $\Gamma \backslash G / K$,
[1.14.1] Corollary: Functions $f \in L^{2}(\Gamma \backslash G / K)$ have expansions

$$
f=\sum_{\operatorname{cfm} F}\langle f, F\rangle \cdot F+\frac{(\ell-1)}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle f, E_{s}\right\rangle \cdot E_{s} d s+(\ell-1) \sum_{s_{o}}\left\langle f, \operatorname{Res}_{s_{o}} E_{s}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s}
$$

and Plancherel

$$
|f|_{L^{2}(\Gamma \backslash X)}^{2}=\sum_{\operatorname{cfm} F}|\langle f, F\rangle|^{2}+\frac{(\ell-1)}{4 \pi} \int_{-\infty}^{\infty}\left|\left\langle f, E_{s}\right\rangle\right|^{2} d t+(\ell-1) \sum_{s_{o}}\left|\left\langle f, \operatorname{Res}_{s_{o}} E_{s}\right\rangle\right|^{2}
$$

where integrals involving Eisenstein series are isometric extensions, as in the previous section.
Again, although the discrete part of the expansion converges in a straightforward $L^{2}$ sense, the continuous/Eisenstein part only makes sense as an isometric extension of literal integrals. Nevertheless, the Plancherel theorem is a literal equality.

The factor of $(\ell-1)$ is purely artifactual, and could be normalized away.

### 1.15 Exotic eigenfunctions, discreteness of pseudo-cuspforms

An important variant approach both to the discrete decomposition of the space of cuspforms as above, and to the meromorphic continuation of Eisenstein series, as in [11.5], is the decomposition of spaces of pseudo-cuspforms

$$
\left.L_{b}^{2}(\Gamma \backslash G / K)=\left\{f \in L^{2}(\Gamma \backslash G / K): c_{P} f\left(a_{y}\right)=0 \text { for } y>b\right\} \quad \text { (for fixed } b \geq 1\right)
$$

with respect to a self-adjoint operator ${ }^{[16]} \widetilde{\Delta}_{b}$ closedly related to $\Delta$, but subtly different. For any $b>0$, the corresponding space of pseudo-cuspforms contains the space of genuine cuspforms $L_{o}^{2}(\Gamma \backslash G / K)$. This operator $\widetilde{\Delta}_{b}$ is a pseudo-Laplacian. The basic, surprising result is
[1.15.1] Theorem: $L_{b}^{2}(\Gamma \backslash \underset{\sim}{~} G / K)$ is a direct sum of eigenspaces for the pseudo-Laplacian $\widetilde{\Delta}_{b}$, each of finite dimension. In particular, $\widetilde{\Delta}_{b}$ has compact resolvent. (Proof in [10.3].)

Without further information, this does not instantly prove that the subspace consisting of genuine cuspforms decomposes discretely for $\widetilde{\Delta}_{b}$, because the description [9.2] of the domain $\widetilde{\Delta}_{b}$ puts technical requirements on possible eigenfunctions.

In any case, for $b \gg 1$, the space $L_{b}^{2}(\Gamma \backslash G / K)$ contains many functions not in the space of genuine cuspforms, for example, pseudo-Eisenstein series $\Psi_{\varphi}$ with data $\varphi$ supported in the interval [1, b]. As in [1.12], these are expressible as integrals of genuine Eisenstein series. However, by the theorem, apparently these pseudo-Eisenstein series are (infinite) sums of $L^{2}$-eigenfunctions for $\widetilde{\Delta}_{b}$ orthogonal to cuspforms. Further, truncations $\wedge^{b} E_{s_{o}}$ of genuine Eisenstein series are square-integrable, by [1.11.3] or [1.11.4], for any $s_{o} \in \mathbb{C}$ away from the poles of $s \rightarrow E_{s}$. By [1.12], these truncations are expressible as integrals of genuine Eisenstein series, but, by the theorem here, are also infinite sums of $L^{2}$-eigenfunctions for $\widetilde{\Delta}_{b}$. Thus, evidently, there are many exotic eigenfunctions for $\widetilde{\Delta}_{b}$, pseudo-cuspforms in a strong sense. Indeed,
[1.15.2] Corollary: The eigenfunctions for $\widetilde{\Delta}_{b}$ in $L_{b}^{2}(\Gamma \backslash G / K)$ with eigenvalues $\lambda=s(s-1)<-1 / 4$ are exactly the truncated Eisenstein series $\wedge^{b} E_{s}$ with $c_{P} E\left(a_{b}\right)=0$. (Proof in [10.4].)

These particular truncated Eisenstein series are indeed not smooth. The slightly non-intuitive nature of the operator $\widetilde{\Delta}_{b}$ explains the situation, in [10.4]. For example, in addition to meeting the Gelfand condition of constant-term vanishing about height $b$ :
[1.15.3] Corollary: An $L^{2}$-eigenfunction $u$ for $\widetilde{\Delta}_{b}$ with eigenvalue $\lambda$ satisfies $\left(\widetilde{\Delta}_{b}-\lambda\right) u=0$ locally at heights above $b$.

[^2]
## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

1. Groups $K_{v}=G L_{2}\left(\mathfrak{o}_{v}\right) \subset G_{v}=G L_{2}\left(k_{v}\right)$
2. Discrete subgroup $G L_{2}(k) \subset G L_{2}(\mathbb{A})$, reduction theory
3. Invariant measures
4. Hecke operators, integral operators
5. Decomposition by central characters
6. Discrete decomposition of cuspforms
7. Pseudo-Eisenstein series
8. Eisenstein series
9. Meromorphic continuation of Eisenstein series
10. Truncation and Maaß-Selberg relations
11. Decomposition of pseudo-Eisenstein series: level one
12. Decomposition of pseudo-Eisenstein series: higher level
13. Plancherel for continuous/Eisenstein spectrum: level one
14. Spectral expansion, Plancherel theorem: level one
15. Exotic eigenfunctions, discreteness of pseudo-cuspforms

Appendix A: compactness of $\mathbb{J}^{1} / k^{\times}$
Appendix B: meromorphic continuation
Appendix C: Hecke-Maaß periods of Eisenstein series
This chapter treats a slightly less elementary example, automorphic forms on $G L_{2}(k) \backslash G L_{2}(\mathbb{A})$ for a number field $k$. The shape of the group elements is still two-by-two matrices, but the contents are not the purely archimedean $\mathbb{R}, \mathbb{C}, \mathbb{H}$ of the previous chapter, now involving $p$-adic and adelic scalars. Among several advantages, this viewpoint consistently maintains the unicuspidality of the quotient. In contrast, a classical approach to congruence subgroups of $S L_{2}(\mathbb{Z})$ apparently produces an ever-growing number of cusps, and for Hilbert-Blumenthal groups $G L_{2}(\mathfrak{o})$ for rings of integers $\mathfrak{o}$ of totally real ${ }^{[17]}$ number fields $k$, even at level one, the number of cusps is a class number. Miraculously, in the adelic formulation, there is only one cusp, regardless of class number or congruence conditions. That is, a single adelic Siegel set covers the quotient, as below in [2.2].

In fact, very little subtle information about $p$-adic numbers or adeles or ideles is used. For most purposes, it is merely the shape of matrices that matters, or their structural role, so things can be cast in a form that scarcely refers to details of the distinctions. The significant result is the compactness of $\mathbb{J}^{1} / k^{\times}$, in the appendix [2.A]. Earlier in the chapter, we prove $p$-adic and archimedean Iwasawa and Cartan decompositions from the most basic features of completions of number fields, with the incidental goal of practicing the relevant physical intuition and noting the truly relevant aspects.

Another point of this example is to accommodate more complicated data in Eisenstein series. With or without congruence conditions, number fields beyond $\mathbb{Q}$ have non-trivial grossencharacters (Hecke characters), and apart from complex quadratic extensions there are always unramified grossencharacters. For non-trivial ideal class groups, there are non-trivial ideal class characters. Using $G L_{2}(\mathbb{A})$ unifies these seemingly disparate features. Thus, the decomposition [2.11-2.12] of pseudo-Eisenstein series involves not only the continuous parameter $s$, but at least one discrete parameter $\chi$ running through grossencharacters with various constraints on ramification. Further, congruence conditions specify further data in Eisenstein series. The functional equation(s) of Eisenstein series will no longer relate one Eisenstein series to itself under $s \rightarrow 1-s$, but must tell how the further data transforms. Suggested by physical analogues, the description of the transformation of the further data is often called a scattering matrix or scattering operator.

Most of the analytical archimedean issues in later chapters are already well illustrated by the previous chapter. The present chapter illustrates interaction of those archimedean issues with $p$-adic.
[17] A finite extension $k$ of $\mathbb{Q}$ is totally real when all archimedean completions are isomorphic to $\mathbb{R}$, rather than to $\mathbb{C}$.

$$
\text { 2.1 Groups } K_{v}=G L_{2}\left(\mathfrak{o}_{v}\right) \subset G_{v}=G L_{2}\left(k_{v}\right)
$$

Throughout this chapter, $k$ is a number field. ${ }^{[18]}$ Let $\mathfrak{o}$ be its ring of algebraic integers. Denote the various archimedean and $p$-adic (non-archimedean) completions by $k_{v}$, where $v<\infty$ means non-archimedean, and $v \mid \infty$ means archimedean. For non-archimedean $v$, let $\mathfrak{o}_{v}$ be the local integers. Normalize all the norms $|\cdot|_{v}$ so that the product formula $\prod_{v}|t|_{v}=1$ holds for $t \in k^{\times}$, preferably by taking the norm in $k_{v}$ to be the composition of the Galois norm to the corresponding completion $\mathbb{Q}_{p}$ of $\mathbb{Q}$ and then the standard $p$-adic norm on $\mathbb{Q}_{p}$, by $|t|_{v}=\left|N_{k_{v} / \mathbb{Q}_{p}} t\right|_{p}$, and similarly for archimedean places. When useful, $\varpi_{v}$ will be a local parameter at a non-archimedean place $v$, that is, a prime element in $\mathfrak{o}_{v}$. Let $\mathbb{A}, \mathbb{J}$ be the adeles and ideles of $k$.

Let $G_{v}=G L_{2}\left(k_{v}\right)$. Let $Z_{v}$ be the center of $G_{v}$. Temporarily let $r$ be the number of non-isomorphic archimedean completions of $k$, thus not counting a complex completion and its conjugate as 2 , but just 1. That is, $[k: \mathbb{Q}]=r_{1}+2 r_{2}$ where $r_{1}$ is the number of real completions, and $r_{2}$ the number of complex, and $r=r_{1}+r_{2}$. Let $Z^{+}$be the positive real scalar matrices diagonally imbedded across all archimedean $v$, by the map

$$
\delta: t \longrightarrow\left(\ldots, t^{1 / r}, \ldots\right) \quad(\text { for } t>0)
$$

This map $\delta$ gives a section of the idele norm map $|t|=\prod_{v}\left|t_{v}\right|_{v}$, in that $|\delta(t)|=t$. As usual, let

$$
P_{v}=\left\{\left(\begin{array}{cc}
* & * \\
0 & *
\end{array}\right) \in G_{v}\right\} \quad N_{v}=\left\{\left(\begin{array}{cc}
1 & * \\
0 & 1
\end{array}\right) \in G_{v}\right\} \quad M_{v}=\left\{\left(\begin{array}{cc}
* & 0 \\
0 & *
\end{array}\right) \in G_{v}\right\}
$$

We have already noted the compact subgroups $K_{v} \approx S O_{2}(\mathbb{R}) \subset S L_{2}(\mathbb{R})$ and $K_{v} \approx S U(2) \subset S L_{2}(\mathbb{C})$ for archimedean completions $k_{v} \approx \mathbb{R}$ and $k_{v} \approx \mathbb{C}$, and the corresponding Iwasawa decompositions [1.3].
[2.1.1] Claim: For $v<\infty$, the subgroup

$$
K_{v}=G L_{2}\left(\mathfrak{o}_{v}\right)=\left\{\text { two-by-two matrices with entries in } \mathfrak{o}_{v} \text { and determinant in } \mathfrak{o}_{v}^{\times}\right\}
$$

is a compact, open subgroup of $G_{v}=G L_{2}\left(k_{v}\right)$. We have Iwasawa decomposition

$$
G_{v}=P_{v} \cdot K_{v}=N_{v} \cdot M_{v} \cdot K_{v}
$$

and Cartan decomposition

$$
G_{v}=K_{v} \cdot M_{v} \cdot K_{v}
$$

Proof: The local fields $k_{v}$ are finite-dimensional vectorspaces over respective $\mathbb{Q}_{p}$ and $\mathbb{R}$, so are locally compact. For non-archimedean $v$, the local integers are both closed and open:

$$
\mathfrak{o}_{v}=\left\{x \in k_{v}:|x|_{v} \leq 1\right\}=\left\{x \in k_{v}:|x|_{v}<\left|\varpi^{-1}\right|_{v}\right\}
$$

Similarly for the local units:

$$
\mathfrak{o}_{v}^{\times}=\left\{x \in k_{v}:|x|_{v}=1\right\}=\left\{x \in k_{v}:\left|\varpi_{v}\right|_{v}<|x|_{v}<\left|\varpi^{-1}\right|_{v}\right\}
$$

From this, the conditions defining the subgroups $K_{v}$ are both open and closed. Since $K_{v}$ is a closed subset of the compact set

$$
\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right): a, b, c, d \in \mathfrak{o}_{v}\right\} \approx \mathfrak{o}_{v} \times \mathfrak{o}_{v} \times \mathfrak{o}_{v} \times \mathfrak{o}_{v}
$$

it is compact. Given $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G_{v}$, either $|c|_{v} \geq|d|_{v}$ or $|c|_{v} \leq|d|_{v}$, so either $d / c \in \mathfrak{o}_{v}$ or $c / d \in \mathfrak{o}_{v}$, respectively. Thus, either

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left(\begin{array}{cc}
-d / c & 1 \\
1 & 0
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
0 & *
\end{array}\right) \quad\left(\operatorname{when}\left(\begin{array}{cc}
-d / c & 1 \\
1 & 0
\end{array}\right) \in K_{v}\right)
$$

[^3]or
\[

\left($$
\begin{array}{ll}
a & b \\
c & d
\end{array}
$$\right) \cdot\left($$
\begin{array}{cc}
1 & 0 \\
-c / d & 1
\end{array}
$$\right)=\left(\begin{array}{cc}

* \& * <br>
0 \& *
\end{array}\right) \quad\left(when\left($$
\begin{array}{cc}
1 & 0 \\
-c / d & 1
\end{array}
$$\right) \in K_{v}\right)
\]

giving the Iwasawa decomposition.
The Cartan decomposition is a corollary of the structure theorem for finitely-generated modules over a principal ideal domain such as $\mathfrak{o}_{v}$, as follows. Given $g \in G_{v}$, multiply by a scalar matrix (in $M_{v}$ ) so that all entries of the modified $g$ are in $\mathfrak{o}_{v}$. (Of course, this does not at all assure that the determinant is in $\mathfrak{o}_{v}^{\times}$.) The columns of such $g$ generate a rank-two $\mathfrak{o}_{v}$-submodule $V$ of $\mathfrak{o}_{v}^{2}$. By the structure theorem, there is an $\mathfrak{o}_{v}$-basis $f_{1}, f_{2}$ for $\mathfrak{o}_{v}^{2}$ and $d_{1}, d_{2}$ in $\mathfrak{o}_{v}$ such that $V=\mathfrak{o}_{v} d_{1} f_{1}+\mathfrak{o}_{v} d_{2} f_{2}$. Since $\left\{d_{j} f_{j}\right\}$ is another $\mathfrak{o}_{v}$-basis for $V$, there is $h \in K_{v}$ expressing the two columns of $g$ as $\mathfrak{o}_{v}$-linear combinations of $d_{1} f_{1}, d_{2} f_{2}$ (and vice-versa). That is, in terms of matrix multiplication, writing $d_{1} f_{1}, d_{2} f_{2}$ as column vectors,

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=\left(\begin{array}{ll}
d_{1} f_{1} & d_{2} f_{2}
\end{array}\right) \cdot h
$$

At the same time, there is $h^{\prime} \in K_{v}$ such that $h^{\prime} e_{1}=f_{1}$ and $h^{\prime} e_{2}=f_{2}$, where $\left\{e_{j}\right\}$ is the standard $\mathfrak{o}_{v}$-basis for $\mathfrak{o}_{v}^{2}$. Thus,

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)=h^{\prime} \cdot\left(\begin{array}{cc}
d_{1} & 0 \\
0 & d_{2}
\end{array}\right) \cdot h
$$

giving the $p$-adic Cartan decomposition.
Unlike the archimedean situation, the compact $K_{v}$ has substantial intersections with both $N_{v}$ and $M_{v}$, so with $P_{v}$. Indeed, since $k_{v}$ is an ascending union $k_{v}=\bigcup_{\ell \geq 0} \varpi^{-\ell} \cdot \mathfrak{o}_{v}$, the subgroup $N_{v}$ is an ascending union of compact, open subgroups:

$$
N_{v}=\bigcup_{\ell \geq 0}\left(\begin{array}{cc}
1 & \varpi_{v}^{-\ell} \mathfrak{o}_{v} \\
0 & 1
\end{array}\right)
$$

Again unlike the archimedean situation, $K_{v}$ has a neighborhood basis at 1 consisting of compact, open subgroups, namely, the (local) principal congruence subgroups

$$
K_{v, \ell}=\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in G L_{2}\left(\mathfrak{o}_{v}\right): a=1 \bmod \varpi^{\ell}, b=0 \bmod \varpi^{\ell}, c=0 \bmod \varpi^{\ell}, \text { and } d=1 \bmod \varpi^{\ell}\right\}
$$

The corresponding adele group is $G_{\mathbb{A}}=G L_{2}(\mathbb{A})$, meaning two-by-two matrices with entries in $\mathbb{A}$, with determinant in the ideles $\mathbb{J}$. This group is also an ascending union (colimit) of products

$$
G_{S}=\prod_{v \in S} G_{v} \times \prod_{v \notin S} K_{v} \quad(S \text { a finite set of places } v, \text { including archimedean places })
$$

ordering finite sets $S$ (of places $v$ ) by containment. Similarly, $P_{\mathbb{A}}, M_{\mathbb{A}}, N_{\mathbb{A}}$, and $Z_{\mathbb{A}}$ are the adelic forms of those groups, that is, obtained by allowing entries in $\mathbb{A}$, or, equivalently, as colimits of products of local groups. Let $K_{\mathbb{A}}=\prod_{v} K_{v} \subset G_{\mathbb{A}}$. Let $\delta:(0, \infty) \rightarrow \mathbb{J}$ the usual diagonal imbedding of $(0, \infty)$ to the archimedean component of the ideles by $\delta(t)=\left(\ldots, t^{1 / d_{v}}, \ldots\right)$ where $d_{v}$ is the local degree, so that $\delta:(0, \infty) \rightarrow \mathbb{J}$ gives a one-sided inverse to the idele norm $|\alpha|=\prod_{v \leq \infty}|\alpha|_{v}$. Then

$$
Z_{\mathbb{A}} / Z^{+} Z_{k} \approx \mathbb{J} / \delta(0, \infty) \cdot k^{\times} \approx \mathbb{J}^{1} / k^{\times}=\text {compact }
$$

where $Z_{k}$ is invertible scalar matrices with entries in $k$. The compactness is non-trivial [2.A], but standard.

### 2.2 Discrete subgroup $G L_{2}(k) \subset G L_{2}(\mathbb{A})$, reduction theory

Let $G_{k}=G L_{2}(k), P_{k}, M_{k}, N_{k}$ be the groups with entries in $k$. Here we demonstrate that a single (adelic) Siegel set surjects to the quotient $G_{k} \backslash G_{\mathbb{A}}$, and that (adelic) Siegel sets behave well. First,
[2.2.1] Claim: $G_{k}$ is a discrete subgroup of $G_{\mathbb{A}}$.
Proof: To show that a subgroup of a topological group is discrete, it suffices to show that there is a neighborhood of the identity containing no element of the subgroup except 1 , since multiplication $U \rightarrow g U$ is homeomorphism of neighborhoods $U$ of 1 to neighborhoods $g U$ of $g$.

We do this in two steps. First, the subgroup $H=\prod_{v \mid \infty} G_{v} \times \prod_{v<\infty} K_{v}$ is an open neighborhood of $1 \in G_{\mathbb{A}}$, so it suffices to show that the group $G_{k} \cap H$ is discrete. The condition on $H$ is that entries are locally integral at all finite places, and the determinants are local units. An element of $k$ that is a local integer everywhere is an integer, and an element of $k^{\times}$that is a local unit everywhere is a unit in $\mathfrak{o}^{\times}$. Thus, $G_{k} \cap H=G L_{2}(\mathfrak{o})$, and it suffices to show that the projection of $G L_{2}(\mathfrak{o})$ to $G_{\infty}=\prod_{v \mid \infty} G_{v}$ is discrete in the latter. The topology on $G_{\infty}$ is the subspace topology from the real vectorspace of 2-by-2 matrices with entries in $k_{\infty}=\prod_{v \mid \infty} k_{v}$, which itself has the product topology. From classical algebraic number theory, o is discrete in $k_{\infty}$, giving the discreteness.

On some occasions, one uses

$$
G^{1}=\left\{g \in G_{\mathbb{A}}:|\operatorname{det} g|=1\right\}
$$

noting that $G_{\mathbb{A}}=Z^{+} \times G^{1}$. The product formula $\prod_{v \leq \infty}|t|_{v}=1$ for $t \in k^{\times}$shows that $G_{k} \subset G^{1}$. In particular, $G_{k}$ is still discrete in $Z^{+} \backslash G_{\mathbb{A}} \approx G^{1}$.

Now define local and global height functions $h_{v}$ and $h$. For $v$-adic completion $k_{v} \approx \mathbb{R}$, let $h_{v}$ be the usual real Hilbert-space norm on $k_{v}^{2} \approx \mathbb{R}^{2}$. To accommodate the product formula, for $k_{v} \approx \mathbb{C}$, let $h_{v}$ be the square of the usual complex Hilbert-space norm on $k_{v}^{2} \approx \mathbb{C}^{2}$. For $k_{v}$ non-archimedean, let $h_{v}\left(x_{1}, x_{2}\right)=\max \left\{\left|x_{1}\right|_{v},\left|x_{2}\right|_{v}\right\}$. Put $h(x)=\prod_{v \leq \infty} h_{v}(x)$. There is good behavior under scalar multiplication, via the product formula: for $t \in k^{\times}$,

$$
h(t \cdot x)=\prod_{v \leq \infty} h_{v}(t \cdot x)=\prod_{v \leq \infty}|t|_{v} \cdot h_{v}(x)=\prod_{v \leq \infty}|t|_{v} \cdot \prod_{v \leq \infty} h_{v}(x)=1 \cdot \prod_{v \leq \infty} h_{v}(x)
$$

Sufficient conditions are given below for finiteness of the product. By design, the isometry groups of the height functions $h_{v}$ are the compact subgroups $K_{v}$ already specified.

For each prime $v$ the group $K_{v}$ is transitive on the collection of vectors in $k_{v}^{2}$ with given norm: at archimedean places, this is because all vectors of a given length can be rotated to each other, while at non-archimedean places the suitable analogue of length is greatest common divisor.

Let $G_{k}, P_{k}, N_{k}, M_{k}, Z_{k}$ be the corresponding groups of matrices with entries in $k$, and use subscript $\mathbb{A}$ to denote the adelic points.

Now we identify a class of vectors with finite height. First, given $x \in k^{2}-\{0\}$, for all but finitely-many $v$ the components of the vector $x$ are all $v$-integral, and generate the local integers $\mathfrak{o}_{v}$. In particular, for all but finitely-many $v$ the $v^{t h}$ local height $h_{v}(x)$ of $x \in k^{r}$ is 1 , and the infinite product for $h(x)$ is a finite product. Write vectors as row vectors, and let $G_{\mathbb{A}}=G L_{2}(\mathbb{A})$ act on the right by matrix multiplication. A non-zero vector $x \in \mathbb{A}^{2}$ is primitive when $x \in\left(k^{2}-\{0\}\right) \cdot G_{\mathbb{A}}$.
[2.2.2] Theorem: For fixed $g \in G_{\mathbb{A}}$ and for fixed $c>0$,

$$
\operatorname{card}\left(k^{\times} \backslash\left\{x \in k^{2}-\{0\}: h(x \cdot g)<c\right\}\right)<\infty
$$

For compact $C \subset G_{\mathbb{A}}$ there are positive implied constants such that

$$
\left.h(x)<_{C} h(x \cdot g)<_{C} h(x) \quad \text { (for all } g \in C, \text { for all primitive } x\right)
$$

Proof: Fix $g \in G_{\mathbb{A}}$. Since $K=K_{\mathbb{A}}=\prod_{v} K_{v}$ preserves heights, via Iwasawa decompositions locally everywhere, we may suppose that $g$ is in the group $P_{\mathbb{A}}$ of upper triangular matrices in $G_{\mathbb{A}}$. Choose
representatives $x=\left(x_{1}, x_{2}\right)$ for non-zero vectors in $k^{\times} \backslash k^{2}$ either of the form $x=\left(1, x_{2}\right)$ or $x=(0,1)$. There is just one vector of the latter shape, so we consider $x=\left(1, x_{2}\right)$ :

$$
x \cdot g=\left(\begin{array}{ll}
1 & x_{2}
\end{array}\right)\left(\begin{array}{ll}
a & b \\
0 & d
\end{array}\right)=\left(\begin{array}{ll}
a & b+x_{2} d
\end{array}\right) \quad\left(\text { for } g=\left(\begin{array}{cc}
a & b \\
0 & d
\end{array}\right)\right)
$$

At each place $v$, including archimedean ones, $\max \left(|a|_{v},\left|b+x_{2} d\right|_{v}\right) \leq h_{v}(x g)$, so

$$
\left|b+x_{2} d\right|_{v} \prod_{w \neq v}|a|_{w} \leq \prod_{\text {all } w} h_{w}(x g)=h(x g)
$$

Since $g$ is fixed, $a$ is fixed, and at almost all places $|a|_{w}=1$. Thus, for $h(x g)<c$,

$$
\left|b+x_{2} d\right|_{v}<c \cdot\left(\prod_{w \neq v}|a|_{w}\right)^{-1}<_{g, c} 1 \quad \quad \text { (for fixed } g, \text { for } x \text { with } h(x g)<c, \text { for all places } v \text { ) }
$$

For the product formula to hold we are using the normalization of norms that $\left|\varpi_{v}\right|=q_{v}^{-1}$, where $\varpi_{v}$ is a local parameter at $v$ and $q_{v}$ is the residue field cardinality at $v$. There are only finitely-many places $v$ with residue field cardinality less than a given constant, so, in fact, $\left|b+x_{2} d\right|_{v} \leq 1$ for $v$ outside a finite set depending on $g$ and $c$. Therefore, $b+x_{2} d$ lies in a compact subset $\Omega$ of $\mathbb{A}$ depending on $g, c$. Since $b, d$ are fixed, and since $k$ is discrete (and closed by [1.5.3]) in $\mathbb{A}$, the collection of images $\left\{b+x_{2} d: x_{2} \in k\right\}$ is discrete in $\mathbb{A}$. The intersection of a closed, discrete set and a compact set is finite, so collection of $x_{2} \in k$ so that $b+x_{2} d$ lies in $\Omega$ is finite, proving the first assertion.

For the last assertion, use the Cartan decompositions $G_{v}=K_{v} \cdot M_{v} \cdot K_{v}$ from [2.1]. The map $\theta_{1} \times m \times \theta_{2} \longrightarrow \theta_{1} m \theta_{2}$, with $\theta_{1}, \theta_{2} \in K_{v}$ and $m \in M_{v}$, is not an injection, so we cannot immediately infer that for a given compact $C \subset G_{v}$ the set

$$
\left\{m \in M_{v}: \text { for some } c \in C, c \in K_{v} m K_{v}\right\}
$$

is compact. Since $K_{v}$ is compact, $C^{\prime}=K_{v} \cdot C \cdot K_{v}$ is compact, and now $\theta_{1} m \theta_{2} \in C^{\prime}$ with $\theta_{i} \in K_{v}$ implies $m \in C^{\prime} \cap M_{v}$, which is compact. Thus, any compact subset of $G_{\mathbb{A}}$ is contained in a set $\left\{\theta_{1} m \theta_{2}: \theta_{1}, \theta_{2} \in K, m \in C_{M}\right\}$ with compact $C_{M} \subset M_{\mathbb{A}}$. Since $K$ preserves heights and since the set of primitive vectors is stable under $K$, the set of values $\{h(x g) / h(x): x$ primitive, $g \in C\}$ is contained in a set

$$
\left\{\frac{h(x \delta)}{h(x)}: x \text { primitive, } m \in C_{M}\right\}
$$

for some compact $C_{M} \subset M$. Letting the diagonal entries of $m$ be $m_{i}$,

$$
0<\inf _{m \in C_{M}} \inf _{i}\left|m_{i}\right| \leq \frac{h(x m)}{h(x)} \leq \sup _{m \in C_{M}} \sup _{i}\left|m_{i}\right|<+\infty
$$

This gives the desired bound.
To compare to the purely archimedean height functions $\eta$ used in the four earlier examples, for $g$ uppertriangular,

For example, with $k_{v} \approx \mathbb{R}$, for $g=n_{x} a_{y}=\left(\begin{array}{cc}1 & x \\ 0 & 1\end{array}\right)\left(\begin{array}{cc}\sqrt{y} & 0 \\ 0 & 1 / \sqrt{y}\end{array}\right) \in S L_{2}(\mathbb{R})$,

$$
h_{v}\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot g\right)=h_{v}\left(\begin{array}{ll}
0 & \left.\frac{1}{\sqrt{y}}\right)=\frac{1}{\sqrt{y}}
\end{array}\right.
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

However, we want local height functions which are right $K_{v}$-invariant and

$$
\eta_{v}\left(\begin{array}{ll}
a & b \\
0 & d
\end{array}\right)=\left|\frac{a}{d}\right|_{v}
$$

so put

$$
\eta_{v}(g)=|\operatorname{det} g|_{v} \cdot h_{v}\left(\left(\begin{array}{ll}
0 & 1) \cdot g)^{-2} \quad\left(\left(\text { for } g \in G_{v}\right)\right)
\end{array}\right.\right.
$$

and $\eta(g)=\prod_{v} \eta_{v}(g)$. This matches earlier use for $S L_{2}(\mathbb{R})$. Left multiplication by $\gamma \in G_{k}$ does not change $|\operatorname{det} g|$ (with idele norm), because of the product formula:

$$
|\operatorname{det}(\gamma g)|=|\operatorname{det} \gamma \cdot \operatorname{det} g|=|\operatorname{det} \gamma| \cdot|\operatorname{det} g|=1 \cdot|\operatorname{det} g|
$$

[2.2.3] Lemma: $\eta$ is left $P_{k}$-invariant, and $Z_{\mathbb{A}}$-invariant.
Proof: Via the product formula, with $p=\left(\begin{array}{cc}a & b \\ 0 & d\end{array}\right) \in P_{k}$,

$$
\begin{aligned}
& \eta(p \cdot g)=|\operatorname{det} p g| \cdot h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot p g\right)^{-2}=|\operatorname{det} p g| \cdot h\left(\left(\begin{array}{ll}
0 & d
\end{array}\right) \cdot g\right)^{-2} \\
& =|\operatorname{det} p| \cdot|\operatorname{det} g| \cdot|d| \cdot h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot g\right)^{-2}=|\operatorname{det} g| \cdot h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot g\right)^{-2}
\end{aligned}
$$

For the center-invariance, with $z=\left(\begin{array}{cc}t & 0 \\ 0 & t\end{array}\right) \in Z_{\mathbb{A}}$,

$$
\begin{aligned}
& \eta(z \cdot g)=|\operatorname{det} z g| \cdot h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot z g\right)^{-2}=|\operatorname{det} z g| \cdot h\left(\left(\begin{array}{ll}
0 & t
\end{array}\right) \cdot g\right)^{-2} \\
& =\left|t^{2}\right| \cdot|\operatorname{det} g| \cdot|t|^{-2} \cdot h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot g\right)^{-2}=|\operatorname{det} g| \cdot h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot g\right)^{-2}
\end{aligned}
$$

as claimed.
[2.2.4] Corollary: For fixed $g \in G_{\mathbb{A}}$, there are finitely-many $\gamma \in P_{k} \backslash G_{k}$ such that $\eta(\gamma \cdot g)>\eta(g)$.
Proof: There is a natural bijection

$$
k^{\times} \backslash\left(k^{2}-\{0\}\right) \longleftrightarrow P_{k} \backslash G_{k} \quad \text { by } \quad k^{\times} \cdot\left(\begin{array}{ll}
c & d
\end{array}\right) \longleftrightarrow P_{k} \cdot\left(\begin{array}{cc}
* & * \\
c & d
\end{array}\right)
$$

for any invertible matrix with bottom row $\left(\begin{array}{cc}c & d\end{array}\right)$. Indeed, $G_{k}$ is transitive on non-zero vectors, and $P_{k}$ is the stabilizer, acting on the right, of the line (minus a point) $\left(\begin{array}{ll}0 & *\end{array}\right)=k^{\times} \cdot\left\{\left(\begin{array}{ll}0 & 1\end{array}\right)\right\}$. The theorem shows that there are finitely-many $x \in k^{\times} \backslash\left(k^{2}-\{0\}\right)$ such that $h(x g)<c$, that is, such that $h(x g)^{-1}>c^{-1}$. Since $|\operatorname{det} g|$ is $G_{k}$-invariant, the bijection just demonstrated gives the assertion of the corollary.
[2.2.5] Corollary: $\sup _{\gamma \in G_{k}} \eta(\gamma \cdot g)<\infty$, and the sup is attained, and

$$
G_{k} \cdot\left\{g \in G_{\mathbb{A}}: \eta(g) \geq \eta(\gamma \cdot g) \text { for all } \gamma \in G_{k}\right\}=G_{\mathbb{A}}
$$

Proof: The previous corollary shows that the sup is finite and that the sup is attained. Thus, the indicated set is a (possibly redundant) collection of representatives for all orbits, by choosing group elements attaining the sup in their $G_{k}$-orbit.

Critical in legitimizing treatment of truncated Eisenstein series:
[2.2.6] Corollary: Given $t_{o}>0$, there is $t_{1} \gg 1$ such that, for $\eta(g) \geq t_{1}$, if $\eta(\gamma \cdot g) \geq t_{o}$ then $\gamma \in P_{k}$.
Proof: It suffices to take $g=n m$ since $\eta$ is right $K_{\mathbb{A}}$-invariant, invoking Iwasawa. Since $\eta$ is $Z_{\mathbb{A}}$-invariant it suffices to consider $m=\left(\begin{array}{cc}m_{1} & 0 \\ 0 & 1\end{array}\right)$. Adjusting on the left by $M_{k}$, by the compactness lemma [2.A], we can take $m_{1}$ of the special form $m_{1}=m_{o} \cdot \delta(t)$ for $t>0$ and $m_{o}$ in a sufficiently large compact subset of
$\mathbb{J}$, where $\delta:(0,+\infty) \rightarrow \mathbb{J}$ imbeds the ray $(0, \infty)$ at archimedean places. Take compact $C \subset N_{\mathbb{A}}$ such that $N_{k} \cdot C=N_{\mathbb{A}}$. For $v \in k^{2}-\{0\}$,

$$
h(v \cdot n m)=h\left(v \cdot m \cdot m^{-1} n m\right)
$$

For $m$ of the special sort indicated, given $t_{1}>0$, there is compact $C^{\prime} \subset N_{\mathbb{A}}$ such that if $\eta(m)=\left|m_{1} / m_{2}\right| \geq t_{1}$, then $m^{-1} C m \subset C^{\prime}$. Let $(c d) \in k^{2}-\{0\}$. From the second assertion of the theorem, there are constants depending only on $C^{\prime}$ such that, for all (primitive) $x=v \cdot m$,

$$
h(v \cdot m)=h(x)<_{C^{\prime}} h(x \cdot n)<_{C^{\prime}} h(x)=h(v \cdot m) \quad\left(\text { for all } n \in C^{\prime}\right)
$$

Thus, it suffices to treat simply $g=m$. In that case, with $v=(c d)$ with $c \neq 0$,

$$
h(v \cdot m)=h((c d) \cdot m)=h\left(c m_{1} \quad d\right) \geq\left|c m_{1}\right|=|c| \cdot\left|m_{1}\right|=\left|m_{1}\right|
$$

by the product rule, since $c \in k^{\times}$. Thus, with $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G_{k}$, but not in $P_{k}$,

$$
\eta(\gamma \cdot m)=\frac{|\operatorname{det} \gamma m|}{h((c d) \cdot m)^{2}}=\frac{\left|m_{1}\right|}{h\left(\left(\begin{array}{cc}
c & d
\end{array}\right) \cdot m\right)^{2}} \leq \frac{\left|m_{1}\right|}{\left|m_{1}\right|^{2}}=\frac{1}{\left|m_{1}\right|}
$$

With whatever constants are implied in the simplifications in the first part of the proof, a sufficiently high lower bound $\eta(m)=\left|m_{1}\right| \geq t_{1}$ assures that $\eta(\gamma \cdot m)$ is below $t_{o}$.

An element $g \in G_{\mathbb{A}}$ such that $\eta(g) \geq \eta(\gamma \cdot g)$ for all $\gamma \in G_{k}$ is reduced. Given the above preparation, as an application of Dirichlet's pigeon-hole principle, after Minkowski, we can prove
[2.2.7] Theorem: There is a constant $t_{o}>0$ depending on $k$ such that $\eta(g) \geq t_{o}$ for reduced $g \in G_{\mathbb{A}}$,
Proof: Since heights are right $K$-invariant, take $g=n m$ with $n=n_{x} \in N_{\mathbb{A}}$ and $m \in M_{\mathbb{A}}$. Adjusting by the center, take

$$
m=\left(\begin{array}{cc}
y & 0 \\
0 & 1
\end{array}\right) \quad n=n_{x}=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)
$$

with $y \in \mathbb{J}, x \in \mathbb{A}$. Let $\mathbb{J}^{1}$ be the ideles of idele-norm 1 , and let $\delta:(0,+\infty) \rightarrow \mathbb{J}$ by

$$
\delta\left(y_{\infty}\right)=\left(y_{\infty}^{\frac{1}{n}}, \ldots, y_{\infty}^{\frac{1}{n}}, 1,1,1, \ldots\right) \quad(\text { where }, \text { temporarily, } n=[k: \mathbb{Q}])
$$

with non-trivial values at the archimedean components. Then $\mathbb{J}=\mathbb{J}^{1} \times \delta(0,+\infty)$. Let $U=\prod_{v \mid \infty} k_{v}^{\times} \times$ $\prod_{v<\infty} \mathfrak{o}_{v}^{\times}$. The quotient $k^{\times} \backslash \mathbb{J}^{1}$ is compact, by [2.A], so $k^{\times} \backslash \mathbb{J}^{1} / U$ is finite.

Thus, adjusting on the left by $\left\{\left(\begin{array}{cc}m_{1} & 0 \\ 0 & 1\end{array}\right): m_{1} \in k^{\times}\right\}$and on the right by $\left\{\left(\begin{array}{cc}m_{1} & 0 \\ 0 & 1\end{array}\right): m_{1} \in U\right\}$, we can suppose that $y=\delta\left(y_{\infty}\right) \cdot \theta$ with $y_{\infty} \in(0,+\infty)$ and $\theta$ in a finite list $\Theta$ of finite ideles, essentially representatives for the ideal class group. We can take $\theta \in \Theta$ everywhere locally integral at finite places. Write

$$
m=m_{\infty} \cdot m_{\text {fin }} \quad \text { with } \quad m_{\infty}=\left(\begin{array}{cc}
\delta\left(y_{\infty}\right) & 0 \\
0 & 1
\end{array}\right) \quad m_{\text {fin }}=\left(\begin{array}{cc}
\theta & 0 \\
0 & 1
\end{array}\right)
$$

For fixed $\theta \in \Theta$, with

$$
V_{\theta}=\theta\left(\prod_{v<\infty} \mathfrak{o}_{v}\right) \theta^{-1} \quad U_{\theta}=m_{\mathrm{fin}}\left(N_{\mathbb{A}} \cap \prod_{v<\infty} K_{v}\right) m_{\mathrm{fin}}^{-1}
$$

we have

$$
U_{\theta}=\left\{\left(\begin{array}{ll}
1 & u \\
0 & 1
\end{array}\right): u \in V_{\theta}\right\}
$$

Let $\mathbb{A}_{\infty}=k \otimes_{\mathbb{Q}} \mathbb{R}$ be the archimedean component of the adeles. For each fixed $\theta \in \Theta$, acting on $n_{x}$ on the left by $N_{k}$ is equivalent to adjusting $x \in \mathbb{A}$ by $k$. By additive approximation, we can adjust $x$ by $k$ to be in

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

$\mathbb{A}_{\infty}+V_{\theta}$. Right multiplication of $n_{x} m$ by $N_{\mathbb{A}} \cap \prod_{v<\infty} K_{v}$ is equivalent to adjusting $x$ by $V_{\theta}$. Thus, without loss of generality, $x \in \mathbb{A}_{\infty}$. None of these adjustments changed the height $\eta(n m)$.

The inequality

$$
h\left(\left(\begin{array}{ll}
0 & 1
\end{array}\right) \cdot n m\right) \leq h\left(\left(\begin{array}{ll}
1 & -\alpha
\end{array}\right) \cdot n m\right)
$$

holds for every $\alpha \in k$, by the reduced property of $g=n m$. Letting $h_{\mathrm{fin}}=\prod_{v<\infty} h_{v}$ and $h_{\infty}=\prod_{v \mid \infty} h_{v}$, since $\left(\begin{array}{ll}0 & 1\end{array}\right)$ is fixed by $\left(\begin{array}{ll}* & * \\ 0 & 1\end{array}\right)$,

$$
1=h\left(\left(\begin{array}{ll}
0 & 1) \cdot n m) \leq h_{\infty}\left(\delta\left(y_{\infty}\right), x-\alpha_{\infty}\right) \cdot h_{\mathrm{fin}}\left(\theta,-\alpha_{\mathrm{fin}}\right)
\end{array}\right.\right.
$$

where $\alpha_{\infty}$ and $\alpha_{\text {fin }}$ are the projections to the archimedean and finite components of the adeles. This is

$$
\frac{1}{h_{\mathrm{fin}}\left(\theta,-\alpha_{\mathrm{fin}}\right)} \leq \prod_{v \mid \infty}\left(\left|y_{\infty}\right|_{v}^{2 / d_{v}}+\left|x_{v}-\alpha_{v}\right|_{v}^{2 / d_{v}}\right)^{d_{v} / 2}
$$

where $d_{v}$ is the local degree at the $v^{t h}$ archimedean place. We want to use Dirichlet's principle to choose $\alpha \in k$ so that $\left|x_{v}-\alpha_{v}\right|_{v}$ is much smaller than $h_{\text {fin }}\left(\theta,-\alpha_{\mathrm{fin}}\right)$, thereby to give a lower bound on $y_{\infty}$.

Choose a $\mathbb{Z}$-basis $\omega_{1}, \omega_{2}, \ldots$ for $\mathfrak{o}$, and put

$$
F=\left\{\sum_{j} r_{j} \cdot \omega_{j}: \text { each } 0 \leq r_{j}<1\right\} \subset \mathbb{A}_{\infty}
$$

Thus, given $x \in \mathbb{A}_{\infty}$, there is $\beta \in \mathfrak{o}$ such that $x-\beta \in F$. For fixed large $1 \leq \ell \in \mathbb{Z}$, for each integer $1 \leq a \leq \ell^{[k: \mathbb{Q}]}+1$, let $b=b_{a} \in \mathfrak{o}$ such that $a x-b \in F$. Since $F$ is a disjoint union of $\ell^{[k: \mathbb{Q}]}$ translates of $\ell^{-1} F$, by the pigeon-hole principle there are $a, b$ and $a^{\prime}, b^{\prime}$ such that $(a x-b)-\left(a^{\prime} x-b^{\prime}\right) \in \ell^{-1} F$. Thus,

$$
x-\frac{b-b^{\prime}}{a-a^{\prime}} \in \frac{1}{\ell\left(a-a^{\prime}\right)} \cdot F
$$

Put $p=b-b^{\prime} \in \mathfrak{o}, q=a-a^{\prime} \in \mathbb{Z}$, and $\alpha=p / q$. Without loss of generality, $q>0$. With $\mu=\sup _{x \in F, v \mid \infty}\left|x_{v}\right|_{v}^{2 / d_{v}}$, we have

$$
\frac{1}{h_{\mathrm{fin}}\left(\theta,-\frac{p}{q}\right)} \leq \prod_{v \mid \infty}\left(y_{\infty}^{2}+\frac{\mu}{(\ell \cdot q)^{2}}\right)^{d_{v} / 2}=\left(y_{\infty}^{2}+\frac{\mu}{(\ell \cdot q)^{2}}\right)^{[k: \mathbb{Q}] / 2}
$$

Now

$$
h_{\mathrm{fin}}\left(\theta,-\frac{p}{q}\right) \leq \prod_{v<\infty} \max \left\{\left|\theta_{v}\right|_{v},\left|\frac{p}{q}\right|_{v}\right\} \leq \prod_{v<\infty} \max \left\{\left|\theta_{v}\right|_{v}, 1\right\} \cdot \prod_{v<\infty} \max \left\{1,\left|\frac{1}{q}\right|_{v}\right\}=q^{n}
$$

since $\theta \in \Theta$ is everywhere locally integral. Then

$$
\frac{1}{q^{n}} \leq\left(y_{\infty}^{2}+\frac{\mu}{(\ell \cdot q)^{2}}\right)^{n / 2}
$$

or

$$
\frac{1}{q^{2}} \leq y_{\infty}^{2}+\frac{\mu}{(\ell \cdot q)^{2}}
$$

Since $1 \leq q \leq \ell^{n}$, this implies

$$
\frac{1}{\ell^{2 n}} \cdot\left(1-\frac{\mu}{\ell^{2}}\right) \leq \frac{1}{q^{2}} \cdot\left(1-\frac{\mu}{\ell^{2}}\right) \leq y_{\infty}^{2}
$$

Taking $\ell^{2} \geq 2 \mu$ gives a uniform positive lower bound $y_{\infty} \geq t_{1}=\frac{1}{2 \ell^{2}}>0$. For each of the finitely-many $\theta \in \Theta$,

$$
\eta(m)=\eta\left(m_{\infty} \cdot m_{\mathrm{fin}}\right)=\eta\left(\delta\left(y_{\infty}\right)\right) \cdot \eta(\theta)=y_{\infty}^{n} \cdot \eta(\theta) \geq t_{1}^{n} \cdot \min _{\theta \in \Theta} \eta(\theta)
$$

That is, every reduced $g \in G_{\mathbb{A}}$ has $\eta(g)$ bounded from below by that (positive) quantity.
For compact $C \subset N_{\mathbb{A}}$ and $t>0$, the corresponding Siegel set is

$$
\mathfrak{S}_{C, t}=\left\{n m k: n \in C, k \in K_{\mathbb{A}},\left|m_{1} / m_{2}\right| \geq t\right\} \quad\left(\text { where } m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right)
$$

[2.2.8] Corollary: Let $C$ be any compact subset of $N_{\mathbb{A}}$ sufficiently large so that $N_{k} \cdot C=N_{\mathbb{A}}$. With $t_{o}$ as in the theorem, $\mathfrak{S}_{C, t_{o}}$ surjects to the quotient $G_{k} \backslash G_{\mathbb{A}}$. That is, $G_{k} \cdot \mathfrak{S}_{C, t_{o}}=G_{\mathbb{A}}$.
Proof: The theorem asserts that $S=\left\{g: \eta(g) \geq t_{o}\right\}$ surjects to $G_{k} \backslash G_{\mathbb{A}}$. The set $S$ is left $P_{k}$-invariant and left $N_{\mathbb{A}}$-invariant. Thus, we can certainly adjust on the left by $N_{k}$ so that with $g \in n m K$ in Iwasawa coordinates $n \in C$.

### 2.3 Invariant measures

We seldom need explicit formulaic evaluation of integrals on groups $G_{v}=G L_{2}\left(k_{v}\right)$ or their subgroups. Rather, qualitative features of the invariant integrals, such as uniqueness and unwinding properties, play the main roles.

Locally, from [14.4], up to scalar multiples there is a unique right $G_{v^{-}}$invariant measure on $G_{v}$, left $P_{v^{-}}$ invariant measure on $P_{v}$, and (left and right) $K_{v}$-invariant measure on $K_{v}$, for each place $v$. Even though $P_{v} \cap K_{v}$ is non-trivial, given any two of the scalar multiple choices, the third is determined, so that

$$
\int_{G_{v}} f=\int_{P_{v}} \int_{K_{v}} f(p h) d p d h
$$

The idea of the proof from [5.2] and [14.4] is that the group $H=P_{v} \times K_{v}$ acts transitively on $G_{v}$ by $(p \times k)(g)=p^{-1} g k$, with isotropy group $P_{v} \times K_{v}$ at $1 \in G_{v}$. Since the modular function of $P_{v} \times K_{v}$ is inevitably trivial on the compact $P_{v} \cap K_{v}$, there is a unique $H$-invariant measure on $G_{v} \approx\left(P_{v} \cap K_{v}\right) \backslash\left(P_{v} \times K_{v}\right)$. Since the (left and right) $G_{v}$-invariant measure is such, these must be the same, by uniqueness. For example, for $f$ right $K_{v}$-invariant,

$$
\int_{G_{v}} f=\int_{P_{v}} f(p) d p \quad\left(f \text { right } K_{v} \text {-invariant, left } P_{v} \text {-invariant measure } d p \text { on } P_{v}\right)
$$

Even more simply, $P_{v}=N_{v} M_{v} \approx N_{v} \times M_{v}$ has a left (or right) invariant measure given by the product of the invariant measures on $N_{v}$ and $M_{v}$. Archimedean examples were already considered in [1.6], and $p$-adic examples below.

Similarly, globally, there is a unique right $G_{\mathbb{A}}$-invariant measure on $G_{\mathbb{A}}, Z_{\mathbb{A}} \backslash G_{\mathbb{A}}$, and $Z^{+} \backslash G_{\mathbb{A}}$. Given these, there are unique right $G_{\mathbb{A}}$-invariant measures on $G_{k} \backslash G_{\mathbb{A}}, Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$, and $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ such that the corresponding unwindings are correct: for example,

$$
\int_{Z^{+} G_{\mathbb{A}}} f=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}}\left(\sum_{\gamma \in G_{k}} f \circ \gamma\right) \quad\left(\text { for every } f \in C_{c}^{o}\left(Z^{+} G_{\mathbb{A}}\right)\right)
$$

Comparisons between global integrals and products of local integrals are as expected: for $f(g)=\prod_{v} f_{v}\left(g_{v}\right)$ in $C_{c}^{o}\left(G_{\mathbb{A}}\right)$ expressible as a product of functions $f_{v} \in C_{c}^{o}\left(G_{v}\right)$, up to a scalar depending on all the normalizations,

$$
\int_{G_{\mathbb{A}}} f=\prod_{v} \int_{G_{v}} f_{v}
$$

despite the fact that the adele group $G_{\mathbb{A}}$ is not the product of the local groups $G_{v}$, but only the colimit of the products $G_{S}=\prod_{v \in S} G_{v} \times \prod_{v \notin S} K_{v}$. Indeed, in practice $f_{v}$ will be $K_{v}$-invariant for all but finitely-many $v$, so

$$
\int_{G_{\mathbb{A}}} f=\lim _{S} \int_{G_{S}} f=\lim _{S}\left(\prod_{v \in S} \int_{G_{v}} f_{v} \cdot \prod_{v \notin S} \int_{K_{v}} f_{v}\right)=\lim _{S}\left(\prod_{v \in S} \int_{G_{v}} f_{v} \cdot 1\right)=\prod_{v \in S} \int_{G_{v}} f_{v}
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

Nevertheless, on some occasions explicit computations are useful or necessary. Measures and integrals on $\mathbb{R}$ and $\mathbb{C}$ are familiar. On $\mathbb{Q}_{p}$ and its finite extensions, somewhat less so. However, the totally-disconnected nature of $\mathbb{Q}_{p}$ and finite extensions makes measure and integration simpler, at least for nice functions. We treat $\mathbb{Q}_{p}$, and every non-archimedean $k_{v}$ is a finite cartesian product of such. We do not need to prove uniqueness, since this follows for general reasons [14.4].

To give a (regular, Borel) measure it suffices to tell the measure of every open. Since $\mathbb{Q}_{p}$ is a group whose group operation and inversion are continuous, for an invariant measure it suffices to tell the measure of a local basis at 0 , since every translate (coset!) of a given basis element must have the same measure. Such a basis is $p^{\ell} \mathbb{Z}_{p}$. Since these are subgroups, we can easily compare them: for $1 \leq \ell<\ell^{\prime}$ the subgroup $p^{\ell^{\prime}} \mathbb{Z}_{p}$ is of index $p^{\ell^{\prime}-\ell}$ in $p^{\ell} \mathbb{Z}_{p}$. The ratio of measures must be the index. Thus, normalizing everything by taking the measure of $\mathbb{Z}_{p}$ to be 1 , the measure of $p^{\ell} \mathbb{Z}_{p}$ is its index in $\mathbb{Z}_{p}$, namely, $p^{-\ell}$. Larger opens are unions of translates of sets $p^{\ell} \mathbb{Z}_{p}$. This gives the standard additive Haar measure on $\mathbb{Q}_{p}$ for $p<\infty$.

On finite extensions $k_{v}$ of $\mathbb{Q}_{p}$, the same process produces an additive Haar measure giving $\mathfrak{o}_{v}$ total measure 1. For $k_{v} / \mathbb{Q}_{p}$ unramified, this is almost always a good normalization. However, for $k_{v} / \mathbb{Q}_{p}$ ramified, other choices may have advantages, for example with respect to local Fourier transforms.

A multiplicative Haar measure $d^{\times} x$ on $\mathbb{Q}_{p}^{\times}$can be arranged from the additive $d^{+} x$, much as for $\mathbb{R}^{\times}$or $\mathbb{C}^{\times}$, namely, $d^{\times} x=d^{+} x /|x|_{v}$. However, this gives the local units $\mathbb{Z}_{p}^{\times}$measure $\frac{p-1}{p}$, not 1 . Since $\mathbb{Z}_{p}^{\times}$is the unique maximal compact subgroup of $\mathbb{Q}_{p}^{\times}$, we might prefer to give the local units measure 1. A similar device applies to $k_{v}$ for $v<\infty$. In practice, the superscripts are not used, because context explains and determines which measure is meant.

Since $N_{v} \approx k_{v}$, the invariant measure on $N_{v}$ is just the additive Haar measure from $k_{v}$. Since $M_{v} \approx k_{v}^{\times} \times k_{v}^{\times}$, a product of multiplicative Haar measures is the invariant measure.

Much as in the archimedean cases considered earlier, a left-invariant measure on $P_{v}=N_{v} M_{v}$ is $d(n m)=d n d m /|\alpha(m)|_{v}$, where $\alpha\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right)=m_{1} / m_{2}$. That is,

$$
d\left(\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right)=\frac{d^{+} x d^{\times} m_{1} d^{\times} m_{2}}{\left|m_{1} / m_{2}\right|_{v}}=\frac{d^{+} x d^{+} m_{1} d^{+} m_{2}}{\left|m_{1}\right|_{v}^{2}}
$$

[2.3.1] Claim: Quotients $Z^{+} M_{k} \backslash \mathfrak{S}_{C, t}$ of Siegel sets have finite volume.
Proof: The notation has compact $C \subset N_{\mathbb{A}}$ and $t>0$. Letting $K_{\mathbb{A}}=\prod_{v \leq \infty} K_{v}$, up to normalization,

$$
\int_{Z^{+} M_{k} \backslash \mathfrak{S}_{C, t}} 1 d g=\int_{C} 1 d n \cdot \int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} 1 \frac{d m}{|\alpha(m)|} \cdot \int_{K} 1 d k \asymp \int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} 1 \frac{d m}{|\alpha(m)|}
$$

Further, $M_{k} \backslash M_{\mathbb{A}} \approx\left(k^{\times} \backslash \mathbb{J}\right) \times\left(k^{\times} \backslash \mathbb{J}\right)$, and the integrand is $M_{\mathbb{A}} \cap K_{\mathbb{A}}$-invariant. By [2.A], the group $k^{\times} \backslash \mathbb{J}$ has compact subgroup $k^{\times} \backslash \mathbb{J}^{1}$, on which $|\alpha(m)|$ is trivial, and $k^{\times} \backslash \mathbb{J} \approx \delta(0, \infty) \times k^{\times} \backslash \mathbb{J}^{1}$. For brevity, write $\mathbb{R}^{+}=\delta(0, \infty)$. In effect, $Z^{+}$is the diagonal copy of $\mathbb{R}^{+}$in $\mathbb{J} \times \mathbb{J}$. Thus,

$$
Z^{+} M_{k} \backslash M_{\mathbb{A}} \approx Z^{+} \backslash\left(\left(\mathbb{R}^{+} \times k^{\times} \backslash \mathbb{J}^{1}\right) \times\left(\mathbb{R}^{+} \times k^{\times} \backslash \mathbb{J}^{1}\right)\right)
$$

so, the further quotient by the kernel of $m \rightarrow|\alpha(m)|$ has representatives $a_{y}=\left(\begin{array}{cc}\delta(y) & 0 \\ 0 & 1\end{array}\right)$ for $y>0$. We have $\left|\alpha\left(a_{y}\right)\right|=y^{[k: \mathbb{Q}]}$ for $y>0$. Thus, up to normalization, the integral is

$$
\int_{y \geq t} 1 \frac{d y / y}{y^{[k: \mathbb{Q}]}}<\infty
$$

The quotient $M_{k} \backslash \mathfrak{S}_{C, t}$ without that further quotient by $Z^{+}$will not have finite volume, because $\mathbb{J} / k^{\times}$is non-compact.

Thus,
[2.3.2] Corollary: The quotient $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ has finite volume.

### 2.4 Hecke operators, integral operators

The simplest non-archimedean analogues of the differential operators on $G_{v}$ for archimedean $v$ are integral operators of the form

$$
\varphi \cdot f=\int_{G_{v}} \varphi(g) g \cdot f d g \quad\left(\text { for } \varphi \in C_{c}^{o}\left(G_{v}\right) \text { and } f \in V\right)
$$

for any continuous action $G_{v} \times V \rightarrow V$ on a quasi-complete, locally convex topological vectorspace $V$. The integrand is a continuous, compactly-supported $V$-valued function, so has a Gelfand-Pettis integral [14.1]. Thus, for $f \in V=L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, with $G_{v}$ acting by right translation, pointwise we have

$$
(\varphi \cdot f)(x)=\int_{G_{v}} \varphi(g)(g \cdot f)(x) d g=\int_{G_{v}} \varphi(g) f(x g) d g \quad\left(\text { for } \varphi \in C_{c}^{o}\left(G_{v}\right) \text { and } f \in V\right)
$$

at least almost-everywhere. Better, for general reasons [6.1] the right-translation action
$G_{v} \times L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous, so the integral converges as an $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$-valued integral, and concern about pointwise values is unnecessary. The composition of two such operators is readily described as the operator attached to the convolution: for $\varphi, \psi \in C_{c}^{o}\left(G_{\mathbb{A}}\right)$,

$$
\varphi \cdot(\psi \cdot f)=\int_{G_{\mathbb{A}}} \varphi(g) g \cdot\left(\int_{G_{\mathbb{A}}} \psi(h) h \cdot f d h\right) d g=\int_{G_{\mathbb{A}}} \int_{G_{\mathbb{A}}} \varphi(g) \psi(h)(g h \cdot f) d h d g
$$

because the operation of $\varphi$ moves inside the Gelfand-Pettis integral. Replacing $h$ by $g^{-1} h$ gives

$$
\int_{G_{\mathbb{A}}} \int_{G_{\mathbb{A}}} \varphi(g) \psi\left(g^{-1} h\right) h \cdot f d h d g=\int_{G_{\mathbb{A}}}\left(\int_{G_{\mathbb{A}}} \varphi(g) \psi\left(g^{-1} h\right) d g\right) h \cdot f d h
$$

by changing the order of integration. The inner integral is one expression for the convolution $\varphi * \psi$.
[2.4.1] Lemma: The adjoint to the action of $\varphi \in C_{c}^{o}\left(G_{\mathbb{A}}\right)$ on $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is given by the action of $\check{\varphi} \in C_{c}^{o}\left(G_{\mathbb{A}}\right)$, where $\check{\varphi}(g)=\overline{\varphi\left(g^{-1}\right)}$.
Proof: This is a direct computation: for $f, F \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, by properties of Gelfand-Pettis integrals,

$$
\langle\varphi \cdot f, F\rangle=\left\langle\int_{G_{\mathbb{A}}} \varphi(g) g \cdot f d g, F\right\rangle=\int_{G_{\mathbb{A}}} \varphi(g)\langle g \cdot f, F\rangle d g=\int_{G_{\mathbb{A}}} \varphi(g)\left\langle f, g^{-1} \cdot F\right\rangle d g
$$

because the right translation action of $G_{\mathbb{A}}$ is unitary:

$$
\langle g \cdot f, F\rangle=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f(x g) \overline{F(x)} d x=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f(x) \overline{F\left(x g^{-1}\right)} d x=\left\langle f, g^{-1} \cdot F\right\rangle
$$

by changing variables. This gives

$$
\langle\varphi \cdot f, F\rangle=\left\langle f, \int_{G_{\mathbb{A}}} \overline{\varphi(g)} g^{-1} \cdot F d g\right\rangle=\left\langle f, \int_{G_{\mathbb{A}}} \overline{\varphi\left(g^{-1}\right)} g \cdot F d g\right\rangle=\langle f, \check{\varphi} \cdot F\rangle
$$

by replacing $g$ by $g^{-1}$.
In the four earlier purely archimedean examples, we only considered automorphic forms invariant under right translation by the standard compact subgroups. It is reasonable to consider comparable requirements here, for simplicity possibly requiring right $K_{v}$-invariance for all places $v$. It is also reasonable to relax this condition to requiring right $K_{v}$-invariance almost everywhere, that is, at all but finitely-many places.

A somewhat relaxed version of $K_{\mathbb{A}}$-invariance, to cope with the finitely-many places where right $K_{v^{-}}$ invariance is not required, is $K$-finiteness of a function $f$ on $G_{\mathbb{A}}$ or $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ or other quotients of $G_{\mathbb{A}}$,

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namely, the requirement that the vectorspace of functions spanned by $\left\{x \rightarrow f(x h): h \in K_{\mathbb{A}}\right\}$ is finitedimensional. At the extreme of $K_{\mathbb{A}}$-invariant $f$, this space is one-dimensional. [19]
[2.4.2] Lemma: For $v$ non-archimedean, $K_{v}$-finiteness is equivalent to invariance under some finite-index subgroup $K^{\prime} \subset K_{v}$.
Proof: Let $f$ be in a topological vectorspace $V$ on which $G_{v}$ acts continuously. For $f$ invariant under $K^{\prime}$, the collection of translates of $f$ under $K_{v}$ is finite, given (with possible redundancy) by $g \cdot f$ for representatives $g$ for $K_{v} / K^{\prime}$. On the other hand, when the collection of all right translates of $f$ by $K_{v}$ is a finite-dimensional (complex) vectorspace $F \subset V$, the map $K_{v} \rightarrow \operatorname{Aut}_{\mathbb{C}}(F)$ is a continuous group homomorphism $\rho$ to some $G L_{n}(\mathbb{C})$. Given a neighborhood $U$ of $1 \in G L_{n}(\mathbb{C})$, there is a small-enough neighborhood $U^{\prime}$ of $1 \in K_{v}$ such that $\rho\left(U^{\prime}\right)$ is inside $U$. In fact, we can take $U^{\prime}$ to be a subgroup, for example, $\left\{g=1_{2} \bmod \varpi_{v}^{n}\right\}$ for varying $n$. Then $\rho\left(U^{\prime}\right)$ is a subgroup of $G L_{n}(\mathbb{C})$ inside $U$. Granting for a moment the no small subgroups property of real or complex Lie groups, that a sufficiently small neighborhood of 1 contains no subgroups except $\{1\}$, it must be that $\rho\left(U^{\prime}\right)=\{1\}$. Since $K_{v}$ is compact and $U^{\prime}$ is open, the cover of $K_{v}$ by cosets of $U^{\prime}$ has a finite subcover, so $U^{\prime}$ is of finite index in $K_{v}$. The proof is complete upon proof of the no small subgroups property, following.
[2.4.3] Claim: $G L_{n}(\mathbb{C})$ has the no small subgroup property, that a sufficiently small neighborhood of 1 contains no subgroup larger than $\{1\}$.
Proof: For an $n$-by- $n$ complex matrix $x$, let $\|x\|$ be the operator norm

$$
\|x\|=\sup _{v \in \mathbb{C}^{n}:|v| \leq 1}|x \cdot v| \quad\left(\text { where }\left|\left(v_{1}, \ldots, v_{n}\right)\right|=\sqrt{\left|v_{1}\right|^{2}+\ldots+\left|v_{n}\right|^{2}}\right)
$$

With $r>0$ small enough so that $\sum_{\ell>2} r^{\ell} / \ell!<1+r$, the matrix exponential $x \rightarrow e^{x}$ is a bijection from $E_{r}=\{x:\|x\|<r\}$ to a neighborhood of $1 \in G L_{n}(\mathbb{C})$. We claim that $U=\left\{e^{x}: x \in E_{r / 2}\right\}$ contains no subgroup other than $\{1\}$. Given $0 \neq x \in E_{r / 2}$, there is $1 \leq \ell \in \mathbb{Z}$ such that $\ell \cdot x \in E_{r / 2}$ but $(\ell+1) \cdot x \notin E_{r / 2}$. Then $\ell \cdot x \in E_{r / 2}$, but $(\ell+1) \cdot x \notin E_{r / 2}$. Still, $(\ell+1) \cdot x \in E_{r}$, so by the injectivity of the exponential on $E_{r}, e^{x} \notin U$.

Unsurprisingly, it turns out that $K$-finite functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ are better behaved than arbitrary functions. Of course, most $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ are not $K$-finite.

For non-archimedean $v$, the spherical Hecke operators for $G_{v}$ are the integral operators given by left-andright $K_{v}$-invariant $\varphi \in C_{c}^{o}\left(G_{v}\right)$, also denoted $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$. Since $K_{v}$ is open, such functions are locally constant: given $x \in G_{v}, \varphi(x h)=\varphi(x)$ for all $h \in K_{v}$, but $x K_{v}$ is a neighborhood of $x$. Then the compact support implies that such $\varphi$ takes only finitely-many distinct values. Thus, the associated integral operator is really a finite sum. Nevertheless, expression as integral operators seems to explain the behavior well.
[2.4.4] Claim: The action of spherical Hecke operators attached to $\varphi \in G_{v}$ stabilizes $K_{v}$-invariant vectors $f$ in any continuous group action $G_{v} \times V \rightarrow V$ for quasi-complete, locally convex $V$.
Proof: Granting properties of Gelfand-Pettis integrals, this is a direct computation: for $f \in V$ and $h \in K_{v}$,

$$
h \cdot(\varphi \cdot f)=h \cdot \int_{G_{v}} \varphi(g) g \cdot f d g=\int_{G_{v}} h \cdot(\varphi(g) g \cdot f) d g=\int_{G_{v}} \varphi(g) h g \cdot f d g=\int_{G_{v}} \varphi\left(h^{-1} g\right) g \cdot f d g
$$

by replacing $g$ by $h^{-1} g$. Since $\varphi$ is left $K_{v}$-invariant, this is just $\varphi \cdot f$ again.
[2.4.5] Claim: For $v$ archimedean or non-archimedean, the spherical Hecke algebra $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ with convolution is commutative.
Proof: Gelfand's trick is to find an involutive anti-automorphism $\sigma$ of $G_{v}$, that is, $g \rightarrow g^{\sigma}$ such that $(g h)^{\sigma}=h^{\sigma} g^{\sigma}$ and $\left(g^{\sigma}\right)^{\sigma}=g$, stabilizing double cosets for $K_{v}$, that is, using the Cartan decomposition [2.1], such that $\left(K_{v} m K_{v}\right)^{\sigma}=K_{v} m K_{v}$ for all $m \in M_{v}$. Here, transpose $g^{\sigma}=g^{\top}$ is such an anti-automorphism, since we have a Cartan decomposition $G_{v}=K_{v} M_{v} K_{v}, K_{v}$ is stabilized by transpose, and the diagonal
[19] In the simplest example, Fourier series on the circle $\mathbb{T}$, smoothness is equivalent to rapid decay of Fourier coefficients, while $\mathbb{T}$-finiteness is equivalent to having only finitely-many non-zero Fourier coefficients.
subgroup $M_{v}$ of $G_{v}$ is fixed pointwise by transpose. Then for $\varphi$ in the spherical Hecke algebra, with $g=k_{1} m k_{2}$ in Cartan decomposition,

$$
\varphi\left(g^{\sigma}\right)=\varphi\left(k_{2}^{\sigma} m^{\sigma} k_{1}^{\sigma}\right)=\varphi\left(m^{\sigma}\right)=\varphi(m)=\varphi\left(k_{1} m k_{2}\right)=\varphi(g)
$$

Then the commutativity is a direct computation:

$$
\begin{gathered}
(\varphi * \psi)(x)=(\varphi * \psi)\left(x^{\sigma}\right)=\int_{G_{v}} \varphi(g) \psi\left(g^{-1} x^{\sigma}\right) d g=\int_{G_{v}} \varphi\left(g^{\sigma}\right) \psi\left(\left(g^{-1} x^{\sigma}\right)^{\sigma}\right) d g \\
=\int_{G_{v}} \varphi\left(g^{\sigma}\right) \psi\left(x\left(g^{\sigma}\right)^{-1}\right) d g=\int_{G_{v}} \varphi(g) \psi\left(x g^{-1}\right) d g
\end{gathered}
$$

by replacing $g$ by $g^{\sigma}$. Then replace $g$ by $g^{-1} x$, and then by $g^{-1}$, to obtain

$$
(\varphi * \psi)(x)=\int_{G_{v}} \varphi(g x) \psi\left(g^{-1}\right) d g=\int_{G_{v}} \varphi\left(g^{-1} x\right) \psi(g) d g=(\psi * \varphi)(x)
$$

as claimed.
It is easy to see that the spherical Hecke algebra is stable under adjoints. Thus, it is plausible to ask for simultaneous eigenvectors for the spherical Hecke algebra. That is, for $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, we might additionally try to require that $f$ be a spherical Hecke eigenfunction at almost all non-archimedean $v$, and be an eigenfunction for invariant Laplacians or Casimir at archimedean places. However, in infinite-dimensional Hilbert spaces there is no general promise of existence of such simultaneous eigenvectors.

### 2.5 Decomposition by central characters

We have seen that $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ has finite invariant volume, while $G_{k} \backslash G_{\mathbb{A}}$ does not. The further quotient $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ certainly has finite invariant volume.

Functions on $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ are automorphic forms (or automorphic functions) with trivial central character, since they are invariant under the center $Z_{\mathbb{A}}$ of $G_{\mathbb{A}}$. Such automorphic forms give a reasonable class to consider, but we can treat a larger class with little further effort. Namely, the compact abelian group $Z_{\mathbb{A}} / Z^{+} Z_{k} \approx \mathbb{J}^{1} / k^{\times}$, being a quotient of the center $Z_{\mathbb{A}}$ of $G_{\mathbb{A}}$, acts on functions on $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ in a fashion that commutes with right translation by $G_{\mathbb{A}}$. In particular, the action of $Z_{\mathbb{A}} / Z^{+} Z_{k}$ commutes with the integral operators on $G_{v}$ for $v<\infty$, and with the Casimir or Laplacians on $G_{v}$ at archimedean places.

Thus, for each central character $\omega$ of $Z_{\mathbb{A}} / Z^{+} Z_{k}$, we can consider the space $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ of all left $Z^{+} G_{k}$-invariant $f$ on $G_{\mathbb{A}}$ such that $|f| \in L^{2}\left(Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}\right)$ and $f(z g)=\omega(a) \cdot f(g)$ for all $z \in Z_{\mathbb{A}}$.
[2.5.1] Claim: $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ decomposes by central characters:

$$
L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\text { completion of } \bigoplus_{\omega} L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)
$$

Proof: The argument applies to any compact abelian group $A$ acting on a Hilbert space $V$ by unitary operators, meaning $\langle a \cdot v, a \cdot w\rangle=\langle v, w\rangle$ for all $a \in A$ and $v, w \in V$. For a character $\omega$ of $A$, let $V_{\omega}$ be the $\omega$-eigenspace:

$$
V_{\omega}=\{v \in V: a \cdot v=\omega(a) \cdot v, \text { for all } a \in A\}
$$

For $\omega \neq \omega^{\prime}, V_{\omega}$ and $V_{\omega^{\prime}}$ are orthogonal: with $a \in A$ such that $\omega(a) \neq \omega^{\prime}(a)$ and $v \in V_{\omega}, v^{\prime} \in V_{\omega^{\prime}}$,

$$
\begin{gathered}
\left\langle v, v^{\prime}\right\rangle=\frac{1}{\omega(a)}\left\langle a \cdot v, v^{\prime}\right\rangle=\frac{1}{\omega(a)}\left\langle v, a^{-1} v^{\prime}\right\rangle=\frac{1}{\omega(a)}\left\langle v, \omega^{\prime}\left(a^{-1}\right) v^{\prime}\right\rangle \\
=\frac{\overline{\omega^{\prime}\left(a^{-1}\right)}}{\omega(a)}\left\langle v, \omega^{\prime}\left(a^{-1}\right) v^{\prime}\right\rangle=\frac{\omega^{\prime}(a)}{\omega(a)}\left\langle v, v^{\prime}\right\rangle
\end{gathered}
$$

giving orthogonality.

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Give $A$ an invariant measure with total measure 1. First, $\int_{A} \omega(a)^{-1} a \cdot v d a$ exists as a Gelfand-Pettis $V$-valued integral, so maps $V \rightarrow V$ continuously, and in fact maps to $V_{\omega}$ : using the commutativity of the integral with continuous maps, for $b \in A$,

$$
b \cdot \int_{A} \omega(a)^{-1} a \cdot v d a=\int_{A} \omega(a)^{-1} b a \cdot v d a=\int_{A} \omega\left(b^{-1} a\right)^{-1} a \cdot v d a=\omega(b) \cdot \int_{A} \omega(a)^{-1} a \cdot v d a
$$

Take $v \neq 0$ in $V$. The scalar-valued function $a \rightarrow\langle a \cdot v, v\rangle$ is continuous on $A$, and, since $\langle 1 \cdot v, v\rangle=|v|^{2} \neq 0$, is not identically 0 . By [6.11], $L^{2}(A)$ is the completion of the direct sum of the one-dimensional spaces of functions $\mathbb{C} \cdot \omega$ as $\omega$ ranges over characters. Thus, in $L^{2}(A)$,

$$
0 \neq\langle a v, v\rangle=\sum_{\omega} \int_{A} \omega(b)^{-1}\langle b v, v\rangle d b \cdot \omega(a)=\sum_{\omega}\left\langle\int_{A} \omega(b)^{-1} b v d a, v\right\rangle \cdot \omega(a)
$$

Thus, not all the coefficients on the right-hand side can be 0 , so the projection of non-zero $v \in V$ to some $V_{\omega}$ must be non-zero. Thus, the completion of the sum of the $V_{\omega}$ is all of $V$.

### 2.6 Discrete decomposition of cuspforms

Automorphic forms or automorphic functions are functions of various sorts on $G_{k} \backslash G_{\mathbb{A}}$, with $G_{k}=G L_{2}(k)$, $G_{\mathbb{A}}=G L_{2}(\mathbb{A})$. Here, because $G_{k} \backslash G_{\mathbb{A}}$ has infinite volume, it is reasonable to look at the further quotient $Z^{+} G_{k} \backslash G_{\mathbb{A}}$, for example. Naturally $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is the space of square-integrable automorphic forms. The constant term of an automorphic form $f$ is

$$
c_{P} f(g)=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n g) d n
$$

[2.6.1] Claim: Constant terms are functions on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$.
Proof: By changing variables, we can see that $g \rightarrow c_{P} f(g)$ is a left $N_{\mathbb{A}}$-invariant function on $G_{\mathbb{A}}$ :

$$
c_{P} f\left(n^{\prime} x\right)=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(n \cdot n^{\prime} x\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(\left(n n^{\prime}\right) \cdot x\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n \cdot x) d n \quad\left(\text { for } n^{\prime} \in N_{\mathbb{A}}\right)
$$

Similarly, for $m \in M_{k}$,

$$
c_{P} f(m x)=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n \cdot m x) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(m \cdot m^{-1} n m \cdot x\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(m^{-1} n m \cdot x\right) d n
$$

since $f$ itself is left $M_{k}$-invariant. Then replacing $n$ by $m n m^{-1}$ gives the expression for $c_{P} f(g)$, noting that conjugation by $m \in M_{k}$ stabilizes $N_{k}$, and by the product formula the change of measure on $N_{\mathbb{A}}$ is trivial. Invariance under $Z^{+}$is even easier.

A cuspform is a function $f$ on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ meeting Gelfand's condition ${ }^{[20]} c_{P} f=0$. When $f$ is merely measurable, so does not have well-defined pointwise values everywhere, this condition is best interpreted
[20] In fact, the Gelfand condition for $f$ on $G_{k} \backslash G_{\mathbb{A}}$ to be a cuspform is that $\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} f(n g) d n=0$ as a function of $g \in G_{\mathbb{A}}$ for the unipotent radical $N^{Q}$ of every parabolic $Q$. For $G L_{2}(k)$, proper parabolic subgroups can be characterized as stabilizers of lines in $k^{2}$, and their unipotent radicals as pointwise-fixers of lines. Since $G L_{2}(k)$ is transitive on lines, all proper parabolics (and their unipotent radicals) are conjugate. Thus, vanishing of one constant term (as a function on $G_{\mathbb{A}}$ ) implies vanishing of every constant term, by a change of variables in the integral: for $h \in G L_{2}(k)$,

$$
\int_{h N_{k} h^{-1} \backslash h N_{\mathbb{A}} h^{-1}} f(n g) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(h n h^{-1} g\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(n \cdot h^{-1} g\right) d n=0
$$

using the left $G L_{2}(k)$-invariance. Thus, vanishing of the constant term along $N$ implies vanishing along every conjugate of $N$.
distributionally, as is clarified in the next section, using pseudo-Eisenstein series. The space of squareintegrable cuspforms is

$$
L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\left\{f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right): c_{P} f=0\right\}
$$

The fundamental theorem proven in [7.1-7.7] is the discrete decomposition of spaces of cuspforms. A simple version addresses the space

$$
L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)=\left\{\text { right- } K_{\mathbb{A}} \text {-invariant square-integrable cuspforms with central character } \omega\right\}
$$

where $K_{\mathbb{A}}=\prod_{v \leq \infty} K_{v}$. This space is $\{0\}$ unless $\omega$ is unramified, that is, is trivial on $Z_{\mathbb{A}} \cap K_{\mathbb{A}}$, since $K_{\mathbb{A}}$-invariance implies $Z_{\mathbb{A}} \cap K_{\mathbb{A}}$-invariance, and we also require $Z_{\mathbb{A}}$, $\omega$-equivariance.

Since the spherical Hecke algebras act by right translation, and the Gelfand condition is an integral on the left, spaces of cuspforms are stable under all these integral operators. It is less clear a priori how they behave with respect to the invariant Laplacians [4.2].
[2.6.2] Theorem: $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ has an orthonormal basis of simultaneous eigenfunctions for invariant Laplacians $\Delta_{v}$ at archimedean places, and for spherical Hecke algebras $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ at non-archimedean places. Each simultaneous eigenspace occurs with finite multiplicity, that is, is finitedimensional. (Proof in [7.1-7.7].)

In contrast, the full spaces $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ do not have bases of simultaneous $L^{2}$-eigenfunctions: as in [2.11-2.12], the orthogonal complement of cuspforms in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ mostly consists of integrals of non- $L^{2}$ eigenfunctions for the Laplacians and Hecke operators, the Eisenstein series $E_{s}$ introduced just below in [2.8].

For spaces of automorphic forms more complicated than being right $K_{v}$-invariant for every place $v$, there is generally no decomposition in terms of simultaneous eigenspaces for commuting operators. The decomposition argument in [7.7] directly uses the non-commutative algebras of test functions on the groups $G_{v}$ :

$$
C_{c}^{\infty}\left(G_{v}\right)=\left\{\begin{array}{r}
\text { compactly-supported smooth functions for } v \text { archimedean } \\
\text { compactly-supported locally-constant functions for } v \text { non-archimedean }
\end{array}\right.
$$

Both cases are called smooth. Letting right translation be $R_{g} f(x)=f(x g)$ for $x, g \in G_{\mathbb{A}}$, the action of $\varphi \in C_{c}^{\infty}\left(G_{v}\right)$ on functions $f$ on $G_{k} \backslash G_{\mathbb{A}}$ is

$$
\varphi \cdot f=\int_{G_{v}} \varphi(g) R_{g} f d g
$$

This makes sense not just as a pointwise-value integral, but as a Gelfand-Pettis integral when $f$ lies in any quasi-complete, locally convex topological vectorspace $V$ on which $G_{v}$ acts so that $G_{v} \times V \rightarrow V$ is continuous. Such $V$ is a representation of $G_{v}$. The multiplication in $C_{c}^{\infty}\left(G_{v}\right)$ compatible with such actions is convolution: associativity $\varphi \cdot(\psi \cdot f)=(\varphi * \psi) \cdot f$.

Here, we are mostly interested in actions $G_{v} \times X \rightarrow X$ on Hilbert-spaces $X$. Such a representation is (topologically) irreducible when $X$ has no closed, $G_{v}$-stable subspace. The convolution algebras $C_{c}^{\infty}\left(G_{v}\right)$ are not commutative, so, unlike the commutative case, few irreducible representations are one-dimensional. In fact, typical irreducible representations of $C_{c}^{\infty}\left(G_{v}\right)$ turn out to be infinite-dimensional. Fortunately, there is no mandate to attempt to classify these irreducibles. Indeed, the spectral theory of compact self-adjoint operators still proves [7.7] discrete decomposition with finite multiplicities, for example, as follows.

For every place $v$, let $K_{v}^{\prime}$ be a compact subgroup of $G_{v}$, and for all but a finite set $S$ of places require that $K_{v}^{\prime}=K_{v}$, the standard compact subgroup. For simplicity, we still assume $K_{v}^{\prime}=K_{v}$ at archimedean places. Put $K^{\prime}=\prod_{v} K_{v}^{\prime}$. Let $\omega$ be a central character trivial on $Z_{\mathbb{A}} \cap K^{\prime}$, so that the space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$ of right $K^{\prime}$-invariant cuspforms with central character $\omega$ is not $\{0\}$ for trivial reasons. For $v \in S$, we have a subalgebra $C_{v}^{\infty}\left(K_{v}^{\prime} \backslash G_{v} / K_{v}^{\prime}\right)$ of the convolution algebra of test functions at $v$, stabilizing $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$. [2.6.3] Theorem: $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$ is the completion of the orthogonal direct sum of subspaces, each consisting of simultaneous eigenfunctions for invariant Laplacians $\Delta_{v}$ at archimedean places, of simultaneous eigenfunctions for spherical Hecke algebras $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ at non-archimedean places $v \notin S$, and irreducible $C_{v}^{\infty}\left(K_{v}^{\prime} \backslash G_{v} / K_{v}^{\prime}\right)$-representations at $v \in S$. Each occurs with finite multiplicity. (Proof in [7.1-7.7].)

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The technical features of decomposition with respect to non-commutative rings of operators certainly bear amplification, postponed to [7.2.18] and [7.7]. The notion of multiplicity is made precise in [9.D.14]. In anticipation,
[2.6.4] Theorem: $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ is the completion of the orthogonal direct sum of irreducibles $V$ for the simultaneous action of all algebras $C_{c}^{\infty}\left(G_{v}\right)$. Each irreducible occurs with finite multiplicity. (Proof in [7.7].)
[2.6.5] Corollary: $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ is the completion of the orthogonal direct sum of subspaces, each consisting of simultaneous eigenfunctions for invariant Laplacians $\Delta_{v}$ at archimedean places, of simultaneous eigenfunctions for spherical Hecke algebras $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ at non-archimedean places $v \notin S$, and irreducible $C_{v}^{\infty}\left(K_{v}^{\prime} \backslash G_{v} / K_{v}^{\prime}\right)$-representations at $v \in S$. Each occurs with finite multiplicity.

Again, the various sorts of orthogonal complements to spaces of cuspforms are mostly not direct sums of irreducibles, but are integrals of Eisenstein series, as we see below.

### 2.7 Pseudo-Eisenstein series

Returning to the larger spaces $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ or $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ or $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}\right.$, $\omega$ ), we want to express the orthogonal complement of cuspforms in terms of simultaneous eigenfunctions for invariant Laplacians at archimedean places, and for spherical Hecke algebras at finite places when possible. To consider larger, non-commutative algebras of operators, the more complicated notion of irreducible representation must replace the notion of simultaneous eigenvector. Therefore, we emphasize the commutating operators. As it happens, the pseudo-Eisenstein series here and the genuine Eisenstein series in the next section avoid some of the subtleties that cuspforms may require.

To exhibit explicit $L^{2}$ functions demonstrably spanning the orthogonal complement to cuspforms, we will recast the Gelfand condition that the constant term vanish as a requirement of vanishing as a distribution on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$, and give an equivalent distributional vanishing condition on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.

Vanishing as a distribution is that

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=0 \quad\left(\text { for all } \varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

where $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ consists of compactly-supported functions on that quotient which are smooth in the archimedean coordinates and locally constant in the non-archimedean coordinates. Smoothness of such $\varphi$ can be described more precisely in a fashion that makes clearer the non-interaction of this property with taking a quotient on the left. Namely, smoothness for archimedean places should mean indefinite differentiability on the right with respect to the differential operators coming from the Lie algebra, as in [4.1], and, given the compact support, (uniform) smoothness for non-archimedean places should mean that there exists a compact, open subgroup $K^{\prime}$ of $\prod_{v<\infty} K_{v}$ under which $\varphi$ is right invariant.

As mentioned briefly in the previous section, the nature of $c_{P} f$ for $f$ merely $L^{2}$ is potentially obscure. For example, it is not likely that $c_{P} f \in L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. Instead, for general reasons [6.1], $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is dense in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ in the $L^{2}$ topology, and for general reasons [6.1] the left action of $N_{k} \backslash N_{\mathbb{A}}$ on the Fréchet space $C^{o}\left(Z^{+} N_{k} \backslash G_{\mathbb{A}}\right)$ is a continuous map $\left(N_{k} \backslash N_{\mathbb{A}}\right) \times C^{o}\left(Z^{+} N_{k} \backslash G_{\mathbb{A}}\right) \rightarrow C^{o}\left(Z^{+} N_{\mathbb{A}} \backslash G_{\mathbb{A}}\right)$, so $c_{P} f$ exists as a $C^{o}\left(Z^{+} N_{\mathbb{A}} \backslash G_{\mathbb{A}}\right)$-valued Gelfand-Pettis integral [14.1]. Then one sees directly that $c_{P} f$ is left $M_{k}$-invariant. For such $f$, the integral of $c_{P} f$ against $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ is the integral of a compactly-supported, continuous function. There is no immediate necessity of elaborating a general notion of distribution on $p$-adic groups or adele groups, since cuspforms are ordinary functions, essentially having pointwise values.

For $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, the corresponding pseudo-Eisenstein series is

$$
\Psi_{\varphi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma \cdot g)
$$

Convergence is good:
[2.7.1] Claim: The series for a pseudo-Eisenstein series $\Psi_{\varphi}$ is locally finite, meaning that for $g$ in a fixed compact in $G_{\mathbb{A}}$, there are only finitely-many non-zero summands in $\Psi_{\varphi}(g)=\sum_{\gamma} \varphi(\gamma g)$. Thus, $\Psi_{\varphi} \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$.
Proof: Grant for a moment that there is compact $C \subset G_{\mathbb{A}}$ such that the image of $C$ in the quotient contains the (compact) support of $\varphi$. Fix compact $C_{o} \subset G$ in which $g$ is constrained to lie. A summand $\varphi(\gamma g)$ is nonzero only if $\gamma g \in Z^{+} N_{\mathbb{A}} M_{k} \cdot C$, which holds only if $\gamma \in Z^{+} N_{\mathbb{A}} M_{k} \cdot C \cdot g^{-1}$, so $\gamma \in G_{k} \cap\left(Z^{+} N_{\mathbb{A}} M_{k} \cdot C \cdot C_{o}^{-1}\right)$.

In the quotient $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$, the image of $G_{k}$ is $P_{k} \backslash G_{k}$, is closed and discrete [1.5.3], while the continuous image of the compact set $C \cdot C_{o}^{-1}$ is compact. Thus, left modulo $Z^{+} N_{\mathbb{A}} M_{k}$, that intersection is the intersection of a closed discrete set and a compact set, so finite, as in [1.5.3]. Therefore, the series is locally finite, and defines a smooth function on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. Summing over left translates certainly retains right $K_{\mathbb{A}}$-invariance.

Similarly, $\Psi_{\varphi}$ has compact support in $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ : for a summand $\varphi(\gamma g)$ to be non-zero, it must be that $g \in G_{k} \cdot C$. The image $G_{k} \backslash\left(G_{k} \cdot C\right)$ is compact, being the continuous image of the compact set $C$.

To prove the existence of $C$, let $q: G \rightarrow Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{A}$ be the quotient map. Let $U$ be a neighborhood of $1 \in G_{\mathbb{A}}$ having compact closure $\bar{U}$. For each $g \in G_{\mathbb{A}}, g U$ is a neighborhood of $g$. The images $q(g U)$ are open, by the characterization of the quotient topology. The support $\operatorname{spt}(\varphi)$ is covered by the opens $q(g U)$, and admits a finite subcover $q\left(g_{1} U\right), \ldots, q\left(g_{n} U\right)$. The set $C=g_{1} \bar{U} \cup \ldots \cup g_{n} \bar{U}$ is compact, and its image covers the support of $\varphi$.
[2.7.2] Claim: Square-integrable cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ are the orthogonal complement in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to the subspace spanned by the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. In particular, the pseudo-Eisenstein series $\Psi_{\varphi}$ fit into an adjunction

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot f \quad\left(\text { for } f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

Proof: As noted above, for general reasons [6.1] $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is dense in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, and we consider $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. This allows unwinding as in [5.2]:

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi(g)\left(\int_{N_{k} \backslash N_{\mathbb{A}}} f(n g) d n\right) d g=\int_{Z^{+} N_{k} M_{k} \backslash G_{\mathbb{A}}} \varphi(g) f(g) d g
$$

Winding up, using the left $G_{k}$-invariance of $f$ and $N_{k} M_{k}=P_{k}$,

$$
\int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}} f(g) \varphi(g) d g=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \sum_{\gamma \in P_{k} \backslash G_{k}} f(\gamma \cdot g) \varphi(\gamma \cdot g) d g=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f(g)\left(\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma g)\right) d g
$$

The inner sum in the last integral is the pseudo-Eisenstein series attached to $\varphi$. By Cauchy-SchwarzBunyakowsky,

$$
\left|\int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}} f \varphi\right|=\left|\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \Psi_{\varphi}\right| \leq|f|_{L^{2}} \cdot\left|\Psi_{\varphi}\right|_{L^{2}}
$$

which proves that the functional $f \rightarrow \int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}} f \varphi$ on $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous in the $L^{2}$ topology, so extends by continuity to a continuous linear functional on $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Indeed, this inequality asserts continuity of $f \rightarrow c_{P} f$ as a linear map from $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to distributions on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$ with the weak dual topology as in [13.14].

Similarly, with

$$
C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}, \omega\right)=\left\{\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right): \varphi(z g)=\omega(z) \cdot \varphi(g), \text { for all } z \in Z_{\mathbb{A}}, g \in G\right\}
$$

we have the comparable assertion, now keeping track of complex conjugations:
[2.7.3] Claim: Square-integrable cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ with central character $\omega$ are the orthogonal complement in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ to the subspace spanned by the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in$ $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}, \omega\right)$. The pseudo-Eisenstein series $\Psi_{\varphi}$ fit into an adjunction

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \bar{\varphi} \cdot c_{P} f=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \overline{\Psi_{\varphi}} \cdot f \quad\left(\text { for } f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)\right)
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

(Formation of pseudo-Eisenstein series respects central characters.)
It is useful to understand simpler sub-families of pseudo-Eisenstein series, toward their spectral decomposition in terms of genuine Eisenstein series below in [2.11,2.12,2.13].

With

$$
M^{1}=\left\{\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right): m_{1}, m_{2} \in \mathbb{J},\left|m_{1}\right|=1=\left|m_{2}\right|\right\}
$$

the group $M_{k} \backslash M^{1}$ is compact, because $\mathbb{J}^{1} / k^{\times}$is compact [2.A]. Certainly $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ is inside $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, so such functions $\varphi$ admit decompositions in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ by characters $\chi$ of the compact abelian group $M_{k} \backslash M^{1}$ acting on the left, as in [6.11]. The integral expressing the $\chi^{\text {th }}$ component

$$
\varphi^{\chi}(g)=\int_{M_{k} \backslash M^{1}} \chi(m)^{-1} \varphi(m g) d m
$$

is a Gelfand-Pettis integral converging in $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ for any quasi-complete [14.7] locally convex [13.11] topology on this space. That is, the Fourier components $\varphi^{\chi}$ of a compactly-supported smooth function along $M_{k} \backslash M^{1}$ are again compactly-supported smooth, and their sum converges to the original in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, at least. The support of $\varphi^{\chi}$ is worst $\left(M_{k} \backslash M^{1}\right) \times \operatorname{spt} \varphi$.
[2.7.4] Lemma: A function $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ has constant term $c_{P} f$ integrating to 0 against $\varphi$ in $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ if and only if $c_{P} f$ integrates to 0 against every $M_{k} \backslash M^{1}$-component $\varphi^{\chi}$ of $\varphi$.
Proof: The technicality is that there is no claim that constant terms of functions in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ are in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. Fortunately, this is not an obstacle: as earlier, it suffices to consider $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, so $c_{P} f \in C^{o}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. With $u$ the characteristic function of $\left(M_{k} \backslash M^{1}\right) \times \operatorname{spt} \varphi$, the truncation $u \cdot c_{P} f$ is in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, and truncation does not alter the integrals against $\varphi^{\chi}$ or $\varphi$. Letting $\langle$,$\rangle be the inner$ product in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, since $\varphi=\sum_{\chi} \varphi^{\chi}$ in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$,

$$
\left\langle c_{P} f, \varphi\right\rangle=\left\langle u \cdot c_{P} f, \varphi\right\rangle=\sum_{\chi}\left\langle u \cdot c_{P} f, \varphi^{\chi}\right\rangle=\sum_{\chi}\left\langle c_{P} f, \varphi^{\chi}\right\rangle
$$

giving the assertion.
For central character $\omega$ and character $\chi$ extending $\omega$ to $M_{k} \backslash M^{1}$, define a space of functions on $G_{\mathbb{A}}$ by ${ }^{[21]}$

$$
J_{\chi}=\left\{\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right): \varphi(m g)=\chi(m) \cdot \varphi(g) \text { for all } m \in M^{1}, g \in G_{\mathbb{A}}\right\}
$$

[2.7.5] Remark: In [2.13.5] we will show that pseudo-Eisenstein series made from $J_{\chi}$ and $J_{\chi^{\prime}}$ with distinct characters $\chi^{\prime} \neq \chi$ and $\chi^{\prime} \neq \chi^{w}$ are mutually orthogonal.
[2.7.6] Corollary: Square-integrable cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ with central character $\omega$ are the orthogonal complement in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ to the subspace spanned by the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in J_{\chi}$, as $\chi$ ranges over characters of $M^{1}$ extending $\omega$.
Proof: The lemma shows that it suffices to form pseudo-Eisenstein series from the $M_{k} \backslash M^{1}$-components $\varphi^{\chi}$, and each $\varphi^{\chi}$ is in $J_{\chi}$.
[2.7.7] Claim: For any compact subgroup $K^{\prime} \subset K_{\mathbb{A}}$, right $K^{\prime}$-invariant square-integrable cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}\right)$ are the orthogonal complement in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}\right)$ to the subspace spanned by the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}\left(N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}} / K^{\prime}\right)$.
Proof: The point is that for $f$ right $K_{v}^{\prime}$-invariant, $c_{P} f$ remains $K_{v}^{\prime}$-invariant, so we need only test against test functions $\varphi$ with the same right $K_{v}^{\prime}$ invariance as $f$, at all places $v$, because integration against more general $\varphi$ has the same effect as integrating against right $K_{v}^{\prime}$-invariant ones: giving $K_{v}^{\prime}$ total measure one for convenience,

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi(g)\left(\int_{K_{v}^{\prime}} c_{P} f(g h) d h\right) d g
$$

[^4]\[

$$
\begin{gathered}
=\int_{K_{v}^{\prime}} \int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} f(g h) \varphi(g) d g d h=\int_{K_{v}^{\prime}} \int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} f(g) \varphi\left(g h^{-1}\right) d g d h \\
=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} f(g)\left(\int_{K_{v}^{\prime}} \varphi\left(g h^{-1}\right) d h\right) d g
\end{gathered}
$$
\]

as claimed.
Right $K_{v}$-invariance requires that $\left.\chi\right|_{M_{v}}$ be right $\left(M_{v} \cap K_{v}\right)$-invariant, so $\chi$ is unramified at $v$, as is $\omega$. That is, the set of right $K_{v}$-invariant elements of $J_{\chi}$ is just $\{0\}$ unless $\chi$ is unramified at $v$.
[2.7.8] Claim: Fix a central character $\omega$, and character $\chi$ of $M_{k} \backslash M^{1}$ extending $\omega$. Fix a place $v$. The space of right $K_{v}$-invariant pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in J_{\chi}$ is stable under the invariant Laplacians for archimedean $v$, or under spherical Hecke operators for non-archimedean places $v: \Delta_{v} \Psi_{\varphi}=E_{\Delta_{v} \varphi}$ for archimedean $v$ and $\eta \cdot \Psi_{\varphi}=E_{\eta \cdot \varphi}$ for $\eta \in C_{c}^{\infty}\left(K_{v} \backslash G_{v} / K_{v}\right)$ for non-archimedean $v$.
Proof: Since the Laplacians $\Delta_{v}$ commute with the group action, the effect of $\Delta_{v}$ on a pseudo-Eisenstein series is reflected entirely in its effect on the data: the sum is locally finite, so interchange of the operator and the sum is easy, giving

$$
\Delta_{v} \Psi_{\varphi}=\Delta_{v} \sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi \circ \gamma=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \Delta_{v}(\varphi \circ \gamma)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}\left(\Delta_{v} \varphi\right) \circ \gamma=E_{\Delta_{v} \varphi}
$$

Similarly, the action of the spherical Hecke algebra is on the right, while the winding-up to form a pseudoEisenstein series is on the left:

$$
\eta \cdot \Psi_{\varphi}=\eta \cdot \sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \varphi \circ \gamma=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \eta \cdot(\varphi \circ \gamma)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}(\eta \cdot \varphi) \circ \gamma=E_{\eta \cdot \varphi}
$$

as claimed.
As a simple special situation, consider cuspforms $f$ right invariant under the standard compact subgroup $K_{v}$ for all $v$. Thus, we can invoke the Iwasawa decomposition $G_{v}=P_{v} K_{v}$ everywhere locally, and the constant term $c_{P} f$ is a function on

$$
Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}} \approx Z^{+} N_{\mathbb{A}} M_{k} \backslash N_{\mathbb{A}} M_{\mathbb{A}} K_{\mathbb{A}} / K_{\mathbb{A}} \approx Z^{+} M_{k} \backslash M_{\mathbb{A}} /\left(M_{\mathbb{A}} \cap K_{\mathbb{A}}\right)
$$

The quotient $Z^{+} M_{k} \backslash M_{\mathbb{A}}$ is the quotient of $k^{\times} \backslash \mathbb{J} \times k^{\times} \backslash \mathbb{J}$ by a diagonal copy of the ray $\mathbb{R}^{+}=\delta(0,+\infty)$, as above, thus, with representatives of the form

$$
\left(\begin{array}{cc}
\mathbb{R}^{+} \times \mathbb{J}^{1} / k^{\times} & 0 \\
0 & \mathbb{J}^{1} / k^{\times}
\end{array}\right)
$$

Thus, for fixed central character $\omega$ and character $\chi$ on $M_{k} \backslash M^{1}$ extending $\omega$, a test function $\varphi$ on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}} \approx Z^{+} M_{k} \backslash M_{\mathbb{A}}$ that is in $J_{\chi}$ is entirely specified by a test function $\varphi_{\infty}$ on the ray $\delta(0, \infty)$ :

$$
\varphi\left(a_{y}\right) \cdot m=\varphi_{\infty}(y) \cdot \chi(m) \quad\left(\text { for } a_{y}=\left(\begin{array}{cc}
\delta(y) & 0 \\
0 & 1
\end{array}\right), \text { and } m \in M^{1}, y>0\right)
$$

### 2.8 Eisenstein series

We can attempt to make a pseudo-Eisenstein series $\Psi_{\varphi}$ which is an eigenfunction for an invariant Laplacian $\Delta_{v}$ (or Casimir operator) at archimedean $v$, or for Hecke operators at non-archimedean $v$, by using a right $K_{v}$-invariant $\varphi$ which is such an eigenfunction. However, we already saw in [1.9] that left $N_{v}$-invariant right $K_{v}$-invariant eigenfunctions on $G_{v}$ with trivial central character are

$$
Z_{v} N_{v}\left(\begin{array}{cc}
y & 0 \\
0 & 1
\end{array}\right) K_{v} \longrightarrow y^{s} \quad(\text { for } y>0, \text { for suitable } s \in \mathbb{C})
$$

with eigenvalues $s\left(s-1\right.$ ) (up to normalization). That is, these are characters on $M_{v}$, and are not compactly supported modulo $Z_{v}$. At non-archimedean places, a parallel computation, but now of the effect of spherical Hecke operators, gives a parallel result, illustrating the constraints on eigenfunctions for spherical Hecke algebras:
[2.8.1] Claim: Let $f$ be a function on $N_{v} \backslash G_{v} / K_{v}$, with (unramified) central character $\omega_{v}$, which is an eigenfunction for the spherical Hecke algebra $C_{c}^{\infty}\left(K_{v} \backslash G_{v} / K_{v}\right)$. Then there is a character $\chi_{v}$ on $M_{v}$ extending $\omega_{v}$ such that $f$ is a linear combination of two Hecke eigenfunctions of the special form $f_{1}(n m k)=\chi_{v}(m)$ and $f_{2}(n m k)=\chi_{v}(m)^{-1} \cdot \omega(m)\left|m_{1} / m_{2}\right|_{v}$ for $n \in N_{v}, m \in M_{v}, k \in K_{v}$ and character $\chi_{v}$ on $M_{v}$ extending $\omega_{v}$ on $Z_{v}$.
Proof: By the Iwasawa decomposition, the right $K_{v}$-invariance and left $N_{v}$-invariance of $f$, and central character, determine $f$ completely by its values on elements $g^{\ell}=\left(\begin{array}{cc}\varpi^{\ell} & 0 \\ 0 & 1\end{array}\right)$. We need just a single Hecke operator, the one attached to the characteristic function $\eta$ of the set $C=K_{v}\left(\begin{array}{cc}\varpi & 0 \\ 0 & 1\end{array}\right) K_{v}$. Let $\eta \cdot f=\lambda f$ for $\lambda \in \mathbb{C}$. Then

$$
\lambda \cdot f(g)=(\eta \cdot f)(g)=\int_{G_{v}} \eta(h) f(g h) d h=\int_{C} f(g h) d h
$$

By the $p$-adic Cartan decomposition [2.1], $C$ is exactly the collection of two-by-two matrices with entries in $\mathfrak{o}_{v}$ and determinant in $\varpi \mathfrak{o}_{v}^{\times}$with local parameter $\varpi=\varpi_{v}$. By $p$-adic Iwasawa decomposition [2.1], $C$ is the disjoint union of right $K_{v}$-cosets

$$
C=\left(\begin{array}{cc}
1 & 0 \\
0 & \varpi
\end{array}\right) K_{v} \cup \bigcup_{b \bmod \varpi}\left(\begin{array}{cc}
\varpi & b \\
0 & 1
\end{array}\right) K_{v}
$$

Giving $K_{v}$ measure 1 and letting $q_{v}$ be the residue field cardinality,

$$
\begin{gathered}
\lambda \cdot f\left(g^{\ell}\right)=\int_{C} f\left(g^{\ell} h\right) d h=f\left(\left(\begin{array}{cc}
\varpi^{\ell} & 0 \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
1 & 0 \\
0 & \varpi
\end{array}\right)\right)+\sum_{b} f\left(\left(\begin{array}{cc}
\varpi^{\ell} & 0 \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\varpi & b \\
0 & 1
\end{array}\right)\right) \\
=f\left(\begin{array}{cc}
\varpi^{\ell} & 0 \\
0 & \varpi
\end{array}\right)+\sum_{b} f\left(\begin{array}{cc}
\varpi^{\ell+1} & b \\
0 & 1
\end{array}\right)=\omega(\varpi) f\left(\begin{array}{cc}
\varpi^{\ell-1} & 0 \\
0 & 1
\end{array}\right)+q_{v} \cdot f\left(\begin{array}{cc}
\varpi^{\ell+1} & 0 \\
0 & 1
\end{array}\right) \\
=\omega(\varpi) \cdot f\left(g^{\ell-1}\right)+q_{v} \cdot f\left(g^{\ell+1}\right)
\end{gathered}
$$

This gives the recursion

$$
\binom{f\left(g^{\ell+1}\right)}{f\left(g^{\ell}\right)}=\left(\begin{array}{cc}
\lambda / q_{v} & -\omega(\varpi) / q_{v} \\
1 & 0
\end{array}\right)\binom{f\left(g^{\ell}\right)}{f\left(g^{\ell-1}\right)}
$$

The eigenvalues of that two-by-two matrix are

$$
\{\alpha, \beta\}=\left\{\frac{\lambda \pm \sqrt{\lambda^{2}-4 q_{v}}}{2 q_{v}}\right\}
$$

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with eigenvectors $\binom{\alpha}{1}$ and $\binom{\beta}{1}$. Thus, there are two such eigenfunctions, both with value 1 at 1 :

$$
f_{1}\left(g^{\ell}\right)=\alpha^{\ell} \quad f_{2}\left(g^{\ell}\right)=\beta^{\ell}
$$

extended to $M_{\mathbb{A}}$ to have central character $\omega$. That is, on $M_{\mathbb{A}}$, these two functions are characters. Since $\alpha \cdot \beta=\omega(\varpi) / q$, the two characters are related as asserted.

This last claim shows the impossibility of making Hecke eigenfunction pseudo-Eisenstein series with $\varphi$ in $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. However, it does illustrate a systematic device to make Hecke eigenfunctions:
[2.8.2] Claim: For non-archimedean $v$, any function $f$ on $G_{v}$ of the form $f(n m k)=\chi(m)$ for unramified character $\chi$ on $M_{v}$ is an eigenfunction for the spherical Hecke algebra.
Proof: Let $I_{\chi}$ be the space of smooth functions $f$ on $G_{v}$ with the property $f(n m k)=\chi(m) \cdot f(k)$ for all $n \in N_{v}, m \in M_{v}$, and $k \in K_{v}$. ${ }^{[22]}$ Here the smoothness means that, for each $f$, there is a compact open subgroup $K^{\prime} \subset K_{v}$ such that $f$ is right $K^{\prime}$-invariant. Thus, by $p$-adic Iwasawa decomposition, $I_{\chi}$ is a colimit of finite-dimensional spaces (compare [13.8]). The action of $G_{v}$ on $I_{\chi}$ by right translation $(g \cdot f)(h)=f(h g)$ is continuous, so $\eta$ in the spherical Hecke algebra $C_{c}^{\infty}\left(K_{v} \backslash G / K_{v}\right)$ acts by the integrated version of the action:

$$
(\eta \cdot f)(h)=\int_{G_{v}} \eta(g) f(h g) d g
$$

By changing variables in the integral, the action of such $\eta$ preserves right $K_{v}$-invariance. By $p$-adic Iwasawa decomposition, the subspace of $I_{\chi}$ of right $K_{v}$-invariant functions is one-dimensional, spanned by $f(n m k)=\chi(m)$ itself. Since this one-dimensional space is stabilized by the spherical Hecke algebra, this $f$ is inevitably an eigenfunction for the Hecke algebra.

We wish to decompose pseudo-Eisenstein series $\Psi_{\varphi}$ into $\Delta_{v}$-eigenfunctions and spherical Hecke algebra eigenfunctions to the extent possible. We have already seen that we can take $\varphi$ in the spaces $J_{\chi}$ of [2.7], for $\chi$ a character on $M_{k} \backslash M^{1}$. The previous two claims suggest taking this further: every character on $Z^{+} M_{k} \backslash M_{\mathbb{A}}$ can be written in the form

$$
\nu^{s} \chi: a_{y} \cdot m \longrightarrow|y|^{s} \cdot \chi(m) \quad\left(\text { for } a_{y}=\left(\begin{array}{cc}
\delta(y) & 0 \\
0 & 1
\end{array}\right), \text { and } m \in M^{1}, y>0\right)
$$

for suitable complex $s$ and character $\chi$ on $M_{k} \backslash M^{1}$. Let

$$
I_{s, \chi}=\left\{f \in C^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right): f(n m g)=\left(\nu^{s} \chi\right)(m) \cdot f(g) \text { for all } n \in N_{\mathbb{A}}, m \in M_{\mathbb{A}}\right\}
$$

A genuine Eisenstein series $E_{f}$ for $f \in I_{s, \chi}$ is

$$
E_{f}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} f(\gamma \cdot g) \quad\left(\text { for } f \in I_{s, \chi}\right)
$$

One immediate issue is convergence: unlike pseudo-Eisenstein series $\Psi_{\varphi}$ where $\varphi$ has controlled support, the sum for $E_{f}$ is not locally finite. Ignoring convergence for a moment, $E_{f}$ is genuine in the sense that it is a spherical Hecke algebra eigenfunction at all but (at worst) finitely-many non-archimedean places, since smoothness at finite places requires that $f$ is right $K^{\prime}$-invariant for some compact open subgroup $K^{\prime}=\prod_{v<\infty} K_{v}^{\prime}$ of $\prod_{v<\infty} K_{v}$, and the product topology requires that $K_{v}^{\prime}=K_{v}$ for all but finitely-many $v<\infty$.
[22] This space $I_{\chi}$ is an example of an unramified principal series representation of $G_{v}$, meaning that it is induced from $\tilde{\chi}(n m)=\chi(m)$ on $P_{v}$, with $\chi$ unramified. The previous two claims touch on the importance of principal series representations in the application of representation theory of $p$-adic groups to automorphic forms. A strong form of the generalization of these claims to a wide class of $p$-adic groups is in [Borel 1976], [Matsumoto 1977], [Casselman 1980].

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

The extreme, simplest case is that $f \in I_{s, \chi}$ is right $K_{v}$-invariant at all places $v$, that is, is spherical everywhere locally. From all the Iwasawa decompositions for groups $G_{v}$, up to a scalar there is a unique such $f$, namely, $f(n m k)=f(m)=\left(\nu^{s} \chi\right)(m)$. The everywhere spherical Eisenstein series attached to an unramified grossencharakter $\chi$ is

$$
E_{s, \chi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} f(\gamma \cdot g) \quad\left(\text { for } f(n m k)=f(m)=\left(\nu^{s} \chi\right)(m)\right)
$$

[2.8.3] Claim: Assuming the series expression for the everywhere-spherical Eisenstein series $E_{s, \chi}$ is convergent, it is an eigenfunction for the invariant Laplacians at archimedean places, and for the spherical Hecke algebras at non-archimedean places.
Proof: Assuming convergence, the invariance of Laplacians and spherical Hecke operators under left translation implies that we need merely check that the function $f(n m k)=f(m)=\widetilde{\chi}(m)$ itself is an eigenfunction. In [1.9] we saw the archimedean assertion, and the claim above proves the non-archimedean assertion.
[2.8.4] Claim: Assuming the series expression for the everywhere-spherical Eisenstein series $E_{s, \chi}$ is convergent, its constant term is

$$
c_{P} E_{s, \chi}(z n m k)=\left(\nu^{s} \chi\right)(m)+c_{s, \chi} \cdot\left(\nu^{1-s} \chi^{w}\right)(m) \quad\left(\text { for } z \in Z^{+}, n \in N_{\mathbb{A}}, m \in M_{\mathbb{A}}, k \in K_{\mathbb{A}}\right)
$$

where $\chi^{w}(m)=\chi\left(w m w^{-1}\right)$ with long Weyl element $w=\left(\begin{array}{ll}0 & 1 \\ 1 & 0\end{array}\right)$, and

$$
c_{s, \chi}=\frac{\Lambda\left(2 s-1, \chi_{1} / \chi_{2}\right)}{\Lambda\left(2 s, \chi_{1} / \chi_{2}\right)} \quad\left(\text { with } \chi\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)=\chi_{1}\left(m_{1}\right) \chi_{2}\left(m_{2}\right)\right)
$$

where $\Lambda\left(s, \chi_{1} / \chi_{2}\right)$ is the Hecke grossencharacter $L$-function completed by multiplying by the appropriate Gamma factors.
Proof: Via the Bruhat decomposition $G_{k}=P_{k} \sqcup P_{k} w N_{k}$,

$$
P_{k} \backslash G_{k}=P_{k} \backslash P_{k} \sqcup P_{k} \backslash P_{k} w N_{k} \approx\{1\} \sqcup w N_{k}
$$

The small Bruhat cell $P_{k}$ produces the first summand in the constant term:

$$
\int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\gamma \in P_{k} \backslash P_{k}} f(\gamma n g) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n g) d n=f(g) \cdot \int_{N_{k} \backslash N_{\mathbb{A}}} 1 d n
$$

The large Bruhat cell $P_{k} w N_{k}$ gives

$$
\int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\gamma \in w N_{k}} f(\gamma n g) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\gamma \in N_{k}} f(w \gamma n g) d n=\int_{N_{\mathbb{A}}} f(w n g) d n
$$

by unwinding, as in [5.2]. Since $c_{P} f$ will be left $N_{\mathbb{A}}$-invariant and right $K_{\mathbb{A}}$-invariant, it suffices to evaluate this integral on $g=m \in M_{\mathbb{A}}$. Then

$$
\int_{N_{\mathbb{A}}} f(w n m) d n=\int_{N_{\mathbb{A}}} f\left(w m \cdot m^{-1} n m\right) d n=\int_{N_{\mathbb{A}}} f(w m \cdot n) \nu(m) d n
$$

by replacing $n$ by $m n m^{-1}$, with $\nu(m)$ resulting from change-of-measure. This is

$$
\begin{gathered}
\int_{N_{\mathbb{A}}} f\left(w m w^{-1} \cdot w n\right) \\
=\nu(m) d n=\nu(m) \nu^{s}\left(w m w^{-1}\right) \chi\left(w m w^{-1}\right) \cdot \int_{N_{\mathbb{A}}} f(w n) d n \\
=\nu(m)^{1-s} \chi\left(w m w^{-1}\right) \cdot \int_{N_{\mathbb{A}}} f(w n) d n
\end{gathered}
$$

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Since the right $K_{\mathbb{A}}$-invariance is preserved by integrating on the left, this is (unique up to constant) the spherical function in $I_{1-s, \chi^{w}}$. That normalization constant is very significant, being a ratio of $L$-functions, as follows.

Let $f_{v}^{o}(n m k)=\left(\nu_{v}^{s} \chi_{v}\right)(m)$ be the normalized spherical vector on $G_{v}$, where $\nu_{v}^{s} \chi_{v}$ is the $v^{t h}$ local factor of $\nu^{s} \chi$,

$$
\left(\begin{array}{lll}
m_{1} & 0 & m_{2}
\end{array}\right) \longrightarrow\left|m_{1} / m_{2}\right|^{s} \cdot \chi_{1}\left(m_{1}\right) \chi_{2}\left(m_{2}\right) \quad\left(\text { on } M_{v}\right)
$$

The integral giving the normalizing constant factors over primes:

$$
\int_{N_{\mathbb{A}}} f(w n) d n=\prod_{v \leq \infty} \int_{N_{v}} f_{v}^{o}(w n) d n
$$

To evaluate the $v^{t h}$ factor, we must determine the local Iwasawa decomposition of $w n$ for $n \in N_{v}$. At $k_{v} \approx \mathbb{R}$, as in [1.3]

$$
w\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)=\left(\begin{array}{ll}
0 & 1 \\
1 & x
\end{array}\right)=\left(\begin{array}{cc}
\frac{-1}{\sqrt{1+x^{2}}} & * \\
0 & \sqrt{1+x^{2}}
\end{array}\right) \cdot\left(\begin{array}{cc}
\frac{x}{\sqrt{1+x^{2}}} & \frac{-1}{\sqrt{1+x^{2}}} \\
\frac{1}{\sqrt{1+x^{2}}} & \frac{x}{\sqrt{1+x^{2}}}
\end{array}\right)
$$

with that last matrix in $S O_{2}(\mathbb{R})$. Unramified unitary characters on $k_{v}^{\times} \approx \mathbb{R}^{\times}$are of the form $\alpha_{v}(y)=|y|^{i t_{v}}$ for some purely imaginary $i t_{v}$. With $|y|^{i t_{v}}=\chi_{1}(y) / \chi_{2}(y)$, the corresponding local integral is evaluated via the standard trick $\int_{0}^{\infty} e^{-t y} t^{s} d t / t=y^{-s} \Gamma(s)$ : first, with $i t_{v}=0$,

$$
\begin{gathered}
\int_{\mathbb{R}} f_{v}^{o}(w n) d n=\int_{\mathbb{R}} \frac{1}{\left(1+x^{2}\right)^{s}} d x=\frac{1}{\pi^{-s} \Gamma(s)} \int_{0}^{\infty} \int_{\mathbb{R}} e^{-\pi t\left(1+x^{2}\right)} t^{s} \frac{d t}{t} \\
=\frac{1}{\pi^{-s} \Gamma(s)} \int_{0}^{\infty} \int_{\mathbb{R}} e^{-\pi t-\pi x^{2}} t^{s-\frac{1}{2}} \frac{d t}{t}=\frac{\pi^{-\left(s-\frac{1}{2}\right)} \Gamma\left(s-\frac{1}{2}\right)}{\pi^{-s} \Gamma(s)}
\end{gathered}
$$

Replacing $s$ by $s-i t_{v}$, the general unramified case is

$$
\int_{\mathbb{R}} f_{v}^{o}(w n) d n=\int_{\mathbb{R}} \frac{1}{\left(1+x^{2}\right)^{s-i t_{v}}} d x=\frac{\pi^{-\left(s-i t_{v}-\frac{1}{2}\right)} \Gamma\left(s-i t_{v}-\frac{1}{2}\right)}{\pi^{-\left(s-i t_{v}\right)} \Gamma\left(s-i t_{v}\right)}
$$

Similarly, at $k_{v} \approx \mathbb{C}$, with trivial $\chi_{1}, \chi_{2}$,

$$
w\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)=\left(\begin{array}{ll}
0 & 1 \\
1 & x
\end{array}\right)=\left(\begin{array}{cc}
\frac{-1}{\sqrt{1+|x|^{2}}} & * \\
0 & \sqrt{1+|x|^{2}}
\end{array}\right) \cdot\left(\begin{array}{cc}
\frac{x}{\sqrt{1+|x|^{2}}} & \frac{-1}{\sqrt{1+|x|^{2}}} \\
\frac{1}{\sqrt{1+|x|^{2}}} & \frac{\bar{x}}{\sqrt{1+|x|^{2}}}
\end{array}\right)
$$

With the normalization of local norms $|t|_{\mathbb{C}}=\left|N_{\mathbb{R}}^{\mathbb{C}} t\right|$ for the product formula, up to measure constants the local integral is

$$
\begin{gathered}
\int_{\mathbb{C}} f_{v}^{o}(w n) d n=\int_{\mathbb{C}} \frac{1}{\left(1+|x|^{2}\right)^{2 s}} d x=\frac{1}{\pi^{-2 s} \Gamma(2 s)} \int_{0}^{\infty} \int_{\mathbb{R}} e^{-\pi t\left(1+|x|^{2}\right)} t^{2 s} \frac{d t}{t} \\
=\frac{1}{\pi^{-2 s} \Gamma(2 s)} \int_{0}^{\infty} \int_{\mathbb{R}} e^{\left.-\pi t-\pi|x|^{2}\right)} t^{2 s-1} \frac{d t}{t}=\frac{\pi^{-(2 s-1)} \Gamma(2 s-1)}{\pi^{-2 s} \Gamma(2 s)}
\end{gathered}
$$

In the general unramified case, with $\chi_{1}(t) / \chi_{2}(t)=|t|_{\mathbb{C}}^{i t_{v}}=|t|^{2 i t_{v}}$, again there is a shift $s \rightarrow s-i t_{v}$.
At non-archimedean places, the Iwasawa decomposition has a different nature:

$$
w\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)=\left(\begin{array}{ll}
0 & 1 \\
1 & x
\end{array}\right)= \begin{cases}\left(\begin{array}{cc}
\frac{-1}{x} & 1 \\
0 & x
\end{array}\right)\left(\begin{array}{ll}
1 & 0 \\
\frac{1}{x} & 1
\end{array}\right) & \text { for }|x|_{v} \geq 1 \\
\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)\left(\begin{array}{ll}
0 & 1 \\
1 & x
\end{array}\right) & \text { for }|x|_{v} \leq 1\end{cases}
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

There is a non-archimedean analogue of the Gamma trick: with ch the characteristic function of $\mathfrak{o}_{v}$, with multiplicative Haar measure giving $\mathfrak{o}_{v}^{\times}$total measure 1,

## [2.8.5] Lemma:

$$
\int_{k_{v}^{\times}} \operatorname{ch}(t y) \operatorname{ch}(t)|t|_{v}^{s} d t= \begin{cases}\frac{1}{1-q_{v}^{-s}} & \left(\text { for } y \in \mathfrak{o}_{v}\right) \\ \frac{|y|^{-s}}{1-q_{v}^{-s}} & \left(\text { for } y \notin \mathfrak{o}_{v}\right)\end{cases}
$$

Proof: For $y \in \mathfrak{o}_{v}$, the integral becomes an Iwasawa-Tate local zeta integral

$$
\int_{k_{v}^{\times}} \operatorname{ch}(t)|t|^{s} d t=\sum_{\ell \geq 0} \int_{\varpi^{\ell} \mathfrak{o}_{v}^{\times}} q_{v}^{-\ell s}=\sum_{\ell \geq 0} 1 \cdot q_{v}^{-\ell s}=\frac{1}{1-q_{v}^{-s}}
$$

For $y \notin \mathfrak{o}_{v}$, replace $t$ by $t / y$ in the integral, producing the $|y|_{v}^{-s}$ factor and then the integral just evaluated.

Returning to the evaluation of the non-archimedean local factor in the constant term, let

$$
\gamma(s, y)=\int_{k_{v}^{\times}} \operatorname{ch}(t y) \operatorname{ch}(t)|t|_{v}^{s} d^{\times} t
$$

emphasizing that it is multiplicative Haar measure. With trivial $\chi_{1}, \chi_{2}$, the lemma gives

$$
\begin{gathered}
\int_{N_{v}} f_{v}^{o}(w n) d n=\left(1-q_{v}^{-2 s}\right) \int_{k_{v}} \gamma(2 s, x) d x=\left(1-q_{v}^{-2 s}\right) \int_{k_{v}^{\times}} \int_{k_{v}} \operatorname{ch}(t x) \operatorname{ch}(t)|t|_{v}^{2 s} d x d^{\times} t \\
=\left(1-q_{v}^{-2 s}\right) \int_{k_{v}^{\times}} \int_{k_{v}} \operatorname{ch}(x) \operatorname{ch}(t)|t|_{v}^{2 s-1} d x d^{\times} t
\end{gathered}
$$

by replacing $x$ by $x / t$. The integral in $x$ is just meas $\left(\mathfrak{o}_{v}\right)$, which at $k_{v}$ unramified over the corresponding $\mathbb{Q}_{p}$ is reasonably taken to be 1 . Thus,

$$
\int_{N_{v}} f_{v}^{o}(w n) d n=\left(1-q_{v}^{-2 s}\right) \int_{k_{v}^{\times}} \operatorname{ch}(t)|t|_{v}^{2 s-1} d^{\times} t=\frac{1-q_{v}^{-2 s}}{1-q_{v}^{-(2 s-1)}}=\frac{1 /\left(1-q_{v}^{-(2 s-1)}\right)}{1 /\left(1-q_{v}^{-2 s}\right)}
$$

The adjustment for a general unramified character again shifts $s$ to $s-i t_{v}$. The products over all places give the indicated ratios of completed $L$-functions, apart from ratios of powers of the conductor, which correspond to the additive measure normalization at ramified places. A ratio of a value and a shift only leaves a constant, not immediately important here.
[2.8.6] Claim: For $\operatorname{Re}(s)>1$, the series expression for $E_{f}$ with (continuous) $f \in I_{s, \chi}$ converges absolutely and uniformly on compacts, to a continuous function on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.
Proof: The function $f$ is dominated by the spherical vector, since $|f(z n m k)|=\left|\left(\nu^{s} \chi\right)(m)\right| \cdot|f(k)|$ and the continuous function $f$ is bounded on the compact $K_{\mathbb{A}}$. Also, $\chi$ has absolute value 1 , so we may as well take $\chi$ trivial. And it suffices to treat $s=\sigma \in \mathbb{R}$. Use the height functions $h_{v}$ on $k_{v}^{2}$ and $h$ on $\mathbb{A}^{2}$, and $\eta(g)=|\operatorname{det} g| / h\left(v_{o} g\right)^{2}$. In particular, $\eta(z n m k)=\left|m_{1} / m_{2}\right|$ for $m=\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right)$. Also, $\nu(m)^{-1} d n d m$ is left Haar measure $d p$ on $P_{\mathbb{A}}$. Thus, it suffices to prove convergence of

$$
\sum_{\gamma \in P_{k} \backslash G_{k}} \eta(\gamma g)^{\sigma}=\sum_{\gamma \in P_{k} \backslash G_{k}}|\operatorname{det} \gamma g|^{\sigma} \cdot h\left(v_{o} \gamma g\right)^{-2 \sigma}=|\operatorname{det} g|^{\sigma} \sum_{\gamma \in P_{k} \backslash G_{k}} h\left(v_{o} \gamma g\right)^{-2 \sigma}
$$

By reduction theory [2.2], for compact $C \subset G_{\mathbb{A}}$, there are constants $0<c \leq c^{\prime}<+\infty$ such that

$$
c \cdot h(v) \leq h(v g) \leq c^{\prime} \cdot h(v) \quad\left(\text { for all } g \in C, \text { for all primitive } v \in \mathbb{A}^{2}\right)
$$

so

$$
h\left(v_{o} \gamma g\right) \leq c^{\prime} \cdot h\left(v_{o} \gamma\right) \leq \frac{c^{\prime}}{c} \cdot c \cdot h\left(v_{o} \gamma g^{\prime}\right) \quad\left(\text { for all } g, g^{\prime} \in C, \text { for all } \gamma \in G_{k}\right)
$$

Thus, convergence of the series is equivalent to convergence of an averaged integral $\int_{C} E_{\sigma}$. By discreteness of $G_{k}$ in $G_{\mathbb{A}}$, we can shrink $C$ so that, for $\gamma$ in $G_{k}$, if $\gamma C \cap C \neq \phi$ then $\gamma=1$. Then

$$
\int_{C} E_{\sigma}=\int_{C} \sum_{\gamma \in P_{k} \backslash G_{k}}|\operatorname{det} \gamma g|^{\sigma} h\left(v_{o} \gamma g\right)^{-2 \sigma} d g=|\operatorname{det} g|^{\sigma} \int_{P_{k} \backslash G_{k} \cdot C} h\left(v_{o} g\right)^{-2 \sigma} d g
$$

Let $\mu$ be the infimum of $h(v)$ over non-zero primitive $v$ in $\mathbb{A}^{2}$. From reduction theory [2.2] this infimum is attained, so $\mu>0$, and $c \cdot \mu \leq h\left(v_{o} \gamma g\right)$ for all $g \in C$ and $\gamma \in G_{k}$, and $G_{k} \cdot C$ is contained in a set

$$
Y=\left\{g \in G_{\mathbb{A}}: h\left(v_{o} g\right) \geq c \cdot \mu \text { and } c_{1} \leq|\operatorname{det} g| \leq c_{2}\right\} \quad\left(\text { with } 0<c_{1} \text { and } c_{2}<+\infty\right)
$$

The set $Y$ is right $K_{\mathbb{A}}$-stable, since $h$ is $K_{\mathbb{A}}$-invariant. Using Iwasawa decompositions, with left Haar measure $d p$ on $P_{\mathbb{A}}$,

$$
\int_{P_{k} \backslash G_{k} \cdot C}|\operatorname{det} g|^{\sigma} h\left(v_{o} g\right)^{-2 \sigma} d g \leq \int_{P_{k} \backslash Y}|\operatorname{det} g|^{\sigma} h\left(v_{o} g\right)^{-2 \sigma} d g=\int_{P_{k} \backslash\left(P_{\mathbb{A}} \cap Y\right)}|\operatorname{det} p|^{\sigma} h\left(v_{o} p\right)^{-2 \sigma} d p
$$

The set $Y$ is left $N_{\mathbb{A}}$-stable, and the induced measure on the compact quotient $N_{k} \backslash N_{\mathbb{A}}$ is finite, so up to a constant the integral is

$$
\int_{M_{k} \backslash\left(M_{\mathbb{A}} \cap Y\right)}|\operatorname{det} m|^{\sigma} h\left(v_{o} m\right)^{-2 \sigma} \nu(m)^{-1} d m=\int_{M_{k} \backslash\left(M_{\mathbb{A}} \cap Y\right)} \nu(m)^{\sigma-1} d m
$$

From

$$
M_{\mathbb{A}} \cap Y \subset\left\{\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right):\left|m_{1} / m_{2}\right| \geq c \mu, c_{1} \leq\left|m_{1} m_{2}\right| \leq c_{2}\right\}
$$

and compactness of $\mathbb{J}^{1} / k^{\times}$,

$$
M_{k} \backslash\left(M_{\mathbb{A}} \cap Y\right)=\text { compact } \times\left\{\left(\begin{array}{cc}
\delta\left(y_{1}\right) & 0 \\
0 & \delta\left(y_{2}\right)
\end{array}\right)\right\} \quad\left(\text { with } y_{2} \geq c \mu \text { and } c_{1} \leq y_{1} y_{2} \leq c_{2}\right)
$$

Thus,

$$
\int_{M_{k} \backslash\left(M_{\mathbb{A}} \cap Y\right)} \nu(m)^{\sigma-1} d m=\int_{M_{k} \backslash\left(M^{1} \cap Y\right)} d m \cdot \int_{c \mu}^{\infty} \int_{c_{1} / y_{2}}^{c_{2} / y_{2}}\left(\frac{y_{1}}{y_{2}}\right)^{\sigma-1} \frac{d y_{1}}{y_{1}} \frac{d y_{2}}{y_{2}}
$$

Replacing $y_{1}$ by $y_{1} / y_{2}$, the latter elementary integral becomes

$$
\int_{c \mu}^{\infty} \int_{c_{1}}^{c_{2}}\left(\frac{y_{1}}{y_{2}^{2}}\right)^{\sigma-1} \frac{d y_{1}}{y_{1}} \frac{d y_{2}}{y_{2}}=(\text { constant }) \cdot \int_{c \mu}^{\infty}\left(y_{2}^{2}\right)^{1-\sigma} \frac{d y_{2}}{y_{2}}
$$

which converges for $\sigma>1$. This also proves the uniform convergence on compacts.
We also want moderate growth on Siegel sets: for $n \in N_{\mathbb{A}}, k \in K_{\mathbb{A}}, z \in Z^{+}$, and $m=a_{y} \cdot m^{\prime}$ with $m^{\prime} \in M^{1}$,

$$
\left|E_{f}(z n m k)\right|<_{t, C} y^{\operatorname{Re}(s)} \quad\left(\text { on } \mathfrak{S}_{t, C}, \text { implied constant depending on } t, C\right)
$$

And we want convergence to a smooth function:
[2.8.7] Claim: The series for $E_{s}$ converges in the $C^{\infty}$ topology for $\operatorname{Re}(s)>1$, and produces a $C^{\infty}$ moderategrowth function on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. (Proofs in [2.B], [11.5].)

As in [13.5], the idea of the archimedean aspect of the $C^{\infty}$ topology is that it is given by the collection of seminorms given by sups on compacts of all derivatives, for example left- $G$-invariant derivatives on

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$G_{\infty}=\prod_{v \mid \text { infty }} G_{v}$ from the Lie algebra, preserving left $G_{k}$-invariance, and stabilizing a useful class of Eisenstein series. The non-archimedean smoothness is simpler, being right invariance under some compact open subgroup of $K^{\prime} \subset \prod_{v<\infty} K_{v}$, which leads to taking an ascending union (colimit) over such $K^{\prime}$.

The everywhere spherical computation of constant terms applies to computation of local components of $c_{P} f$ at good primes for general $f \in I_{s, \chi}$, that is, places $v$ where $f$ is right $K_{v}$-invariant and $\chi$ is unramified. However, at the other, bad, primes for $f \in I_{s, \chi}$, where $f$ is right $K_{v}^{\prime}$-invariant only for $K_{v}^{\prime} \subset K_{v}$ of high index, the local integrals

$$
f \longrightarrow\left(g \longrightarrow \int_{N_{v}} f(w n g) d n\right)
$$

are naturally more complicated. Still, these maps visibly commute with the right translation action of $G_{v}$, and have predictable left-equivariance under $M_{v}$ for the same reason as in the simpler computation:
[2.8.8] Claim: The constant term of the Eisenstein series $E_{f}$ for $f \in I_{s, \chi}$ is

$$
c_{P} E_{f}=f+C_{s, \chi}(f) \quad\left(\text { with } C_{s, \chi} f(g)=\int_{N_{\mathbb{A}}} f(w n g) d n\right)
$$

The map $C_{s, \chi}$ is a $G_{\mathbb{A}}$-map in the sense that $g \cdot\left(C_{s, \chi} f\right)=C_{s, \chi}(g \cdot f)$ where $g \cdot f$ is right translation.
Proof: Integration on the left certainly commutes with right translation. As in the earlier, simpler, case, the small Bruhat cell gives

$$
\int_{N_{k} \backslash N_{\mathbb{A}}} f(n g) d n=f(g) \cdot \int_{N_{k} \backslash N_{\mathbb{A}}} 1 d n
$$

and the volume of $N_{k} \backslash N_{\mathbb{A}}$ is reasonably normalized to 1 . The big Bruhat cell integral unwinds:

$$
\int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\gamma \in N_{k}} f(w \gamma n g) d n=\int_{N_{\mathbb{A}}} f(w n g) d n
$$

For $m \in M_{\mathbb{A}}$,

$$
\begin{aligned}
\int_{N_{\mathbb{A}}} f(w n m g) d n= & \int_{N_{\mathbb{A}}} f\left(w m \cdot m^{-1} n m g\right) d n=\int_{N_{\mathbb{A}}} \nu(m) f\left(w m w^{-1} \cdot w n g\right) d n \\
& =\nu(m)^{1-s} \chi\left(w m w^{-1}\right) \cdot \int_{N_{\mathbb{A}}} f(w n g) d n
\end{aligned}
$$

showing that this part of the constant term is in $I_{1-s, \chi^{w}}$.
The scattering matrix/operator is the map ${ }^{[23]}$

$$
f \longrightarrow\left(g \longrightarrow \int_{N_{\mathbb{A}}} f(w n g) d n\right) \quad\left(\text { from } I_{s, \chi} \text { to } I_{1-s, \chi^{w}}\right)
$$

Since the unwound integral over $N_{\mathbb{A}}$ is a (limit of) product(s) of integrals over $N_{v}$, it is a (tensor) product of local maps $C_{s, \chi, v}$ among corresponding local spaces

$$
I_{s, \chi, v}=\left\{f \in C^{\infty}\left(G_{v}\right): f(n m g)=\left(\nu^{s} \chi\right)(m) \cdot f(g), \text { for all } n \in N_{v}, m \in M_{v}, g \in G_{v}\right\}
$$

in the natural way. For example, for monomial $f \in I_{s, \chi}$, that is, of the form $f(g)=\prod_{v \leq \infty} f_{v}\left(g_{v}\right)$, with $f_{v}$ the spherical vector for $v$ outside a finite set $S$ of places,

$$
\int_{N_{\mathbb{A}}} f(w n g) d n=\int_{N_{\mathbb{A}}} \prod_{v} f_{v}(w n g) d n=\prod_{v} \int_{N_{v}} f_{v}(w n g) d n
$$

[23] Since this map respects the right translation action of $G_{v}$ and/or of $G_{\mathbb{A}}$ on functions, it is an instance of an intertwining operator among representations of $G_{v}$ and/or $G_{\mathbb{A}}$.

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The earlier, simple constant term computation shows that for places $v \notin S$, the local operator sends the spherical vector in $I_{s, \chi, v}$ to the spherical vector in $I_{1-s, \chi^{w}, v}$, multiplied by the $v^{t h}$ Euler factor of $\Lambda\left(2 s-1, \chi^{w}\right) / \Lambda(2 s, \chi)$.
[2.8.9] Remark: The space of functions $I_{s, \chi}$ has central character

$$
\left(\begin{array}{cc}
z & 0 \\
0 & z
\end{array}\right) \longrightarrow \nu^{s}\left(\begin{array}{cc}
z & 0 \\
0 & z
\end{array}\right) \cdot \chi_{1}(z) \chi_{2}(z)=\chi_{1}(z) \chi_{2}(z)
$$

The condition $\chi^{w}=\chi$ is that $\chi_{1}=\chi_{2}$. Thus, for $f \in I_{s, \chi}$, the normalized function $g \rightarrow \chi_{1}(\operatorname{det} g)^{-1} f(g)$ has trivial central character, and, further, is in $I_{s, 1}$.

### 2.9 Meromorphic continuation of Eisenstein series

This is an issue of showing that a family of Eisenstein series $E_{f_{s}}$ with $f_{s} \in I_{s, \chi}$ has a meromorphic continuation in $s$ beyond the range of convergence $\operatorname{Re}(s)>1$. So certainly $E_{f_{s}}$ must be holomorphic in $\operatorname{Re}(s)>1$, and the dependence of $f_{s}$ on $s$ must be constrained for this to be plausible. The simplest example is to take $f_{s}$ to be the everywhere-spherical vector in $I_{s, \chi}$, with $\chi$ unramified everywhere. As special argument applicable to $G L_{2}$ is in [2.B], and instantiation of a more general approach is in [11.5]. The basic theorem is
[2.9.1] Theorem: The everywhere-spherical Eisenstein series $E_{s, \chi}$ has a meromorphic continuation in $s \in \mathbb{C}$, as a smooth function of moderate growth on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. As a function of $s, E_{s, \chi}(g)$ is of at most polynomial growth vertically, uniformly in bounded strips, uniformly for $g$ in compacts. (Proofs in [2.B] and [11.5].)
[2.9.2] Corollary: At archimedean $v$, let $t_{v} \in \mathbb{R}$ be associated to the character $\chi$, as in the proof of [2.8.4], so that $E_{s, \chi}$ is an eigenfunction for the $v^{t h}$ invariant Laplacian $\Delta_{v}$, with eigenvalue $\lambda_{s, \chi}=\left(s-i t_{v}\right)\left(s-i t_{v}-1\right)$. This eigenfunction property persists under meromorphic continuation.
Proof: Both $\Delta_{v} E_{s, \chi}$ and $\lambda_{s, \chi} \cdot E_{s, \chi}$ are holomorphic function-valued functions of $s$, taking values in the topological vector space of smooth moderate-growth functions. They agree in the region of convergence $\operatorname{Re}(s)>1$, then apply the vector-valued form [15.2] of the Identity Principle from complex analysis. ///
[2.9.3] Corollary: The meromorphic continuation of $E_{s, \chi}$ implies the meromorphic continuation of the constant term $c_{P} E_{s, \chi}(m)=\left(\nu^{s} \chi\right)(m)+c_{s, \chi} \cdot\left(\nu^{1-s} \chi^{w}\right)(m)$, and, in particular, of the function

$$
c_{s, \chi}=\frac{\Lambda\left(2 s-1, \chi_{1} / \chi_{2}\right)}{\Lambda\left(2 s, \chi_{1} / \chi_{2}\right)} \quad\left(\text { with } \chi\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)=\chi_{1}\left(m_{1}\right) \chi_{2}\left(m_{2}\right)\right)
$$

Proof: Since $E_{s, \chi}$ meromorphically continues at least as a smooth function, the integral over the compact set $N_{k} \backslash N_{\mathbb{A}}$ expressing a pointwise value $c_{P} E_{s, \chi}(g)$ of the constant term certainly converges absolutely. In fact, the integral converges as a continuous-function-valued function $n \rightarrow\left(g \rightarrow E_{s, \chi}(n m)\right)$, so has a continuous-function-valued Gelfand-Pettis integral $m \rightarrow c_{P} E_{s, \chi}(m)$. In brief, the constant term has a meromorphic continuation. Then the vector-valued form of the Identity Principle from complex analysis implies that the form of the constant term persists outside the region of convergence $\operatorname{Re}(s)>1$. In particular, this gives the meromorphic continuation of $c_{s, \chi}$.

The theory of the constant term in [8.2] yields
[2.9.4] Claim: For every $s$ away from poles of $s \rightarrow E_{s, \chi}$, in a fixed Siegel set $\mathfrak{S}_{t, C}$,

$$
E_{s, \chi}(n m k)-\left(\left(\nu^{s} \chi\right)(m)+c_{s, \chi}\left(\nu^{1-s} \chi^{w}\right)(m)\right)<_{B} \quad \nu(m)^{-B}
$$

That is, $E_{s, \chi}-c_{P} E_{s, \chi}$ is rapidly decreasing in standard Siegel sets.
Proof: Since $E_{s, \chi}$ is an eigenfunction for invariant Laplacians and for spherical Hecke algebras, and is of moderate growth, the theory of the constant term [8.2] exactly assures that $E_{s, \chi}$ is asymptotic to its constant term, in the sense of the assertion.
[2.9.5] Corollary: The poles of $E_{s, \chi}$ are exactly the poles of the constant term $c_{s, \chi}$. ///

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Granting the meromorphic continuation and the asymptotic estimation of the Eisenstein series by its constant term, the functional equation is determined by its constant term:
[2.9.6] Corollary: $E_{s, \chi}$ has the functional equation $E_{1-s, \chi}=c_{1-s, \chi} E_{s, \chi}$, and $c_{s, \chi} \cdot c_{1-s, \chi^{w}}=1$.
[2.9.7] Remark: Note that the functional equation does not generally relate $E_{1-s, \chi}$ to $E_{s, \chi}$, but to $E_{s, \chi^{w}}$.
Proof: Take $\operatorname{Re}(s)>\frac{1}{2}$ and $s$ off the real line. The function $f=E_{1-s, \chi}-c_{1-s, \chi} E_{s, \chi^{w}}$ has constant term

$$
\begin{aligned}
c_{P} f(m)=\left(\left(\nu^{1-s} \chi\right)(m)+\right. & \left.c_{1-s, \chi} \cdot\left(\nu^{s} \chi^{w}\right)(m)\right)-c_{1-s, \chi} \cdot\left(\left(\nu^{s} \chi\right)(m)+c_{s, \chi^{w}} \cdot\left(\nu^{1-s} \chi\right)(m)\right) \\
& =\nu^{1-s}(m) \cdot\left(1-c_{1-s, \chi} c_{s, \chi}\right) \cdot \chi(m)
\end{aligned}
$$

For $\sigma=\operatorname{Re}(s)>\frac{1}{2}, \nu^{1-s}$ is square-integrable on $\mathfrak{S}_{t, C}$ : via an Iwasawa decomposition, noting that $\nu(m)^{-1} d n d m$ is left Haar measure on $P_{\mathbb{A}}$,

$$
\int_{\mathfrak{S}_{t, C}}\left|\nu^{1-s}\right|^{2}=\int_{\mathfrak{S}_{t, C}} \nu^{2-2 \sigma}(m) \nu(m)^{-1} d n d m d k \ll \int_{m \in Z^{+} M_{k} \backslash M_{\mathbb{A}}: \nu(m) \geq t} \nu(m)^{1-2 \sigma} d m \ll \int_{t}^{\infty} y^{1-2 \sigma} \frac{d y}{y}
$$

By the theory of the constant term [8.2], on a standard Siegel set

$$
\left.f=c_{P} f+(\text { rapidly decreasing }) \ll_{s} \eta^{1-\sigma}+\text { (rapidly decreasing }\right)
$$

Thus, on $\mathfrak{S}_{t, C}$,

$$
|f|^{2} \ll \mid \nu^{1-\sigma}+\left.(\text { rapidly decreasing })\right|^{2}
$$

$$
=\nu^{2(1-\sigma)}+2 \cdot \eta^{1-\sigma} \cdot(\text { rapidly decreasing })+(\text { rapidly decreasing })^{2}=\nu^{2(1-\sigma)}+(\text { rapidly decreasing })
$$

Thus, $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. For archimedean $v, f$ is a $\Delta_{v}$-eigenfunction, with eigenvalue of the form $\lambda=\left(s-i t_{v}\right)\left(s-i t_{v}-1\right)$ for $i t_{v}$ purely imaginary, depending on $\chi$. This eigenvalue is not real for $\operatorname{Re}(s)>\frac{1}{2}$ and $s \notin \mathbb{R}$. But

$$
\lambda \cdot\langle f, f\rangle=\langle\lambda f, f\rangle=\left\langle\Delta_{v} f, f\right\rangle=\overline{\left\langle f, \Delta_{v} f\right\rangle}=\overline{\langle f, \lambda f\rangle}=\bar{\lambda} \cdot \overline{\langle f, f\rangle}=\bar{\lambda} \cdot\langle f, f\rangle
$$

We did not use symmetry properties of $\Delta_{v}$, but only that $\langle f, F\rangle=\overline{\langle F, f\rangle}$. Necessarily $E_{1-s, \chi}-c_{1-s, \chi} E_{s, \chi^{w}}=$ 0 for such $s$. For all $g \in G_{\mathbb{A}}$, by the Identity Principle applied to the $\mathbb{C}$-valued meromorphic functions $s \longrightarrow\left(E_{1-s, \chi}(g)-c_{1-s, \chi} E_{s, \chi^{w}}(g)\right)$, the same identity applies for all $s$ away from poles. Since the constant term is identically 0 , necessarily $c_{1-s, \chi} c_{s, \chi^{w}}=1$.

The more general scenario needs some restrictions to stay near enough to the simple case to apply the same causal mechanisms. In particular, generalizing right $K_{\mathbb{A}^{-}}$invariance, the function $f \in I_{s, \chi}$ must be right $K_{\mathbb{A}}$-finite, in the sense that the collection of right translates of $f$ by $K_{\mathbb{A}}$ spans a finite-dimensional space of functions. For non-archimedean places, this is equivalent to being fixed by a finite-index subgroup in $\prod_{v<\infty} K_{v}$, but for archimedean places there is no such equivalence. Also, unsurprisingly, the dependence of $f$ on the complex parameter $s$ must also be controlled: take $f(n m k)=\left(\nu^{s} \chi\right)(m) f_{o}(k)$ with the function $f_{o}$ on $K_{\mathbb{A}}$ independent of $s$, and right $K_{\mathbb{A}}$-finite, and write $E\left(s, \chi, f_{o}\right)=E_{f}$. Of course, to avoid the potential ambiguity due to the non-triviality of $M_{\mathbb{A}} \cap K_{\mathbb{A}}$, it must be that $\chi$ is trivial on $M^{1} \cap K_{\mathbb{A}}$, and $f_{o}(m k)=\chi(m) f_{o}(k)$ for $m \in M^{1}$ and $k \in K_{\mathbb{A}}$, or else $f=0$. The scattering operator $C_{s, \chi}$ not only flips $s \rightarrow 1-s$ and $\chi \rightarrow \chi^{w}$, but also acts on the function $f_{o}$ by (possibly a meromorphic continuation of)

$$
\left(C_{s, \chi} f_{o}\right)(k)=\int_{N_{\mathbb{A}}} f(w n k) d k \quad\left(\text { for } k \in K_{\mathbb{A}}\right)
$$

The constant term becomes

$$
c_{P} E_{f}=c_{P} E\left(s, \chi, f_{o}\right)=\nu^{s} \chi \otimes f_{o}+\nu^{1-s} \chi^{w} \otimes C_{s, \chi}\left(f_{o}\right)
$$

[2.9.8] Theorem: $E\left(s, \chi, f_{o}\right)$ has a meromorphic continuation in $s \in \mathbb{C}$, as a smooth function of moderate growth on $\Gamma \backslash G$. As a function of $s, E\left(s, \chi, f_{o}\right)(g)$ is of at most polynomial growth vertically, uniformly in bounded strips, uniformly for $g$ in compacts. (Proofs in [2.B] and [11.5].)

The general analogue of the argument in the special case proves meromorphic continuation of scattering matrix/operators, with the qualification that they be restricted to $K_{\mathbb{A}}-f i n i t e$ functions $I_{s, \chi}^{\mathrm{fin}}$ in $I_{s \chi}$, commensurate with the conditions for meromorphic continuation of Eisenstein series.
[2.9.9] Corollary: The scattering matrix/operator $C_{s, \chi}$ restricted to a map $C_{s, \chi}^{\mathrm{fin}}: I_{s, \chi}^{\mathrm{fin}} \rightarrow I_{1-s, \chi^{w}}^{\mathrm{fin}}$, has a meromorphic continuation.
Proof: The appropriate sense of meromorphic continuation is that $C_{s, \chi} f$ has a meromorphic continuation as a $I_{1-s, \chi^{w}}^{\mathrm{fin}}$-valued function for every $f \in I_{s, \chi}^{\mathrm{fin}}$. The meromorphic continuation of $E_{f}$ gives the meromorphic continuation of $c_{P} E_{f}=f+C_{s, \chi} f$, and the special form of $f$ assures that $s \rightarrow f$ is entire, so $C_{s, \chi} f$ has a meromorphic continuation.
///
[2.9.10] Corollary: For $f_{o}$ as in the theorem, the functional equation $E\left(1-s, \chi, f_{o}\right)=E\left(s, \chi^{w}, C_{1-s, \chi^{w}} f_{o}\right)$ holds, and $C_{1-s, \chi^{w}} \circ C_{s, \chi}=1$. The operator $C_{s, \chi}^{\mathrm{fin}}$ has poles exactly where $E_{f}$ has a pole for some $f \in I_{s, \chi}^{\mathrm{fin}}$. Proof: Arranging to cancel the $\nu^{s}$ part of the constant terms,

$$
\begin{gathered}
c_{P}\left(E\left(1-s, \chi, f_{o}\right)-E\left(s, \chi^{w}, C_{1-s, \chi^{w}} f_{o}\right)\right) \\
=\left(\nu^{1-s} \chi \otimes f_{o}+\nu^{s} \chi^{w} \otimes C_{1-s, \chi} f_{o}\right)-\left(\nu^{s} \chi^{w} \otimes C_{1-s, \chi} f_{o}+\nu^{1-s} \chi \otimes C_{1-s, \chi^{w}} C_{1-s, \chi} f_{o}\right) \\
=\nu^{1-s} \chi \otimes\left(f_{o}-C_{1-s, \chi^{w}} C_{1-s, \chi} f_{o}\right)
\end{gathered}
$$

The theory of the constant term [8.2] implies that Eisenstein series $E\left(s, \chi, f_{o}\right)$ are asymptotic to their constant terms. In $\operatorname{Re}(s)>\frac{1}{2}$, the function $\nu^{1-s}$ is in $L^{2}$ on Siegel sets, so $E\left(1-s, \chi, f_{o}\right)-E\left(s, \chi^{w}, C_{1-s, \chi^{w}} f_{o}\right)$ is in $L^{2}$. However, the eigenvalues of the invariant Casimir operators $\Omega_{v}$ at archimedean places are not real in $\operatorname{Re}(s)>\frac{1}{2}$ off the real line, so this difference must be 0 . This holds for all $f_{o}$.

### 2.10 Truncation and Maaß-Selberg relations

The genuine Eisenstein series are not in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, but from the theory of the constant term [8.2] the only obstruction is the constant term, which is sufficiently altered by truncation. The Maaß-Selberg relations are computation of the $L^{2}$ inner products of the resulting truncated Eisenstein series.

The truncation operators $\wedge^{T}$ for large positive real $T$ act on an automorphic form $f$ by killing off $f$ 's constant term on $g=n m k$ for large $\nu(m)$. Thus, for a right $K_{\mathbb{A}}$-invariant function, one might imagine that

$$
\text { (naive } T \text {-truncation of } f)(n m k)= \begin{cases}f(g) & \text { for } \nu(m) \leq T \\ f(g)-c_{P} f(g) & \text { for } \nu(m)>T\end{cases}
$$

This is flawed. On a standard Siegel set $\mathfrak{S}_{t, C}$ this description is good, but it fails to describe the truncated function on the whole group $G_{\mathbb{A}}$, in the sense that this failed truncation is not an automorphic form, that is, as a left $Z^{+} G_{k}$-invariant function. Truncation should produce automorphic forms. For sufficiently large $T$ the same effect is achieved by first defining the tail $c_{P}^{T} f$ of the constant term $c_{P} f$ of $f$ :

$$
c_{P}^{T} f(n m k)=\left\{\begin{array}{cl}
0 & (\text { for } \nu(m) \leq T) \\
c_{P} f(n m k) & (\text { for } \nu(m)>T)
\end{array}\right.
$$

Although $c_{P}^{T} f$ need not be smooth, nor compactly supported, by design, for $T$ large, its support is sufficiently high to control analytical issues: writing $\Psi(\varphi)=\Psi_{\varphi}$ for legibility,
[2.10.1] Claim: For $T$ sufficiently large, the pseudo-Eisenstein series $\Psi\left(c_{P}^{T} f\right)$ is a locally finite sum, hence, uniformly convergent on compacts.

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Proof: The tail $c_{P}^{T} f$ is left $N_{\mathbb{A}}$-invariant. The reduction theory of [2.2] shows that a set $\left\{n m k: \nu(m) \geq t_{o}\right\}$ does not meet $\gamma \cdot\{n m k: \nu(m) y \geq t\}$ for $\gamma \in G_{k}$ unless $\gamma \in P_{k}$, for large-enough $t$ depending on $t_{o}$. Thus, for large-enough $T$, the set $S=\{n m k: \nu(m) \geq T\}$ does not meet its translate $\gamma \cdot S$ unless $\gamma \in P_{k}$. Thus, $\gamma_{1} \cdot S$ does not meet $\gamma_{2} \cdot S$ unless $\gamma_{1} P_{k}=\gamma_{2} P_{k}$.

Similarly,
[2.10.2] Claim: On a standard Siegel set $\mathfrak{S}_{t, C}, \Psi\left(c_{P}^{T} f\right)=c_{P}^{T} f$ for all $T$ sufficiently large depending on $t$.
Proof: By reduction theory [2.2], for large-enough $T$ depending on $t_{o}$, a set $\left\{n m k: \nu(m)>t_{o}\right\}$ does not meet $\gamma \cdot\{n m k: \nu(m)>T\}$ unless $\gamma \in P_{k}$. Thus, for large-enough $T,\{n a m k: \nu(m)>T\}$ does not meet $\mathfrak{S}_{t_{o}, C}$ unless $\gamma \in P_{k}$, and the sole non-zero summand is $c_{P}^{T} f$.

A proper definition of the truncation operator $\wedge^{T}$ is

$$
\wedge^{T} f=f-\Psi\left(c_{P}^{T} f\right)
$$

The critical effect of the truncation procedure is to have
[2.10.3] Corollary: For $K_{\mathbb{A}}$-finite $f \in I_{s, \chi}$, for $s$ away from poles, the truncated Eisenstein series $\wedge^{T} E_{f}$ is of rapid decay in all Siegel sets.

Proof: By the previous claim and by the theory of the constant term [8.2], $E_{f}-c_{P} E_{f}$ is of rapid decay in standard Siegel sets. (Meromorphic continuation uses $K_{\mathbb{A}}$-finiteness.).

Surprisingly, inner products of truncated Eisenstein series have a useful explication. Let

$$
X^{-}=\left\{g \in Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}: \eta(g)<T\right\} \quad X^{+}=\left\{g \in Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}: \eta(g) \geq T\right\}
$$

[2.10.4] Theorem: (Maaß-Selberg relation) Given $\chi, \chi^{\prime}$ characters of $M_{k} \backslash M_{\mathbb{A}}$ and $f \in I_{s, \chi}$ and $f^{\prime} \in I_{r, \chi^{\prime}}$

$$
\begin{gathered}
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \wedge^{T} E_{f} \cdot \overline{\wedge^{T} E_{f^{\prime}}}=\int_{X^{-}} f \cdot \overline{f^{\prime}}+\int_{X^{-}} f \cdot \overline{C_{r, \chi^{\prime}}\left(f^{\prime}\right)}-\int_{X^{+}} C_{s, \chi}(f) \cdot \overline{f^{\prime}}-\int_{X^{+}} C_{s, \chi}(f) \cdot \overline{C_{r, \chi^{\prime}}\left(f^{\prime}\right)} \\
=\frac{T^{s+\bar{r}-1}}{s+\bar{r}-1} \int_{M_{k} \backslash M^{1}} \chi \overline{\chi^{\prime}} \int_{K_{\mathbb{A}}} f \cdot \overline{f^{\prime}}+\frac{T^{(1-s)+\bar{r}-1}}{(1-s)+\bar{r}-1} \int_{M_{k} \backslash M^{1}} \chi^{w} \overline{\chi^{\prime}} \int_{K_{\mathbb{A}}} C_{s, \chi}(f) \cdot \overline{f^{\prime}} \\
+\frac{T^{s+(1-\bar{r})-1}}{s+(1-\bar{r})-1} \int_{M_{k} \backslash M^{1}} \chi \overline{\chi^{\prime w}} \int_{K_{\mathbb{A}}} f \cdot \overline{C_{r, \chi^{\prime}}\left(f^{\prime}\right)}+\frac{T^{(1-s)+(1-\bar{r})-1}}{(1-s)+(1-\bar{r})-1} \int_{M_{k} \backslash M^{1}} \chi^{w} \overline{\chi^{\prime w}} \int_{K_{\mathbb{A}}} C_{s, \chi}(f) \cdot \overline{C_{r, \chi^{\prime}}\left(f^{\prime}\right)}
\end{gathered}
$$

[2.10.5] Remark: The integrals over $M_{k} \backslash M^{1}$ are 0 unless the integrand is the trivial character on $M^{1}$. Proof: Because the tail of the constant term of $E_{f^{\prime}}$ is orthogonal to the truncation $\wedge^{T} E_{f}$ of $E_{f}$,

$$
\int_{Z+G_{k} \backslash G_{\mathbb{A}}} \wedge^{T} E_{f} \cdot \overline{\wedge^{T} E_{f^{\prime}}}=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \wedge^{T} E_{f} \cdot \overline{E_{f^{\prime}}}
$$

This is

$$
\int_{Z+G_{k} \backslash G_{\mathbb{A}}}\left(E_{f}-\Psi\left(\begin{array}{ll}
0 & (\text { for } \eta<T) \\
f+C_{s, \chi}(f) & (\text { for } \eta \geq T)
\end{array}\right)\right) \cdot \overline{E_{f^{\prime}}}=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi\left(\begin{array}{ll}
f & (\text { for } \eta<T) \\
-C_{s, \chi}(f) & (\text { for } \eta \geq T)
\end{array}\right) \cdot \overline{E_{f^{\prime}}}
$$

Unwinding the awkward pseudo-Eisenstein series gives

$$
\begin{aligned}
& \int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \int_{N_{k} \backslash N_{\mathbb{A}}}\left\{\begin{array}{ll}
f & (\text { for } \eta<T) \\
-C_{s, \chi}(f) & (\text { for } \eta \geq T)
\end{array} \cdot \overline{E_{f^{\prime}}}\right. \\
& =\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{A}}\left\{\begin{array}{ll}
f & (\text { for } \eta<T) \\
-C_{s, \chi}(f) & (\text { for } \eta \geq T)
\end{array} \cdot\left(\int_{N_{k} \backslash N_{\mathbb{A}}} \overline{E_{f^{\prime}}}(n g) d n\right) d g\right. \\
& =\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{A}}\left\{\begin{array}{ll}
f & (\text { for } \eta<T) \\
-C_{s, \chi}(f) & (\text { for } \eta \geq T)
\end{array} \cdot \overline{\left(f^{\prime}+C_{r, \chi^{\prime}}\left(f^{\prime}\right)\right)}\right. \\
& =\int_{X^{-}} f \cdot \overline{f^{\prime}}+\int_{X^{-}} f \cdot \overline{C_{r, \chi^{\prime}}\left(f^{\prime}\right)}-\int_{X^{+}} C_{s, \chi}(f) \cdot \overline{f^{\prime}}-\int_{X^{+}} C_{s, \chi}(f) \cdot \overline{C_{r, \chi^{\prime}}\left(f^{\prime}\right)}
\end{aligned}
$$

The sets $X^{ \pm}$are stable under the left action of $M_{k} \backslash M^{1}$. Since $f$ is left $\chi$-equivariant for $M_{k} \backslash M^{1}$ and $\overline{f^{\prime}}$ is left $\overline{\chi^{\prime}}$-equivariant, via the Iwasawa decomposition, noting that $\nu(m)^{-1} d n d m$ is left Haar measure on $P_{\mathbb{A}}$,

$$
\begin{aligned}
& \int_{X^{-}} f \cdot \overline{f^{\prime}}=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}: \nu<T} \int_{K_{\mathbb{A}}} f(m k) \cdot \overline{f^{\prime}(m k)} \nu^{-1}(m) d m d k \\
= & \int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}: \nu<T}\left(\nu^{s} \chi\right)(m) \cdot \overline{\left(\nu^{r} \chi\right)(m)} \nu^{-1}(m) d m \cdot \int_{K_{\mathbb{A}}} f(k) \cdot \overline{f^{\prime}(k)} d k
\end{aligned}
$$

The left-most integral is left $M_{k} \backslash M^{1}$-equivariant by $\chi \overline{\chi^{\prime}}$. When this is a non-trivial character the integral is 0 , by the usual cancellation trick: with $m_{o} \in M^{1}$ such that $\chi\left(m_{o}\right) \neq \chi^{\prime}\left(m_{o}\right)$, by replacing $m$ by $m_{o} m$ in the integral,

$$
\begin{gathered}
\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}: \nu<T}\left(\nu^{s} \chi\right)(m) \cdot \overline{\chi^{\prime}(m)} d m=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}: \nu<T}\left(\nu^{s} \chi\right)\left(m_{o} m\right) \cdot \overline{\chi^{\prime}\left(m_{o} m\right)} d m \\
=\chi \overline{\chi^{\prime}}\left(m_{o}\right) \int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}: \nu<T}\left(\nu^{s} \chi\right)(m) \cdot \overline{\chi^{\prime}(m)} d m
\end{gathered}
$$

For $\chi^{\prime}=\chi$, the integral over $M_{k} \backslash M^{1}$ gives a volume. What remains is the integral over the image of the fragment $(0, T)$ of the ray $(0, \infty)$, giving

$$
\int_{0}^{T} y^{s+\bar{r}-1} \frac{d y}{y}=\frac{T^{s+\bar{r}-1}}{s+\bar{r}-1}
$$

The other three summands are similarly evaluated.
[2.10.6] Corollary: Unless $\chi^{\prime}=\chi$ or $\chi^{\prime}=\chi^{w}$, for $f \in I_{s, \chi}$ and $f^{\prime} \in I_{r, \chi^{\prime}}$, the corresponding truncated Eisenstein series $\wedge^{T} E_{f}$ and $\wedge^{T} E_{f^{\prime}}$ are orthogonal.
Proof: For $\chi^{\prime} \neq \chi$ and $\chi^{\prime} \neq \chi^{w}$, all four integrals over $M_{k} \backslash M^{1}$ vanish.
The situation $\chi^{w}=\chi$ can be adjusted, by multiplying by $\chi_{1}(\operatorname{det} g)^{-1}$, to have trivial central character. Thus, the following corollary refers essentially to the case of trivial central character:
[2.10.7] Corollary: For characters $\chi^{\prime}=\chi=\chi^{w}$ of $M_{k} \backslash M_{\mathbb{A}}$, and simplest Eisenstein series $E_{s, \chi}, E_{r, \chi}$ attached to the everywhere-spherical elements of $I_{s, \chi}$ and $I_{r, \chi}$,

$$
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \wedge^{T} E_{s, \chi} \overline{\wedge^{T} E_{r, \chi}}=\frac{T^{s+\bar{r}-1}}{s+\bar{r}-1}+\frac{c_{s, \chi} T^{(1-s)+\bar{r}-1}}{(1-s)+\bar{r}-1}+\frac{c_{\bar{r}, \bar{\chi}} T^{s+(1-\bar{r})-1}}{s+(1-\bar{r})-1}+\frac{c_{s, \chi} c_{\bar{r}, \bar{\chi}} T^{(1-s)+(1-\bar{r})-1}}{(1-s)+(1-\bar{r})-1}
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

The following result has a special, more direct argument, but the proof mechanism used here is more broadly applicable.
[2.10.8] Corollary: For $f \in I_{s, \chi}$ of the form $f(n m k)=\left(\nu^{s} \chi\right)(m) \cdot f(k)$ for $n \in N_{\mathbb{A}}, m \in M_{\mathbb{A}}$, and $k \in K_{\mathbb{A}}$, with $\left.f\right|_{K_{\mathbb{A}}}$ independent of $s$, neither the Eisenstein series $E_{f}$ nor the scattering operator $C_{s, \chi}$ has any poles in $\operatorname{Re}(s) \geq \frac{1}{2}$ off the interval $\left(\frac{1}{2}, 1\right]$. The poles on $\left(\frac{1}{2}, 1\right]$, if any, are simple. When $\chi \neq \chi^{w}$, there are no poles on $\left(\frac{1}{2}, 1\right]$. Any residues in $\operatorname{Re}(s)>\frac{1}{2}$ are square-integrable.
Proof: Suppose $E_{f}$ has a pole $s_{o}=\sigma_{o}+i t_{o}$ of order $\ell>0$ with $t_{o} \neq 0$ and $\sigma_{o}>\frac{1}{2}$. Certainly the order of pole of the constant term can be no greater than that of $E_{f}$, so the second summand $C_{s, \chi}(f)$ has a pole of order at most $\ell$ at $s=s_{o}$. The first summand, $f$ itself, as a function of $s$ is entire, by the assumptions about the dependence of $f$ on $s$. Take $r=s=\sigma_{o}+i t$ in the theorem, giving an equality of the form

$$
\int_{Z+G_{k} \backslash G_{\mathbb{A}}}\left|\wedge^{T} E_{f}\right|^{2}=\frac{T^{2 \sigma_{o}-1}}{2 \sigma_{o}-1} A_{1}+\frac{T^{-2 i t}}{-2 i t} A_{2}+\frac{T^{2 i t}}{2 i t} A_{3}+\frac{T^{1-2 \sigma_{o}}}{1-2 \sigma_{o}} A_{4}
$$

The left-hand side of the Maaß-Selberg relation blows up like $\left(t-t_{o}\right)^{-2 \ell}$ as $t \rightarrow t_{o}$ on $\mathbb{R}$. The second and third terms blow up at most like $C_{s, \chi}(f)$ does, which is at worst $\left(t-t_{o}\right)^{-\ell}$. The fourth term blows up at worst like $\left|C_{s, \chi}(f)\right|^{2}$, which is at worst $\left(t-t_{o}\right)^{-2 \ell}$. Thus, as $t \rightarrow t_{o}$, the left-hand side and the fourth term on the right dominate. However, the left-hand side is positive, while the fourth term is negatrive, since $1-2 \sigma<0$. That is, there can be no such pole.

Next, let $s_{o}=\sigma_{o}$ be a pole of $E_{f}$ of order $\ell \geq 1$ on $\left(\frac{1}{2}, 1\right]$. Looking at the same expression, again, $A_{1}$ does not blow up as $t \rightarrow t_{o}=0$, unlike the previous case the second and third terms blow up at most like $t^{-(\ell+1)}$ since $t_{o}=0$, and the fourth again at most like $t^{-2 \ell}$. Again, the fourth term is negative, and if $\ell>1$ dominates the right-hand side as $t \rightarrow 0$, contradicting the positivity of the left-hand side. Thus, $\ell=1$, in which case the second and third terms' blow-up may be the same order as the left-hand side, and as the fourth term on the right-hand side. This proves that any pole on $\left(\frac{1}{2}, 1\right]$ is simple. Further, when $\chi \neq \chi^{w}$, the second and third terms are identically 0 , so there can be no pole on $\left(\frac{1}{2}, 1\right]$ in that case.

To prove square-integrability of a residue at $\sigma_{o} \in\left(\frac{1}{2}, 1\right]$, treat the Eisenstein series as a meromorphic function-valued function, as in [15.2]. Its Laurent coefficients coefficients are functions in the same topological vector space, by the vector-valued form of Cauchy's formulas [15.2]. From the Maaß-Selberg expression again, at $r=s=\sigma_{o}+i t$, multiplying through by $t^{2}$ and letting $t \rightarrow 0$, the first term on the right-hand side disappears, the powers of $T$ in the second and third terms become $T^{0}$, giving

$$
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}}\left|\operatorname{Res}_{\sigma_{o}} E_{f}^{T}\right|^{2}=\frac{\operatorname{Res}_{\sigma_{o}} A_{2}}{2}+\frac{\operatorname{Res}_{\sigma_{o}} A_{3}}{2}+\frac{T^{1-2 \sigma_{o}}}{1-2 \sigma_{o}} \lim _{t \rightarrow t_{o}} t^{2} A_{4}
$$

Since $1-2 \sigma_{o}<0$, the limit of the last term is 0 as $T \rightarrow+\infty$, given the square-integrability of the residue. Properties of meromorphic vector-valued functions [15.2] and Gelfand-Pettis integrals [14.1] assure that taking residues commutes with taking the limit as $T \rightarrow \infty$. The two remaining terms are equal, since the pole is on the real line.

Suppose $s_{o}=\frac{1}{2}+i t_{o}$ is a pole of $E_{f}$ of order $\ell \geq 1$ with $t_{o} \neq 0$. Take $r=s=\sigma+i t_{o}$ with $\sigma>\frac{1}{2}$ in the theorem, giving an equality of the form

$$
\int_{Z+G_{k} \backslash G_{\mathbb{A}}}\left|\wedge^{T} E_{f}\right|^{2}=\frac{T^{2 \sigma-1}}{2 \sigma-1} A_{1}+\frac{T^{-2 i t_{o}}}{-2 i t_{o}} A_{2}+\frac{T^{2 i t_{o}}}{2 i t_{o}} A_{3}+\frac{T^{1-2 \sigma}}{1-2 \sigma} A_{4}
$$

The left-hand side is positive, and blows up like $\left(\sigma-\frac{1}{2}\right)^{-2 \ell}$ as $\sigma \rightarrow \frac{1}{2}^{+}$, while the first three terms on the right blow up with orders at most 1 , and the fourth term is negative, impossible.

For a possible pole at $s_{o}=\frac{1}{2}$, take $r=s=\frac{1}{2}+\frac{\varepsilon}{2}(1+i)$ with $\varepsilon>0$, giving an equality of the form

$$
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}}\left|\wedge^{T} E_{f}\right|^{2}=\frac{T^{\varepsilon}}{\varepsilon} A_{1}+\frac{T^{-i \varepsilon}}{-i \varepsilon} A_{2}+\frac{T^{i \varepsilon}}{i \varepsilon} A_{3}+\frac{T^{-\varepsilon}}{-\varepsilon} A_{4}
$$

with the second and third terms absent unless $\chi^{w}=\bar{\chi}$. Thus, for $\chi^{w} \neq \bar{\chi}$, the left-hand side is positive and blows up like $\varepsilon^{-2 \ell}$ as $\varepsilon \rightarrow 0^{+}$, while the first term on the right blows up like $\varepsilon^{-1}$, and the fourth term
is negative, so this is impossible. For $\chi^{w}=\bar{\chi}$, for $\operatorname{Re}(s)=\frac{1}{2}$, the functional equation $1=C_{s, \chi}^{\mathrm{fin}} \circ C_{1-s, \chi^{w}}^{\mathrm{fin}}$ becomes

$$
1=C_{s, \chi}^{\mathrm{fin}} \circ C_{1-s, \chi^{w}}^{\mathrm{fin}}=C_{s, \chi}^{\mathrm{fin}} \circ C_{\bar{s}, \bar{\chi}}^{\mathrm{fin}}
$$

Thus, $\left(C_{\frac{1}{2}, \chi}^{\mathrm{fin}}\right)^{2}=1$, so $C_{s, \chi}^{\mathrm{fin}}$ has neither pole nor zero at $s=\frac{1}{2}$. Thus, the first three terms on the right blow up like $\varepsilon^{-1}$, while the last is negative, impossible.

### 2.11 Decomposition of pseudo-Eisenstein series: level one

From [2.7], the pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in J_{\chi}$ and varying $\chi$ generate the orthogonal complement to cuspforms in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Thus, the orthogonal complement to cuspforms is the $L^{2}$-closure of the set of these pseudo-Eisenstein series. For this section, we take trivial central character, and consider only the simplest case, right $K_{\mathbb{A}}$-invariant pseudo-Eisenstein series. These are everywhere spherical case, or level one. As earlier,

$$
\chi\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)=\chi_{1}\left(m_{1}\right) \cdot \chi_{2}\left(m_{2}\right) \quad\left(\text { for } m_{1}, m_{2} \in \mathbb{J}^{1}\right)
$$

so $\chi_{2}=\chi_{1}^{-1}$, and $\chi^{w}=\chi^{-1}$. Since $M_{k} \backslash M^{1}$ is compact, $\bar{\chi}=\chi^{-1}$. Thus, $\bar{\chi}=\chi^{-1}=\chi^{w}$. The potential ambiguity in the decomposition $g=n m k$ must be accommodated in $\chi$, or else $f=0$, so $\chi$ is unramified everywhere locally. Thus, for genuine Eisenstein series $E_{f}$, we take $f \in I_{s, \chi}$ right $K_{\mathbb{A}}$-invariant, so necessarily of the form $f(n m k)=\left(\nu^{s} \chi\right)(m)$ for $n \in N_{\mathbb{A}}, m \in M^{1}$, and $k \in K_{\mathbb{A}}$, up to a constant multiple. In principle, the constant multiple could depend on $s$, but we want $\left.f\right|_{K_{\mathbb{A}}}$ to be independent of $s$, for meromorphic continuation. Thus, take $\left.f\right|_{K_{\mathbb{A}}}=1$.

The essential harmonic analysis is Fourier transform on the real line, as Mellin transform on functions on the ray $(0,+\infty)$.

From [2.7], pseudo-Eisenstein series $\Psi_{\varphi}$ are in $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, so their integrals against genuine Eisenstein series $E_{f}$ converge absolutely, since $E_{f}$ is continuous, even after meromorphic continuation. Thus, even though this $\langle$,$\rangle cannot be the L^{2}$ pairing, since $E_{f} \notin L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, write

$$
\left\langle\Psi_{\varphi}, E_{s}\right\rangle=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot \overline{E_{f}}
$$

First consider $\chi^{w}=\chi$.
[2.11.1] Theorem: Fix unramified $\chi$ with $\chi^{w}=\chi$. Let $\varphi \in J_{\chi}$ be right $K_{\mathbb{A}}$-invariant, with trivial central character. Let $s_{o}$ run over poles of $E_{s, \chi}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$. The pseudo-Eisenstein series $\Psi_{\varphi}$ is expressible in terms of genuine Eisenstein series $E_{s, \chi}$ by an integral converging absolutely and uniformly on compacts in $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ : pointwise, uniformly on compacts,

$$
\begin{aligned}
\Psi_{\varphi}= & \frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\sum_{\operatorname{Re}\left(s_{o}\right) \geq \frac{1}{2}}\left\langle\Psi_{\varphi}, \operatorname{Res}_{s_{o}} E_{s, \chi}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s, \chi}(g) \\
& =\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \frac{\chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}
\end{aligned}
$$

[2.11.2] Remark: By various devices, for example the Poisson summation argument of [2.B], the only possible pole in $\operatorname{Re}(s) \geq \frac{1}{2}$ is at $s_{o}=1$, with residue a constant multiple of $\chi_{1} \circ$ det. However, the general pattern of argument does not depend on our fortuitous knowledge of these further details.
Proof: By an easy part of the Paley-Wiener theorem, the Mellin transform of $\varphi_{\infty} \in C_{c}^{\infty}(0, \infty)$ is entire, and has rapid decay vertically, and Mellin inversion is

$$
\varphi_{\infty}(y)=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left(\int_{0}^{\infty} \varphi_{\infty}(r) r^{-s} \frac{d r}{r}\right) y^{s} d s \quad \text { (for any real } \sigma \text { ) }
$$

$$
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$$

With

$$
\varphi\left(z a_{y} n m k\right)=\varphi_{\infty}(y) \chi(m) \quad\left(\text { for } z \in Z^{+}, a_{y}=\left(\begin{array}{cc}
\delta(y) & 0 \\
0 & 1
\end{array}\right), y>0, n \in N_{\mathbb{A}}, m \in M^{1}, k \in K_{\mathbb{A}}\right)
$$

define a kind of Mellin transform $J_{\chi} \rightarrow I_{s, \chi}$ by

$$
\mathcal{M} \varphi(s)(g)=\int_{0}^{\infty} r^{-s} \varphi\left(a_{r} g\right) \frac{d r}{r} \quad\left(\text { with } a_{r}=\left(\begin{array}{cc}
\delta(r) & 0 \\
0 & 1
\end{array}\right) \text { with } r>0\right)
$$

We decompose the pseudo-Eisenstein series $\Psi_{\varphi}$ along the ray $(0, \infty)$ :

$$
\Psi_{\varphi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma g)=\frac{1}{2 \pi i} \sum_{\gamma \in P_{k} \backslash G_{k}} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s)(\gamma g) d s
$$

Since $\chi$ is specified, and $\varphi$ is right $K_{\mathbb{A}}$-invariant, in fact

$$
\mathcal{M} \varphi(s)\left(z a_{y} n m k\right)=\mathcal{M} \varphi(s)(1) \cdot y^{s} \chi(m)
$$

Thus, although $\mathcal{M} \varphi(s)$ is a function on $G_{\mathbb{A}}$ for each $s$, it is simply a scalar multiple of the everywherespherical function in $I_{s, \chi}$. Thus, for subsequent computations, suppress the argument $g \in G_{\mathbb{A}}$, and just write

$$
\mathcal{M} \varphi(s)=\int_{0}^{\infty} r^{-s} \varphi\left(a_{r}\right) \frac{d r}{r} \quad \text { (with } a_{r}=\left(\begin{array}{cc}
\delta(r) & 0 \\
0 & 1
\end{array}\right) \text { with } r>0 \text { ) }
$$

and, commensurately,

$$
\Psi_{\varphi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma g)=\frac{1}{2 \pi i} \sum_{\gamma \in P_{k} \backslash G_{k}} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot y_{\gamma g}^{s} \chi\left(m_{\gamma g}\right) d s
$$

Taking $\sigma=0$ would be natural, but with $\sigma=0$ the double integral (sum and integral) is not absolutely convergent, and the two integrals cannot be interchanged. For $\sigma>1$, the Eisenstein series is absolutely convergent, so the rapid vertical decrease of $\mathcal{M} \varphi$ in $s$ makes the double integral absolutely convergent, and by Fubini the two integrals can be interchanged:

$$
\Psi_{\varphi}(g)=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot\left(\sum_{\gamma \in P_{k} \backslash G_{k}} y_{\gamma g}^{s} \chi\left(m_{\gamma g}\right)\right) d s \quad \quad(\text { with } \sigma>1)
$$

The inner sum is the everywhere spherical Eisenstein series $E_{s, \chi}$, so, pointwise in $g \in G_{\mathbb{A}}$,

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi} d s \quad(\text { for } \sigma>1)
$$

Although this does express $\Psi_{\varphi}$ as a superposition of eigenfunctions $E_{s, \chi}$ for invariant Laplacians and for spherical Hecke operators, it is unsatisfactory, because it should not refer to $\mathcal{M} \varphi$, but to $\Psi_{\varphi}$, to have an intrinsic integral formula. Elimination of this issue is the remainder of the argument.

We move the line of integration in

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi} d s \quad(\text { for } \sigma>1)
$$

to the left, to $\sigma=1 / 2$, which is stabilized by the functional equation of $E_{s, \chi}$. From the corollary [2.10.8] to the Maaß-Selberg relations, there are only finitely-many poles of $E_{s}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$, removing one possible

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obstacle to the contour move. From the theorem [2.B], [11.5] on meromorphic continuation, we know that even the meromorphically continued $E_{s, \chi}$ is of polynomial growth vertically in $s$, uniformly in bounded strips in $s$, uniformly for $g$ in compacts. Thus, we may move the contour, picking up finitely-many residues:

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi} d s+\sum_{s_{o}} \mathcal{M} \varphi\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}} E_{s, \chi}
$$

since the poles of $E_{s, \chi}$ in $\operatorname{Re}(s)>\frac{1}{2}$ are simple and $\mathcal{M} \varphi(s)$ is entire in $s$. The $1 / 2 \pi i$ from inversion cancels the $2 \pi i$ in the residue formula. The integral in the expression of $\Psi_{\varphi}$ in terms of $E_{s, \chi}$ can be folded in half, integrating from $\frac{1}{2}+i 0$ to $\frac{1}{2}+i \infty$ rather than from $\frac{1}{2}-i \infty$ to $\frac{1}{2}+i \infty$ :

$$
\Psi_{\varphi}-(\text { residual part })=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi}(g) d s=\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s, \chi}+\mathcal{M} \varphi(1-s) E_{1-s, \chi} d s
$$

The functional equation is $E_{1-s, \chi}=c_{1-s, \chi} E_{s, \chi^{w}}$, and we are assuming $\chi^{w}=\chi$, so

$$
\Psi_{\varphi}-(\text { residual })=\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s, \chi}+\mathcal{M} \varphi(1-s) c_{1-s, \chi} E_{s, \chi} d s
$$

To rewrite this in terms of $\Psi_{\varphi}$, use the adjunction/unwinding property of $\Psi_{\varphi}$ :

$$
\begin{aligned}
&\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle= \int_{Z+N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot \overline{c_{P} E_{s, \chi}}=\int_{Z+N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi(g) \cdot \overline{\left(y_{g}^{s} \chi\left(m_{g}\right)+c_{s, \chi} y_{g}^{1-s} \chi^{w}\left(m_{g}\right)\right)} d g \\
&=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash P_{\mathbb{A}}} \int_{K_{\mathbb{A}}} \varphi(p k) \cdot \overline{\left(y_{p k}^{s} \chi\left(m_{p k}+c_{s, \chi} y_{p k}^{1-s} \chi^{w}\left(m_{p k}\right)\right)\right.} d p d k \\
&= \int_{K_{\mathbb{A}}} 1 d k \cdot \int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} \varphi\left(a_{y}\right) \chi(m) \cdot \overline{\left(y^{s} \chi(m)+c_{s, \chi} y^{1-s} \chi^{w}(m)\right)} d m \frac{d y}{y^{2}} \\
&=\int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} \varphi\left(a_{y}\right) \chi(m) \cdot \overline{\left(y^{s} \chi(m)+c_{s, \chi} y^{1-s} \chi^{w}(m)\right)} d m \frac{d y}{y^{2}}=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot \overline{\left(y^{s}+c_{s, \chi} y^{1-s}\right)} \frac{d y}{y^{2}}
\end{aligned}
$$

by using the Iwasawa decomposition, the right $K_{\mathbb{A}^{-}}$-invariance, and $\chi^{w}=\chi$. On $\operatorname{Re}(s)=\frac{1}{2}$, this is

$$
\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot\left(y^{1-s}+c_{1-s, \bar{\chi}} y^{s}\right) \frac{d y}{y^{2}}=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot\left(y^{-s}+c_{1-s, \bar{\chi}} y^{-(1-s)}\right) \frac{d y}{y}=\mathcal{M} \varphi(s)+c_{1-s, \bar{\chi}} \mathcal{M} \varphi(1-s)
$$

Using $\bar{\chi}=\chi$,

$$
\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\mathcal{M} \varphi(s)+c_{1-s, \chi} \mathcal{M} \varphi(1-s)
$$

Thus,

$$
\Psi_{\varphi}-(\text { residual })=\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s
$$

The integral can be restored to be over the whole line $\operatorname{Re}(s)=\frac{1}{2}$, since the integrand is invariant under $s \rightarrow 1-s$ : by the functional equations of $E_{s, \chi}$ and $c_{s, \chi}$,

$$
\left\langle\Psi_{\varphi}, E_{1-s, \chi}\right\rangle \cdot E_{1-s, \chi}=\left\langle\Psi_{\varphi}, c_{1-s, \chi} E_{s, \chi}\right\rangle \cdot c_{1-s, \chi} E_{s, \chi}=\overline{c_{1-s, \chi}} c_{1-s, \chi}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi}=\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi}
$$

$$
\text { 2. The quotient } Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})
$$

Thus, dividing by 2 ,

$$
\Psi_{\varphi}-(\text { residual part })=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s
$$

It remains to explicate the finitely-many residual contributions $\mathcal{M} \varphi\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}} E_{s, \chi}$. In fact, by [2.10.8] or [2.B], there are no poles unless $\chi_{1} / \chi_{2}=1$ on $\mathbb{J}^{1}$, and then the only pole in that region is at $s_{o}=1$, with residue $\chi_{1} \circ$ det. However, we want to illustrate more widely applicable methods, as follows.

As do the pseudo-Eisenstein series, $E_{s, \bar{\chi}}$ fits into an adjunction

$$
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \cdot E_{s, \bar{\chi}}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} f(g) \cdot y_{g}^{s} \bar{\chi}\left(m_{g}\right) d g \quad\left(\text { for } f \text { on } Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)
$$

whenever the implied integrals converge absolutely. By the identity principle from complex analysis, the same formula holds for the meromorphic continuation of $E_{s, \bar{\chi}}$ for $s$ away from poles. For right $K_{\mathbb{A}}$-invariant $f$, via Iwasawa decomposition,

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} f(g) \cdot y_{g}^{s} \bar{\chi}\left(m_{g}\right) d g=\int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} c_{P} f\left(a_{y} m\right) \cdot y^{s} \bar{\chi}(m) \frac{d y}{y^{2}}
$$

even though $P_{\mathbb{A}} \cap K_{\mathbb{A}}$ is not simply $\{1\}$. The integration over the compact group $M_{k} \backslash M^{1}$ computes the $\chi$-component $\left(c_{P} f\right)^{\chi}$ of $c_{P} f$ with respect to the left action of $M^{1}$ :

$$
\int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} c_{P} f\left(a_{y} m\right) \cdot y^{s} \bar{\chi}(m) \frac{d y}{y^{2}}=\int_{0}^{\infty}\left(c_{P} f\right)^{\chi}\left(a_{y}\right) \cdot y^{s-1} \frac{d y}{y}=\mathcal{M}\left(c_{P} f\right)^{\chi}(1-s)
$$

On $\operatorname{Re}(s)=\frac{1}{2}$, where $\bar{s}=1-s$, using $1-(1-s)=s$,

$$
\left\langle f, E_{s, \chi}\right\rangle=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \cdot \overline{E_{s, \chi}}=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \cdot E_{\bar{s}, \bar{\chi}}=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \cdot E_{1-s, \bar{\chi}}=\mathcal{M}\left(c_{P} f\right)^{\chi}(s)
$$

Taking $f$ to be the pseudo-Eisenstein series $\Psi_{\varphi}$,

$$
\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\mathcal{M}\left(c_{P} \Psi_{\varphi}\right)^{\chi}(s) \quad\left(\text { on } \operatorname{Re}(s)=\frac{1}{2}\right)
$$

At a pole $s_{o}$ of $E_{s, \chi}$ in $\operatorname{Re}(s) \geq \frac{1}{2}, c_{s, \chi}$ also has a pole of the same order. Since $c_{s, \chi} \cdot c_{1-s, \chi}=1$ for $\chi^{w}=\chi$, necessarily $c_{1-s, \chi}$ has a zero at $s_{o}$. Thus, from

$$
\mathcal{M} c_{P} \Psi_{\varphi}(s)=\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\mathcal{M} \varphi(s)+c_{1-s, \chi} \mathcal{M} \varphi(1-s)
$$

at a pole $s_{o}$ of $E_{s}$

$$
\mathcal{M} c_{P} \Psi_{\varphi}\left(s_{o}\right)=\mathcal{M} \varphi\left(s_{o}\right)+c_{1-s_{o}, \chi} \mathcal{M} \varphi\left(1-s_{o}\right)=\mathcal{M} \varphi\left(s_{o}\right)+0 \cdot \mathcal{M} \varphi\left(1-s_{o}\right)=\mathcal{M} \varphi\left(s_{o}\right)
$$

That is, the value $\mathcal{M} c_{P} \Psi_{\varphi}$ at $s_{o}$ is just the value of $\mathcal{M} \varphi$ there, so the coefficients appearing in the decomposition of $\Psi_{\varphi}$ are intrinsic. Thus, the decomposition above has an intrinsic form as in the statement of the theorem. This completes the argument for the decomposition of right $K_{\mathbb{A}}$-invariant pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in J_{\chi}$, with trivial central character, and $\chi^{w} \neq \chi$. In fact, the residues at poles are constant multiples of $\chi_{1}(\operatorname{det} g)$, from [2.B].

Still with trivial central character and right $K_{\mathbb{A}}$-invariance, consider $\chi^{w} \neq \chi$ :
[2.11.3] Theorem: Fix unramified $\chi$ with $\chi^{w} \neq \chi$. Let $\varphi \in J_{\chi}$ be right $K_{\mathbb{A}}$-invariant, with trivial central character. The pseudo-Eisenstein series $\Psi_{\varphi}$ is expressible in terms of genuine Eisenstein series $E_{s, \chi}$ by an integral converging absolutely and uniformly on compacts in $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ :

$$
\Psi_{\varphi}(g)=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s
$$

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[2.11.4] Remark: As in the corollary [2.10.8] to the Maaß-Selberg relation, there is no pole at all unless $\chi_{1} / \chi_{2}=1$ on $\mathbb{J}^{1}$. This absence of poles is also visible by the Poisson summation argument [2.B].
Proof: As in the situation of the previous theorem, pointwise in $g \in G_{\mathbb{A}}$,

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi} d s \quad(\text { for } \sigma>1)
$$

Move the line of integration to the left, to $\sigma=1 / 2$, using the lack of poles for these Eisenstein series in $\operatorname{Re}(s)$,

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi} d s
$$

To rewrite this in terms of $\Psi_{\varphi}$, by the adjunction/unwinding property of $\Psi_{\varphi}$,

$$
\begin{gathered}
\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot \overline{c_{P} E_{s, \chi}}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi(g) \cdot \overline{\left(y_{g}^{s} \chi\left(m_{g}\right)+c_{s, \chi} y_{g}^{1-s} \chi^{w}\left(m_{g}\right)\right)} d g \\
\quad=\int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} \varphi\left(a_{y}\right) \chi(m) \cdot \overline{\left(y^{s} \chi(m)+c_{s, \chi} y^{1-s} \chi^{w}(m)\right)} d m \frac{d y}{y^{2}}=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot \overline{y^{s}} \frac{d y}{y^{2}}
\end{gathered}
$$

by using the Iwasawa decomposition and the right $K_{\mathbb{A}}$-invariance, since $\varphi$ is left $\chi$-equivariant under $M^{1}$ and $\chi^{w} \neq \chi$. On $\operatorname{Re}(s)=\frac{1}{2}$, this is

$$
\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot y^{1-s} \frac{d y}{y^{2}}=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot y^{-s} \frac{d y}{y}=\mathcal{M} \varphi(s)
$$

Since the functional equation $E_{1-s, \chi}=c_{1-s, \chi} E_{s, \chi^{w}}$ involves $E_{s, \chi^{w}}$ and $\chi^{w} \neq \chi$, we anticipate needing the complementary computation

$$
\begin{aligned}
& \left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot \overline{c_{P} E_{s, \chi^{w}}}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi(g) \cdot \overline{\left(y_{g}^{s} \chi^{w}\left(m_{g}\right)+c_{s, \chi^{w}} y_{g}^{1-s} \chi\left(m_{g}\right)\right)} \\
& \quad=\int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} \varphi\left(a_{y}\right) \chi(m) \cdot \overline{\left(y^{s} \chi^{w}(m)+c_{s, \chi^{w}} y^{1-s} \chi(m)\right)} d m \frac{d y}{y^{2}}=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot \overline{c_{s, \chi^{w}} y^{1-s}} \frac{d y}{y^{2}}
\end{aligned}
$$

using the Iwasawa decomposition and the right $K_{\mathbb{A}}$-invariance, since $\varphi$ is left $\chi$-equivariant under $M^{1}$ and $\chi^{w} \neq \chi$. On $\operatorname{Re}(s)=\frac{1}{2}$, using $\chi^{w}=\bar{\chi}$ due to the trivial central character, this is

$$
\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot c_{1-s, \chi} y^{s} \frac{d y}{y^{2}}=\int_{0}^{\infty} \varphi\left(a_{y}\right) \cdot c_{1-s, \chi} y^{-(1-s)} \frac{d y}{y}=c_{1-s, \chi} \mathcal{M} \varphi(1-s)
$$

The integral in the expression of $\Psi_{\varphi}$ in terms of $E_{s, \chi}$ can be folded in half, integrating from $\frac{1}{2}+i 0$ to $\frac{1}{2}+i \infty$ rather than from $\frac{1}{2}-i \infty$ to $\frac{1}{2}+i \infty$ :

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi}(g) d s=\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s, \chi}+\mathcal{M} \varphi(1-s) E_{1-s, \chi} d s
$$

The functional equation is $E_{1-s, \chi}=c_{1-s, \chi} E_{s, \chi^{w}}$, and $\chi^{w} \neq \chi$, so

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) E_{s, \chi}+\mathcal{M} \varphi(1-s) c_{1-s, \chi} E_{s, \chi^{w}} d s
$$

$$
\begin{gathered}
\text { 2. The quotient } Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A}) \\
=\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\frac{1}{2 \pi i} \int_{\frac{1}{2}+i 0}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}} d s
\end{gathered}
$$

using $\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\mathcal{M} \varphi(s)$ and $\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle=c_{1-s, \chi} \mathcal{M} \varphi(1-s)$. The integrals can be restored to be over the whole line $\operatorname{Re}(s)=\frac{1}{2}$, since the two integrals are interchanged under $s \rightarrow 1-s$ :

$$
\begin{gathered}
\left\langle\Psi_{\varphi}, E_{1-s, \chi}\right\rangle \cdot E_{1-s, \chi}=\left\langle\Psi_{\varphi}, c_{1-s, \chi} E_{s, \chi^{w}}\right\rangle \cdot c_{1-s, \chi} E_{s, \chi}=\overline{c_{1-s, \chi}} c_{1-s, \chi}\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi} \\
=c_{s, \chi^{w}} c_{1-s, \chi}\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}}=1 \cdot\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}}
\end{gathered}
$$

and similarly for $E_{s, \chi^{w}}$. Dividing by 2 ,

$$
\Psi_{\varphi}=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}} d s
$$

This completes the argument for the decomposition of right $K_{\mathbb{A}}$-invariant pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in J_{\chi}$, with trivial central character, and $\chi^{w} \neq \chi$.

### 2.12 Decomposition of pseudo-Eisenstein series: higher level

A similar argument applies to decomposition of pseudo-Eisenstein series without the everywhere-spherical constraint, necessarily tracking the additional information about the restriction of $f \in I_{s, \chi}$ to $\left(P_{\mathbb{A}} \cap K_{\mathbb{A}}\right) \backslash K_{\mathbb{A}}$. We retain the trivial central character condition, for simplicity. For meromorphic continuation and analytical properties of Eisenstein series, $\left.f\right|_{K_{\mathbb{A}}}$ should be $K_{\mathbb{A}}$-finite and $f_{K_{\mathbb{A}}}$ should not depend on $s$. For simplicity, still require invariance under $K_{\infty}=\prod_{v \mid \infty} K_{v}$, so the relaxation of conditions will be at non-archimedean places. Putting $K_{\mathrm{fin}}=\prod_{v<\infty} K_{v}$, we require $K_{\text {fin }}$-finiteness.

Let $\Theta_{v}=P_{v} \cap K_{v}$, and $\Theta_{\text {fin }}=\prod_{v<\infty} \Theta_{v}$. With fixed $\chi$, restrictions of $\varphi \in J_{\chi}$ or $f \in I_{s, \chi}$ to $K_{\text {fin }}$ necessarily lie in ${ }^{[24]}$

$$
\Phi=\left\{u \in C^{\infty}\left(K_{\mathrm{fin}}\right): u(\theta k)=\chi(\theta) \cdot u(k), \text { for } \theta \in \Theta_{\mathrm{fin}}, k \in K_{\mathrm{fin}}\right\}
$$

where smooth in this context means means locally constant. There is a natural right $K_{\text {fin }}$-invariant inner product on $\Phi$ by

$$
\left\langle u_{1}, u_{2}\right\rangle=\int_{K_{\mathbb{A}}} u_{1} \cdot \overline{u_{2}}
$$

For each irreducible representation $\rho$ of the compact group $K_{\text {fin }}$, let $\Phi^{\rho}$ be the $\rho$ isotypic component in $\Phi$, namely, the sum of all isomorphic copies of $\rho$ inside $\Phi$ [9.D.14]. The dimension of the space $\operatorname{Hom}_{K_{\text {fin }}}(\rho, \Phi)$ of $K_{\mathrm{fin}}$-homomorphisms of $\rho$ to $\Phi$ is the multiplicity of $\rho$ in $\Phi_{\mathrm{fin}}$. It is not obvious that the following claim is true, nor that it will be needed in the proof:
[2.12.1] Claim: $\Phi$ is a direct sum of irreducible representations of $K_{\text {fin }}$, and is multiplicity-free, in the sense that multiplicity of any irreducible $\rho$ in $\Phi$ is at most 1. (Proof after proof of the theorem.)

Similarly, let $J_{\chi}^{\rho}$ be the elements of $J_{\chi}$ which restrict to $\Phi_{\text {fin }}^{\rho}$ on $K_{\text {fin }}$, and are right $K_{\infty}$-invariant. Let $I_{s, \chi}^{\rho}$ be the collection of elements in $I_{s, \chi}$ which restrict to $\Phi_{\text {fin }}^{\rho}$ on $K_{\text {fin }}$, and are right $K_{\infty}$-invariant. Indeed, for each fixed $s, I_{s, \chi}^{\rho}$ is in bijection with the space $\Phi_{\text {fin }}^{\rho}$ by extending $u \in \Phi_{\text {fin }}^{\rho}$ by

$$
u_{s, \chi}\left(z a_{y} m k\right)=y^{s} \chi(m) u\left(k_{\text {fin }}\right) \quad\left(\text { for } z \in Z^{+}, n \in N_{\mathbb{A}}, y>0, m \in M^{1}, k_{o} \in K_{\text {fin }}, k_{\infty} \in K_{\infty}\right)
$$

with corresponding Eisenstein series

$$
E(s, \chi, u)(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} u_{s, \chi}(\gamma \cdot g)
$$

[24] Again, this $\Phi$ is an induced representation, but we do not immediately need any properties of such.

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Let $u_{1}, \ldots, u_{\ell}$ be an orthonormal basis for $\Phi^{\rho}$, and let $u_{j, s, \chi} \in I_{s, \chi}^{\rho}$ be the corresponding extensions. We use the fact from [2.B] that the only poles of Eisenstein series are at $s=1$, and the residues are constant multiples of $\chi_{1} \circ$ det.
[2.12.2] Theorem: For $\varphi \in J_{\chi}^{\rho}$ and $\chi^{w}=\chi$, the pseudo-Eisenstein series $\Psi_{\varphi}$ is expressible in terms of genuine Eisenstein series $E\left(s, \chi, u_{j}\right)$ by an integral converging absolutely and uniformly on compacts in $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ : pointwise, uniformly on compacts,

$$
\Psi_{\varphi}=\sum_{j} \frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E\left(s, \chi, u_{j}\right)\right\rangle \cdot E\left(s, \chi, u_{j}\right) d s+\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \frac{\chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}
$$

Proof: To track the dependence on $u \in \Phi$, modify the earlier notation slightly: let

$$
\mathcal{M}_{s} \varphi(g)=\int_{0}^{\infty} r^{-s} \varphi\left(a_{r} g\right) \frac{d r}{r} \quad \text { (with } a_{r}=\left(\begin{array}{cc}
\delta(r) & 0 \\
0 & 1
\end{array}\right) \text { with } r>0 \text { ) }
$$

Thus, $\mathcal{M}_{s}$ is a map $J_{\chi}^{\rho} \rightarrow I_{s, \chi}^{\rho}$. By Mellin inversion,

$$
\varphi(g)=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M}_{s} \varphi(s)(g) d s \quad\left(\text { for } \varphi \in J_{\chi}^{\rho}\right)
$$

Further, $\mathcal{M}_{s} \varphi=\sum_{j=1}^{\ell}\left\langle\left.\mathcal{M}_{s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle \cdot u_{j, s, \chi}$, so

$$
\begin{gathered}
\Psi_{\varphi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma \cdot g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M}_{s} \varphi(\gamma g) d s \\
=\sum_{\gamma \in P_{k} \backslash G_{k}} \sum_{j=1}^{\ell} \frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left\langle\left.\mathcal{M}_{s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle \cdot u_{j, s, \chi}(\gamma g) d s
\end{gathered}
$$

For $\sigma>1$, the integral and the infinite sum can be interchanged, giving

$$
\sum_{j=1}^{\ell} \frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left\langle\left.\mathcal{M} \varphi(s)\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle \cdot\left(\sum_{\gamma \in P_{k} \backslash G_{k}} u_{j, s, \chi}(\gamma g)\right) d s=\sum_{j=1}^{\ell} \frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty}\left\langle\left.\mathcal{M}_{s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle \cdot E\left(s, \chi, u_{j}\right) d s
$$

As in the simpler cases, fold up the integral:

$$
\int_{\sigma+i 0}^{\sigma+i \infty}\left\langle\left.\mathcal{M} \varphi(s)\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle E\left(s, \chi, u_{j}\right)+\left\langle\left.\mathcal{M} \varphi(1-s)\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle E\left(1-s, \chi, u_{j}\right) d s
$$

The functional equations of such Eisenstein series can be written

$$
E(1-s, \chi, u)=E\left(s, \chi^{w},\left.C_{1-s, \chi}\left(u_{1-s, \chi}\right)\right|_{K_{\mathbb{A}}}\right)
$$

so

$$
\Psi_{\varphi}=\sum_{j} \int_{\sigma+i 0}^{\sigma+i \infty}\left\langle\left.\mathcal{M}_{s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle \cdot E\left(s, \chi, u_{j}\right)+\left\langle\left.\mathcal{M}_{1-s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle \cdot E\left(s, \chi^{w},\left.C_{1-s, \chi}\left(u_{j, 1-s, \chi}\right)\right|_{K_{\mathbb{A}}}\right) d s
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

Here $\chi^{w}=\chi$, but it is not obvious that $\left.C_{1-s, \chi}\left(u_{j, 1-s, \chi}\right)\right|_{K_{\mathbb{A}}}$ is simply related to $u_{j}$, unlike the simple case of right $K_{\mathbb{A}^{-}}$invariant functions. The maps

$$
\left.u \longrightarrow u_{1-s, \chi} \longrightarrow C_{1-s, \chi}\left(u_{1-s, \chi}\right) \longrightarrow C_{1-s, \chi}\left(u_{1-s, \chi}\right)\right|_{K_{\mathbb{A}}} \in \Phi \quad\left(\text { for } u \in \Phi^{\rho}\right)
$$

all respect the right translation action of $K_{\mathbb{A}}$, so the image is again in $\Phi^{\rho}$. Now we use the claim: $\Phi^{\rho}$ consists of a single copy of $\rho$, so the composition of these maps is an automorphism of $\rho$ respecting the action of $K_{\text {fin }}$. Irreducible Hilbert-space representations of compact groups are finite-dimensional [9.C.7], so by a suitable form [9.D.12] of Schur's lemma, $\left.u \rightarrow C_{1-s, \chi}\left(u_{1-s, \chi}\right)\right|_{K_{\text {fin }}}$ is a scalar $c_{1-s, \chi}^{\rho}$ depending on $s, \chi$, and $\rho$, but not on $u \in \Phi^{\rho}$ :

$$
\Psi_{\varphi}=\sum_{j} \int_{\sigma+i 0}^{\sigma+i \infty}\left(\left\langle\left.\mathcal{M}_{s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle+\left\langle\left.\mathcal{M}_{1-s} \varphi\right|_{K_{\mathbb{A}}}, u_{j}\right\rangle c_{1-s, \chi}^{\rho}\right) \cdot E\left(s, \chi, u_{j}\right) d s
$$

As in the simpler cases, to express this in terms of $\Psi_{\varphi}$ itself, not $\varphi$, unwind and use the Iwasawa decomposition:

$$
\begin{aligned}
&\left\langle\Psi_{\varphi}, E(s, \chi, u)\right\rangle=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot \overline{c_{P} E(s, \chi, u)}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot \overline{\left(u_{s, \chi}+C_{s, \chi} u_{s, \chi}\right)} \\
&=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} \int_{K_{\mathbb{A}}} \varphi(p k) \cdot \overline{\left(u_{s, \chi}(p k)+C_{s, \chi} u_{s, \chi}(p k)\right.} d p d k \\
&=\int_{0}^{\infty} \int_{M_{k} \backslash M^{1}} \int_{K_{\mathbb{A}}} \chi(m) \varphi\left(a_{y} k\right) \cdot \overline{\left(y^{s} \chi(m) u_{s, \chi}(k)+y^{1-s} \chi^{w}(m) C_{s, \chi}\left(u_{s, \chi}\right)(k)\right)} \frac{d y}{y^{2}} d m d k
\end{aligned}
$$

On $\operatorname{Re}(s)=\frac{1}{2}$, and with $\chi^{w}=\chi=\bar{\chi}$, we have $\overline{C_{s, \chi}}=C_{1-s, \bar{\chi}}=C_{1-s, \chi}$, and $\overline{u_{s, \chi}}=\bar{u}_{1-s, \bar{\chi}}=\bar{u}_{1-s, \chi}$. Also, of course, $u_{s, \chi}$ is just $u$ itself on $K_{\mathbb{A}}$. Again use the fact that $\left.u \rightarrow C_{1-s, \chi}\left(u_{1-s, \chi}\right)\right|_{K_{\text {fin }}}$ is a scalar $c_{1-s, \chi}^{\rho}$, so this is

$$
\begin{gathered}
\int_{0}^{\infty} \int_{K_{\mathbb{A}}} \varphi\left(a_{y} k\right) \cdot y^{-s} \bar{u}(k)+\varphi\left(a_{y} k\right) \cdot y^{-(1-s)} C_{1-s, \chi}\left(\bar{u}_{1-s, \bar{\chi}}\right)(k) \frac{d y}{y} d k \\
=\int_{K_{\mathbb{A}}} \mathcal{M}_{s} \varphi(k) \cdot \bar{u}(k) d k+\int_{K_{\mathbb{A}}} \mathcal{M}_{1-s} \varphi \cdot C_{1-s, \chi}\left(\bar{u}_{1-s, \bar{\chi}}\right)(k) d k \\
=\int_{K_{\mathbb{A}}} \mathcal{M}_{s} \varphi(k) \cdot \bar{u}(k) d k+\int_{K_{\mathbb{A}}} \mathcal{M}_{1-s} \varphi(k) \cdot c_{1-s, \chi}^{\rho} \cdot \bar{u}(k) d k=\left\langle\left.\mathcal{M}_{s} \varphi\right|_{K_{\mathbb{A}}}+\left.\mathcal{M}_{1-s} \varphi\right|_{K_{\mathbb{A}}} \cdot c_{1-s, \chi}^{\rho}, u\right\rangle
\end{gathered}
$$

Thus, the coefficients in the expression for $\Psi_{\varphi}$ are these inner products, apart from the residue picked up by moving the contour from $\sigma>1$ to $\sigma=\frac{1}{2}$. Regardless of choice of the basis $u_{j}$, the residues are all constant multiples of $\chi_{1} \circ$ det, so their sum must be as indicated.
[2.12.3] Theorem: For $\varphi \in J_{\chi}^{\rho}$ and $\chi^{w} \neq \chi$, the pseudo-Eisenstein series $\Psi_{\varphi}$ is expressible in terms of genuine Eisenstein series $E\left(s, \chi, u_{j}\right)$ by an integral converging absolutely and uniformly on compacts in $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ : pointwise, uniformly on compacts,

$$
\Psi_{\varphi}=\sum_{j} \frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E\left(s, \chi, u_{j}\right)\right\rangle \cdot E\left(s, \chi, u_{j}\right) d s+\left\langle\Psi_{\varphi}, E\left(s, \chi^{w}, u_{j}\right)\right\rangle \cdot E\left(s, \chi^{w}, u_{j}\right) d s
$$

(The proof combines the argument from the level-one analogue with the multiplicity-free claim.
Now we prove the multiplicity-free property $\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{K_{\text {fin }}}(\rho, \Phi) \leq 1$ : First, continuity of $\rho$ requires that it restricts to 1 on all but finitely-many $K_{v}$. Irreducibles of finite products of compact groups are (external) tensor products of irreducibles of the factors [9.C.8]. Thus, it suffices to prove a local fact, that

$$
\Phi_{v}=\left\{u \in C^{\infty}\left(K_{v}\right): u(\theta k)=\chi(\theta) \cdot u(k), \text { for } \theta \in \Theta_{v}, k \in K_{v}\right\}
$$

is multiplicity-free. By the Gelfand-Kazhdan criterion [6.11], it suffices to find an involutive antiautomorphism $\sigma$ of $K_{v}$ such that every left and right $\Theta_{v}$-invariant distribution $u$ on $K_{v}$ is invariant under $\sigma$, $u^{\sigma}=u$, where $u^{\sigma}(\varphi)=u\left(\varphi^{\sigma}\right)$, for all $\varphi \in C_{c}^{\infty}\left(K_{v}\right)$, where $(\varphi \circ \sigma)(k)=u\left(k^{\sigma}\right)$. We will use $g^{\sigma}=w\left(g^{\top}\right) w^{-1}$ with $w$ the Weyl element, which stabilizes $\Theta_{v}$. We find representatives for $\Theta_{v} \backslash K_{v} / \Theta_{v}$, for $v$ non-archimedean. Suppress the subscript $v$ in what follows.

Given $g=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in K$, for $c=0$, we have representative 1 . For $c \in \mathfrak{o}^{\times}$, left multiplication by $\Theta$ can make $c=1$, and also $a=0$ by subtracting an integer multiple of the lower row from the upper. Then (the modifed version of) $b$ is in $\mathfrak{o}^{\times}$, so can be made 1 , and right multiplication by $\Theta$ makes $d=0$, by subtracting an integer multiple of the left column from the right. This gives representatives $w$. For the intermediate cases $0<\operatorname{ord} c=\ell<\infty$, both $a, d$ must be units for the determinant to be a unit. Right multiplication by $\Theta$ makes $b=0$, by subtracting an integer of the left column from the right, and then $a=d=1$ by left or right multiplication by $\Theta$. Then left and right multiplication by $\left(\begin{array}{ll}u & 0 \\ 0 & 1\end{array}\right)$ with $u \in \mathfrak{o}^{\times}$does not disturb $a=d=1$, and makes $c=\varpi^{\ell}$ with chosen local parameter $\varpi$. Thus, there are representatives

$$
r_{\infty}=\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right), r_{o}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right), r_{1}=\left(\begin{array}{cc}
1 & 0 \\
\varpi & 1
\end{array}\right), r_{2}=\left(\begin{array}{cc}
1 & 0 \\
\varpi^{2} & 1
\end{array}\right), r_{3}=\left(\begin{array}{cc}
1 & 0 \\
\varpi^{3} & 1
\end{array}\right), \ldots
$$

Each representative is fixed under $\sigma$, so the double cosets are stabilized by $\sigma$. Every double coset $\Theta r \Theta$ is closed, being continuous images of compacts. The double coset $\Theta r_{\infty} \Theta=\Theta$ is not open, but all other double cosets $\Theta r_{\ell} \Theta$ are open, being defined by the open condition $\left|\varpi^{\ell+1}\right|<|c|<\left|\varpi^{\ell-1}\right|$. The characteristic function $\varphi_{\ell}$ of $\Theta r_{\ell} \Theta$ is a test function for $\ell<\infty$. The integrations

$$
u_{\infty}(\varphi)=\int_{\Theta} \varphi \quad u_{o}(\varphi)=\int_{\Theta r_{o} \Theta} \varphi \quad u_{1}(\varphi)=\int_{\Theta r_{1} \Theta} \varphi \quad \ldots \quad \quad\left(\text { for } \varphi \in C_{c}^{\infty}(K)\right)
$$

are left and right $\Theta$-invariant, and $\sigma$-invariant.
The uniqueness of invariant functionals [14.4] shows that $u_{\ell}$ is the unique $\Theta \times \Theta$-invariant distribution on $K$ supported on the compact, open set $\Theta r_{\ell} \Theta$ for $\ell<\infty$, up to constants. Left and right $\Theta$-invariant distributions factor through the two-sided averaging map $\varphi \rightarrow \int_{\Theta \times \Theta} \varphi\left(\theta k \theta^{\prime}\right) d \theta d \theta^{\prime}$. The space $\mathcal{D}=C_{c}^{\infty}(K)$ is the colimit over compact open subgroups $H$ of the finite-dimensional spaces $\mathcal{D}^{H}$ of test functions left and right $H$-invariant: indeed, by smoothness, $u \in \mathcal{D}$ is left $H_{1}$-invariant and right $H_{2}$-invariant for some compact-open subgroups $H_{i}$, and take $H=H_{1} \cap H_{2}$. For $\varphi \in \mathcal{D}^{H}$, the $\Theta \times \Theta$-averaged $\varphi$ is constant on $\Theta H \Theta$. The representatives $r_{\ell}$ approach $r_{\infty}=1_{2} \in K$. Thus, $\Theta r_{\ell} \Theta$ for every $\ell \geq \ell_{o}$ with $\ell_{o}=\ell_{o}(H)$ depending on $H$. Letting $\operatorname{ch}_{\Theta H \Theta}$ be the characteristic function of $\Theta H \Theta$,

$$
u(\varphi)=u\left(\varphi-\varphi(0) \cdot \operatorname{ch}_{\Theta H \Theta}\right)+\varphi(0) \cdot u\left(\operatorname{ch}_{\Theta H \Theta}\right)=u\left(\varphi-\varphi(0) \cdot \operatorname{ch}_{\Theta H \Theta}\right)+u_{\infty}(\varphi) \cdot u\left(\operatorname{ch}_{\Theta H \Theta}\right)
$$

The test function $\varphi-\varphi(0) \cdot \operatorname{ch}_{\Theta H \Theta}$ ) is supported on the finitely-many double cosets $\Theta r_{\ell} \Theta$ with $0 \leq \ell<\ell_{o}(H)$. The restriction of $u$ to test functions supported on this finite union of compact, open double cosets is a linear combination of $u_{0}, u_{1}, \ldots, u_{\ell_{o}-1}$, and the constant $u\left(\operatorname{ch}_{\Theta H \Theta}\right)$ does not depend on the individual $\varphi \in \mathcal{D}^{H}$. Thus, the restriction of $u$ to $\mathcal{D}^{H}$ is $\sigma$-invariant. Thus, $u$ is $\sigma$-invariant on colim ${ }_{H} \mathcal{D}^{H}$, the ascending union of the spaces $\mathcal{D}^{H}$. This verifies the hypothesis for application of the Gelfand-Kazhdan criterion, so $\Phi_{v}$ is multiplicity-free.

### 2.13 Plancherel for pseudo-Eisenstein series: level one

The previous decompositions can be refined to prove convergence of the integral as a $C^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ valued integral, from a corresponding result for behavior of Fourier inversion integrals. This refinement gives a Plancherel theorem for pseudo-Eisenstein series. For simplicity, we treat trivial central character. This entails $\chi_{2}=\chi_{1}^{-1}$, so $\chi^{w}=\bar{\chi}=\chi^{-1}$. More significantly, we restrict our attention to level one, that is, right $K_{\mathbb{A}}$-invariant, pseudo-Eisenstein series $\Psi_{\varphi}$ in this section. One corollary, awkward to obtain otherwise, is the mutual orthogonality of pairs of pseudo-Eisenstein series made from data in $J_{\chi}$ and $J_{\chi^{\prime}}$ with $\chi^{\prime} \neq \chi$ and $\chi^{\prime} \neq \chi^{w}$ on $M_{k} \backslash M^{1}$.
[2.13.1] Claim: With $\chi^{w}=\chi$, the integral in

$$
\Psi_{\varphi}=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \frac{\chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}
$$

converges as a vector-valued integral, taking values in the Fréchet space $C^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ of continuous functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.
Proof: Let $\psi_{\xi}(x)=e^{i \xi x}$ on $\mathbb{R}$. From [14.3], and as already applied in [1.13], the integral expressing Fourier inversion for Schwartz functions $f$ on the real line

$$
f(x)=\frac{1}{2 \pi} \int_{-\infty}^{\infty}\left(\int_{-\infty}^{\infty} f(u) \bar{\psi}_{\xi}(u) d u\right) \psi_{\xi}(x) d \xi=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \psi_{\xi}(x) \cdot \widehat{f}(\xi) d \xi
$$

converges as a Gelfand-Pettis integral with values in the Fréchet space $C^{o}(\mathbb{R})$. Changing coordinates, Mellin inversion gives convergence as Gelfand-Pettis integral with values in smooth functions $C^{o}(0,+\infty)$. With

$$
\varphi\left(z a_{y} n m k\right)=\varphi_{\infty}(y) \chi(m) \quad\left(\text { for } z \in Z^{+}, a_{y}=\left(\begin{array}{cc}
\delta(y) & 0 \\
0 & 1
\end{array}\right), y>0, n \in N_{\mathbb{A}}, m \in M^{1}, k \in K_{\mathbb{A}}\right)
$$

define a transform $J_{\chi} \rightarrow I_{s, \chi}$ by

$$
\left.\mathcal{M} \varphi(s)(g)=\int_{0}^{\infty} r^{-s} \varphi\left(a_{r} g\right) \frac{d r}{r} \quad \text { (with } a_{r}=\left(\begin{array}{cc}
\delta(r) & 0 \\
0 & 1
\end{array}\right) \text { with } r>0\right)
$$

Because $\varphi$ is completely determined except as a function on the ray $(0,+\infty)$, the inversion integral

$$
\varphi(g)=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s)(g) d s
$$

converges as a vector-valued integral with values in the Fréchet space $C^{o}\left(G_{\mathbb{A}}\right)$. From [6.1], left and right translation by $G_{\mathbb{A}}$ are continuous maps on $C^{o}\left(G_{\mathbb{A}}\right)$, so the linear operators of left translation by $G_{k}$ commute with the integral, and in the region of convergence, the expression of the pseudo-Eisenstein series $\Psi_{\varphi}$

$$
\Psi_{\varphi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma g)=\frac{1}{2 \pi i} \sum_{\gamma \in P_{k} \backslash G_{k}} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \varphi(s)(\gamma g) d s
$$

converges as a vector-valued integral with values in that Fréchet space. By the same steps as in the proof of the numerical form of the theorem,

$$
\Psi_{\varphi}=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle E_{s, \chi} d s+\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \frac{\chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}
$$

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as a $C^{o}\left(G_{\mathbb{A}}\right)$-valued Gelfand-Pettis integral. Since the integrand is in $C^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$-valued and the topology on this subspace is the restriction of that from $C^{o}\left(G_{\mathbb{A}}\right)$, the integral converges in $C^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. $\quad / / /$
[2.13.2] Corollary: For $\varphi, \psi$ right $K_{\mathbb{A}}$-invariant functions in $J_{\chi}$ with $\chi^{w}=\chi$,

$$
\left\langle\Psi_{\varphi}, \Psi_{\psi}\right\rangle=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \overline{\left\langle\Psi_{\psi}, E_{s, \chi}\right\rangle} d s+\frac{\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \overline{\left\langle\Psi_{\psi}, \chi_{1} \circ \operatorname{det}\right\rangle}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}
$$

Proof: For $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, the map $F \rightarrow \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} F \cdot \bar{f}$ is a continuous linear functional on $F \in C^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, so the Gelfand-Pettis property legitimizes the obvious interchange:

$$
\begin{gathered}
\left\langle\Psi_{\varphi}, f\right\rangle=\left\langle\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle E_{s, \chi} d s+\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \frac{\chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}, f\right\rangle \\
=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle\left\langle E_{s, \chi}, f\right\rangle d s+\frac{\left\langle\Psi_{\varphi}, \chi_{1} \circ \operatorname{det}\right\rangle \cdot\left\langle\chi_{1} \circ \operatorname{det}, f\right\rangle}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}
\end{gathered}
$$

where $\left\langle E_{s, \chi}, f\right\rangle$ converges because $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Taking $f=E_{\psi}$ for $\psi$ right $K_{\mathbb{A}}$ invariant in $J_{\chi}$, this gives the asserted isometry.

The discussion for trivial central character, $\chi^{w} \neq \chi$, and right $K_{\mathbb{A}}$-invariance proceeds along similar lines: [2.13.3] Claim: With $\chi^{w} \neq \chi$, the integral

$$
\Psi_{\varphi}=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi}+\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}} d s
$$

converges as a vector-valued integral, taking values in the Fréchet space $C^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ of continuous functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.
[2.13.4] Corollary: For $\varphi, \psi$ right $K_{\mathbb{A}}$-invariant functions in $J_{\chi}$ with $\chi^{w} \neq \chi$,

$$
\left\langle\Psi_{\varphi}, \Psi_{\psi}\right\rangle=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \overline{\left\langle\Psi_{\psi}, E_{s, \chi}\right\rangle}+\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \overline{\left\langle\Psi_{\psi}, E_{s, \chi^{w}}\right\rangle} d s
$$

These decomposition formulas facilitate comparison of pseudo-Eisenstein series:
[2.13.5] Corollary: For $\chi^{\prime} \neq \chi$ and $\chi^{\prime} \neq \chi^{w}$, pseudo-Eisenstein series made from $J_{\chi}$ are orthogonal to those made from $J_{\chi^{\prime}}$.
Proof: For $\varphi \in J_{\chi}$, for $\chi^{w}=\chi$ or not, we have a convergent $C^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$-valued integral

$$
\Psi_{\varphi}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot E_{s, \chi} d s+\sum_{s_{o}} \mathcal{M} \varphi\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}} E_{s, \chi}
$$

where $s_{o}$ runs over poles of $E_{s, \chi}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$. Inner product with the compactly-supported $\Psi_{\psi}$ is a continuous functional, so this inner product passes inside the integral by [14.1], giving

$$
\left\langle\Psi_{\varphi}, \Psi_{\psi}\right\rangle=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot\left\langle E_{s, \chi}, \Psi_{\psi}\right\rangle d s+\sum_{s_{o}} \mathcal{M} \varphi\left(s_{o}\right) \cdot\left\langle\operatorname{Res}_{s_{o}} E_{s, \chi}, \Psi_{\psi}\right\rangle
$$

Similarly, from [15.2], taking residues commutes with application of the functional. Unwinding,

$$
\left\langle E_{s, \chi}, \Psi_{\psi}\right\rangle=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} E_{s, \chi} \cdot \bar{\psi}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}}\left(\varphi_{s, \chi}+c_{s, \chi} \varphi_{1-s, \chi^{w}}\right) \cdot \bar{\psi}
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

Under left multiplication by $M^{1}$, the function $\varphi_{s, \chi}$ is equivariant by $\chi$, and $\varphi_{1-s, \chi^{w}}$ by $\chi^{w}$, while $\psi$ is left equivariant by $\chi^{\prime}$. Thus, the integral is 0 .

As in [1.13], these decomposition formulas suggest the form of a Plancherel theorem for the $\chi^{\text {th }}$ fragment of the complement to cuspforms. As in [1.13], for each fixed $\chi$ we must identify the closure of the image of

$$
\Psi_{\varphi} \longrightarrow \begin{cases}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \oplus\left\langle\Psi_{\varphi}, \operatorname{Res}_{s=1} E_{s, \chi}\right\rangle & \left(\text { when } \chi^{w}=\chi, \text { with } K_{\mathbb{A}} \text {-invariant } \varphi \in J_{\chi}\right) \\ \left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \oplus\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle & \left(\text { when } \chi^{w} \neq \chi, \text { with } K_{\mathbb{A}} \text {-invariant } \varphi \in J_{\chi}\right)\end{cases}
$$

and certify that residues behave compatibly with the simplest outcome. The functional equation of $E_{s, \chi}$ constrains the functions $s \rightarrow\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle$ and $s \rightarrow\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle$ :

$$
\left\langle\Psi_{\varphi}, E_{1-s, \chi}\right\rangle=\left\langle\Psi_{\varphi}, c_{1-s, \chi} E_{s, \bar{\chi}}\right\rangle=\overline{c_{1-s, \chi}} \cdot\left\langle\Psi_{\varphi}, E_{s, \bar{\chi}}\right\rangle=c_{s, \chi^{w}} \cdot\left\langle\Psi_{\varphi}, E_{s, \chi}{ }^{w}\right\rangle
$$

We wil show that the $L^{2}$ closure of the set of images is as large as possible, given these obvious constraints. In both cases, the map $\varphi \rightarrow \mathcal{M} \varphi$ is essentially Fourier transform, so maps test functions to a space of functions dense in the Schwartz functions on $L^{2}\left(\frac{1}{2}+i \mathbb{R}\right)$. Then we proceed differently depending on whether $\chi^{w}=\chi$ or not. The case $\chi^{w}=\chi$ is much like [1.13]: for $\chi^{w}=\chi$, let $V_{\chi}$ be the subspace of $L^{2}\left(\frac{1}{2}+i \mathbb{R}\right)$ functions meeting $f(1-s)=c_{s, \chi} \cdot f(s)$.
[2.13.6] Claim: With fixed $\chi^{w}=\chi$, the map

$$
\Psi_{\varphi} \longrightarrow\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \oplus\left\langle\Psi_{\varphi}, \operatorname{Res}_{s=1} E_{s, \chi}\right\rangle
$$

has dense image in $V \oplus \mathbb{C}$, and is an $L^{2}$-isometry.
Proof: In this case, $\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\mathcal{M} c_{P} \Psi_{\varphi}(s)=\mathcal{M} \eta(s)+c_{1-s, \chi} \mathcal{M} \eta(1-s)$. For $F$ in the Schwartz space on $\frac{1}{2}+i \mathbb{R}$, the averaging $F(s)+c_{1-s, \chi} F(1-s)$ maps to a dense subspace of $V$. Thus, ignoring for a moment the residual summand, the images $\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle$ are dense in $V$, as desired.

The residue is $\chi_{1} \circ$ det, and this should be orthogonal to $\Psi_{\psi}$ with $\varphi^{\prime} \in J_{\chi^{\prime}}$ and $\chi^{\prime} \neq \chi$. Indeed, unwinding the pseudo-Eisenstein series and using Iwasawa,

$$
\left\langle\Psi_{\psi}, \chi_{1} \circ \operatorname{det}\right\rangle=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \psi(g) \cdot c_{P} \bar{\chi}_{1}(\operatorname{det}(g)) d g=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} \psi(m) \cdot \bar{\chi}_{1}(\operatorname{det}(m)) \frac{d m}{\delta(m)}
$$

where $\delta$ is the modular function of $P_{\mathbb{A}}$. Let $r$ be the number of isomorphism classes of archimedean completions of $k$, and let

$$
\left.A^{+}=\left\{\left(\begin{array}{cc}
t^{1 / r} & 0 \\
0 & 1
\end{array}\right): t>0\right\} \quad \text { (on the diagonal in } M_{\infty}=\prod_{v \mid \infty} M_{v}\right)
$$

Using $Z^{+} M_{k} \backslash M_{\mathbb{A}} \approx A^{+} \times M_{k} \backslash M^{1}$, the integrand is equivariant by a non-trivial character of $M_{k} \backslash M^{1}$, so is 0 . Even more simply, the various functions $\chi_{1} \circ$ det are mutually orthogonal.

Since $\chi_{1} \circ$ det is in the orthogonal complement to cuspforms, it is in the closure of the space generated by pseudo-Eisenstein series. We have just shown that it is orthogonal to all of these except those with data from $J_{\chi}$, so it must be in the closure of the images from $J_{\chi}$ alone. By subtraction, the integral part of the decomposition is also in the closure of the pseudo-Eisenstein series, so the images are $L^{2}$ dense in $V \oplus \mathbb{C}$, as claimed.

Then the spectral-coefficient map extends by continuity, to give an $L^{2}$ isometry, the statement of a Plancherel theorem for this fragment of $L^{2}$.

For the $\chi^{w} \neq \chi$ case, let

$$
V=\left\{f=f_{1} \oplus f_{2} \in L^{2}\left(\frac{1}{2}+i \mathbb{R}\right) \oplus L^{2}\left(\frac{1}{2}+i \mathbb{R}\right): f_{1}(1-s)=f_{2}(s)\right\} \subset L^{2}\left(\frac{1}{2}+i \mathbb{R}\right) \oplus L^{2}\left(\frac{1}{2}+i \mathbb{R}\right)
$$

[2.13.7] Claim: With fixed $\chi^{w} \neq \chi$, the map

$$
\Psi_{\varphi} \longrightarrow\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \oplus\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle
$$

has dense image in $V$, and is an $L^{2}$-isometry.
Proof: In this case, $\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle=\mathcal{M} c_{P} \Psi_{\varphi}=\mathcal{M} \eta(s)$, and this is in the Schwartz space on $\frac{1}{2}+i \mathbb{R}$, which is dense in $L^{2}$. The functional equation relating $E_{s, \chi}$ determines $\left\langle\Psi_{\varphi}, E_{1-s, \chi^{w}}\right\rangle$. Thus, the images are dense in $V$, as desired. Then the spectral-coefficient map extends by continuity, to give an $L^{2}$ isometry, the statement of a Plancherel theorem for this fragment of $L^{2}$.

### 2.14 Spectral expansion, Plancherel theorem: level one

From [2.7], the collection of right $K_{\mathbb{A}}$-invariant pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in J_{\chi}$ and $\chi$ running over pairs $\chi$ of unramified characters $\chi_{1}, \chi_{2}$ with $\chi_{2}=\chi^{-1}$ (due to trivial central character) is dense in the orthogonal complement in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to right $K_{\mathbb{A}}$-invariant cuspforms.

For unramified $\chi_{1}$ on $\mathbb{J}$, let $\chi\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right)=\chi_{1}\left(m_{1} / m_{2}\right)$. For $\chi_{1}^{\prime} \neq \chi_{1}^{ \pm 1}$, the adjunctions [2.7] satisfied by pseudo-Eisenstein series show that $\left\langle\Psi_{\varphi}, E_{\varphi^{\prime}}\right\rangle=0$ for $\varphi \in J_{\chi}$ and $\varphi^{\prime} \in J_{\chi^{\prime}}$. Lettin $F$ run over an orthonormal basis for the space of cuspforms on $Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}$ with trivial central character, we have an automorphic Plancherel theorem at level one:
[2.14.1] Theorem: With trivial central character, for $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$,

$$
\begin{aligned}
f=\sum_{\mathrm{cfm} F}\langle f, F\rangle & \cdot F+\sum_{\chi_{1}}\left(\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle f, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\frac{\left\langle f, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}\right) \\
& +\sum_{\chi: \chi^{w} \neq \chi}\left(\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi}+\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}} d s\right) \quad \quad \text { (in an } L^{2} \text { sense) }
\end{aligned}
$$

and

$$
\begin{aligned}
|f|_{L^{2}}^{2}=\sum_{\operatorname{cfm} F}|\langle f, F\rangle|^{2}+ & \sum_{\chi: \chi^{w}=\chi}\left(\frac{1}{4 \pi} \int_{-\infty}^{\infty}\left|\left\langle f, E_{\frac{1}{2}+i t, \chi}\right\rangle\right|^{2} d t+\frac{\left|\left\langle f, \chi_{1} \circ \operatorname{det}\right\rangle\right|^{2}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}\right) \\
& +\sum_{\chi: \chi^{w} \neq \chi}\left(\frac{1}{4 \pi} \int_{-\infty}^{\infty}\left|\left\langle\Psi_{\varphi}, E_{\frac{1}{2}+i t, \chi}\right\rangle\right|^{2}+\left|\left\langle\Psi_{\varphi}, E_{\frac{1}{2}+i t, \chi^{w}}\right\rangle\right|^{2} d t\right)
\end{aligned}
$$

The integrals suggested by the notation are not literal integrals, but are the extension-by-continuity of the corresponding literal integrals, as with Fourier transform and Fourier inversion on $L^{2}(\mathbb{R})$.

Combining the decomposition of right $K_{\mathbb{A}}$-invariant $L^{2}$ cuspforms (with trivial central character) and the decomposition of their orthogonal complement:
[2.14.2] Corollary: Functions $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ with trivial central character have $L^{2}$ expansions

$$
\begin{aligned}
f=\sum_{\operatorname{cfm} F}\langle f, F\rangle \cdot F & +\sum_{\chi: \chi^{w}=\chi}\left(\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle f, E_{s, \chi}\right\rangle \cdot E_{s, \chi} d s+\frac{\left\langle f, \chi_{1} \circ \operatorname{det}\right\rangle \cdot \chi_{1} \circ \operatorname{det}}{\left|\chi_{1} \circ \operatorname{det}\right|_{L^{2}}^{2}}\right) \\
& +\sum_{\chi: \chi^{w} \neq \chi}\left(\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\varphi}, E_{s, \chi}\right\rangle \cdot E_{s, \chi}+\left\langle\Psi_{\varphi}, E_{s, \chi^{w}}\right\rangle \cdot E_{s, \chi^{w}} d s\right) \quad \text { (in an } L^{2} \text { sense) }
\end{aligned}
$$

with corresponding equality of $L^{2}$ norms, where integrals involving Eisenstein series are isometric extensions, as in the previous section

### 2.15 Exotic eigenfunctions, discreteness of pseudo-cuspforms

An important variant approach both to the discrete decomposition of the space of cuspforms [1.7] and [2.6], and to the meromorphic continuation of Eisenstein series as in [11.5], is the notion of pseudo-cuspform. The largest space of pseudo-cuspforms with cut-off height $b \geq 0$ is

$$
L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\left\{f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right): c_{P} f\left(m^{\prime} a_{y} k\right)=0 \text { for } m^{\prime} \in M^{1}, b<y \in \mathbb{R}, k \in K\right\}
$$

The idea is that the constant terms vanish above height $b$. With $b=0$, this is the space of squareintegrable cuspforms. More precisely, via the adjunction [1.7.3], $L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is the orthogonal complement in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to all pseudo-Eisenstein series $\Psi_{\varphi}$ with data $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ supported on

$$
Z^{+} N_{\mathbb{A}} M_{k} \backslash\left\{z n m a_{y} k: z \in Z^{+}, n \in N_{\mathbb{A}}, m \in M^{1}, b<y \in \mathbb{R}, k \in K_{\mathbb{A}}\right\}
$$

That is, these are pseudo-Eisenstein series $\Psi_{\varphi}$ with data $\varphi$ supported above height $y=b$.
However, as throughout this chapter, right $K_{\mathbb{A}}$-finiteness assumptions avoid some relatively uninteresting secondary complications. Thus, for simplicity, we consider only right $K_{\infty} K^{\prime}$-fixed functions for $K^{\prime}$ a fixed open subgroup of $K_{\text {fin }}$.

The pseudo-Laplacian $\widetilde{\Delta}_{b}$ is the Friedrichs self-adjoint extension [9.2] of the sum $\Delta=\sum_{\left.v\right|_{\infty}} \Delta_{v}$ of the invariant Laplacians on the archimedean quotients $G_{v} / K_{v}$, restricted to the test functions in the space $L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K+\infty K^{\prime}\right)$ of pseudo-cuspforms. For any $b>0$, the corresponding space of square-integrable pseudo-cuspforms contains the space of genuine cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. The basic, unexpected result is [2.15.1] Theorem: $L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is a direct sum of eigenspaces for $\widetilde{\Delta}_{b}$, each of finite dimension. In particular, $\widetilde{\Delta}_{b}$ has compact resolvent. (Proof in [10.3].)

Without further information, this does not immediately prove that the subspace consisting of genuine cuspforms decomposes discretely for $\widetilde{\Delta}_{b}$, because the description [9.2] of $\widetilde{\Delta}_{b}$ imposes technical conditions on possible eigenfunctions, and one should check that the smooth cuspforms are dense in the space of $L^{2}$ cuspforms.

In any case, for $b \gg 1$, the space $L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ contains many functions not in the space of genuine cuspforms, for example, pseudo-Eisenstein series $\Psi_{\varphi}$ with data $\varphi$ supported in the interval [1, b]. As in [2.11] and [2.12], these are expressible as integrals of genuine Eisenstein series. However, by the theorem, apparently these pseudo-Eisenstein series are (infinite) sums of $L^{2}$-eigenfunctions for $\widetilde{\Delta}_{b}$ orthogonal. Similarly, by [2.10.3] and [2.10.4], truncated Eisenstein series $\wedge^{b} E_{f}$ are in $L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Because they are in the span of pseudoEisenstein series with compactly supported data, by [2.11] and [2.12] they are integrals of genuine Eisenstein series. Again, however, by the theorem, they are also (infinite) sums of $\widetilde{\Delta}_{b}$-eigenfunctions. Evidently, there are many exotic eigenfunctions for $\widetilde{\Delta}_{b}$. Indeed,
[2.15.2] Corollary: The eigenfunctions for $\widetilde{\Delta}_{b}$ in $L_{b}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\text {fin }} K^{\prime}\right)$ with eigenvalues $\lambda=s(s-1)<$ $-1 / 4$ are exactly the truncated Eisenstein series $\wedge^{b} E_{f}$ with $c_{P} E_{f}\left(a_{b}\right)=0$, for right $K_{\infty}$-invariant right $K^{\prime}$-right-invariant $f \in I_{s, \chi}$, for $s(s-1)<0$ and character $\chi$ on $M_{k} \backslash M^{1} /\left(M^{1} \cap K_{\infty} K^{\prime}\right)$. (Proof in [10.4].).

These truncated Eisenstein series are not smooth. The slightly non-intuitive nature of the operator $\widetilde{\Delta}_{b}$ explains the situation, in [10.4]. For example, in addition to meeting the Gelfand condition of constant-term vanishing about height $b$, eigenfunctions of the pseudo-Laplacian $\widetilde{\Delta}_{b}$ are pseudo-cuspforms in a stronger sense:
[2.15.3] Corollary: An $L^{2}$-eigenfunction $u$ for $\widetilde{\Delta}_{b}$ with eigenvalue $\lambda$ satisfies $\left(\widetilde{\Delta}_{b}-\lambda\right) u=0$ locally at heights above $b$.

## 2.A Appendix: compactness of $\mathbb{J}^{1} / k^{\times}$

The following compactness result has both finiteness of class numbers and Dirichlet's units theorem as corollaries. Indeed, the compactness can be proven as a consequence of these two results. However, the compactness can be proven directly, and is what proves useful here.
[2.A.1] Theorem: $\mathbb{J}^{1} / k^{\times}$is compact.
Proof: As in [5.2], Haar measure on $\mathbb{A}=\mathbb{A}_{k}$ and Haar measure on the (topological group) quotient $\mathbb{A} / k$ are inter-related in the sense that choice of one uniquely determines the other by the relation

$$
\int_{\mathbb{A}} f(x) d x=\int_{\mathbb{A} / k} \sum_{\gamma \in k} f(\gamma+x) d x \quad\left(\text { for } f \in C_{c}^{o}(\mathbb{A})\right)
$$

Normalize the measure on $\mathbb{A}$ so that, mediated by this relation, $\mathbb{A} / k$ has measure 1 . We have a Minkowskilike claim, a measure-theoretic pigeon-hole principle, that a compact subset $C$ of $\mathbb{A}$ with measure greater than 1 cannot inject to the quotient $\mathbb{A} / k$. Suppose, to the contrary, that $C$ injects to the quotient. With $f$ the characteristic function of $C$,

$$
1<\int_{\mathbb{A}} f(x) d x=\int_{\mathbb{A} / k} \sum_{\gamma \in k} f(\gamma+x) d x \leq \int_{\mathbb{A} / k} 1 d x=1
$$

with the last inequality by injectivity. Contradiction. For idele $\alpha$, the change-of-measure on $\mathbb{A}$ is given conveniently by

$$
\frac{\text { meas }(\alpha E)}{\text { meas }(E)}=|\alpha| \quad \text { (for measurable } E \subset \mathbb{A} \text { ) }
$$

Given $\alpha \in \mathbb{J}^{1}$, we will adjust $\alpha$ by $k^{\times}$to lie in a compact subset of $\mathbb{J}^{1}$. Fix compact $C \subset \mathbb{A}$ with measure $>1$. The topology on $\mathbb{J}$ is strictly finer than the subspace topology with $\mathbb{J} \subset \mathbb{A}$ : the genuine topology is by imbedding $\mathbb{J} \rightarrow \mathbb{A} \times \mathbb{A}$ by $\alpha \rightarrow\left(\alpha, \alpha^{-1}\right)$. For $\alpha \in \mathbb{J}^{1}$, both $\alpha C$ and $\alpha^{-1} C$ have measure $>1$, neither injects to the quotient $k \backslash \mathbb{A}$. So there are $x \neq y$ in $k$ so that $x+\alpha C=y+\alpha C$. Subtracting,

$$
0 \neq a=x-y \in \alpha(C-C) \cap k
$$

That is, $a \cdot \alpha^{-1} \in C-C$. Likewise, there is $0 \neq b \in \alpha^{-1}(C-C) \cap k$, and $b \cdot \alpha \in C-C$. There is an obvious constraint

$$
\left.a b=(a \cdot \alpha)\left(b \cdot \alpha^{-1}\right) \in(C-C)^{2} \cap k^{\times}=\text {compact } \cap \text { (discrete subgroup }\right)=\text { finite }
$$

as in [1.5.3]. Let $\Xi=(C-C)^{2} \cap k^{\times}$be this finite set. Paraphrasing: given $\alpha \in \mathbb{J}^{1}$, there are $a \in k^{\times}$and $\xi \in \Xi(\xi=a b$ above $)$ such that $\left(a \cdot \alpha^{-1},\left(a \cdot \alpha^{-1}\right)^{-1}\right) \in(C-C) \times \xi^{-1}(C-C)$. That is, $\alpha^{-1}$ can be adjusted by $a \in k^{\times}$to be in the compact $C-C$, and, simultaneously, for one of the finitely-many $\xi \in \Xi$, $\left(a \cdot \alpha^{-1}\right)^{-1} \in \xi \cdot(C-C)$. In the topology on $\mathbb{J}$, for each $\xi \in \Xi$,

$$
\left((C-C) \times \xi^{-1}(C-C)\right) \cap \mathbb{J}=\text { compact in } \mathbb{J}
$$

The continuous image in $\mathbb{J} / k^{\times}$of each of these finitely-many compacts is compact. Their union covers the closed subset $\mathbb{J}^{1} / k^{\times}$, so the latter is compact.
2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

## 2.B Appendix: meromorphic continuation

A somewhat special argument gives precise information about the meromorphic continuation of certain Eisenstein series for $G L_{2}$, in particular about possible poles. Analogous arguments are possible in a few other situations. As above, with $\chi$ denoting a pair of characters $\chi_{1}, \chi_{2}$ on $\mathbb{J}^{1}$ as above, take

$$
f\left(z n a_{r} m \cdot k\right)=|r|^{s} \cdot \chi_{1}\left(m_{1}\right) \chi_{2}\left(m_{2}\right) \cdot f_{o}(k) \quad\left(\text { with } m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) \in M^{1}\right)
$$

with $f_{o}$ independent of $s$. In particular, this argument for meromorphic continuation uses the following expression for $f$.

With diagonal map $\delta:(0,+\infty) \rightarrow \mathbb{J}$ as earlier, abuse notation by extending $\chi_{1}, \chi_{2}$ to characters on $\mathbb{J}$ by extending trivially on $\delta(0,+\infty)$ :

$$
\chi_{j}(\delta(r) \cdot \theta)=\chi_{j}(\theta) \quad\left(\text { for } r>0 \text { and } \theta \in \mathbb{J}^{1}, j=1,2\right)
$$

By changing variables in the integral, one finds that any function $f$ expressed as

$$
f(g)=\frac{|\operatorname{det} g|^{s} \chi_{1}(\operatorname{det} g)}{\zeta\left(2 s, \frac{\chi_{1}}{\chi_{2}}, \Phi(0,-)\right)} \int_{\mathbb{J}}|t|^{2 s} \frac{\chi_{1}}{\chi_{2}}(t) \cdot \Phi(t \cdot e \cdot g) d t \quad \quad\left(\text { with } e=\left(\begin{array}{ll}
0 & 1
\end{array}\right)\right)
$$

with Schwartz function $\Phi$ on $\mathbb{A}^{2}$ and global Iwasawa-Tate zeta integral

$$
\zeta\left(2 s, \frac{\chi_{1}}{\chi_{2}}, \Phi(0,-)\right)=\int_{\mathbb{J}}|t|^{2 s} \frac{\chi_{1}}{\chi_{2}}(t) \Phi(0, t) d t
$$

is in $I_{s, \chi}$. The division by $\zeta\left(2 s, \frac{\chi_{1}}{\chi_{2}}, \Phi(0,-)\right)$ normalizes $f(1)=1$. We do not consider the issue of exactly which $f \in I_{s, \chi}$ can be expressed in this form.

Every Schwartz function $\Phi$ is a finite sum of monomial functions $\Phi=\bigotimes_{v} \Phi_{v}$, so permissible functions $f$ are finite sums of monomial functions $f=\bigotimes_{v} f_{v}$ with the local functions $f_{v}$ on $G_{v}$ expressible as

$$
f_{v}(g)=\frac{|\operatorname{det} g|^{s} \chi_{1}(\operatorname{det} g)}{\zeta_{v}\left(2 s, \frac{\chi_{1}}{\chi_{2}}, \Phi_{v}(0,-)\right)} \int_{k_{v}^{\times}}|t|_{v}^{2 s} \frac{\chi_{1}}{\chi_{2}}(t) \cdot \Phi_{v}(t \cdot e \cdot g) d t \quad \quad\left(\text { with } e=\left(\begin{array}{ll}
0 & 1
\end{array}\right)\right)
$$

The product formula makes $f$ left $P_{k}$-invariant. The corresponding Eisenstein series is

$$
E_{f}(g)=E(s, \chi, \Phi)(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} f(\gamma \cdot g)
$$

Let $w=\left(\begin{array}{rr}0 & -1 \\ 1 & 0\end{array}\right)$.
[2.B.1] Theorem: These Eisenstein series admit meromorphic continuations to $\mathbb{C}$, and have no poles in $\operatorname{Re}(s) \geq \frac{1}{2}$ unless $\chi_{1} / \chi_{2}=1$, in which case there is a unique pole in $\operatorname{Re}(s) \geq \frac{1}{2}$, at $s=1$, with residue a constant multiple of $\chi_{1}(\operatorname{det} g)$. The functional equation is

$$
\frac{\zeta(2 s, \rho, \Phi(0,-))}{\chi_{1} \circ \operatorname{det}} \cdot E(s, \chi, \Phi)=\frac{\zeta\left(2-2 s, \rho^{-1},(\widehat{\Phi} \circ w)(0,-)\right)}{\chi_{2} \circ \operatorname{det}} \cdot E\left(1-s, \chi^{w}, \widehat{\Phi} \circ w\right)
$$

[2.B.2] Corollary: For trivial central character $\omega$, these Eisenstein series $E_{f}$ have no poles in $\operatorname{Re}(s) \geq \frac{1}{2}$ unless $\chi_{1}^{2}=1$, in which case there is a unique pole in $\operatorname{Re}(s) \geq \frac{1}{2}$, at $s=1$, with residue a constant multiple of $\chi_{1}(\operatorname{det} g)$.
Proof: (of Corollary) With trivial central character, $\chi_{2}=\chi_{1}^{-1}$.

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Proof: To isolate the Poisson summation effect, for unitary character $\rho=\chi_{1} / \chi_{2}$ on $\mathbb{J} / k^{\times}$, let

$$
\mathcal{E}(s, \rho, \Phi)(g)=\frac{\zeta(2 s, \rho, \Phi(0,-))}{|\operatorname{det} g|^{s} \chi_{1}(\operatorname{det} g)} \cdot E(s, \chi, \Phi)(g)=\int_{\mathbb{J}}|t|^{2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \Phi(t \cdot x \cdot g) d t
$$

With $\mathbb{J}^{+}=\{t \in \mathbb{J}:|t| \geq 1\}$ and $\mathbb{J}^{-}=\{t \in \mathbb{J}:|t| \leq 1\}$, we follow Riemann, Hecke, Iwasawa, and Tate, and break the integral into two pieces:

$$
\mathcal{E}(s, \rho, \Phi)(g)=\int_{\mathbb{J}+}|t|^{2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \Phi(t \cdot x \cdot g) d t+\int_{\mathbb{J}-}|t|^{2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \Phi(t \cdot x \cdot g) d t
$$

By the first lemma below, the integral over $\mathbb{J}^{+}$is absolutely convergent for all $s \in \mathbb{C}$, so is entire. Adelic Poisson summation converts the integral over $\mathbb{J}^{-}$to an integral over $\mathbb{J}^{+}$, plus two elementary terms: first,

$$
\begin{gathered}
\sum_{0 \neq x \in k^{2}} \Phi(t x g)=\sum_{x \in k^{2}} \Phi(t x g)-\Phi(0)=|t|^{-2}|\operatorname{det} g|^{-1} \sum_{x \in k^{2}} \widehat{\Phi}\left(t^{-1} x\left(g^{\top}\right)^{-1}\right)-\Phi(0) \\
=|t|^{-2}|\operatorname{det} g|^{-1} \cdot \sum_{0 \neq x \in k^{2}} \widehat{\Phi}\left(t^{-1} x\left(g^{\top}\right)^{-1}\right)-\Phi(0)+|t|^{-2}|\operatorname{det} g|^{-1} \cdot \widehat{\Phi}(0)
\end{gathered}
$$

Thus, inverting $t$ to replace the integral over $\mathbb{J}^{-}$by one over $\mathbb{J}^{+}$,

$$
\begin{gathered}
\int_{\mathbb{J}-}|t|^{2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \Phi(t x g) d t \\
=\int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t)\left(|t|^{-2}|\operatorname{det} g|^{-1} \sum_{0 \neq x \in k^{2}} \widehat{\Phi}\left(t^{-1} x\left(g^{\top}\right)^{-1}\right)-\Phi(0)+|t|^{-2}|\operatorname{det} g|^{-1} \cdot \widehat{\Phi}(0)\right) d t \\
=|\operatorname{det} g|^{-1}\left(\int_{\mathbb{J}^{+}}|t|^{2-2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \widehat{\Phi}\left(t x\left(g^{\top}\right)^{-1}\right) d t+\widehat{\Phi}(0) \int_{\mathbb{J}^{-}}|t|^{2 s} \rho^{-1}(t)|t|^{-2} d t\right)-\Phi(0) \int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t) d t
\end{gathered}
$$

Altogether,

$$
\begin{gathered}
\mathcal{E}(s, \rho, \Phi)(g)=\int_{\mathbb{J}^{+}}|t|^{2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \Phi(t \cdot x \cdot g) d t-\Phi(0) \int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t) d t \\
+|\operatorname{det} g|^{-1}\left(\int_{\mathbb{J}^{+}}|t|^{2-2 s} \rho^{-1}(t) \sum_{0 \neq x \in k^{2}} \widehat{\Phi}\left(t x\left(g^{\top}\right)^{-1}\right) d t+\widehat{\Phi}(0) \int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t)|t|^{-2} d t\right)
\end{gathered}
$$

Multiplying through by $|\operatorname{det} g|^{\frac{1}{2}}$ symmetrizes this:

$$
\begin{gathered}
|\operatorname{det} g|^{\frac{1}{2}} \cdot \mathcal{E}(s, \rho, \Phi)(g)=|\operatorname{det} g|^{\frac{1}{2}} \cdot \int_{\mathbb{J}+}|t|^{2 s} \rho(t) \sum_{0 \neq x \in k^{2}} \Phi(t \cdot x \cdot g) d t-|\operatorname{det} g|^{\frac{1}{2}} \cdot \Phi(0) \int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t) d t \\
+\left\lvert\, \operatorname{det}\left(g^{\top}\right)^{\frac{1}{2}}\left(\int_{\mathbb{J}^{+}}|t|^{2-2 s} \rho^{-1}(t) \sum_{0 \neq x \in k^{2}} \widehat{\Phi}\left(t x\left(g^{\top}\right)^{-1}\right) d t+\left.\left|\operatorname{det}\left(g^{\top}\right)^{\frac{1}{2}} \widehat{\Phi}(0) \int_{\mathbb{J}-}\right| t\right|^{2 s} \rho(t)|t|^{-2} d t\right)\right. \\
=\left|\operatorname{det}\left(g^{\top}\right)^{-1}\right|^{\frac{1}{2}} \cdot \mathcal{E}\left(1-s, \rho^{-1}, \widehat{\Phi}\right)\left(\left(g^{\top}\right)^{-1}\right)
\end{gathered}
$$

We examine the two elementary integrals which, if non-zero, give the poles. If $\rho(\tau) \neq 1$ for some $\tau \in \mathbb{J}^{1}$, then by changing variables,

$$
\int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t) d t=\int_{\mathbb{J}^{-}}|\tau t|^{2 s} \rho(\tau t) d t=\rho(\tau) \int_{\mathbb{J}^{-}}|t|^{2 s} \rho(t) d t
$$

## 2. The quotient $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$

so the integral must vanish. On the other hand, when $\rho(\tau)=1$ on $\mathbb{J}^{1}$, we give the compact group $\mathbb{J}^{1} / k^{\times}$ measure 1, and

$$
\int_{\mathbb{J}-}|t|^{2 s} d t=\int_{0}^{1} y^{2 s} \frac{d y}{y}=\frac{1}{2 s} \quad \text { and } \quad \int_{\mathbb{J}^{-}}|t|^{2 s-2} d t=\int_{0}^{1} y^{2 s-2} \frac{d y}{y}=\frac{1}{2 s-2}
$$

Thus, when $\chi_{1} / \chi_{2}$ is not identically 1 on $\mathbb{J}^{1}$, there are no poles. When $\rho=\chi_{1} / \chi_{2}$ is identically 1 on $\mathbb{J}^{1}$, there are polar terms in $|\operatorname{det} g|^{\frac{1}{2}} \cdot \mathcal{E}(s, \chi, \Phi)(g)$, and they are symmetrical:

$$
-\frac{\Phi(0) \cdot|\operatorname{det} g|^{\frac{1}{2}}}{2 s}-\frac{\widehat{\Phi}(0) \cdot\left|\operatorname{det}\left(g^{T}\right)^{-1}\right|^{\frac{1}{2}}}{2(1-s)}
$$

Thus, the preliminary form of the functional equation:

$$
|\operatorname{det} g|^{\frac{1}{2}} \cdot \mathcal{E}(s, \chi, \Phi)(g)=\left|\operatorname{det}\left(g^{\top}\right)^{-1}\right|^{\frac{1}{2}} \cdot \mathcal{E}\left(1-s, \chi^{w}, \widehat{\Phi}\right)\left(\left(g^{\top}\right)^{-1}\right)
$$

We would prefer not to have a relationship involving $\left(g^{\top}\right)^{-1}$. Fortunately, $w^{-1}\left(g^{\top}\right)^{-1} w=g /(\operatorname{det} g)$. Thus, replacing $x$ by $x w^{-1}$ in the sum, replacing $\widehat{\Phi}$ by $\widehat{\Phi} \circ w$, and replacing $t$ by $t \cdot \operatorname{det} g$, in the region for $\operatorname{Re}(1-s)>1$ for convergence,

$$
\left|\operatorname{det}\left(g^{\top}\right)^{-1}\right|^{\frac{1}{2}} \cdot \mathcal{E}\left(1-s, \chi^{w}, \widehat{\Phi}\right)\left(\left(g^{\top}\right)^{-1}\right)=|\operatorname{det} g|^{-\frac{1}{2}} \cdot|\operatorname{det} g|^{2-2 s} \frac{\chi_{2}}{\chi_{1}}(\operatorname{det} g) \cdot \mathcal{E}\left(1-s, \chi^{w}, \widehat{\Phi} \circ w\right)(g)
$$

Thus, by the identity principle,

$$
|\operatorname{det} g|^{\frac{1}{2}} \cdot \mathcal{E}(s, \chi, \Phi)(g)=|\operatorname{det} g|^{-\frac{1}{2}} \cdot|\operatorname{det} g|^{2-2 s} \frac{\chi_{2}}{\chi_{1}}(\operatorname{det} g) \cdot \mathcal{E}\left(1-s, \chi^{w}, \widehat{\Phi} \circ w\right)(g)
$$

and

$$
\begin{aligned}
& |\operatorname{det} g|^{\frac{1}{2}} \cdot \frac{\zeta(2 s, \rho, \Phi(0,-))}{|\operatorname{det} g|^{s} \chi_{1}(\operatorname{det} g)} \cdot E(s, \chi, \Phi)(g) \\
& =|\operatorname{det} g|^{-\frac{1}{2}} \cdot|\operatorname{det} g|^{2-2 s} \frac{\chi_{2}}{\chi_{1}}(\operatorname{det} g) \cdot \frac{\zeta\left(2-2 s, \rho^{-1}, \widehat{\Phi} \circ w(0,-)\right)}{|\operatorname{det} g|^{1-s} \chi_{2}(\operatorname{det} g)} \cdot E\left(1-s, \chi^{w}, \widehat{\Phi} \circ w\right)(g)
\end{aligned}
$$

simplifying to

$$
\frac{\zeta(2 s, \rho, \Phi(0,-))}{\chi_{1}(\operatorname{det} g)} \cdot E(s, \chi, \Phi)(g)=\frac{\zeta\left(2-2 s, \rho^{-1}, \widehat{\Phi} \circ w(0,-)\right)}{\chi_{2}(\operatorname{det} g)} \cdot E\left(1-s, \chi^{w}, \widehat{\Phi} \circ w\right)(g)
$$

[2.B.3] Lemma: Half-zeta integrals over $\mathbb{J}^{+}$are absolutely convergent for all $s \in \mathbb{C}$.
Proof: Fix $g \in G_{\mathbb{A}}$, let $\varphi(t)=\Phi(t e g)$ and $\varphi_{v}(t)=\Phi_{v}(t e g)$. By the lemma below,

$$
|\varphi(t)|<_{N} \prod_{v} \sup \left(\left|t_{v}\right|_{v}, 1\right)^{-2 N} \quad\left(\text { for adele } t=\left\{t_{v}\right\}, \text { for all } N\right)
$$

For idele $t$ let $\nu(t)=\prod_{v} \sup \left(\left|t_{v}\right|_{v},\left|t_{v}\right|_{v}^{-1}\right)$. Almost all factors on the right-hand side are 1 , so there is no problem with convergence. Apply

$$
(\sup (a, 1))^{2}=\sup \left(a^{2}, 1\right)=a \cdot \sup \left(a, \frac{1}{a}\right) \quad(\text { for } a>0)
$$

to every factor:

$$
\prod_{v} \sup \left(\left|t_{v}\right|_{v}, 1\right)^{-2 N}=|t|^{-N} \prod_{v} \sup \left(\left|t_{v}\right|_{v},\left|t_{v}^{-1}\right|_{v}\right)^{-N}=|t|^{-N} \nu(t)^{-N}
$$

Thus, on $\mathbb{J}^{+}$,

$$
\prod_{v} \sup \left(\left|t_{v}\right|_{v}, 1\right)^{-2 N}=|t|^{-N} \nu(t)^{-N} \leq \nu(t)^{-N} \quad\left(\text { when } t \in \mathbb{J}^{+}, N \geq 0\right)
$$

With $\sigma=\operatorname{Re} s$, for every $N \geq 0$

$$
\left.\left.\left|\int_{\mathbb{J}^{1}}\right| y\right|^{s} \varphi(t) d t\left|\ll \int_{\mathbb{J}^{1}}\right| t\right|^{\sigma} \nu(t)^{-N} d t \ll \int_{\mathbb{J}}|t|^{\sigma} \nu(t)^{-N} d t=\prod_{v}\left(\int_{k_{v}^{\times}}|t|^{\sigma} \sup \left(|t|, \frac{1}{|t|^{-N}} d t\right)\right.
$$

For $N>|\sigma|$, the non-archimedean local integrals are absolutely convergent:
$\int_{k_{v}^{\times}}|t|^{\sigma} \sup \left(|t|, \frac{1}{|t|}\right)^{-N} d t=\sum_{\ell=0}^{\infty} q_{v}^{-\sigma-N}+\sum_{\ell=1}^{\infty} q_{v}^{\sigma-N}=\frac{1}{1-q^{-\sigma-N}}+\frac{q^{\sigma-N}}{1-q^{\sigma-N}}=\frac{1-q^{-2 N}}{\left(1-q^{-\sigma-N}\right)\left(1-q^{\sigma-N}\right)}$
The archimedean local integrals are convergent for similar reasons. For $2 N>1$ and $N>|\sigma|+1$, the Euler product is dominated by the Euler product for the expression $\zeta_{k}(N+\sigma) \zeta_{k}(N-\sigma) / \zeta_{k}(2 N)$ in terms of the zeta function $\zeta_{k}(s)$ of $k$, which converges absolutely.
[2.B.4] Lemma: For all $N$, a Schwartz function $\varphi$ on $\mathbb{A}$ satisfies

$$
|\varphi(t)| \ll{ }_{\varphi, N} \prod_{v} \sup \left(\left|t_{v}\right|_{v}, 1\right)^{-2 N} \quad(\text { for } t \in \mathbb{A})
$$

Proof: By definition, a Schwartz function $\varphi$ on $\mathbb{A}$ is a finite sum of monomials $\varphi=\bigotimes_{v} \varphi_{v}$. Thus, it suffices to consider monomial $\varphi$, and to prove the local assertion that for $\varphi_{v} \in \mathscr{S}\left(k_{v}\right)$

$$
\left|\varphi_{v}(t)\right| \ll N_{N} \sup \left(\left|t_{v}\right|_{v}, 1\right)^{-2 N} \quad\left(\text { for } t \in k_{v}\right)
$$

At archimedean places, the definition of the Schwartz space requires that

$$
\left.\sup _{t \in k_{v}}\left(1+|t|_{v}\right)^{N} \cdot\left|\varphi_{v}(t)\right|<\infty \quad \quad \quad \text { archimedean } k_{v}, \text { for all } N\right)
$$

so

$$
\left|\varphi_{v}(t)\right| \ll N_{N}\left(1+|t|_{v}\right)^{-2 N} \leq \sup \left(|t|_{v}, 1\right)^{-2 N}
$$

Almost everywhere, $\varphi_{v}$ is the characteristic function of the local integers. At such places,

$$
\left|\varphi_{v}(t)\right|=\left\{\begin{array}{ll}
1 & \left(\text { for }|t|_{v} \leq 1\right) \\
0 & \left(\text { for }|t|_{v}>1\right)
\end{array}\right\} \leq \sup \left(|t|_{v}, 1\right)^{-2 N} \quad(\text { for all } N)
$$

At the remaining bad finite primes, $\varphi_{v} \in \mathscr{S}\left(k_{v}\right)$ is compactly supported and locally compact. Then, similar to the good prime case,

$$
\left|\varphi_{v}(t)\right|<_{\varphi_{v}}\left\{\begin{array}{ll}
1 & \left(t \in \operatorname{spt} \varphi_{v}\right) \\
0 & \left(t \notin \operatorname{spt} \varphi_{v}\right)
\end{array}\right\}<_{\varphi_{v}, N} \sup \left(|t|_{v}, 1\right)^{-2 N} \quad(\text { for all } N)
$$

This proves the lemma.

## 2.C Appendix: Hecke-Maaß periods of Eisenstein series

These examples, essentially due to Hecke and Maaß, include as special cases both sums of values at Heegner points, and integrals over hyperbolic geodesics. The set-up of the previous appendix allows a simple computation for $G L_{2}$ over a number field $k$.

Let $\ell$ be a quadratic field extension of a number field $k$. Let $G=G L_{2}(k)$, and $H \subset G$ a copy of $\ell^{\times}$inside $G$, by specifying the isomorphism in

$$
\ell^{\times} \subset \operatorname{Aut}_{k}(\ell) \approx \operatorname{Aut}_{k}\left(k^{2}\right)=G L_{2}(k)
$$

Let $P$ be the standard parabolic of upper-triangular elements in $G$. Factor the idele class group $\mathbb{J}_{\ell} / \ell^{\times}=$ $(0, \infty) \times \mathbb{J}_{\ell}^{1} / \ell^{\times}$where the ray $(0, \infty)$ is imbedded on the diagonal in the archimedean factors. Let $\chi$ be a character on $\mathbb{J}_{k} / k^{\times}$trivial on the ray $(0, \infty)$, and define a character on $P_{\mathbb{A}}$ by

$$
\chi_{s}\left(\begin{array}{ll}
a & * \\
0 & d
\end{array}\right)=\left|\frac{a}{d}\right|^{s} \cdot \chi(a) \cdot \chi^{-1}(d)
$$

Let $\varphi_{s, \chi}$ be a left $\chi_{s}$-equivariant function on $G_{\mathbb{A}}: \varphi_{s, \chi} h(p g)=\chi_{s}(p) \cdot \varphi_{s, \chi}(g)$ for all $p \in P_{\mathbb{A}}$ and $g \in G_{\mathbb{A}}$. At places $v$ where $\chi$ is unramified, we may as well take $\varphi_{s, \chi}$ to be right $K_{v}$-invariant, where $K_{v}$ is the standard maximal compact in $G_{v}$. This function has trivial central character. Ignoring the ambiguity at bad primes, put

$$
E_{s, \chi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s, \chi}(\gamma \cdot g)
$$

Let $Z$ be the center of $G$, and $\omega$ the quadratic character of $\ell / k$. Let $S$ be a (finite) set of places including those ramified in $\ell / k$ or ramified for $\chi$.

## [2.C.1] Theorem:

$$
\int_{Z_{\mathbb{A}} H_{k} \backslash H_{\mathbb{A}}} E_{s, \chi}=\frac{L^{S}(s, \chi) \cdot L^{S}(s, \chi \cdot \omega)}{L^{S}\left(2 s, \chi^{2}\right)} \times(\text { bad prime factors })
$$

where $L^{S}(s, \alpha)$ denotes the $L$-function attached to a Hecke character $\alpha$ over $k$ dropping the local factors attached to places $v \in S$.
[2.C.2] Remark: In fact, as in the proof, the numerator arises as an $L$-function over the quadratic extension $\ell$, namely $L_{\ell}^{S}(s, \chi \circ N)$, where $N$ is the norm $\ell \rightarrow k$. Then quadratic reciprocity gives the indicated factorization into $L$-functions over the base field $k$.
Proof: The subgroup $P_{k}$ is the isotropy group of a $k$-line $k \cdot e$ for a fixed non-zero $e \in k^{2} \approx \ell$. The group $G_{k}$ is transitive on these $k$-lines, so $P_{k} \backslash G_{k} \approx\{k$-lines $\}$. The critical-but-trivial point is that the action of $\ell^{\times}$ on $\ell$ is transitive on non-zero elements. Thus, $P_{k} \cdot H_{k}=G L_{2}(k)$. That is, the period integral unwinds

$$
\int_{Z_{\mathbb{A}} H_{k} \backslash H_{\mathbb{A}}} E_{s, \chi}=\int_{Z_{\mathbb{A}}\left(P_{k} \cap H_{k}\right) \backslash H_{\mathbb{A}}} \varphi_{s, \chi}=\int_{Z_{\mathbb{A}} \backslash H_{\mathbb{A}}} \varphi_{s, \chi}
$$

since $H \cap P=Z$. With $\varphi_{s, \chi}$ chosen to factor over primes $\varphi_{s, \chi}=\bigotimes_{v} \varphi_{s, \chi, v}$, the unwound period integral likewise factors over primes

$$
\int_{Z_{\mathbb{A}} H_{k} \backslash H_{\mathbb{A}}} E_{s, \chi}=\int_{Z_{\mathbb{A}} \backslash H_{\mathbb{A}}} \varphi_{s, \chi}=\prod_{v} \int_{Z_{v} \backslash H_{v}} \varphi_{s, \chi, v}
$$

A graceful way to evaluate the local integrals is to use an integral representation of the local vectors $\varphi_{s, \chi, v}$ akin to well-known archimedean devices involving the Gamma function. That is, express $\varphi_{v}$ in terms of Iwasawa-Tate local zeta integrals

$$
L_{k, v}\left(2 s, \chi^{2}\right)=\int_{k_{v}^{\times}}|t|_{v}^{2 s} \chi^{2}(t) \Phi_{v}(t e) d t
$$

as

$$
\varphi_{s, \chi, v}(g)=\frac{1}{L_{k, v}\left(2 s, \chi_{v}^{2}\right)} \cdot|\operatorname{det} g|_{v}^{s} \chi_{v}(\operatorname{det} g) \cdot \int_{k_{v}^{\times}}|t|_{v}^{2 s} \chi_{v}^{2}(t) \cdot \Phi_{v}(t \cdot e \cdot g) d t
$$

for suitable Schwartz functions $\Phi_{v}$ on $k_{v}^{2} \approx \ell_{v}=\ell \otimes_{k} k_{v}$. The leading local zeta integral factor gives the normalization $\varphi_{v}(1)=1$ at $g=1$. Then

$$
\int_{Z_{v} \backslash H_{v}} \varphi_{s, \chi, v}=\frac{1}{L_{k, v}\left(2 s, \chi_{v}^{2}\right)} \cdot \int_{k_{v}^{\times} \backslash \ell_{v}^{\times}}|\operatorname{det} h|_{v}^{s} \chi_{v}^{2}(\operatorname{det} h) \cdot \int_{k_{v}^{\times}}|t|_{v}^{2 s} \chi_{v}^{2}(t) \cdot \Phi_{v}(t \cdot e \cdot h) d t d h
$$

Let $N_{v}$ be the $k_{v}$-extension $\ell \otimes_{k} k_{v} \rightarrow k_{v}$ of the norm map $\ell \rightarrow k$. Since

$$
|\operatorname{det} h|_{k_{v}}=|N h|_{k_{v}}=|h|_{\ell_{v}}
$$

and since $\chi_{v}(\operatorname{det} h)=\chi_{v}(N h)$, the local factor of the period becomes

$$
\begin{gathered}
\frac{1}{L_{k, v}(2 s)} \cdot \int_{k_{v}^{\times} \backslash \ell_{v}^{\times}}|h|_{\ell_{v}}^{s} \chi_{v}(N h) \cdot \int_{k_{v}^{\times}}|t|_{\ell_{v}}^{s} \chi_{v}^{2}(t) \cdot \Phi_{v}(t \cdot e \cdot h) d t d h \\
=\frac{1}{L_{k, v}\left(2 s, \chi_{v}^{2}\right)} \cdot \int_{k_{v}^{\times} \backslash \ell_{v}^{\times}} \int_{k_{v}^{\times}}|t \cdot h|_{\ell_{v}}^{s} \chi_{v}(N(t h)) \cdot \Phi_{v}(t \cdot e \cdot h) d t d h \\
=\frac{1}{L_{k, v}(2 s)} \cdot \int_{\ell_{v}^{\times}}|h|_{\ell_{v}}^{s} \chi_{v}(N h) \cdot \Phi_{v}(e \cdot h) d t d h=\frac{1}{L_{k, v}\left(2 s, \chi_{v}^{2}\right)} \cdot L_{\ell, v}\left(s, \chi_{v} \circ N\right)
\end{gathered}
$$

where the local $L$-function $L_{\ell, v}$ is the product of the finitely-many local factors $L_{\ell, w}$ for places $w$ of $\ell$ lying over $v$.

Let $S$ be the collection of ramified primes for $\chi$ and primes ramified in $\ell / k$. Let $\omega$ be the quadratic character attached to $\ell / k$, with local characters $\omega_{v}$ : at finite primes $v$ splitting in $\ell / k$, the character is trivial. Let $q_{v}$ be the residue field cardinality at a finite place $v \notin S$, and $\varpi_{v}$ a local parameter.

At a finite place $v$ of $k$ splitting in $\ell / k$, we immediately have

$$
L_{\ell, v}\left(s, \chi_{v} \circ N\right)=\frac{1}{1-\chi_{v}\left(\varpi_{v}\right) q_{v}^{-s}} \cdot \frac{1}{1-\chi_{v}\left(\varpi_{v}\right) q_{v}^{-s}}=L_{k, v}\left(s, \chi_{v}\right)^{2}
$$

Since local $L$-functions of unramified local characters are determined by their values on local parameters, at a finite place $v$ of $k$ inert in $\ell / k$, in which case $N: \ell \otimes_{k} k_{v} \rightarrow k_{v}$ is the local field norm and $\varpi_{v}$ remains prime in $\ell \otimes_{k} k_{v}$, we similarly have

$$
\begin{aligned}
& L_{\ell, v}\left(s, \chi_{v} \circ N\right)=\frac{1}{1-\chi_{v}\left(N \varpi_{v}\right)\left(q_{v}^{2}\right)^{-s}}=\frac{1}{1-\chi_{v}\left(\varpi_{v}^{2}\right)\left(q_{v}^{2}\right)^{-s}} \\
&= \frac{1}{1-\chi_{v}\left(\varpi_{v}\right) q_{v}^{-s}} \cdot \frac{1}{1+\chi_{v}\left(\varpi_{v}\right) q_{v}^{-s}} \\
&=\frac{1}{1-\chi_{v}\left(\varpi_{v}\right) q_{v}^{-s}} \cdot \frac{1}{1-\chi_{v} \omega_{v}\left(\varpi_{v}\right) q_{v}^{-s}}=L_{k, v}\left(s, \chi_{v}\right) \cdot L_{k, v}\left(s, \chi_{v} \cdot \omega_{v}\right)
\end{aligned}
$$

Thus, the good-prime part of the global $L$-function is

$$
L_{\ell}^{S}(s, \chi \circ N)=\prod_{v \notin S} L_{\ell, v}\left(s, \chi_{v} \circ N\right)=\prod_{v \notin S} L_{k, v}\left(s, \chi_{v}\right) \cdot L_{k, v}\left(s, \chi_{v} \cdot \omega_{v}\right)
$$

Thus,

$$
\int_{Z_{\mathbb{A}} H_{k} \backslash H_{\mathbb{A}}} E_{s, \chi}=\frac{L_{k}^{S}(s, \chi) \cdot L_{k}^{S}(s, \chi \cdot \omega)}{L_{k}^{S}\left(2 s, \chi^{2}\right)} \times(\text { bad prime factors })
$$

as claimed.
[2.C.3] Remark: In this normalization, the unitary line is $\operatorname{Re}(s)=\frac{1}{2}$, and

$$
\int_{Z_{\mathbb{A}} H_{k} \backslash H_{\mathbb{A}}} E_{\frac{1}{2}+i t, \chi}=\frac{L^{S}\left(\frac{1}{2}+i t, \chi\right) \cdot L^{S}\left(\frac{1}{2}+i t, \chi \cdot \omega\right)}{L^{S}\left(1+2 i t, \chi^{2}\right)} \times(\text { bad prime factors })
$$

[2.C.4] Remark: The basic remaining issue about the finitely-many bad-prime local integrals is to be sure that we can choose local data $\varphi_{s, \chi, v}$ so that the local inegrals do not vanish identically. This can be accomplished, for example, by taking the bad-prime local functions $\varphi_{s, \chi, v}$ to be 0 off a very small neighborhood of the local points $P_{v}$ of the parabolic $P$.

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

1. Parabolic subgroups of $G L_{r}$
2. Groups $K_{v}=G L_{r}\left(\mathfrak{o}_{v}\right) \subset G_{v}=G L_{r}\left(k_{v}\right)$
3. Discrete subgroup $G L_{r}(k)$, reduction theory
4. Invariant differential operators and integral operators
5. Hecke operators and integral operators
6. Decomposition by central characters
7. Discrete decomposition of cuspforms
8. Pseudo-Eisenstein series
9. Cuspidal-data pseudo-Eisenstein series
10. Minimal-parabolic Eisenstein series
11. Cuspidal-data Eisenstein series
12. Continuation of minimal-parabolic Eisenstein series
13. Continuation of cuspidal-data Eisenstein series
14. Truncation and Maaß-Selberg relations
15. Minimal-parabolic decomposition
16. Cuspidal-data decomposition
17. Plancherel for pseudo-Eisenstein series
18. Automorphic spectral expansions

Appendix A: Bochner's lemma
Appendix B: Phragment-Lindelöf theorem
We keep most of the conventions and context of the previous chapter, except now $G$ is the group $G L_{r}$ of $r$-by- $r$ matrices. The novelties originate in the greater variety of parabolic subgroups in $G L_{r}$, the latter explicated in the first section. This variety increases the subtlety of the Gelfand condition defining the space of cuspforms, with pursuant proliferation of types of pseudo-Eisenstein series and Eisenstein series on $G L_{r}$ corresponding to the various parabolic subgroups. One new phenomenon is the necessary formation of pseudo-Eisenstein series and Eisenstein series using cuspforms on smaller groups $G L_{r^{\prime}}$.

To narrow somewhat the scope of complications, later in the chapter we mostly treat level one automorphic forms, that is, right $K_{\mathbb{A}}$-invariant ones, for $K_{\mathbb{A}}$ a maximal compact subgroup of $G L_{r}(\mathbb{A})$. This specializes to automorphic forms for $G L_{r}(\mathbb{Z})$ when the ground field is $\mathbb{Q}$.

### 3.1 Parabolic subgroups of $G L_{r}$

It is perhaps impossible to anticipate the significance of these subgroups. Nevertheless, they subsequently prove their importance. ${ }^{[25]}$ Let $G=G L_{r}(k)$ with an arbitrary field $k$, acting on $k^{r}$ by matrix multiplication. A flag $F$ in $k^{n}$ is a nested sequence of one or more non-zero $k$-subspaces (with proper containments) $V_{1} \subset \ldots \subset V_{\ell} \subset k^{r}$. The corresponding parabolic subgroup $P=P^{F}$ is the stabilizer of the flag $F$. The whole group $G$ stabilizes the improper flag $k^{r}$, so is a parabolic subgroup of itself. The proper parabolics are stabilizers of flags $V_{1} \subset \ldots \subset V_{\ell} \subset k^{r}$ with $\ell \geq 1$.

The maximal proper parabolic subgroups are stabilizers $P^{V \subset k^{r}}$ of flags consisting of single proper subspaces $V \subset k^{r}$. Every proper parabolic subgroup $P^{F}$ for flag $F=\left(V_{1} \subset \ldots \subset V_{\ell} \subset k^{r}\right)$ is the intersection of the maximal proper parabolics $P^{V_{i} \subset k^{r}}$. A minimal parabolic, stabilizing a maximal flag, is a Borel subgroup.

With $e_{1}, e_{2}, \ldots, e_{r}$ the standard basis for $k^{r}$, identify $k^{d}=k e_{1}+\ldots+k e_{d}$. By transitivity of $G$ on ordered bases of $k^{r}$, every orbit in the action of $G$ on flags has a unique representative among the standard flags, namely, for some ordered partition $d_{1}+d_{2}+\ldots+d_{\ell}=r$ with $0<d_{j} \in \mathbb{Z}$, the corresponding standard flag is

$$
F^{d_{1}, \ldots, d_{\ell}}=\left(k^{d_{1}} \subset k^{d_{1}+d_{2}} \subset k^{d_{1}+d_{2}+d_{3}} \subset \ldots \subset k^{d_{1}+\ldots+d_{\ell}}\right)
$$

[^5]$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

The stabilizer of $F^{d_{1}, \ldots, d_{\ell}}$ is the standard (proper) parabolic subgroup $P^{d_{1}, \ldots, d_{\ell}}$ of $G$, and is the intersection of the maximal proper parabolics containing it, namely

$$
P^{d_{1}, \ldots, d_{\ell}}=\bigcap_{1 \leq i \leq \ell-1} P^{\left(d_{1}+\ldots+d_{i}\right),\left(d_{i+1}+\ldots+d_{\ell}\right)}
$$

Two standard parabolics $P^{d_{1}, \ldots, d_{\ell}}$ and $P^{d_{1}^{\prime}, \ldots, d_{\ell^{\prime}}^{\prime}}$ are associate when $\ell=\ell^{\prime}$ and the lists $d_{1}, \ldots, d_{\ell}$ and $d_{1}^{\prime}, \ldots, d_{\ell}^{\prime}$ merely differ by being permutations of each other. A parabolic $P^{d_{1}, \ldots, d_{\ell}}$ is self-associate when $d_{i}=d_{j}$ for some $i \neq j$. These notions are important for discussion of constant terms of Eisenstein series, meromorphic continuations, and functional equations. The standard maximal proper parabolics are block-upper-triangular, in the sense

$$
P^{r^{\prime}, r-r^{\prime}}=\left\{\left(\begin{array}{ll}
a & b \\
0 & d
\end{array}\right): a \in G L_{r^{\prime}}, b=r^{\prime} \times\left(r-r^{\prime}\right), d \in G L_{r-r^{\prime}}\right\}
$$

That is, the diagonal blocks are $r^{\prime} \times r^{\prime}$ and $\left(r-r^{\prime}\right) \times\left(r-r^{\prime}\right)$, and the off-diagonal blocks are sized to fit. Next-to-maximal proper parabolics have the shape

$$
P^{r_{1}, r_{2}, r_{3}}=\left\{\left(\begin{array}{ccc}
m_{1} & * & * \\
0 & m_{2} & * \\
0 & 0 & m_{3}
\end{array}\right): m_{1} \in G L_{r_{1}}, m_{2} \in G L_{r_{2}}, m_{3} \in G L_{r_{3}}\right\} \quad\left(\text { for } r_{1}+r_{2}+r_{3}=r\right)
$$

with off-diagonal blocks to fit. The general standard proper parabolic $P^{d_{1}, \ldots, d_{\ell}}$ consists of block-uppertriangular matrices with diagonal blocks of sizes $d_{1} \times d_{1}, d_{2} \times d_{2}, \ldots, d_{\ell} \times d_{\ell}$. The standard Borel subgroup is the subgroup of upper triangular matrices.

The unipotent radical $N^{P}$ of a parabolic $P$ stabilizing a flag $F=\left(V_{1} \subset \ldots \subset V_{\ell} \subset k^{r}\right)$ is the subgroup that fixes the quotients $V_{\ell} / V_{\ell-1}$ pointwise. This characterization shows that $N^{P}$ is a normal subgroup of $P$. For the standard maximal parabolic $P=P^{r^{\prime}, r-r^{\prime}}$, the unipotent radical is

$$
N=N^{P}=N^{r^{\prime}, r-r^{\prime}}=\left\{\left(\begin{array}{cc}
1_{r^{\prime}} & b \\
0 & 1_{r-r^{\prime}}
\end{array}\right): b=r^{\prime} \times\left(r-r^{\prime}\right)\right\}
$$

Containment of parabolics reverses the containment of unipotent radicals: $P \subset Q$ implies $N^{P} \supset N^{Q}$. For example, for next-to-maximal standard proper parabolic $P=P^{r_{1}, r_{2}, r_{3}}$, the unipotent radical is

$$
N=N^{P}=N^{r_{1}, r_{2}, r_{3}}=\left\{\left(\begin{array}{ccc}
1_{r_{1}} & * & * \\
0 & 1_{r_{2}} & * \\
0 & 0 & 1_{r_{3}}
\end{array}\right)\right\}
$$

The standard Levi component (or standard Levi-Malcev component) $M=M^{P}=M^{d_{1}, \ldots, d_{\ell}}$ of the standard parabolic $P=P^{d_{1}, \ldots, d_{\ell}}$ is the subgroup of $P=P^{d_{1}, \ldots, d_{\ell}}$ with all the blocks above the diagonal 0 , namely,

$$
M=M^{P}=M^{d_{1}, \ldots, d_{\ell}}=\left\{\left(\begin{array}{ccccc}
m_{1} & 0 & 0 & \ldots & 0 \\
0 & m_{2} & 0 & \ldots & 0 \\
\vdots & 0 & \ddots & & \vdots \\
& & & \ddots & 0 \\
0 & \ldots & & 0 & m_{\ell}
\end{array}\right): m_{j} \in G L_{d_{j}}\right\}
$$

Unlike the unipotent radical, the standard Levi component is not normal in the standard parabolic. Nevertheless, we have the Levi-Malcev decomposition $P=N^{P} \cdot M^{P}$ for matrices with entries in any field. For the standard parabolics and standard Levi components, this is simply an expression of the behavior of matrix multiplication in block decompositions. For example,

$$
\left(\begin{array}{cc}
a & b  \tag{inblocks}\\
0 & d
\end{array}\right)=\left(\begin{array}{cc}
1_{r^{\prime}} & b d^{-1} \\
0 & 1_{r-r^{\prime}}
\end{array}\right)\left(\begin{array}{ll}
a & 0 \\
0 & d
\end{array}\right)
$$

The standard maximal split torus in $G$ is the subgroup of diagonal matrices, which is also the Levi component $M^{\text {min }}$ of the standard minimal parabolic $P^{\min }=P^{1,1, \ldots, 1}$. The standard Weyl group $W$ can be identified with permutation matrices in $G$, namely, matrices with exactly one non-zero entry in each row and column, and that entry is 1 . The Weyl group normalizes $M^{\text {min }}$. ${ }^{[26]}$ The simplest Bruhat decomposition is [3.1.1] Claim: With $P^{\min }$ the standard minimal parabolic and $N^{\min }$ its unipotent radical, we have a disjoint union

$$
G L_{r}(k)=\bigsqcup_{w \in W} P^{\min } w P^{\min }=\bigsqcup_{w \in W} P^{\min } w N^{\min }
$$

Proof: The second equality follows from the first, by the Levi decomposition: letting $P=P^{\min }$ and $N=N^{P}$ and $M=M^{P}$,

$$
P w P=P w(M N)=P \cdot w M m^{-1} \cdot w N=P M \cdot w N=P \cdot w N
$$

For the first equality, given $g \in G L_{r}(k)$, left multiplication by $N$ can add or subtract multiples of a row of $g$ to or from any higher row. Similarly, right multiplication by $N$ adds or subtracts multiples of a column of $g$ to or from any column farther to the right. Thus, letting $g_{i_{1}, 1}$ be the lowest non-zero entry in the first column (maximal index $i_{1}$ ), left multiplication by $N$ makes all other entries in the first column 0 . Right multiplication by $N$ makes all other entries in the $i_{1}^{t h}$ row 0 . Next, let $g_{i_{2}, 2}$ be the lowest (maximal index $i_{2}$ ) non-zero entry in the (new) second column. Without disturbing the effects of the previous step, all higher entries in the second column, and all entries in the $i_{2}^{t h}$ row to the right, can be made 0 by left and right action of $N$. An induction produces a monomial matrix, that is, one with a single non-zero entry in each row and column. Then left multiplication by $M$ normalizes all non-zero entries to 1 . Thus, $M N g N \in W$.

In fact, the positions of the lowest non-zero entries $g_{i_{j}, j}$ in each column are completely determined by this procedure, and there is no other way to reach a monomial matrix by left and right multiplication by $N$. This explains the disjointness of the decomposition.

Rather than set up notation for the general case, the induction is better illustrated by an example: writing * for unknown entries and $\times$ for non-zero entries, with suitable values in the elements of $N$, the first stage, using the non-zero 2,1 entry, is

$$
\left(\begin{array}{lll}
* & * & * \\
\times & * & * \\
0 & \times & *
\end{array}\right) \rightarrow\left(\begin{array}{lll}
1 & * & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)\left(\begin{array}{ccc}
* & * & * \\
\times & * & * \\
0 & \times & *
\end{array}\right)=\left(\begin{array}{ccc}
0 & * & * \\
\times & * & * \\
0 & \times & *
\end{array}\right) \rightarrow\left(\begin{array}{ccc}
0 & * & * \\
\times & * & * \\
0 & \times & *
\end{array}\right)\left(\begin{array}{lll}
1 & * & * \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right)=\left(\begin{array}{ccc}
0 & * & * \\
\times & 0 & 0 \\
0 & \times & *
\end{array}\right)
$$

The second stage, using the non-zero 3,2 entry, is
$\left(\begin{array}{lll}0 & * & * \\ \times & 0 & 0 \\ 0 & \times & *\end{array}\right) \rightarrow\left(\begin{array}{ccc}1 & 0 & * \\ 0 & 1 & 0 \\ 0 & 0 & 1\end{array}\right)\left(\begin{array}{ccc}0 & * & * \\ \times & 0 & 0 \\ 0 & \times & *\end{array}\right)=\left(\begin{array}{ccc}0 & 0 & * \\ \times & 0 & 0 \\ 0 & \times & *\end{array}\right) \rightarrow\left(\begin{array}{ccc}0 & 0 & * \\ \times & 0 & 0 \\ 0 & \times & *\end{array}\right)\left(\begin{array}{lll}1 & 0 & 0 \\ 0 & 1 & * \\ 0 & 0 & 1\end{array}\right)=\left(\begin{array}{ccc}0 & 0 & * \\ \times & 0 & 0 \\ 0 & \times & 0\end{array}\right)$
The upper-right entry must be invertible, since the original matrix is.
[3.1.2] Corollary: $G=\bigcup_{w \in W} P w Q$ for any standard parabolics $P, Q$.
[3.1.3] Remark: Letting $W^{P}=W \cap P$ and $W^{Q}=W \cap Q$, we have a disjoint union $G=\bigsqcup_{w \in W^{P} \backslash W / W^{Q}} P w Q$. However, except for minimal or maximal-proper parabolics, this requires a subtler proof. Our subsequent examples will not need this precise form of the general case, although it would become relevant in other examples.

Below, to distinguish matrices with entries in $k$ from entries in $k_{v}$, we will write $P_{k}, M_{k}$, and $N_{k}$ in place of the unadorned $P, M, N$ above.

[^6]
### 3.2 Groups $K_{v}=G L_{r}\left(\mathfrak{o}_{v}\right) \subset G_{v}=G L_{r}\left(k_{v}\right)$

Now $k$ is again a number field with integers $\mathfrak{o}$, completions $k_{v}$, and local rings of integers $\mathfrak{o}_{v}$ at nonarchimedean places. Let $G_{v}=G L_{r}\left(k_{v}\right)$, and let $Z_{v}$ be the center of $G_{v}$. At non-archimedean places $v$, let $K_{v}=G L_{r}\left(\mathfrak{o}_{v}\right)$. At real $v$, let $K_{v}$ be the standard orthogonal group $O_{n}(\mathbb{R})=\left\{g \in G L_{n}(\mathbb{R}): g^{\top} g=1_{r}\right\}$, and at complex $v$ let $K_{v}$ be the standard unitary group $U_{n}=\left\{g \in G L_{n}(\mathbb{C}): g^{*} g=1_{r}\right\}$.

Temporarily, let $\ell$ be the number of non-isomorphic archimedean completions of $k$, thus not counting a complex completion and its conjugate as 2 , but just 1 . That is, $[k: \mathbb{Q}]=\ell_{1}+2 \ell_{2}$ where $\ell_{1}$ is the number of real completions, and $\ell_{2}$ the number of complex, and $\ell=\ell_{1}+\ell_{2}$. Let $Z^{+}$be the positive real scalar matrices diagonally imbedded across all archimedean $v$, by the map

$$
t \longrightarrow\left(\ldots, t^{1 / \ell}, \ldots\right) \quad(\text { for } t>0)
$$

This map $\delta$ gives a section of the idele norm map $|t|=\prod_{v}\left|t_{v}\right|_{v}$, in that $|\delta(t)|=t$.
The group $P_{v}$ of $v$-adic points of a standard parabolic $P=P^{F}=P^{d_{1}, \ldots, d_{\ell}}$ is the stabilizer in $G_{v}$ with the same shape as the $k$-rational points $P_{k}$ in the previous section, but with entries in $k_{v}$ rather than $k$. That is, the $v$-adic version of the flag $F=\left(k^{d_{1}} \subset k^{d_{1}+d_{2}} \subset \ldots \subset k^{r}\right)$ is the natural $\left.F_{v}=k_{v}^{d_{1}} \subset k_{v}^{d_{1}+d_{2}} \subset \ldots \subset k_{v}^{r}\right)$, and $P_{v}=P_{v}^{d_{1}, \ldots, d_{\ell}}$ is its stabilizer. [27] Similarly, $N_{v}=N_{v}^{P}$ is the $v$-adic points of the unipotent radical $N=N^{P}$ of $P$, and $M_{v}=M_{v}^{P}$ is the $v$-adic points of the standard Levi component of a standard parabolic $P$. That is, again, the shapes of the matrices are the same as in the previous section, but with entries in $k_{v}$ rather than $k$.

Iwasawa decompositions are analogous to the previous chapter's, with proofs merely iterations of the arguments there:
[3.2.1] Claim: $G_{v}=P_{v} \cdot K_{v}$ for standard minimal parabolic $P$.
Proof: For archimedean $v$, the right action of $K_{v}$ rotates the rows of given $g \in G_{v}$. The bottom row can be rotated to be of the form $(0,0, \ldots, 0,0, *)$. Without disturbing this effect, the second-to-bottom row can be rotated to be of the form $(0,0, \ldots, 0, *, *)$. Continuing with higher rows puts the result in $P_{v}^{\mathrm{min}}$.

For non-archimedean $v$, right multiplication by $K_{v}$ can subtract local-integer multiples of the largest entry in the bottom row from all others, to put the bottom row into the form $(0,0, \ldots, 0, *, 0, \ldots 0)$ with a non-zero entry at just one position. Then a permutation matrix (in $K_{v}$ ) can move the non-zero entry to the far right. Without disturbing this effect, the first $r-1$ entries of the second-to-bottom row can be dealt with similarly, putting it into the form $(0,0, \ldots, 0, *, *)$. Because the determinant is non-zero, the second-to-right entry in the new second-to-bottom row is non-zero. Continuing to modify higher rows puts the result in $P_{v}$.

As in [1.2] and [2.1], Cartan decompositions follow from the spectral theorem for symmetric or hermitian operators at archimedean places, and from the structure theorem for finitely-generated modules at finite places:
[3.2.2] Claim: $G_{v}=K_{v} M_{v} K_{v}$ with $M=M^{\min }$ the standard Levi component of the minimal parabolic.
Proof: For archimedean $v$, letting $g \rightarrow g^{*}$ be either transpose for $k_{v} \approx \mathbb{R}$, or conjugate-transpose for $k_{v} \approx \mathbb{C}$, the matrix $g g^{*}$ is positive-definite symmetric or hermitian-symmetric. The spectral theorem for such operators gives an orthogonal or unitary matrix $k$ such that $k\left(g g^{*}\right) k^{*}=\delta$ is diagonal with strictly positive real diagonal entries. Let $\sqrt{\delta}$ be the positive-definite diagonal square root of $\delta$. Then $h=k^{*} \sqrt{\delta} k$ is a positive-definite hermitian/symmetric square root of $g g^{*}$. Then

$$
\left(h^{-1} \cdot g\right) \cdot\left(h^{-1} \cdot g\right)^{*}=h^{-1} \cdot g g^{*} \cdot h^{-1}=h^{-1} \cdot h^{2} \cdot h^{-1}=1_{r}
$$

Inverting $k^{*} k=1_{r}$ gives $k^{-1}\left(k^{*}\right)^{-1}=1_{r}$, and then $1_{r}=k k^{*}$, so the latter condition also defines $K_{v}$. That is, $h^{-1} g \in K_{v}$, and $g \in k^{*} \sqrt{\delta} k \cdot K_{v} \subset K_{v} \sqrt{\delta} K_{v}$.

For non-archimedean $v$, multiply through by scalar matrix $c \cdot 1_{r}$ so that all entries of $c g$ are in $\mathfrak{o}_{v}$, though of course the determinant may fail to be a local unit. The rows of $R_{1}, \ldots, R_{r}$ of $g \in G_{v}$ are linearly
[27] The $k_{v}$-vectorspace $k_{v}^{r}$ has many subspaces and flags not obtained by extending scalars from subspaces and flags in $k^{r}$, but these will play no role here.
independent, and generate a free $\mathfrak{o}_{v}$-module $F$ of rank $r$ inside $\mathfrak{o}_{v}^{r}$. Observe that $K_{v}=G L_{r}\left(\mathfrak{o}_{v}\right)$ is the stabilizer of $\mathfrak{o}_{v}^{r}$. Since $\mathfrak{o}_{v}$ has a unique non-zero prime ideal, the applicable form of the structure theorem for finitely-generated modules over principal ideal domains is even simpler, and produces an $\mathfrak{o}_{v}$-basis $f_{1}, \ldots, f_{r}$ of $\mathfrak{o}_{v}^{r}$ and elementary divisors $d_{1}|\ldots| d_{r}$ such that $\left\{d_{i} f_{i}\right\}$ is an $\mathfrak{o}_{v}$-basis of $F$. Let $k_{1} \in K_{v}$ be the change-of-basis element such that the $j^{\text {th }}$ row of $k_{1} \cdot g$ is $d_{j} f_{j}$. Let $\left\{e_{i}\right\}$ be the standard basis of $\mathfrak{o}_{v}^{r}$, and $k_{2} \in K_{v}$ such that $f_{j} \cdot g=e_{j}$ for all $j$. Then $\delta=k_{1} g k_{2}$ is diagonal with entries $d_{1}, \ldots, d_{r}$. We can undo the initial multiplication to get $g \in k_{1}^{-1} c^{-1} \delta k_{2}^{-1} \in K_{v} M_{v} K_{v}$.

As in $G L_{2}$, unlike the archimedean situation, for non-archimedean $v$ the compact $K_{v}$ has substantial intersections with $N_{v}^{P}$ and $M_{v}^{P}$ for every standard parabolic $P$. As for $G L_{2}$, unipotent radicals $N_{v}^{P}$ of standard parabolics $P$ are ascending unions of compact, open subgroups:

$$
N_{v}^{P}=\bigcup_{\ell \geq 0}\left\{n \in N_{v}^{P}: n=1_{r} \bmod \varpi^{\ell}\right\}
$$

Again, unlike the archimedean situation, $K_{v}$ has a basis at 1 consisting of compact, open subgroups, namely, the (local) principal congruence subgroups

$$
K_{v, \ell}=\left\{g \in K_{v}=G L_{r}\left(\mathfrak{o}_{v}\right): g=1_{r} \bmod \varpi^{\ell}\right\}
$$

The corresponding adele group is $G_{\mathbb{A}}=G L_{r}(\mathbb{A})$, meaning $r$-by- $r$ matrices with entries in $\mathbb{A}$, with determinant in the ideles $\mathbb{J}$. This group is also an ascending union (colimit) of products

$$
G_{S}=\prod_{v \in S} G_{v} \times \prod_{v \notin S} K_{v} \quad \text { (for } S \text { a finite set of places } v, \text { including archimedean places) }
$$

ordering the finite sets $S$ by containment. Similarly, $P_{\mathbb{A}}, M_{\mathbb{A}}^{P}, N_{\mathbb{A}}^{P}$, and $Z_{\mathbb{A}}$ are the adelic forms of those groups. Let $K_{\mathbb{A}}=\prod_{v} K_{v} \subset G_{\mathbb{A}}$. With the usual one-sided inverse to $\delta:(0, \infty) \rightarrow \mathbb{J}$ the idele norm $|\cdot|: \mathbb{J} \rightarrow(0, \infty)$,

$$
Z_{\mathbb{A}} / Z^{+} Z_{k} \approx \mathbb{J} / \delta(0, \infty) \cdot k^{\times} \approx \mathbb{J}^{1} / k^{\times}=\text {compact }
$$

where $Z_{k}$ is the invertible scalar matrices with entries in $k$, with compactness demonstrated in [2.15].

### 3.3 Discrete subgroup $G_{k}=G L_{r}(k)$, reduction theory

As expected, $G_{k}=G L_{r}(k)$, and $P_{k}, M_{k}^{P}, N_{k}^{P}$ are the corresponding groups with entries in $k$. Proof of the discreteness of $G_{k}$ in $G_{\mathbb{A}}$ is essentially identical to that for $G L_{2}$ in [2.2], and we will not repeat it. Let

$$
G^{1}=\left\{g \in G_{\mathbb{A}}:|\operatorname{det} g|=1\right\}
$$

and $G_{\mathbb{A}}=Z^{+} \times G^{1}$. The product formula $\prod_{v \leq \infty}|t|_{v}=1$ for $t \in k^{\times}$shows that $G_{k} \subset G^{1}$. In particular, $G_{k}$ is still discrete in $Z^{+} \backslash G_{\mathbb{A}} \approx G^{1}$. As in the simpler cases [1.5] and [2.2], reduction theory should show that the quotient $G_{k} \backslash G_{\mathbb{A}}$ is covered by a suitable notion of Siegel set, and that these Siegel sets interact well with each other. We prove that a single Siegel set covers the quotient, but omit the discussion of their interaction.

The notion of (standard) Siegel set becomes somewhat more complicated. The notion of a single numerical height as in [1.5] and [2.2] is replaced by a family: the standard positive simple roots ${ }^{[28]}$ are characters on $M=M^{P}=M^{\min }$ with $P=P^{1, \ldots, 1}=P^{\min }$ the standard minimal parabolic:

$$
\alpha_{i}\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right)=\frac{m_{i}}{m_{i+1}} \quad(\text { for } 1 \leq i<r)
$$

[^7]$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

These simple roots make sense on $M_{k}$ or $M_{v}$ or $M_{\mathbb{A}}$, taking values in $k^{\times}, k_{v}^{\times}$, and $\mathbb{J}$. A standard Siegel set adapted to or aligned with $P=P^{\mathrm{min}}$ is a set of the form

$$
\mathfrak{S}^{P}=\mathfrak{S}_{t, C}^{P}=\left\{g=n m k: n \in C, m \in M_{\mathbb{A}}, k \in K_{\mathbb{A}}, \text { and }\left|\alpha_{i}(m)\right| \geq t \text { for } 1 \leq i<r\right\}
$$

with idele norm $|\cdot|$, for $0<t \in \mathbb{R}$ and compact $C \subset N_{\mathbb{A}}^{P}$. Let $N^{\min }=N^{P}$ with $P=P^{\min }$.
[3.3.1] Theorem: For given $k$, there is $t>0$ and a compact subgroup $C \subset N_{\mathbb{A}}^{\min }$ such that $G_{k} \cdot \mathfrak{S}_{t, C}^{P}=G_{\mathbb{A}}$. That is, $G_{k} \backslash G_{\mathbb{A}}$ is covered by a single, sufficiently large Siegel set.
Proof: We need a notion of height on $\mathbb{A}^{r}$ as in [2.2] for $r=2$. Let $G_{\mathbb{A}}$ act on the right on $\mathbb{A}^{r}$ by matrix multiplication. For real primes $v$ of $k$ the local height function $h_{v}$ on $x=\left(x_{1}, \ldots, x_{r}\right) \in k_{v}^{r}$ is $h_{v}(x)=\sqrt{x_{1}^{2}+\ldots+x_{n}^{2}}$. For complex $v$, take $h_{v}(x)=\left|x_{1}\right|_{\mathbb{C}}+\ldots+\left|x_{r}\right|_{\mathbb{C}}$ with $|z|_{\mathbb{C}}=\left|N_{\mathbb{R}}^{\mathbb{C}} z\right|_{\mathbb{R}}$ to not disturb the product formula. For non-archimedean $v, h_{v}(x)=\sup _{i}\left|x_{i}\right|_{v}$.

A vector $x \in \mathbb{A}^{r}$ is primitive when it is of the form $x_{o} g$ for $g \in G_{\mathbb{A}}$ and $x_{o} \in k^{r}$. For $x=\left(x_{1}, \ldots, x_{r}\right) \in k^{r}$, at almost all non-archimedean primes $v$ the $x_{i}$ 's are in $\mathfrak{o}_{v}$ and have local greatest common divisor 1. Elements of the adele group $g \in G_{\mathbb{A}}$ are in $K_{v}$ almost everywhere, so this is not changed by multiplication by $g$. That is, a primitive vector $x$ has the property that at almost all $v$ the components of $x$ are locally integral and have local greatest common divisor 1. For primitive $x$ the global height is $h(x)=\prod_{v} h_{v}\left(x_{v}\right)$. Since $x$ is primitive, at almost all finite primes the local height is 1 , so this product has only finitely many non- 1 factors. The proof of the following is mostly identical to the $r=2$ case [2.2.2]:
[3.3.2] Claim: For fixed $g \in G L_{r}(\mathbb{A})$ and for fixed $c>0$,

$$
\operatorname{card}\left(k^{\times} \backslash\left\{x \in k^{r}-\{0\}: h(x \cdot g)<c\right\}\right)<\infty
$$

For compact $C \subset G_{\mathbb{A}}$ there are positive implied constants such that

$$
h(x)<_{C} h(x \cdot g)<_{C} h(x) \quad \text { (for all } g \in C, \text { for all primitive } x \text { ) }
$$

Proof: Fix $g \in G_{\mathbb{A}}$. Since $K=K_{\mathbb{A}}=\prod_{v} K_{v}$ preserves heights, via Iwasawa decompositions locally everywhere, we may suppose that $g$ is in the group $P_{\mathbb{A}}$ of upper triangular matrices in $G_{\mathbb{A}}$. Let $g_{i j}$ be the $(i, j)^{t h}$ entry of $g$. Choose representatives $x=\left(x_{1}, \ldots, x_{r}\right)$ for non-zero vectors in $k^{r}$ modulo $k^{\times}$such that, with $\mu$ the first index with $x_{\mu} \neq 0, x_{\mu}=1$. That is, $x$ is of the form $x=\left(0, \ldots, 0,1, x_{\mu+1}, \ldots, x_{r}\right)$.

The more easily written-out case $r=2$ of the first assertion was treated in [2.2.2]. For $x \in k^{r}-\{0\}$ such that $h(x g)<c$, let $\mu-1$ be the least index such that $x_{\mu} \neq 0$. Adjust by $k^{\times}$such that $x_{\mu}=1$. For each $v$, from $h(x g)<c$,

$$
\left|g_{\mu-1, \mu}+x_{\mu} g_{\mu, \mu}\right|_{v} \prod_{w \neq v}\left|g_{\mu-1, \mu-1}\right|_{w} \leq h(g x)<c
$$

For almost all $v$ we have $\left|g_{\mu-1, \mu-1}\right|_{v}=1$, so there is a uniform $c^{\prime}$ such that

$$
\left|g_{\mu-1, \mu}+x_{\mu} g_{\mu, \mu}\right|_{v}<c^{\prime} \quad(\text { for all } v)
$$

For almost all $v$ the residue field cardinality $q_{v}$ is strictly greater than $c^{\prime}$, so for almost all $v$

$$
\left|g_{\mu-1, \mu}+x_{\mu} g_{\mu, \mu}\right|_{v} \leq 1
$$

Therefore, $g_{\mu-1, \mu}+x_{\mu} g_{\mu, \mu}$ lies in a compact subset $C$ of $\mathbb{A}$. Since $k$ is discrete, the collection of $x_{\mu}$ is finite.
Continuing, there are only finitely many choices for the other entries of $x$. Inductively, suppose $x_{i}=0$ for $i<\mu-1$, and $x_{\mu}, \ldots, x_{\nu-1}$ fixed, and show that $x_{\nu}$ has only finitely many possibilities. Looking at the $\nu^{t h}$ component $(x g)_{\nu}$ of $x g$,

$$
\left|g_{\mu-1, \nu}+x_{\mu} g_{\mu, \nu}+\ldots+x_{\nu-1} g_{\nu-1, \nu}+x_{\nu} g_{\nu, \nu}\right|_{v} \prod_{w \neq v}\left|g_{\mu-1, \mu-1}\right|_{w} \leq h(x g) \leq c
$$

For almost all places $v$ we have $\left|g_{\mu-1, \mu-1}\right|_{w}=1$, so there is a uniform $c^{\prime}$ such that for all $v$

$$
\left|(x g)_{\nu}\right|_{v}=\left|g_{\mu-1, \nu}+x_{\mu} g_{\mu, \nu}+\ldots+x_{\nu-1} g_{\nu-1, \nu}+x_{\nu} g_{\nu, \nu}\right|_{v}<c^{\prime}
$$

For almost all $v$ the residue field cardinality $q_{v}$ is strictly greater than $c^{\prime}$, so for almost all $v$

$$
\left|g_{\mu-1, \nu}+x_{\mu} g_{\mu, \nu}+\ldots+x_{\nu-1} g_{\nu-1, \nu}+x_{\nu} g_{\nu, \nu}\right|_{v} \leq 1
$$

Therefore,

$$
g_{\mu-1, \nu}+x_{\mu} g_{\mu, \nu}+\ldots+x_{\nu-1} g_{\nu-1, \nu}+x_{\nu} g_{\nu, \nu}
$$

lies in the intersection of a compact subset $C$ of $\mathbb{A}$ with a closed discrete set, so lies in a finite set. Thus, the number of possibilities for $x_{\nu}$ is finite. Induction gives the first assertion of the claim.

For the second assertion of the claim, let $E$ be a compact subset of $G_{\mathbb{A}}$, and let $K=\prod_{v} K_{v}$. Then $K \cdot E \cdot K$ is compact, being the continuous image of a compact set. So without loss of generality $E$ is left and right $K$-stable. By Cartan decompositions the compact set $E$ of $G_{\mathbb{A}}$ is contained in a set $K C K$ where $C \subset M_{\mathbb{A}}$ is compact. Let $g=\theta_{1} m \theta_{2}$ with $\theta_{i} \in K, m \in C$, and $x$ a primitive vector. By the $K$-invariance of the height,

$$
\frac{h(x g)}{h(x)}=\frac{h\left(x \theta_{1} m \theta_{2}\right)}{h(x)}=\frac{h\left(x \theta_{1} m\right)}{\theta(x)}=\frac{h\left(\left(x \theta_{1}\right) m\right)}{h((x \theta))}
$$

Thus, the set of ratios $h(x g) / h(x)$ for $g$ in a compact set and $x$ ranging over primitive vectors is exactly the set of values $h(x m) / h(x)$ where $m$ ranges over a compact set and $x$ varies over primitives. With diagonal entries $m_{i}$ of $m$, by compactness of $C$,

$$
0<\inf _{m \in C} \inf _{i}\left|m_{i}\right| \leq \frac{h(x m)}{h(x)} \leq \sup _{m \in C} \sup _{i}\left|m_{i}\right|<\infty \quad \text { (for all primitive } x \text { ) }
$$

This proves the second assertion of the claim.
Analogous to [2.2] for $r=2$, we could put $\eta(g)=|\operatorname{det} g| \cdot h\left(e_{r} \cdot g\right)^{-r}$, where $\left\{e_{i}\right\}$ is the standard basis for $k^{r}$. The parabolic $Q=P^{r-1,1}$ is the stabilizer of the line $k \cdot e_{r}$. This modification makes $\eta$ invariant under $Z_{\mathbb{A}}$, as well as left $Q_{k}$-invariant and right $K_{\mathbb{A}}$-invariant.
[3.3.3] Corollary: (of claim) Given $g \in G_{\mathbb{A}}$, there are finitely-many $\gamma \in Q_{k} \backslash G_{k}$ such that $\eta(\gamma \cdot g)>\eta(g)$. Thus, the supremum $\sup _{\gamma} \eta(\gamma \cdot g)$ is attained, and is finite.
Proof: There is a natural bijection $Q_{k} \backslash G_{k} \longleftrightarrow k^{\times} \backslash\left(k^{r}-\{0\}\right)$ mapping a matrix to its bottom row. The claim shows that there are finitely-many $x \in k^{\times} \backslash\left(k^{r}-\{0\}\right)$ such that $h(x g)<c$, that is, such that $h(x g)^{-1}>c^{-1}$. Since $|\operatorname{det} g|$ is $G_{k}$-invariant, the bijection gives the assertion.

Now we prove the theorem by induction on $r$. Given $g \in G_{\mathbb{A}}$, by the corollary there is $x \in k^{r}-\{0\}$ such that $h(x g)>0$ is minimal among values $h\left(x^{\prime} g\right)$ with $x^{\prime} \in k^{r}-\{0\}$. Take $\gamma_{o} \in G_{k}$ so that $e_{r} \gamma_{o}=x$, so $h(x g)=h\left(e_{r} \gamma_{o} g\right)$ is minimal, and $\eta\left(\gamma_{o} g\right)$ is maximal among all values $\eta\left(\gamma \cdot \gamma_{o} g\right)$ for $\gamma \in G_{k}$. By Iwasawa, there is $\theta \in K$ such that $q=\gamma_{o} g \theta \in Q_{\mathbb{A}}$. Then $h(\gamma g \theta)=\left|q_{r r}\right|$ where $q_{i j}$ is the $i j^{t h}$ entry of $q$, and $\eta(q)$ is maximal among all values $\eta(\gamma \cdot q)$ for $\gamma \in G_{k}$. Let $H \subset M^{Q}$ be the subgroup of $G_{\mathbb{A}}$ fixing $e_{r}$ and stabilizing the subspace spanned by $e_{1}, \ldots, e_{r-1}$, so $H \approx G L_{r-1}(\mathbb{A})$. By induction on $r$, beginning at $r=2$ in [2.2], by acting on $q=\gamma g \theta$ on the left by $H_{k}$ and on the right by $H_{\mathbb{A}} \cap K_{\mathbb{A}}$, we can suppose that $q \in P_{\mathbb{A}}^{\min }$ and $\left|q_{i i} / q_{i+1, i+1}\right| \geq t$ for $i<r-1$, without altering $\eta(q)$. The induction step reduces to the case $r=2$. The extremal property $h\left(e_{r} q\right) \leq h\left(x^{\prime} \cdot q\right)$ for all $x^{\prime} \in k^{r}-\{0\}$ certainly implies $h\left(e_{r} q\right) \leq h\left(x^{\prime} \cdot q\right)$ with $x^{\prime}$ ranging over the smaller set of vectors of the form $x^{\prime}=\left(0 \ldots 0 x_{r-1} x_{r}\right)$. Thus, the lower right 2-by- 2 block $q^{\prime}$ of $q$ is reduced as an element of $G L_{2}(\mathbb{A})$. This reduces to the $r=2$ case treated in $[2.2 .7]$, giving $\left|q_{r-1, r-1}\right| /\left|q_{r r}\right| \geq t$ for sufficiently small $t$ and proving the theorem.

### 3.4 Invariant differential operators and integral operators

For archimedean $G_{v}$, for some purposes, such as meromorphic continuation of Eisenstein series [11.5], [11.10], [11.12], the Casimir operator or Laplacian as in [4.2], [4.4] suffices. ${ }^{[29]}$ Beyond that, the tractability of integral operators, as in the rewriting of non-archimedean Hecke operators as such, suggests using integral operators at archimedean places, as well, especially in light of the commutativity result [3.4.3].

As usual, for a continuous action $G_{v} \times V \rightarrow V$ on a quasi-complete, locally convex topological vectorspace $V$, the corresponding integral operators are

$$
\varphi \cdot f=\int_{G_{v}} \varphi(g) g \cdot f d g \quad\left(\text { for } \varphi \in C_{c}^{o}\left(G_{v}\right) \text { and } f \in V\right)
$$

The integrand is a continuous, compactly-supported $V$-valued function, so has a Gelfand-Pettis integral [14.1]. Thus, for $f \in V=L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, with $G_{v}$ acting by right translation, pointwise

$$
(\varphi \cdot f)(x)=\int_{G_{v}} \varphi(g)(g \cdot f)(x) d g=\int_{G_{v}} \varphi(g) f(x g) d g \quad\left(\text { for } \varphi \in C_{c}^{o}\left(G_{v}\right) \text { and } f \in V\right)
$$

In fact, for general reasons [6.1] the right-translation action $G_{v} \times L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous, so the integral converges as an $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathrm{A}}\right)$-valued integral, obviating concern about pointwise values. The composition of two such operators is the operator attached to the convolution: for $\varphi, \psi \in C_{c}^{o}\left(G_{\mathbb{A}}\right)$, by the same computation as in [2.4] and [3.4],

$$
\varphi \cdot(\psi \cdot f)=\int_{G_{v}} \varphi(g) g \cdot\left(\int_{G_{v}} \psi(h) h \cdot f d h\right) d g=\int_{G_{v}} \int_{G_{v}} \varphi(g) \psi(h)(g h \cdot f) d h d g
$$

because the operation of $\varphi$ moves inside the Gelfand-Pettis integral. Replacing $h$ by $g^{-1} h$ gives

$$
\int_{G_{v}} \int_{G_{v}} \varphi(g) \psi\left(g^{-1} h\right) h \cdot f d h d g=\int_{G_{v}}\left(\int_{G_{v}} \varphi(g) \psi\left(g^{-1} h\right) d g\right) h \cdot f d h
$$

by changing the order of integration.
[3.4.1] Lemma: The adjoint to the action of $\varphi \in C_{c}^{o}\left(G_{v}\right)$ on $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is the action of $\check{\varphi} \in C_{c}^{o}\left(G_{v}\right)$, where $\check{\varphi}(g)=\overline{\varphi\left(g^{-1}\right)}$. (Proof identical to [2.4.1].) //I
For simplicity of discussion, we restrict attention to functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ right $K_{v}$-invariant for archimedean $v$. In that situation, for archimedean $v$, the integral operators given by left-and-right $K_{v^{-}}$ invariant $\varphi \in C_{c}^{o}\left(G_{v}\right)$, also denoted $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$, act on right $K_{v}$-invariant functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.
[3.4.2] Claim: The action of integral operators attached to $\varphi \in C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ stabilizes $K_{v}$-invariant vectors $f$ in any continuous group action $G_{v} \times V \rightarrow V$ for quasi-complete, locally convex $V$. (Proof identical to [2.4.3].)

Invoking Gelfand's trick [2.4.5],
[3.4.3] Claim: The action of integral operators attached to $\varphi \in C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ with convolution is commutative, for both non-archimedean and archimedean $v$.
Proof: To apply [2.4.5], we need an involutive anti-automorphism $\sigma$ of $G_{v}$, that is, $g \rightarrow g^{\sigma}$ such that $(g h)^{\sigma}=h^{\sigma} g^{\sigma}$ and $\left(g^{\sigma}\right)^{\sigma}=g$, stabilizing $K_{v}$ and acting trivially on representatives for double cosets $K_{v} \backslash G_{v} / K_{v}$. Use the Cartan decomposition [3.2] $G_{v}=K_{v} M_{v} K_{v}$ and use transpose $g^{\sigma}=g^{\top}$. Transpose stabilizes $K_{v}$, and acts trivially on $M_{v}$.

As in other cases, the algebra of integral operators attached to $\varphi \in C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ is stable under adjoints. Thus, it is plausible to ask for simultaneous eigenvectors for this commutative algebra of integral operators.
[29] For $S L_{r}(\mathbb{R})$ or $S L_{r}(\mathbb{C})$, with $r \geq 3$, the center of the universal enveloping algebra $U \mathfrak{g}$ of the corresponding algebra $\mathfrak{g}$ is generated by $r-1$ commuting operators.

### 3.5 Hecke operators and integral operators

For non-archimedean $G_{v}$, for any continuous action $G_{v} \times V \rightarrow V$ on a quasi-complete, locally convex topological vectorspace $V$, the corresponding integral operators are

$$
\varphi \cdot f=\int_{G_{v}} \varphi(g) g \cdot f d g \quad\left(\text { for } \varphi \in C_{c}^{o}\left(G_{v}\right) \text { and } f \in V\right)
$$

The integrand is a continuous, compactly-supported $V$-valued function, so has a Gelfand-Pettis integral [14.1]. Thus, for $f \in V=L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, with $G_{v}$ acting by right translation,

$$
(\varphi \cdot f)(x)=\int_{G_{v}} \varphi(g)(g \cdot f)(x) d g=\int_{G_{v}} \varphi(g) f(x g) d g \quad\left(\text { for } \varphi \in C_{c}^{o}\left(G_{v}\right) \text { and } f \in V\right)
$$

and for general reasons [6.1] the right-translation action $G_{v} \times L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous, so the integral converges as an $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$-valued integral. The composition of two such operators is the operator attached to the convolution: as in [2.4] and [3.5], for $\varphi, \psi \in C_{c}^{o}\left(G_{\mathbb{A}}\right)$,

$$
\varphi \cdot(\psi \cdot f)=\int_{G_{v}}\left(\int_{G_{v}} \varphi(g) \psi\left(g^{-1} h\right) d g\right) h \cdot f d h
$$

[3.5.1] Lemma: The adjoint to the action of $\varphi \in C_{c}^{o}\left(G_{v}\right)$ on $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is the action of $\check{\varphi} \in C_{c}^{o}\left(G_{v}\right)$, where $\check{\varphi}(g)=\overline{\varphi\left(g^{-1}\right)}$. (Proof identical to [2.4.1].)

It is reasonable to restrict attention to functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ right $K_{v}$-invariant for all $v$. But it is also reasonable to relax this condition to requiring right $K_{v}$-invariance almost everywhere, that is, at all but finitely-many places. A variant of $K_{\mathbb{A}}$-invariance, to cope with the finitely-many places where right $K_{v}$-invariance is not required, is $K_{\mathbb{A}}$-finiteness of a function $f$ on $G_{\mathbb{A}}$ or $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ or other quotients of $G_{\mathbb{A}}$, namely, the requirement that the vectorspace of functions spanned by $\left\{x \rightarrow f(x h): h \in K_{\mathbb{A}}\right\}$ is finite-dimensional. At the extreme of $K_{\mathbb{A}}$-invariant $f$, this space is one-dimensional.
[3.5.2] Lemma: For $v$ non-archimedean, $K_{v}$-finiteness is equivalent to invariance under some finite-index subgroup $K^{\prime} \subset K_{v}$. (Proof identical to [2.4.3].)

Unsurprisingly, it turns out that $K_{\mathbb{A}^{-}}$finite functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ are better behaved than arbitrary functions. Of course, most $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ are not $K_{\mathbb{A}}$-finite.

For non-archimedean $v$, the spherical Hecke operators for $G_{v}$ are the integral operators given by left-andright $K_{v}$-invariant $\varphi \in C_{c}^{o}\left(G_{v}\right)$, also denoted $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$. Since $K_{v}$ is open, such functions are locally constant: given $x \in G_{v}, \varphi(x h)=\varphi(x)$ for all $h \in K_{v}$, but $x K_{v}$ is a neighborhood of $x$. Then the compact support implies that such $\varphi$ takes only finitely-many distinct values. Thus, the associated integral operator is really a finite sum. Nevertheless, expression as integral operators explains the behavior well.
[3.5.3] Claim: The action of spherical Hecke operators attached to $\varphi \in G_{v}$ stabilizes $K_{v}$-invariant vectors $f$ in any continuous group action $G_{v} \times V \rightarrow V$ for quasi-complete, locally convex $V$. (Proof identical to [2.4.3].)
[3.5.4] Claim: For non-archimedean $v$, the spherical Hecke algebra $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ with convolution is commutative. (Again, Gelfand's trick [2.4.5].)

It is easy to see that the spherical Hecke algebra is stable under adjoints. Thus, it is plausible to ask for simultaneous eigenvectors for the spherical Hecke algebra. That is, for $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, we might try to require that $f$ be a spherical Hecke eigenfunction at almost all non-archimedean $v$, in addition to conditions at archimedean places. However, it bears repeating that, in infinite-dimensional Hilbert spaces, there is no general promise of existence of such simultaneous eigenvectors.

### 3.6 Decomposition by central characters

While $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ has finite invariant volume, $G_{k} \backslash G_{\mathbb{A}}$ does not. The further quotient $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ certainly has finite invariant volume. Functions on $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ are automorphic forms (or automorphic functions) with trivial central character, since they are invariant under the center $Z_{\mathbb{A}}$ of $G_{\mathbb{A}}$. We can treat a larger class with little further effort. Namely, the compact abelian group $Z_{\mathbb{A}} / Z^{+} Z_{k} \approx \mathbb{J}^{1} / k^{\times}$, being a quotient of the center $Z_{\mathbb{A}}$ of $G_{\mathbb{A}}$, acts on functions on $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ and commutes with right translation by $G_{\mathbb{A}}$. In particular, the action of $Z_{\mathbb{A}} / Z^{+} Z_{k}$ commutes with the integral operators on $G_{v}$ for $v<\infty$, and with differential operators coming from the Lie algebra $\mathfrak{g}_{v}$ of $G_{v}$ at archimedean places. Thus, for this chapter, an automorphic form or automorphic function is a function on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. For each character $\omega$ of $Z_{\mathbb{A}} / Z^{+} Z_{k}$, the space $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$ of all left $Z^{+} G_{k}$-invariant $f$ on $G_{\mathbb{A}}$ such that $|f| \in L^{2}\left(Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}\right)$ and $f(z g)=\omega(a) \cdot f(g)$ for all $z \in Z_{\mathbb{A}}$ is the space of $L^{2}$ automorphic forms with central character $\omega$.
[3.6.1] Claim: $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ decomposes by central characters:

$$
L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\text { completion of } \bigoplus_{\omega} L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)
$$

(Proof identical to that in [2.5].)

### 3.7 Discrete decomposition of cuspforms

For a standard parabolic subgroup $P$ of $G$ with unipotent radical $N=N^{P}$, the constant term of an automorphic form $f$ along $P$ is

$$
c_{P} f(g)=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n g) d n
$$

For general reasons [6.1], the group $N_{\mathbb{A}}$ acts continuously on the Fréchet space $C^{o}\left(Z^{+} N_{k} \backslash G_{\mathbb{A}}\right)$, and $N_{k} \backslash N_{\mathbb{A}}$ is compact, so for $f \in C^{o}\left(Z^{+} N_{k} \backslash G_{\mathbb{A}}\right)$ the constant-term integral exists as a Gelfand-Pettis integral, and is a continuous function.
[3.7.1] Claim: Constant terms are functions on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$.
Proof: By changing variables, $g \rightarrow c_{P} f(g)$ is a left $N_{\mathbb{A}}$-invariant function on $G_{\mathbb{A}}$ :

$$
c_{P} f\left(n^{\prime} x\right)=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(n \cdot n^{\prime} x\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(\left(n n^{\prime}\right) \cdot x\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n \cdot x) d n \quad\left(\text { for } n^{\prime} \in N_{\mathbb{A}}\right)
$$

Similarly, for $m \in M_{k}$,

$$
c_{P} f(m x)=\int_{N_{k} \backslash N_{\mathbb{A}}} f(n \cdot m x) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(m \cdot m^{-1} n m \cdot x\right) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} f\left(m^{-1} n m \cdot x\right) d n
$$

since $f$ itself is left $M_{k}$-invariant. Replacing $n$ by $m n m^{-1}$ gives the expression for $c_{P} f(g)$, noting that conjugation by $m \in M_{k}$ stabilizes $N_{k}$, and by the product formula the change of measure on $N_{\mathbb{A}}$ is trivial. Invariance under $Z^{+}$is even easier.

A cuspform is a function $f$ on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ meeting Gelfand's condition $c_{P} f=0$ for every standard parabolic $P$. When $f$ is merely measurable, so does not have well-defined pointwise values everywhere, this condition is best interpreted distributionally, as in [2.7] for $G L_{2}$, and addressed below in [3.8], in terms of pseudoEisenstein series. The space of square-integrable cuspforms is

$$
L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\left\{f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right): c_{P} f=0, \text { for all } P\right\}
$$

The fundamental theorem proven in [7.1-7.7] is the discrete decomposition of spaces of cuspforms. A simple version addresses the space

$$
L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)=\left\{\text { right- } K_{\mathbb{A}} \text {-invariant square-integrable cuspforms with central character } \omega\right\}
$$

where $K_{\mathbb{A}}=\prod_{v \leq \infty} K_{v}$. This space is $\{0\}$ unless $\omega$ is unramified, that is, is trivial on $Z_{\mathbb{A}} \cap K_{\mathbb{A}}$, since $K_{\mathbb{A}}$-invariance implies $Z_{\mathbb{A}} \cap K_{\mathbb{A}}$-invariance, and we also require $Z_{\mathbb{A}}$, $\omega$-equivariance.

The spherical Hecke algebras $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ act by right translation, and the Gelfand condition is an integral on the left, so spaces of cuspforms are stable under all these integral operators. The everywherespherical form of the decomposition result is
[3.7.2] Theorem: $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ has an orthonormal basis of simultaneous eigenfunctions for spherical Hecke algebras $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$. Each simultaneous eigenspace occurs with finite multiplicity, that is, is finite-dimensional. (Proof in [7.1-7.7].)

In contrast, the full spaces $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ do not have bases of simultaneous $L^{2}$-eigenfunctions. Instead, as in [2.11-2.12] and [3.15-3.16], the orthogonal complement of cuspforms in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ mostly consists of integrals of non- $L^{2}$ eigenfunctions for the Laplacians and Hecke operators, the Eisenstein series introduced below in [3.9].

For spaces of automorphic forms more complicated than being right $K_{v}$-invariant for every place $v$, there is generally no decomposition in terms of simultaneous eigenspaces for commuting operators. The decomposition argument in [7.7] directly uses the larger non-commutative algebras of test functions on the groups $G_{v}$ :

$$
C_{c}^{\infty}\left(G_{v}\right)=\left\{\begin{array}{r}
\text { compactly-supported smooth functions for } v \text { archimedean } \\
\text { compactly-supported locally constant functions for } v \text { non-archimedean }
\end{array}\right.
$$

Both cases are called smooth. With right translation $R_{g} f(x)=f(x g)$ for $x, g \in G_{\mathbb{A}}$, the action of $\varphi \in C_{c}^{\infty}\left(G_{v}\right)$ on functions $f$ on $G_{k} \backslash G_{\mathbb{A}}$ is

$$
\varphi \cdot f=\int_{G_{v}} \varphi(g) R_{g} f d g
$$

This makes sense not just as a pointwise-value integral, but as a Gelfand-Pettis integral [14.1] when $f$ lies in any quasi-complete, locally convex topological vectorspace $V$ on which $G_{v}$ acts so that $G_{v} \times V \rightarrow V$ is continuous. Such $V$ is a representation of $G_{v}$. The multiplication in $C_{c}^{\infty}\left(G_{v}\right)$ compatible with such actions is convolution: associativity $\varphi \cdot(\psi \cdot f)=(\varphi * \psi) \cdot f$.

Here, we are mostly interested in actions $G_{v} \times X \rightarrow X$ on Hilbert-spaces $X$. Such a representation is (topologically) irreducible when $X$ has no closed, $G_{v}$-stable subspace. The convolution algebras $C_{c}^{\infty}\left(G_{v}\right)$ are not commutative, so, unlike the commutative case, few irreducible representations are one-dimensional. In fact, typical irreducible representations of $C_{c}^{\infty}\left(G_{v}\right)$ turn out to be infinite-dimensional. There is no mandate to attempt to classify these irreducibles. Indeed, the spectral theory of compact self-adjoint operators still proves [7.7] discrete decomposition with finite multiplicities, for example, formulated as follows.

For every place $v$, let $K_{v}^{\prime}$ be a compact subgroup of $G_{v}$, and for all but a finite set $S$ of places require that $K_{v}^{\prime}=K_{v}$, the standard compact subgroup. For simplicity, we still assume $K_{v}^{\prime}=K_{v}$ at archimedean places. Put $K^{\prime}=\prod_{v} K_{v}^{\prime}$. Let $\omega$ be a central character trivial on $Z_{\mathbb{A}} \cap K^{\prime}$, so that the space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$ of right $K^{\prime}$-invariant cuspforms with central character $\omega$ is not $\{0\}$ for trivial reasons. For $v \in S$, we have a subalgebra $C_{c}^{\infty}\left(K_{v}^{\prime} \backslash G_{v} / K_{v}^{\prime}\right)$ of the convolution algebra of test functions at $v$, stabilizing $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$. [3.7.3] Theorem: $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$ is the completion of the orthogonal direct sum of subspaces, each consisting of simultaneous eigenfunctions for spherical algebras $C_{c}^{o}\left(K_{v} \backslash G_{v} / K_{v}\right)$ at $v \notin S$, and irreducible $C_{c}^{\infty}\left(K_{v}^{\prime} \backslash G_{v} / K_{v}^{\prime}\right)$-representations at $v \in S$. Each occurs with finite multiplicity. (Proof in [7.1-7.7].)

The technical features of decomposition with respect to non-commutative rings of operators certainly bear amplification, postponed to [7.7]. In anticipation,
[3.7.4] Theorem: (Gelfand and Piatetski-Shapiro) $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ is the completion of the orthogonal direct sum of irreducibles $V$ for the simultaneous action of all algebras $C_{c}^{\infty}\left(G_{v}\right)$. Each irreducible occurs with finite multiplicity. (Proof in [7.7].)

Again, the various sorts of orthogonal complements to spaces of cuspforms are mostly not direct sums of irreducibles, but are integrals of Eisenstein series, as below, with a relatively small number of squareintegrable residues of Eisenstein series. For $G L_{2}$ or $G L_{3}$ the square-integrable residues of Eisenstein series are relatively boring, but for $G L_{4}$ and larger there are highly non-trivial square-integrable residues, namely, the Speh forms, since for $G L_{4}(\mathbb{R})$ the relevant unitary representations appear in [Speh 1981/2]. The
general pattern for residual spectrum for $G L_{n}$ was conjectured in [Jacquet 1982/3] and proven in [MoeglinWaldspurger 1989].

### 3.8 Pseudo-Eisenstein series

We want to express the orthogonal complement of cuspforms in the larger spaces $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ or $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}, \omega\right)$ or $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K^{\prime}, \omega\right)$ in terms of simultaneous eigenfunctions for spherical Hecke algebras almost everywhere. Therefore, we emphasize the commutative algebras of integral operators attached to left-and-right $K_{v}$-invariant test functions on $G_{v}$. To exhibit explicit $L^{2}$ functions demonstrably spanning the orthogonal complement to cuspforms, recast the Gelfand condition that all constant terms $c_{P} f$ vanish as a requirement that $c_{P} f$ vanishes as a distribution on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$, and give an equivalent distributional vanishing condition on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.

For each standard parabolic $P$, with $N=N^{P}$, the condition that $c_{P} f$ vanishes as a distribution is that

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=0 \quad\left(\text { for all } \varphi \in C_{c}^{\infty}\left(N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

where, again, $C_{c}^{\infty}\left(N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ is compactly-supported functions on that quotient, smooth in the archimedean coordinates and locally constant in the non-archimedean coordinates. Smoothness for archimedean places should mean indefinite differentiability on the right with respect to the differential operators coming from the Lie algebra, as in [4.1]. Given the compact support, (uniform) smoothness for non-archimedean places should mean that there exists a compact, open subgroup $K^{\prime}$ of $\prod_{v<\infty} K_{v}$ under which $\varphi$ is right invariant.

Beyond perhaps having pointwise values almost everywhere, the nature of $c_{P} f$ for $f$ merely $L^{2}$ is potentially obscure. For example, it is not likely that $c_{P} f \in L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. Instead, for general reasons [6.1], $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is dense in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ in the $L^{2}$ topology, and for general reasons [6.1] the left action of $N_{k} \backslash N_{\mathbb{A}}$ on $C^{o}\left(Z^{+} P_{k} \backslash G_{\mathbb{A}}\right)$ is a continuous map $N_{k} \times C^{o}\left(Z^{+} P_{k} \backslash G_{\mathbb{A}}\right) \rightarrow C^{o}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, so $c_{P} f$ exists as a $C^{o}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$-valued Gelfand-Pettis integral [14.1]. For such $f$, the integral of $c_{P} f$ against $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ is the integral of a compactly-supported, continuous function.

The simplest type of pseudo-Eisenstein series is

$$
\Psi_{\varphi}(g)=\Psi_{\varphi}^{P}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma \cdot g) \quad\left(\text { for } \varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

Convergence is good:
[3.8.1] Claim: The series for a pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ is locally finite, meaning that for $g$ in a fixed compact in $G_{\mathbb{A}}$, there are only finitely-many non-zero summands in $\Psi_{\varphi}(g)=\sum_{\gamma} \varphi(\gamma g)$. Further, $\Psi_{\varphi} \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, so these pseudo-Eisenstein series are in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{Q}}\right)$. (Proof identical to [2.7.1].)
[3.8.2] Claim: For $f \in L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, for a standard parabolic $P$, pseudo-Eisenstein series $\Psi_{\varphi}=\Psi_{\varphi}^{P}$ with $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ fit into an adjunction

$$
\int_{Z+N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=\int_{Z+G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot f \quad\left(\text { for } f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

In particular, $c_{P} f=0$ if and only if $\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot f=0$ for all $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$.
Proof: The mechanism of the proof is that of [2.7.2]. For general reasons [6.1] $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is dense in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, and we consider $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. This allows unwinding as in [5.2]:

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi \cdot c_{P} f=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \varphi(g)\left(\int_{N_{k} \backslash N_{\mathbb{A}}} f(n g) d n\right) d g=\int_{Z^{+} N_{k} M_{k} \backslash G_{\mathbb{A}}} \varphi(g) f(g) d g
$$

Winding up, using the left $G_{k}$-invariance of $f$ and $N_{k} M_{k}=P_{k}$,

$$
\int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}} f(g) \varphi(g) d g=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \sum_{\gamma \in P_{k} \backslash G_{k}} f(\gamma \cdot g) \varphi(\gamma \cdot g) d g=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f(g)\left(\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma g)\right) d g
$$

The inner sum in the last integral is the pseudo-Eisenstein series attached to $\varphi$. By Cauchy-SchwarzBunyakowsky,

$$
\left|\int_{Z+P_{k} \backslash G_{\mathbb{A}}} f \varphi\right|=\left|\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \Psi_{\varphi}\right| \leq|f|_{L^{2}} \cdot\left|\Psi_{\varphi}\right|_{L^{2}}
$$

which proves that the functional $f \rightarrow \int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}} f \varphi$ on $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous in the $L^{2}$ topology, so extends by continuity to a continuous linear functional on $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Indeed, this inequality asserts continuity of $f \rightarrow c_{P} f$ as a linear map from $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to distributions on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$ with the weak dual topology as in [13.14].

Similarly, with

$$
C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}, \omega\right)=\left\{\varphi \in C_{c}^{\infty}\left(N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right): \varphi(z g)=\omega(z) \cdot \varphi(g), \text { for all } z \in Z_{\mathbb{A}}, g \in G\right\}
$$

analogously, keeping track of complex conjugations:
[3.8.3] Claim: Let $N=N^{P}$. For $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$, with $\varphi$ in $C_{c}^{\infty}\left(N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}, \omega\right)$,

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \bar{\varphi} \cdot c_{P} f=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \overline{\Psi_{\varphi}} \cdot f
$$

Thus, $c_{P} f=0$ if and only if $\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \overline{\Psi_{\varphi}} \cdot f=0$ for all $\varphi \in C_{c}^{\infty}\left(N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}, \omega\right)$.
For $P=P^{\text {min }}$ the minimal standard parabolic, especially for right $K_{\mathbb{A}}$-invariant functions, as in [2.7.4] for $G L_{2}$, minimal-parabolic pseudo-Eisenstein series with test-function data can be broken up into sub-families parametrized by (tuples of) Hecke characters, as follows. With $P=P^{\text {min }}$, let

$$
M_{P}^{1}=\left\{\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right): m_{1}, \ldots, m_{r} \in \mathbb{J},\left|m_{1}\right|=\ldots=\left|m_{r}\right|=1\right\}
$$

The group $M_{k} \backslash M^{1}$ is compact, because $\mathbb{J}^{1} / k^{\times}$is compact [2.A]. Certainly $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ is inside $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, so such functions $\varphi$ admit decompositions in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ by characters $\chi$ of the compact abelian group $M_{k} \backslash M^{1}$ acting on the left, as in [6.11]. The integral expressing the $\chi^{t h}$ component

$$
\varphi^{\chi}(g)=\int_{M_{k} \backslash M^{1}} \chi(m)^{-1} \varphi(m g) d m
$$

is a Gelfand-Pettis integral converging in $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ for any quasi-complete [14.7] locally convex [13.11] topology on this space. That is, the Fourier components $\varphi^{\chi}$ of a compactly-supported smooth function along $M_{k} \backslash M^{1}$ are again compactly-supported smooth, and their sum converges to the original in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, at least. The support of $\varphi^{\chi}$ is worst $\left(M_{k} \backslash M^{1}\right) \times \operatorname{spt} \varphi$.
[3.8.4] Lemma: A function $f \in L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ has constant term $c_{P} f$ integrating to 0 against $\varphi$ in $C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$ if and only if $c_{P} f$ integrates to 0 against every $M_{k} \backslash M^{1}$-component $\varphi^{\chi}$ of $\varphi$.
Proof: The potential pitfall is that there is no claim that constant terms of functions in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ are in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. Fortunately, this is not an obstacle: as earlier, it suffices to consider $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, so $c_{P} f \in C^{o}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$. With $u$ the characteristic function of $\left(M_{k} \backslash M^{1}\right) \times \operatorname{spt} \varphi$, the truncation $u \cdot c_{P} f$ is in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, and truncation does not alter the integrals against $\varphi^{\chi}$ or $\varphi$. Letting $\langle$,$\rangle be the inner$ product in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$, since $\varphi=\sum_{\chi} \varphi^{\chi}$ in $L^{2}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right)$,

$$
\left\langle c_{P} f, \varphi\right\rangle=\left\langle u \cdot c_{P} f, \varphi\right\rangle=\sum_{\chi}\left\langle u \cdot c_{P} f, \varphi^{\chi}\right\rangle=\sum_{\chi}\left\langle c_{P} f, \varphi^{\chi}\right\rangle
$$

giving the assertion.
[3.8.5] Corollary: With $P=P^{\text {min }}$ and $M=M^{P}$, to know $c_{P} f=0$ it suffices to know orthogonality to $\Psi_{\varphi}$ for $\varphi$ in

$$
\left\{\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}\right): \varphi(m g)=\chi(m) \cdot \varphi(g) \text { for all } m \in M^{1}\right\}
$$

with $\chi$ ranging over characters of the compact group $M_{k} \backslash M^{1}$.

### 3.9 Cuspidal-data pseudo-Eisenstein series

The simplest pseudo-Eisenstein series $\Psi_{\varphi}^{P}$, with $\varphi \in C_{c}^{\infty}\left(Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{Q}}\right)$ having compact support on the relevant quotient, behave well, as in [3.8.1]. For minimal-parabolic $P$, such pseudo-Eisenstein series suffice for the corresponding part of spectral theory, as they have good decompositions in terms of corresponding genuine Eisenstein series, much as in [2.11] and [2.12], as below in [3.15]. However, for $r \geq 3$ and for non-minimal $P$, genuine Eisenstein series with best behavior involve cuspforms on the Levi component $M^{P}$. Anticipating this, we want pseudo-Eisenstein series $\Psi_{\varphi}^{P}$ that facilitate this part [3.16] of the spectral decomposition. This entails minor analytical complications, since the data $\varphi$ can no longer have compact support.

Let $\delta:(0,+\infty) \rightarrow \mathbb{J}$ be the usual imbedding of the ray in the archimedean factors of the ideles, so that $|\delta(t)|=t$ for $t>0$. The centers of the factors $G L_{d_{i}}(\mathbb{A})$ of the standard Levi component $M^{P}$ are copies of $\mathbb{J}$. The standard archimedean split component $A_{P}^{+}$of a parabolic $P=P^{d_{1}, \ldots, d_{\ell}}$ is the product of the copies of $\delta(0, \infty)$ in the product of the centers of the factors of $M^{P}(\mathbb{A})$. Another important subgroup of $M_{\mathbb{A}}^{P}$ is

$$
M_{P}^{1}=\left\{\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{\ell}
\end{array}\right):\left|\operatorname{det} m_{1}\right|=\ldots=\left|\operatorname{det} m_{\ell}\right|=1, m_{i} \in G L_{d_{i}}(\mathbb{A})\right\}
$$

By design, $M_{\mathbb{A}}^{P}=A_{P}^{+} \cdot M_{P}^{1}$, and $M_{k}^{P} \subset M_{P}^{1}$. As already in $G L_{2}$, the center of $M^{P}$ is larger than the center of $G$, so $Z_{\mathbb{A}} M_{k} \backslash M_{\mathbb{A}}$ is not quite a product of quotients of the form $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ or $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$. This discrepancy necessitates looking at test functions on the archimedean split components $A_{P}^{+}$or their quotients $Z^{+} \backslash A_{P}^{+}$, in addition to automorphic data on $M_{k} \backslash M_{P}^{1}$.

In brief, the data $\varphi$ on $Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}$ appropriate for spectral decompositions of pseudo-Eisenstein series $\Psi_{\varphi}^{P}$ in terms of genuine Eisenstein series with good behavior, must specify test function data on the split component $A_{P}^{+}$, and cuspforms on the $M_{k} \backslash M_{P}^{1}$. For the minimal parabolic, the cuspidal data is vacuous, since the Levi component is a product of copies of $G L_{1}$, and test function data and specification of character on the compact abelian group $M_{k} \backslash M_{P}^{1} \approx\left(k^{\times} \backslash \mathbb{J}^{1}\right)^{r}$ suffices for the spectral decomposition [3.15] of minimal parabolic pseudo-Eisenstein series in terms of genuine Eisenstein series with analytic continuations and functional equations. In contrast, for a non-minimal parabolic, some factor of the Levi component is $G L_{r^{\prime}}$ with $r^{\prime}>1$, and the cuspform condition is non-vacuous.

Further, we only consider everywhere spherical automorphic forms, that is, right $K_{\mathbb{A}}$-invariant and left $Z_{\mathbb{A}}$-invariant functions. This has the convenient simplification, via Iwasawa decomposition, that constant terms $c_{P} f$ are identifiable with functions on the quotient of the Levi component of $P$ :

$$
Z_{\mathbb{A}} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}=Z_{\mathbb{A}} N_{\mathbb{A}} M_{k} \backslash N_{\mathbb{A}} M_{\mathbb{A}} K_{\mathbb{A}} / K_{\mathbb{A}} \approx Z_{\mathbb{A}} M_{k} \backslash M_{\mathbb{A}} /\left(M_{\mathbb{A}} \cap / K_{\mathbb{A}}\right) \longleftarrow Z_{\mathbb{A}} M_{k} \backslash M_{\mathbb{A}}
$$

This allows easier description of the cuspidal data, as follows. Let $f_{1}, f_{2}$ be cuspforms on $G L_{r_{1}}(\mathbb{A})$ and $G L_{r_{2}}(\mathbb{A})$, right invariant by the standard maximal compacts everywhere, with central characters $\omega_{1}$ and $\omega_{2}$, necessarily unramified. Anticipating the behavior of corresponding genuine Eisenstein series, we require that $f_{1}$ and $f_{2}$ be eigenfunctions for all the spherical Hecke algebras, including the archimedean places. This includes an eigenfunction condition for invariant Laplacians. That is, $f_{1}$ and $f_{2}$ are cuspforms in the strong sense, beyond satisfaction of the Gelfand condition on vanishing of constant terms. The theory of the constant term [8.3] shows that cuspforms in this strong sense are of rapid decay. Then $f=f_{1} \otimes f_{2}$ is a function on $G L_{r_{1}}(\mathbb{A}) \times G L_{r_{2}}(\mathbb{A}) \approx M_{\mathbb{A}}^{P}$. In the extreme cases where $r_{1}=1$ or $r_{2}=1$, the situation degenerates a little: there is no corresponding $f_{j}$, that is, the corresponding $f_{j}$ is simply the identically- 1 function. For a test function $\eta$ on the ray $(0, \infty)$, define

$$
\varphi(z n m k)=\varphi_{\eta, f}(z n m k)=\eta\left(\frac{\left|\operatorname{det} m_{1}\right|^{r_{2}}}{\left|\operatorname{det} m_{2}\right|^{r_{1}}}\right) \cdot f_{1}\left(m_{1}\right) \cdot f_{2}\left(m_{2}\right)
$$

with $m=\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right) \in M_{\mathbb{A}}^{P}, z \in Z^{+}, n \in N_{\mathbb{A}}, k \in K_{\mathbb{A}}$. The possibly counter-intuitive exponents on the idele norms of the determinants make $\varphi$ invariant under $Z_{\mathbb{A}}$. The corresponding pseudo-Eisenstein series is formed as expected,

$$
\Psi_{\varphi}^{P}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma \cdot g)
$$

However, this sum is not locally finite, so convergence is subtler, and needs properties of strong-sensecuspform data. Convergence will follow from comparison to similarly-formed genuine Eisenstein series in their range of absolute convergence, in [3.11.2].
[3.9.1] Remark: The argument of [3.11.3] for orthogonality of genuine Eisenstein series with cuspidal data attached to non-associate parabolics applies to pseudo-Eisenstein series with cuspidal data as well, showing orthogonality of those attached to non-associate parabolics. For associate parabolics $P, Q$, as for $G L_{2}$ in [2.13.5], spectral decompositions of pseudo-Eisenstein series will make clear [3.17.3] that pseudo-Eisenstein series $\Psi_{\eta, f}^{P}$ and $\Psi_{\theta, f^{\prime}}^{Q}$ with test functions $\eta, \theta$ and cuspidal data $f, f^{\prime}$, are orthogonal if $P=Q$ but $\left\langle f, f^{\prime}\right\rangle=0$, or if $M^{P}=w M^{Q} w^{-1}$ but $f^{w} \neq f^{\prime}$.

### 3.10 Minimal-parabolic Eisenstein Series

In the often-treated example of automorphic forms on $G L_{2}$, there are no Eisenstein series made from cuspidal data, because $G L_{2}$ is so small. In contrast, for $G L_{n}$ with $n>2$, cuspidal-data Eisenstein series play an essential role. However, the minimal-parabolic Eisenstein series for $G L_{r}$ involve no cupidal data, because the Levi component is a product of groups $G L_{1}$, where the cuspidal condition is vacuous. Further, especially in the everywhere-spherical case of right $K_{\mathbb{A}}$-invariant minimal-parabolic Eisenstein series, much of the behavior reduces to $G L_{2}$ via Bochner's lemma [3.B], as we will see. Hartogs' lemma on separate analyticity implying joint analyticity [15.C] removes several ambiguities and potential imprecisions in discussion of functions of one complex variable versus several.

With $\delta$ mapping $(0, \infty)$ to the archimedean factors of $\mathbb{J}$ so that $|\delta(t)|=t$, as earlier, describe Hecke characters $\widetilde{\chi}$ as

$$
\widetilde{\chi}\left(\delta(t) \cdot t_{1}\right)=t^{s} \cdot \chi\left(t_{1}\right) \quad\left(\text { with } t>0, t_{1} \in \mathbb{J}^{1}, s \in \mathbb{C}\right)
$$

Given an $r$-tuple of Hecke characters $\widetilde{\chi}_{1}, \ldots, \widetilde{\chi}_{r}$ with the relation $s_{1}+\ldots+s_{r}=0$ among the complex parameters $s=\left(s_{1}, \ldots, s_{r}\right)$, the right $K_{\mathbb{A}}$-invariant, $Z_{\mathbb{A}}$-invariant minimal-parabolic Eisenstein series $E_{s, \chi}=E_{s, \chi}^{\min }$ on $G L_{r}$ is formed as usual:

$$
E_{s, \chi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s, \chi}^{o}(\gamma \cdot g)
$$

where

$$
\varphi_{s, \chi}^{o}(n m k)=\widetilde{\chi}\left(m_{1}\right) \cdot \ldots \cdot \tilde{\chi}_{r}\left(m_{r}\right) \quad\left(\text { for } n \in N_{\mathbb{A}}^{\min }, m=\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right), k \in K_{\mathbb{A}}\right)
$$

For Hecke characters all of the simplest form $\widetilde{\chi}_{j}\left(\delta(t) \cdot t_{1}\right)=t^{s_{j}}$, this is

$$
\varphi_{s}^{o}(n m k)=\varphi_{s, 1}^{o}(n m k)=\left|m_{1}\right|^{s_{1}}\left|m_{2}\right|^{s_{2}} \ldots\left|m_{r}\right|^{s_{r}}
$$

That is, in terms of the parameter $s, E_{s, \chi}$ is a function-valued function of $r-1$ complex variables, but the parameter space is the complex hyperplane $s_{1}+\ldots+s_{r}=0$ in $\mathbb{C}^{r}$, rather than $\mathbb{C}^{r-1}$. In terms of the positive simple roots $\alpha_{i}(m)=m_{i} / m_{i+1}$, using $s_{1}+\ldots+s_{r}=0$,

$$
\varphi_{s, \chi}^{o}(n m k)=\left|\alpha_{1}(m)\right|^{s_{1}} \cdot\left|\alpha_{2}(m)\right|^{s_{1}+s_{2}} \cdot\left|\alpha_{3}(m)\right|^{s_{1}+s_{2}+s_{3}} \cdot \ldots \cdot\left|\alpha_{r-1}(m)\right|^{s_{1}+\ldots+s_{r-1}}
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

[3.10.1] Claim: (In coordinates) The minimal-parabolic Eisenstein series $E_{s, \chi}(g)$ on $G L_{n}$ converges (absolutely and uniformly for $g$ in compacts) for $\frac{\sigma_{j}-\sigma_{j+1}}{2}>1$ for $j=1, \ldots, r-1$, where $s=\left(s_{1}, \ldots, s_{r}\right) \in \mathbb{C}^{r}$ and $\sigma=\left(\operatorname{Re}\left(s_{1}\right), \ldots, \operatorname{Re}\left(s_{r}\right)\right)$. (Proof below.)

The inequalities describing the region of convergence can be rewritten in a more intrinsic form later relevant to functional equations, as follows. Let $\mathfrak{g l}_{r}(\mathbb{R})$ be the Lie algebra of $G L_{r}(\mathbb{R})$, that is, all $r$-by- $r$ real matrices. Let $\mathfrak{a}$ be the Lie algebra of the diagonal matrices in $G L_{r}(\mathbb{R})$. The non-zero eigenvalues (roots) of $\mathfrak{a}$ on $\mathfrak{g l} l_{r}(\mathbb{R})$ are functionals $a \rightarrow a_{i}-a_{j}$ in the dual space $\mathfrak{a}^{*}$. For $i \neq j$, the corresponding eigenspace (rootspace) is matrices with non-zero entry only at the $i j^{t h}$ entry. The standard positive roots and rootspaces are those with $i<j$. Write $\beta>0$ for positive root $\beta$, and $\beta<0$ when $-\beta>0$. The standard simple positive roots are $a \rightarrow a_{i}-a_{i-1}$ The half-sum of positive roots is

$$
\rho(a)=\sum_{i<j}\left(a_{i}-a_{j}\right) \quad(\text { for } a \in \mathfrak{a})
$$

There is a sort of logarithm $\operatorname{map} M_{\mathbb{A}} \rightarrow \mathfrak{a}$ by

$$
\log \left|\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right)\right|=\left(\begin{array}{ccc}
\log \left|m_{1}\right| & & \\
& \ddots & \\
& & \log \left|m_{r}\right|
\end{array}\right)
$$

and then for $m \in M_{\mathbb{A}}$ and $\alpha \in \mathfrak{a}^{*}$, write

$$
m^{\alpha}=e^{\alpha(\log |m|)}
$$

This enables interpretation of the parameter $s$ as lying in the complexification $\mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$ of the dual $\mathfrak{a}^{*}$ of $\mathfrak{a}$. Using $\langle x, y\rangle=\operatorname{tr}(x y)$ on $\mathfrak{a}$, we can identify $\mathfrak{a}$ with $\mathfrak{a}^{*}$, and transport to $\mathfrak{a}^{*}$ the pairing $\langle$,$\rangle .$
[3.10.2] Corollary: (Intrinsic/conceptual version) The minimal-parabolic Eisenstein series $E_{s, \chi}(g)$ on $G L_{r}$ converges (absolutely and uniformly for $g$ in compacts) for $\langle\alpha, \sigma-2 \rho\rangle>0$ for all positive simple roots $\alpha$. (Proof below.)

That is, the Eisenstein series $E_{s, \chi}$ converges absolutely for $\sigma \in \mathfrak{a}$ in the translate by $2 \rho$ of

$$
\text { positive Weyl chamber }=\left\{x \in \mathfrak{a}^{*}:\langle x, \alpha\rangle>0 \text {, for all positive roots } \alpha\right\} \subset \mathfrak{a}^{*}
$$

Proof: (of claim) For convergence, it suffices to treat Hecke characters only of the form $\widetilde{\chi}(y)=|y|^{s}$. With number field $k$, let $h$ be the standard height function on a $k$-vectorspace with specified basis. Let $P=P^{\min }$ be the standard minimal parabolic of $G$. Let $e_{1}, \ldots, e_{r}$ be the standard basis of $k^{r}$. Any exterior power $\wedge^{\ell}\left(k^{r}\right)$ has (unordered) basis of wedges of the $e_{j}$, and an associated height function. Let $v_{o}=e_{j} \wedge \ldots \wedge e_{r}$, and

$$
\eta_{j}(g)=\frac{h\left(v_{o} \cdot \wedge^{r-j+1} g\right)}{h\left(v_{o}\right)} \quad\left(\text { for } g \in G L_{r}(\mathbb{A})\right)
$$

where $\wedge^{\ell} g$ is the natural action of $g$ on $\wedge^{\ell} k^{r}$. The spherical vector $\varphi_{s}=\varphi_{s, 1}^{o}$, from which the $s^{t h}$ minimalparabolic Eisenstein series $E_{s}=E_{s, 1}$ is made, is expressible as

$$
\varphi_{s}^{o}=\eta_{1}^{s_{1}} \eta_{2}^{s_{2}-s_{1}} \eta_{3}^{s_{3}-s_{1}-s_{2}} \ldots \eta_{r}^{s_{r}-s_{1}-s_{2}-\ldots-s_{r-1}} \quad\left(\text { where } s=\left(s_{1}, \ldots, s_{r}\right)\right)
$$

From reduction theory, given compact $C \subset Z_{\mathbb{A}} \backslash G_{\mathbb{A}}$, for some implied constants depending only on $C$,

$$
h(v)<_{C} \quad h(v \cdot g)<_{C} \quad h(v) \quad\left(\text { for all } 0 \neq v \in k^{r} \text { and } g \in C\right)
$$

and similarly for heights on $\wedge^{\ell} k^{r}$. Therefore, convergence of the series defining the Eisenstein series $E_{s}\left(g_{o}\right)$ is equivalent to convergence of

$$
\int_{C} \sum_{\gamma \in Z_{\mathbb{A}} P_{k} \backslash G_{k}} \varphi_{s}^{o}(\gamma g) d g
$$

Shrinking $C$ sufficiently so that $\gamma \cdot C \cap C \neq \phi$ implies $\gamma=1$,

$$
\int_{C} \sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s}^{o}(\gamma g) d g=\int_{Z_{\mathbb{A}} P_{k} \backslash G_{k} \cdot C} \varphi_{s}^{o}(g) d g
$$

From reduction theory, the infimum $\mu$ of $h(v)$ over non-zero primitive $v$ in $\wedge^{\ell}\left(\mathbb{A}^{r}\right)$ is attained, so is positive. In particular, $\mu \leq h\left(v_{o} \gamma g\right)$ for all $g \in C$ and $\gamma \in G_{k}$. Thus, $G_{k} \cdot C$ is contained in a set

$$
Y=\left\{g \in G_{\mathbb{A}}: 1<_{C} \eta_{j}(g) \text { for } j=1, \ldots, r\right\}
$$

Thus, convergence of the Eisenstein series is implied by convergence of

$$
\int_{Z_{\mathbb{A}} P_{k} \backslash Y}\left|\varphi_{s}^{o}(g)\right| d g
$$

The set $Y$ is stable by right multiplication by the maximal compact subgroup $K_{v} \subset G_{v}$ at all places $v$, so by Iwasawa decomposition this integral is

$$
\left.\int_{Z_{\mathbb{A}} P_{k} \backslash\left(Y \cap P_{\mathbb{A}}\right)}\left|\varphi_{s}^{o}(p)\right| d p \quad \text { (left Haar measure on } Z_{\mathbb{A}} \backslash P_{\mathbb{A}}\right)
$$

With $\rho$ the half-sum of positive roots, the left Haar measure on $Z_{\mathbb{A}} P_{\mathbb{A}}$ is $d(n m)=d n d m / m^{2 \rho}$, where $d n$ is Haar measure on the unipotent radical and $d m$ is Haar measure on $Z_{\mathbb{A}} \backslash M_{\mathbb{A}}$. Since $\varphi_{s}^{o}$ is left $N_{\mathbb{A}}$-invariant and $N_{k} \backslash N_{\mathbb{A}}$ is compact, convergence of the latter integral is equivalent to convergence of

$$
\int_{Z_{\mathbb{A}} M_{k} \backslash\left(Y \cap M_{\mathbb{A}}\right)}\left|\varphi_{s}^{o}(a)\right| \frac{d m}{m^{2 \rho}}=\int_{Z_{\mathbb{A}} M_{k} \backslash\left(Y \cap M_{\mathbb{A}}\right)} m^{\sigma-2 \rho} d m \quad\left(\text { where } \sigma=\left(\operatorname{Re} s_{1}, \ldots, \operatorname{Re} s_{r}\right)\right)
$$

The quotient $k^{\times} \backslash \mathbb{J}^{1}$ of norm-one ideles $\mathbb{J}^{1}$ is compact, by [2.A], and the discrepancy between $Z_{\mathbb{A}} \backslash G_{\mathbb{A}}$ and $S L_{r}(\mathbb{A})$ is absorbed by $M_{\mathbb{A}} \cap \prod_{v} K_{v}$. Thus, convergence of the following integral suffices.

Parametrize a subgroup $H$ of $S L_{n}(\mathbb{A})$ by $r-1$ maps from $\left(G L_{1}(\mathbb{A})\right.$, namely,

$$
h_{j}: t \longrightarrow\left(\begin{array}{ccccccc}
1 & & & & & & \\
& \ddots & & & & & \\
\\
& & 1 & & & & \\
\\
& & & t & & & \\
\\
& & & & & t^{-1} & \\
& & & \\
& & & & & & \ddots \\
\\
& & & & & & \\
&
\end{array}\right) \quad\left(\text { at } j^{t h} \text { and }(j+1)^{t h}\right. \text { positions) }
$$

From

$$
\eta_{i}\left(h_{j}(t)\right)=\left\{\begin{array}{cl}
|t|^{-1} & (\text { for } i=j+1) \\
1 & (\text { otherwise })
\end{array}\right.
$$

we have

$$
Y \cap M_{\mathbb{A}} \cap S L_{n}(\mathbb{A})=\left\{\prod_{j} h_{j}\left(t_{j}\right): t_{j} \in \mathbb{J} \text { and }\left|t_{j}^{-1}\right| \gg 1\right\}
$$

Again using compactness of $k^{\times} \backslash \mathbb{J}^{1}$, noting that $h_{j}(t)^{2 \rho}=|t|^{2}$ for all $j$, convergence of the Eisenstein series is implied by convergence of the archimedean integral

$$
\int_{0}^{\ll 1} t^{\sigma_{j}-\sigma_{j+1}-2} \frac{d t}{t} \quad(\text { for } j=1, \ldots, r-1)
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

These integrals are absolutely convergent for $\sigma_{i}-\sigma_{i+1}-2>0$ for all $i$.
Proof: (of corollary) The absolute convergence condition is $\langle\sigma-2 \rho, \alpha\rangle>0$ for all simple roots $\alpha$.
The general shape of the $P=P^{\text {min }}$ constant term of the simplest $P$ Eisenstein series is easily determined:
[3.10.3] Claim: In the region of convergence, for suitable holomorphic functions $s \rightarrow c_{w, s}$, with $c_{1, s}=1$, the constant term is

$$
c_{P} E_{\rho+s}^{P}(m)=m^{\rho+s}+\sum_{1 \neq w \in W} c_{w, s} m^{\rho+w \cdot s} \quad\left(\text { with } m \in M_{\mathbb{A}}^{\min }\right)
$$

[3.10.4] Remark: We could explicitly compute the coefficients $c_{w, s}$ as part of the proof of this claim, but essentially the same computation occurs in the proof of the functional equations [3.12.1] and the corollary [3.12.3]. Quoting [3.12.3] and [3.12.6] to have a more complete statement here:
[3.10.5] Corollary: $c_{\tau, s}=\xi\langle s, \alpha\rangle / \xi(1+\langle s, \alpha\rangle)$ for reflections $\tau$, and the cocycle relation $c_{w^{\prime}, w \cdot s} \cdot c_{w, s}=c_{w w^{\prime}, s}$ holds for $w, w^{\prime} \in W$ and $s \in \mathfrak{a}^{*} \otimes_{R} \mathbb{C}$. We have

$$
c_{w, s}=\prod_{\beta>0: w \cdot \beta<0} \frac{\xi\langle s, \beta\rangle}{\xi(\langle s, \beta\rangle+1)}
$$

Proof: This begins with an archetypical unwinding argument, using the disjointness of $G_{k}=\bigsqcup_{w} P_{k} w N_{k}$ from the Bruhat decomposition [3.1.1].

$$
\begin{gathered}
c_{P} E_{\rho+s}^{P}(m)=\int_{N_{k} \backslash N_{\mathbb{A}}} E_{\rho+s}^{P}(n \cdot m) d n=\int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{\rho+s}^{o}(\gamma \cdot n m) d n \\
=\sum_{w} \int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\gamma \in P_{k} \backslash P_{k} w N_{k}} \varphi_{\rho+s}^{o}(\gamma \cdot n m) d n=\sum_{w} \int_{N_{k} \backslash N_{\mathbb{A}}} \sum_{\beta \in\left(w^{-1} P_{k} w \cap N_{k}\right) \backslash N_{k}} \varphi_{\rho+s}^{o}(w \beta \cdot n m) d n \\
=\sum_{w} \int_{\left(w^{-1} P_{k} w \cap N_{k}\right) \backslash N_{\mathbb{A}}} \varphi_{\rho+s}^{o}(w \cdot n m) d n=\sum_{w} \int_{\left(w^{-1} N_{k} w \cap N_{k}\right) \backslash N_{\mathbb{A}}} \varphi_{\rho+s}^{o}(w \cdot n m) d n
\end{gathered}
$$

Since $\varphi_{\rho+s}^{o}$ is left $N_{\mathbb{A}}$-invariant, $g \rightarrow \varphi_{\rho+s}^{o}(w g)$ is still left invariant by $w^{-1} N_{\mathbb{A}} w \cap N_{\mathbb{A}}$. Thus, with the volume of $\left(w^{-1} N_{k} w \cap N_{k}\right) \backslash\left(w^{-1} N_{\mathbb{A}} w \cap N_{\mathbb{A}}\right)$ normalized to 1 , the constant term is

$$
\sum_{w} \int_{\left(w^{-1} N_{\mathbb{A}} w \cap N_{\mathbb{A}}\right) \backslash N_{\mathbb{A}}} \varphi_{\rho+s}^{o}(w \cdot n m) d n
$$

The case $w=1$ gives $\varphi_{\rho+s}^{o}(m)=m^{\rho+s}$. More generally, there is a convenient complementary subgroup $N^{w}$ to $w^{-1} N_{\mathbb{A}} w \cap N_{\mathbb{A}}$ inside $N_{\mathbb{A}}$ :
[3.10.6] Lemma: Let $N^{\text {opp }}$ be lower-triangular matrices with 1's on the diagonal, and let $N^{w}=N \cap w^{-1} N^{\mathrm{opp}} w$. Then

$$
N^{w} \cap\left(w^{-1} N w \cap N\right)=\{1\} \quad \text { and } \quad N^{w} \cdot\left(w^{-1} N w \cap N\right)=N
$$

Proof: (of Lemma) First, of course,

$$
w^{-1} N^{\mathrm{opp}} w \cap w^{-1} N w=w^{-1}\left(N^{\mathrm{opp}} \cap N\right) w=w^{-1}\{1\} w=\{1\}
$$

For a root $\alpha(m)=m_{i} / m_{j}$ for $i \neq j$ and $m \in M=M^{\mathrm{min}}$, the corresponding root subgroup is

$$
N^{\alpha}=\left\{n=e^{x}: x \in \mathfrak{n}, m x m^{-1}=\alpha(m) \cdot x\right\} \quad(\mathfrak{n}=\text { Lie algebra of } N)
$$

where $x \rightarrow e^{x}$ is the usual matrix exponential. Thus, $N$ is generated by all the $N_{\beta}$ for positive roots $\beta$, and $N^{\text {opp }}$ is generated by the $N_{\beta}$ with negative roots $\beta$. The action of $W$ permutes roots, so it permutes root subgroups. Every root subgroup is inside either $w^{-1} N w$ or $w^{-1} N^{\text {opp }} w$, so the intersections of these with $N$ generate $N$.

Then

$$
c_{P} E_{\rho+s}^{P}(m)=\sum_{w} \int_{N_{\mathbb{A}}^{w}} \varphi_{\rho+s}^{o}(w \cdot n m) d n
$$

Each root subgroup is stable under conjugation by $M$, so any product that is a subgroup of $N$ is stable by $M$. Thus, in the integral, replace $n$ by $m n m^{-1}$ : letting $\delta^{w}(m)$ be the change of measure $d\left(m n m^{-1}\right) / d n$,

$$
\begin{aligned}
& c_{P} E_{\rho+s}^{P}(m)=\sum_{w} \delta^{w}(m) \int_{N_{\mathbb{A}}^{w}} \varphi_{\rho+s}^{o}(w m n) d n=\sum_{w} \delta^{w}(m) \int_{N_{\mathbb{A}}^{w}} \varphi_{\rho+s}^{o}\left(w m w^{-1} \cdot w n\right) d n \\
& =\sum_{w} \delta^{w}(m)\left(w m w^{-1}\right)^{\rho+s} \int_{N_{\mathbb{A}}^{w}} \varphi_{\rho+s}^{o}(w n) d n=\sum_{w} \delta^{w}(m) m^{w^{-1} \cdot(\rho+s)} \int_{N_{\mathbb{A}}^{w}} \varphi_{\rho+s}^{o}(w n) d n
\end{aligned}
$$

As usual, the sign in the exponent of $w$ in the latter expression is necessary for the action of $W$ on $\mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$ to be associative. Thus,

$$
c_{w, s}=\int_{N_{\mathbb{A}}^{w}} \varphi_{\rho+s}^{o}\left(w^{-1} n\right) d n
$$

Optimistically, to understand $\delta^{w^{-1}}(m) m^{w \cdot(\rho+s)}=\delta^{w^{-1}}(m) m^{w \cdot \rho} \cdot m^{w \cdot s}$, apparently
[3.10.7] Lemma: $\delta^{w^{-1}}(m) m^{w \cdot \rho}=m^{\rho}$ for $m \in M$.
Proof: (of Lemma) Write $\beta>0$ or $\beta<0$ as $\beta$ is a positive or negative root. The character $m \rightarrow \delta^{w}(m)$ is the modular function of $N \cap w^{-1} N^{\mathrm{opp}} w$, so $\delta^{w}(m)=m^{\gamma}$ where

$$
\gamma=\sum_{\beta<0: w^{-1} \beta>0} w^{-1} \cdot \beta=\sum_{\beta>0: w^{-1} \beta<0} w^{-1} \cdot(-\beta)=-\sum_{\beta>0: w^{-1} \beta<0} w^{-1} \cdot \beta
$$

Meanwhile,

$$
\begin{gathered}
w^{-1} \cdot 2 \rho=\sum_{\beta>0} w^{-1} \cdot \beta=\sum_{\beta>0: w^{-1} \cdot \beta>0} w^{-1} \cdot \beta+\sum_{\beta>0: w^{-1} \cdot \beta<0} w^{-1} \cdot \beta \\
=\sum_{\beta>0: w^{-1} \cdot \beta>0} w^{-1} \cdot \beta-\sum_{\beta<0: w^{-1} \cdot \beta>0} w^{-1} \cdot \beta
\end{gathered}
$$

Thus, $w^{-1} \cdot \rho+\gamma=\rho$.
Thus, we obtain an expression for $c_{P} E_{s}(m)$ of the asserted form.

### 3.11 Cuspidal-data Eisenstein series

To keep things relatively simple, our examples of cuspidal-data Eisenstein series for non-minimal proper parabolics will include only maximal proper parabolics. In fact, the general case is a combination of the features of the minimal-parabolic and maximal-proper parabolic.

Let $f_{1}, f_{2}$ be cuspforms on $G L_{r_{1}}(\mathbb{A})$ and $G L_{r_{2}}(\mathbb{A})$, right invariant by the standard maximal compacts everywhere, with trivial central characters. We require that $f_{1}$ and $f_{2}$ be eigenfunctions for all the spherical Hecke algebras, including the archimedean places. This includes an eigenfunction condition for invariant Laplacians. That is, $f_{1}$ and $f_{2}$ are cuspforms in a strong sense, beyond satisfaction of the Gelfand condition on vanishing of constant terms.

The corollary [7.3.19] of the discrete decomposition of cuspforms shows that cuspforms in this strong sense are of rapid decay. The cuspidal data $f=f_{1} \otimes f_{2}$ is a function on $G L_{r_{1}}(\mathbb{A}) \times G L_{r_{2}}(\mathbb{A}) \approx M_{\mathbb{A}}^{P}$. In the extreme cases where $r_{1}=1$ or $r_{2}=1$, the situation degenerates: there is no corresponding $f_{j}$, that is, the corresponding $f_{j}$ is simply the identically- 1 function. Let

$$
\varphi(z n m k)=\varphi_{s, f}(z n m k)=\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{s} \cdot f_{1}\left(m_{1}\right) \cdot f_{2}\left(m_{2}\right)
$$

with $m=\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right) \in M_{\mathbb{A}}^{P}, z \in Z^{+}, n \in N_{\mathbb{A}}, k \in K_{\mathbb{A}}$. The exponents on the idele norms of the determinants make $\varphi$ invariant under $Z_{\mathbb{A}}$. The corresponding genuine Eisenstein series is formed as expected:

$$
E_{s, f}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s, f}(\gamma \cdot g)
$$

[3.11.1] Claim: The cuspidal-data Eisenstein series $E_{s, f}(g)$ converges (absolutely and uniformly for $g$ in compacts) for $\operatorname{Re}(s)>1$.
[3.11.2] Corollary: All cuspidal-data pseudo-Eisenstein series converge (absolutely and uniformly on compacts).

Proof: (of corollary) The genuine Eisenstein series with any $\operatorname{Re}(s)>1$ and the same cuspidal data dominates every pseudo-Eisenstein series with that cuspidal data.
Proof: (of claim) As above, hypotheses on the cuspform $f$ assure that it is bounded, so it suffices to prove the claim with $f$ replaced by 1. Then the argument becomes a variant of that of [3.10.1] and [3.10.2], with $\varphi_{s, f}$ replaced by

$$
\varphi_{s}(n m k)=\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{s}
$$

The sum

$$
E_{s}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s}(\gamma \cdot g)
$$

dominates that for $E_{s, f}$. This $E_{s}$ is a degenerate Eisenstein series when either $r_{1}+r_{1}>2$, so-called because it is missing the cuspidal data, and does not play a direct role in spectral theory.

With $e_{1}, \ldots, e_{r}$ the standard basis of $k^{r}$, let $h$ be the standard height function on $k$-vectorspace $\wedge^{r_{2}}\left(k^{r}\right)$ with basis consisting of $r_{2}$-fold exterior products of the $e_{i}$. Put

$$
\eta(g)=\frac{h\left(e_{r_{1}+1} \wedge e_{r_{1}+2} \wedge \ldots \wedge e_{r} \cdot \wedge^{r_{2}} g\right)}{h\left(e_{r_{1}+1} \wedge e_{r_{1}+2} \wedge \ldots \wedge e_{r}\right)} \quad\left(\text { with } g \text { acting on } \wedge^{r_{2}}\left(k^{r}\right) \text { by } \wedge^{r_{2}} g\right)
$$

Note that $\eta$ is right $K_{v}$-invariant at all $v$, left $N_{\mathbb{A}}$-invariant, and $\eta\left(\begin{array}{cc}m_{1} & * \\ 0 & m_{2}\end{array}\right)=\left|\operatorname{det} m_{2}\right|$. Thus,

$$
\left|\varphi_{s}(g)\right|=\left(\frac{|\operatorname{det} g|^{r_{2}}}{\eta(g)^{r}}\right)^{\sigma} \quad(\text { with } \sigma=\operatorname{Re}(s))
$$

## Garrett: Modern Analysis of Automorphic Forms

From reduction theory, given compact $C \subset Z_{\mathbb{A}} \backslash G_{\mathbb{A}}, h(v)<_{C} h(v \cdot g)<_{C} h(v)$ for all $0 \neq v \in k^{r}$ and $g \in C$. Therefore, convergence of the series defining $E_{s}\left(g_{o}\right)$ is equivalent to convergence of

$$
\int_{C} \sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{\sigma}(\gamma g) d g
$$

Shrinking $C$ sufficiently so that $\gamma \cdot C \cap C \neq \phi$ implies $\gamma=1$,

$$
\int_{C} \sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{\sigma}(\gamma g) d g=\int_{Z_{\mathbb{A}} P_{k} \backslash G_{k} \cdot C} \varphi_{\sigma}(g) d g
$$

Let $\mu$ be the infimum of $h(v)$ over non-zero primitive $v$ in $\wedge^{r_{2}}\left(\mathbb{A}^{r}\right)$. From reduction theory, this infimum is attained, so it is $\mu>0$, and $\mu \ll h\left(v_{o} \gamma g\right)$ for all $g \in C$ and $\gamma \in G_{k}$. Thus, $G_{k} \cdot C$ is contained in a set

$$
Y=\left\{g \in G_{\mathbb{A}}:|\operatorname{det} g|^{r_{2}} / \eta(g)^{r}<_{C} 1\right\}
$$

and convergence of the Eisenstein series is implied by convergence of

$$
\int_{Z_{\mathbb{A}} P_{k} \backslash Y} \varphi_{\sigma}(g) d g
$$

The set $Y$ is stable by right multiplication by the maximal compact subgroup $K_{v} \subset G_{v}$ at all places $v$, so via the Iwasawa decomposition this integral is

$$
\int_{Z_{A} P_{k} \backslash\left(Y \cap P_{\mathbb{A}}\right)} \varphi_{\sigma}(p) d p \quad \text { (left Haar measure on } P \text { ) }
$$

The left Haar measure on $P_{\mathbb{A}}$ is

$$
d(n m)=\frac{d n d m}{\left|\operatorname{det} m_{1}\right|^{r_{2}} \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1}}} \quad\left(\text { where } m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right)
$$

where $d n$ is Haar measure on the unipotent radical and $d m$ is Haar measure on $M_{\mathbb{A}}^{P}$. Since $\varphi_{\sigma}$ is left $N_{\mathbb{A}}^{P}$-invariant and $N_{k}^{P} \backslash N_{\mathbb{A}}^{P}$ is compact, convergence of the latter integral is equivalent to convergence of

$$
\int_{Z_{\mathbb{A}} M_{k} \backslash\left(Y \cap M_{\mathbb{A}}\right)} \varphi_{\sigma}(m) \frac{d m}{\left|\operatorname{det} m_{1}\right|^{r_{2}} \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1}}}=\int_{Z_{\mathbb{A}} M_{k} \backslash\left(Y \cap M_{\mathbb{A}}\right)}\left|\operatorname{det} m_{1}\right|^{r_{2}(\sigma-1)}\left|\operatorname{det} m_{2}\right|^{-r_{1}(\sigma-1)} d m
$$

We have

$$
Y \cap M_{\mathbb{A}}=\left\{m \in M_{\mathbb{A}}:|\operatorname{det} m|^{r_{2}} / \eta(m)^{r}<_{C} 1\right\}=\left\{m \in M_{\mathbb{A}}:\left|\operatorname{det} m_{1}\right|^{r_{2}} /\left|\operatorname{det} m_{2}\right|^{r_{1}}<_{C} 1\right\}
$$

By reduction theory, for example, the quotients $G L_{r_{i}}(k) \backslash G L_{r_{i}}(\mathbb{A})^{1}$ have finite total measure, and $\left|\operatorname{det} m_{1}\right|^{r_{2}} /\left|\operatorname{det} m_{2}\right|^{r_{1}}$ is $Z_{\mathbb{A}} M^{1}$-invariant. Let $M^{1}$ be the copy of $G L_{r_{1}}(\mathbb{A})^{1} \times G L_{r_{2}}(\mathbb{A})^{1}$ inside $M_{\mathbb{A}}$. It suffices to prove convergence of

$$
\begin{gathered}
\int_{Z_{A} M^{1} \backslash\left(Y \cap M_{\mathbb{A}}\right)}\left|\operatorname{det} m_{1}\right|^{r_{2} \sigma} \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1} \sigma} \frac{d m}{\left|\operatorname{det} m_{1}\right|^{r_{2}} \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1}}} \\
=\int_{Z_{A} M^{1} \backslash\left(Y \cap M_{\mathbb{A}}\right)}\left|\operatorname{det} m_{1}\right|^{r_{2}(\sigma-1)} \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1}(\sigma-1)} d m
\end{gathered}
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

The map

$$
\Delta(t)=\left(\begin{array}{cc}
t \cdot 1_{r_{1}} & 0 \\
0 & 1_{r_{2}}
\end{array}\right) \quad(\text { for } t>0)
$$

surjects to $Z_{\mathbb{A}} M^{1} \backslash M_{\mathbb{A}}$, so it suffices to prove convergence of

$$
\int_{0}^{\ll 1}\left(\operatorname{det}\left(t \cdot 1_{r_{1}}\right)\right)^{r_{2}(\sigma-1)} \frac{d t}{t}=\int_{0}^{\ll 1} t^{r_{1} r_{2}(\sigma-1)} \frac{d t}{t}
$$

Convergence is implied by $\sigma>1$.
One benefit of cuspidal data for Eisenstein series is that many constant terms vanish for general reasons. For maximal proper parabolics, the outcome is especially clear. Continue to assume that $f_{1}$ and $f_{2}$ are everywhere spherical, for simplicity. The vanishing conclusion in the following follows without assuming much beyond the Gelfand condition on the cuspforms $f_{1}, f_{2}$, and enough decay for convergence of the Eisenstein series in $\operatorname{Re}(s)>1$. However, the explicit computation of constant terms in the non-vanishing case will need more.
[3.11.3] Theorem: Let $P=P^{r_{1}, r_{2}}$, and $f=f_{1} \otimes f_{2}$ cuspform(s) on $M=M^{P}$. Let $Q$ be another parabolic. Then $c_{Q} E_{s, f}^{P}=0$ unless $Q=P$ or $Q=P^{r_{2}, r_{1}}$, the associate of $P$.
Proof: First, since we claim that it suffices to consider maximal proper $Q$, the underlying reason being that all standard parabolics are intersections of maximal proper standard parabolics, and for standard parabolics $N^{Q \cap Q^{\prime}}=N^{Q} \cdot N^{Q^{\prime}}$. Giving $X=\left(N_{k}^{Q} \cap N_{k}^{Q^{\prime}}\right) \backslash\left(N_{\mathbb{A}}^{Q} \cap N_{\mathbb{A}}^{Q^{\prime}}\right)$ measure 1, and noting that the constant-term integrals make sense for any left $P_{k}^{\mathrm{min}}$-invariant functions,

$$
\begin{aligned}
c_{Q \cap Q^{\prime}} f(g) & =\int_{N_{k}^{Q \cap Q^{\prime}} \backslash N_{\mathbb{A}}^{Q \cap Q^{\prime}}} f(n g) d n=\int_{N_{k}^{Q \cap Q^{\prime}} \backslash N_{\mathbb{A}}^{Q \cap Q^{\prime}}} \int_{X} f(n x g) d x d n \\
& =\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \int_{N_{k}^{Q^{\prime}} \backslash N_{\mathbb{A}}^{Q^{\prime}}} \int_{X} f\left(n^{\prime} n g\right) d x d n=c_{Q}\left(c_{Q^{\prime}} f\right)(g)
\end{aligned}
$$

Next, we show that $c_{Q} E_{s, f}^{P}=0$ for maximal proper $Q$, unless $Q=P$ or its associate. Let $Q=P^{r_{1}^{\prime}, r_{2}^{\prime}}$, and write $\varphi=\varphi_{s, f}$. Take $\operatorname{Re}(s)>1$ for convergence of the series expression for $E_{s, f}^{P}$.

$$
\begin{gathered}
c_{Q} E_{s, f}^{P}(g)=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} E_{s, f}^{P}(n g) d n=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma \cdot n g) d n \\
=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\delta \in P_{k} \backslash G_{k} / Q_{k}} \sum_{\gamma \in\left(\delta^{-1} P_{\left.P_{k} \delta \cap Q_{k}\right) \backslash Q_{k}}\right.} \varphi(\delta \gamma \cdot n g) d n \\
=\sum_{\delta \in P_{k} \backslash G_{k} / Q_{k}} \int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \quad \sum_{\gamma \in\left(\delta^{-1} P_{k} \delta \cap Q_{k}\right) \backslash Q_{k}} \varphi(\delta \gamma \cdot n g) d n
\end{gathered}
$$

It certainly suffices to show that the integral vanishes for every $\delta$. The idea is that enough of the unipotent radical $N_{\mathbb{A}}^{\mathscr{Q}}$ conjugates across each $\delta \gamma$ so that the integral vanishes because of the Gelfand property of $f$. We need to understand $\delta^{-1} P_{k} \delta \cap Q_{k}$.

By the Bruhat decomposition [3.1,1], the Weyl group $W$ gives a collection of representatives for $P_{k} \backslash G_{k} / Q_{k}$. Indeed, letting $W^{P}=P_{k} \cap W$ and $W^{Q}=Q_{k} \cap W$, a set of representatives for $W^{P} \backslash W / W^{Q}$ is a set of representatives for $P_{k} \backslash G_{k} / Q_{k}$. In fact, $W^{P} \backslash W / W^{Q}$ is in bijection with $P_{k} \backslash G_{k} / Q_{k}$ by $W^{P} w W^{Q} \longleftrightarrow P_{k} w Q_{k}$, although we only proved this for $P=Q=P^{\mathrm{min}}$ in [3.1]. We determine representatives for $W^{P} \backslash W / W^{Q}$. Write a permutation matrix $w \in W$ in blocks corresponding to the Levi components of $P, Q$ :

$$
w=\left(\begin{array}{cc}
a & b \\
c & d
\end{array}\right) \quad\left(\text { with } a=r_{1} \times r_{1}^{\prime}, b=r_{1} \times r_{2}^{\prime}, c=r_{2} \times r_{1}^{\prime}, d=r_{2} \times r_{2}^{\prime}\right)
$$

The left action of the upper-left $G L_{r_{1}}$ part of $W^{P}$ inside $G L_{r_{1}} \times G L_{r_{2}} \approx M^{P}$, and the action of the upper-left $G L_{r_{1}^{\prime}}$ part of $W^{Q}$ inside $G L_{r_{1}^{\prime}} \times G L_{r_{2}^{\prime}} \approx M^{Q}$ adjust the matrix $a$ to the form $a=\left(\begin{array}{cc}1_{t_{1}} & 0 \\ 0 & 0\end{array}\right)$ for some size $t_{1}$, so $w$ becomes

$$
w=\left(\begin{array}{ccc}
1_{t_{1}} & 0 & 0 \\
0 & 0 & * \\
0 & * & *
\end{array}\right)
$$

The lower-right parts of $W^{P}$ and $W^{Q}$ further adjust the lower right block of $w$ to be of the form $\left(\begin{array}{cc}0 & 0 \\ 0 & 1_{t_{4}}\end{array}\right)$ for some $t_{4}$, putting the permutation matrix $w$ into the form

$$
w=\left(\begin{array}{cccc}
1_{t_{1}} & 0 & 0 & 0 \\
0 & 0 & * & 0 \\
0 & * & 0 & 0 \\
0 & 0 & 0 & 1_{t_{4}}
\end{array}\right)
$$

Necessarily the remaining entries can be adjusted to be identity matrices of suitable sizes. That is, $W^{P} \backslash W / W^{Q}$ has representatives of the form

$$
w=\left(\begin{array}{cccc}
1_{t_{1}} & 0 & 0 & 0 \\
0 & 0 & 1_{t_{3}} & 0 \\
0 & 1_{t_{2}} & 0 & 0 \\
0 & 0 & 0 & 1_{t_{4}}
\end{array}\right)
$$

where $t_{1}+t_{3}=r_{1}, t_{1}+t_{2}=r_{1}^{\prime}$, and so on. With suitable block sizes,

$$
\begin{aligned}
w^{-1} P w \cap Q & =\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 0 & 1
\end{array}\right)\left(\begin{array}{llll}
* & * & * & * \\
* & * & * & * \\
0 & 0 & * & * \\
0 & 0 & * & *
\end{array}\right)\left(\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \cap Q \\
& =\left(\begin{array}{llll}
* & * & * & * \\
0 & * & 0 & * \\
* & * & * & * \\
0 & * & 0 & *
\end{array}\right) \cap Q=\left(\begin{array}{llll}
* & * & * & * \\
0 & * & 0 & * \\
0 & 0 & * & * \\
0 & 0 & 0 & *
\end{array}\right)
\end{aligned}
$$

Thus, writing the sum as an iterated sum and unwinding,

$$
\begin{gathered}
\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in\left(w^{-1} P_{k} w \cap Q_{k}\right) \backslash Q_{k}} \varphi(w \gamma \cdot n g) d n \\
=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in\left(w^{-1} P_{k} w \cap M_{k}^{Q}\right) \backslash M_{k}^{Q}} \sum_{\nu \in\left((w \gamma)^{-1} P_{k} w \gamma \cap N_{k}^{Q}\right) \backslash N_{k}^{Q}} \varphi(w \gamma \nu \cdot n g) d n \\
=\sum_{\gamma \in\left(w^{-1} P_{k} w \cap M_{k}^{Q}\right) \backslash M_{k}^{Q}} \int_{\left.(w \gamma)^{-1} P_{k} w \gamma \cap N_{k}^{Q}\right) \backslash N_{\mathbb{A}}^{Q}} \varphi(w \gamma \cdot n g) d n
\end{gathered}
$$

For fixed $\gamma$, replacing $n$ by $\gamma^{-1} n \gamma$ gives

$$
\int_{\left(w^{-1} P_{k} w \cap N_{k}^{Q}\right) \backslash N_{\mathbb{A}}^{Q}} \varphi(w n \cdot \gamma g) d n=\int_{\left(w^{-1} P_{k} w \cap N_{k}^{Q}\right) \backslash N_{\mathbb{A}}^{Q}} \varphi\left(w n w^{-1} \cdot w \gamma g\right) d n
$$

A similar computation to that above shows that

$$
w N^{Q} w^{-1} \cap P=\left(\begin{array}{cccc}
1_{t_{1}} & * & 0 & * \\
0 & 1_{t_{3}} & 0 & 0 \\
0 & 0 & 1_{t_{2}} & * \\
0 & 0 & 0 & 1_{t_{4}}
\end{array}\right)
$$

This contains the unipotent radical $N^{\prime}$ of the parabolic $P^{\prime}=P^{t_{1}, t_{3}} \times P^{t_{2}, t_{4}}$ of the Levi component $M^{P} \approx G L_{r_{1}} \times G L_{r_{2}}$ of $P$. Unless $\left(r_{1}^{\prime}, r_{2}^{\prime}\right)=\left(r_{1}, r_{2}\right)$ or $\left(r_{1}^{\prime}, r_{2}^{\prime}\right)=\left(r_{2}, r_{1}\right)$, at least one of those parabolic subgroups of $G L_{r_{j}}$ must be a proper parabolic of the corresponding $G L_{r_{j}}$. That is, for each fixed $\gamma$ the integral over $\left(w^{-1} P_{k} w \cap N_{k}^{Q}\right) \backslash N_{\mathbb{A}}^{Q}$ has a subintegral over $N_{k}^{\prime} \backslash N_{\mathbb{A}}^{\prime}$, which computes the $P^{\prime}$ constant term of the cuspidal data $f$, giving 0 .

This almost gives the vanishing assertion of the theorem. One anomalous case remains, namely, $P \cap Q$ when $P=P^{r_{1}, r_{2}}$ and $Q=P^{r_{2}, r_{1}}$ with $r_{1} \neq r_{2}$. Still, use the fact that $c_{Q \cap P}=c_{Q} \circ c_{P}$. Compute the constant term along $P$, using the fact from above that only $w=1$ gives a non-zero outcome. Thus, for non-self-associate proper maximal $P$ and cuspidal data,

$$
\begin{gathered}
c_{P} E_{s, f}^{P}(g)=\int_{\left(P_{k} \cap N_{k}^{Q}\right) \backslash N_{\mathbb{A}}^{Q}} \varphi(n g) d n=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \varphi_{s, f}(n g) d n \\
=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \varphi_{s, f}(g) d n=\varphi_{s, f}(g) \cdot \int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} 1 d n
\end{gathered}
$$

Because $r_{1} \neq r_{2}, N^{Q} \cap N^{P}$ contains a unipotent radical of some proper parabolic in $M^{P}$, so the cuspidality of $f$ means $c_{Q} \varphi_{s, f}=0$. Thus, $c_{Q \cap P} E_{s, f}^{P}=0$.
[3.11.4] Remark: More generally, Eisenstein series with cuspidal data for parabolics $P=P^{r_{1}, \ldots, r_{\ell}}$ have constant term 0 along parabolics $Q$ unless $Q$ contains some associate of $P$, that is, contains some $P^{r_{1}^{\prime}, \ldots, r_{\ell}^{\prime}}$ with the $r_{i}^{\prime}$ 's a permutation of the $r_{i}$ 's.

The same arguments and vanishing conclusions apply to constant terms of pseudo- Eisenstein series with cuspidal data:
[3.11.5] Corollary: For maximal proper $P$ and cuspidal data $f$ on $M^{P}$, for another parabolic $Q, c_{Q} \Psi_{\eta, f}^{P}=0$ unless $Q$ is associate to $P$.

An optimist would have to hope that cuspidal-data Eisenstein series $E_{s, f}^{P}$ formed from spherical Hecke eigenfunction cuspforms $f=f_{1} \otimes f_{2}$ would itself be a spherical Hecke eigenfunction, and that this is so because $\varphi_{s, f}$ is a spherical Hecke eigenfunction for all $s$. Happily, this is nearly true, with a yet-stronger notion of cuspform, as follows. Fix a non-archimedean $v$ and square-integrable right $K_{v}$-invariant cuspform $f$, and consider the space

$$
\pi_{v}=\left\{\text { finite linear combinations of right translates } g \rightarrow f(g h) \text { with } h \in G_{v}\right\}
$$

generated by $f$ under the action of $G_{v}$ by right translation, suitably topologized. The most direct way to begin description of a suitable requirement ${ }^{[30]}$ on $f$ at $v$ is that $\pi_{v}$ be isomorphic as $G_{v}$ representation to a

[^8]$G_{v^{-}}$subrepresentation of an unramified principal series attached to the standard minimal parabolic $B$ :
$$
I_{\chi}^{B}=\left\{F \in C^{\infty}\left(G_{v}\right): F(b g)=\chi(b) \cdot F(g) \text { for all } b \in B_{v}, g \in G_{v}\right\} \quad \text { (for unramified } \chi \text { on } M_{v}^{B} \text { ) }
$$

Since $f$ is $K_{v}$-fixed, its image in $I_{\chi}^{B}$ contains the subspace of $K_{v}$-fixed vectors, which is one-dimensional by the Iwasawa decomposition $G_{v}=B_{v} \cdot K_{v}$. This is the right feature to prove
[3.11.6] Theorem: Fix a non-archimedean place $v$. With $f_{1}$ and $f_{2}$ as just described, that is, under right translation by $G L_{r_{j}}\left(k_{v}\right)$ generating representations isomorphic to subrepresentations of unramified principal series representations of $G L_{r_{j}}\left(k_{v}\right)$, with trivial central characters, the function $\varphi_{s, f}$ is a spherical Heckealgebra eigenfunction for $G_{v}$. In the region of convergence, $E_{s, f}^{P}$ is a spherical Hecke-algebra eigenfunction with the same eigenvalues as $\varphi_{s, f}$. (Proof in [8.5].)
[3.11.7] Remark: Quantitative details about the spherical Hecke eigenvalues of $\varphi_{s, f}$ and $E_{s, f}^{P}$ for $G_{v}$ in terms of the spherical Hecke eigenvalues of $f=f_{1} \otimes f_{2}$ for $M_{v}$ and $s \in \mathbb{C}$ are visible in the proof [8.5].
[3.11.8] Remark: In the cases of $c_{Q} E_{s, f}^{P} \neq 0$, the proof above shows that non-vanishing occurs only in a few cases: $w=1_{r}$ for $Q=P$ always gives the summand $\varphi_{s, f}^{P}$ of the constant term, and $w=\left(\begin{array}{cc}0 & 1_{r_{1}} \\ 1_{r_{2}} & 0\end{array}\right)$ for $Q=P^{r_{2}, r_{1}}$ for both $r_{1}=r_{2}$ and $r_{1} \neq r_{2}$. Happily, in both these cases, $\left.w^{-1} P_{k} w \cap M_{k}^{Q}\right)=M_{k}^{Q}$, so the sum over $\gamma \in\left(w^{-1} P_{k} w \cap M_{k}^{Q}\right) \backslash M_{k}^{Q}$ is trivial. Thus, for $Q=P$ and $w=1_{r}$, that part of the constant term is easily made explicit, as in the proof above:

$$
c_{P} E_{s, f}^{P}(g)=\varphi_{s, f}(g) \cdot \int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} 1 d n
$$

The other part of the constant terms is significantly more complicated, as follows. With or without $r_{1}=r_{2}$, when $Q=P^{r_{2}, r_{1}}$ and $w=\left(\begin{array}{cc}0 & 1_{r_{1}} \\ 1_{r_{2}} & 0\end{array}\right)$, that part of the constant term is unwound completely, since $w^{-1} P_{k} w \cap N_{k}^{Q}=\{1\}$, so

$$
\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} E_{s, f}^{P}(n g) d n=\int_{N_{\mathbb{A}}^{Q}} \varphi_{s, f}(n g) d n
$$

Since we have supposed that $f$ is right $K_{\mathbb{A}}$-invariant, the integral produces a left $Z^{+} N_{\mathbb{A}} M_{k}^{Q}$-invariant, right $K_{\mathbb{A}}$-invariant function, so by Iwasawa is a function on $M_{k}^{Q} \backslash M_{\mathbb{A}}^{Q}$ and right $M_{\mathbb{A}} \cap K_{\mathbb{A}}$-invariant. The behavior under the center of $M_{\mathbb{A}}^{Q}$ is also easy to assess by changing variables in the integer. Thus, it is reasonable to imagine that it is of the form $\varphi_{1-s, f^{\prime}}^{Q}$ for some cuspform(s) on $M^{Q}$. However, we would not want $f^{\prime}$ to depend on $s$, so the dependence on $s$ should be somehow separate, and this integral should be expressible as $c_{s, f} \cdot \varphi_{1-s, f^{\prime}}^{Q}(m)$ with cuspform(s) on $M^{Q}$ independent of $s$, and $m \in M_{\mathbb{A}}^{Q}$.

This conclusion does hold, but only with substantial assumptions on $f \approx f_{1} \otimes f_{2}$, as follows. Continue to assume that $f$ is a spherical Hecke algebra eigenfunction on $M^{P}$ for all non-archimedean $M_{v}^{P}$. The best further simplifying hypothesis ${ }^{[31]}$ is a form of strong multiplicity one, that the only other cuspforms on $M^{P} \approx G L_{r_{1}} \times G L_{r_{2}}$ with the same spherical Hecke eigenvalues at all finite primes are scalar multiples of $f=f_{1} \otimes f_{2}$. Let $f^{w}=\left(f_{1} \otimes f_{2}\right)^{w}=f_{2} \otimes f_{1}$.
[3.11.9] Theorem: In the non-vanishing cases, with maximal proper $P$, and $Q=P$ or its associate, with the strong multiplicity one assumption above,

$$
\left\{\begin{aligned}
c_{P} E_{s, f}^{P} & =\varphi_{s, f}^{P} & & \left(\text { for } r_{1} \neq r_{2}(\text { not self-associate })\right) \\
c_{P} E_{s, f}^{P} & =\varphi_{s, f}^{P}+c_{s, f}^{P} \varphi_{1-s, f^{w}}^{P} & & \left(\text { for } r_{1}=r_{2}(\text { self-associate }), \text { meromorphic } c_{s, f}^{P}\right) \\
c_{Q} E_{s, f}^{P} & =c_{s, f}^{Q} \cdot \varphi_{1-s, f}^{Q} & & \left(\text { for } r_{1} \neq r_{2}, Q=P^{r_{2}, r_{1}}, \text { meromorphic } c_{s, f}^{Q}\right)
\end{aligned}\right.
$$

[31] The strong multiplicity one assumption convenient for $G L_{r}$ is a theorem of [Shalika 1974]. See also [PiatetskiShapiro 1977]. It apparently does not hold for most other groups. That is, in general, only more complicated conclusions can be reached about unwound integrals appearing in constant term computations.

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

Proof: From the proof of [3.11.3], due to the cuspidal data $f$, most summands of $c_{Q} E_{s, f}^{P}$ corresponding to double cosets $P_{k} w Q_{k} \in P_{k} \backslash G_{k} / Q_{k}$ are 0. In all cases, for $P=Q$, the small Bruhat cell $P=P \cdot 1 \cdot Q$ gives contribution

$$
\begin{aligned}
& \int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in\left(w^{-1} P_{k} w \cap Q_{k}\right) \backslash Q_{k}} \varphi(w \gamma \cdot n g) d n=\int_{N_{k}^{P} \backslash N_{\mathbb{A}}^{P}} \varphi(n g) d n \\
& =\varphi(g) \int_{N_{k}^{P} \backslash N_{\mathbb{A}}^{P}} 1 d n \quad \text { (with } w=1 \text { and } P=Q \text { ) }
\end{aligned}
$$

since $\varphi$ is left $N_{\mathbb{A}}$-invariant.
Now consider $P=P^{r_{1}, r_{2}}$ and $Q=P^{r_{2}, r_{1}}$. From the end of the proof of [3.11.3], a double coset $P w Q$ can give a non-zero contribution to the constant term only if $w N^{Q} w^{-1} \cap P$ contains no unipotent radical of a proper parabolic of the Levi component $M^{P}$ of $P$. As in the proof of [3.11.3], $P_{k} \backslash G_{k} / Q_{k} \approx W^{P} \backslash W / W^{Q}$ has representatives

$$
w=\left(\begin{array}{cccc}
1_{t_{1}} & 0 & 0 & 0 \\
0 & 0 & 1_{t_{3}} & 0 \\
0 & 1_{t_{2}} & 0 & 0 \\
0 & 0 & 0 & 1_{t_{4}}
\end{array}\right)
$$

where $t_{1}+t_{3}=r_{1}, t_{1}+t_{2}=r_{1}$, and so on. The only case other than $w=1$ meeting the condition is with $t_{1}=0=t_{4}$ and $t_{3}=r_{1}, t_{2}=r_{2}$. This has the effect that $w^{-1} P w \cap Q=M^{Q}$. This summand in the constant term unwinds completely:

$$
\begin{gathered}
\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in w^{-1}} \varphi(w \gamma \cdot n g) d n=\int_{P_{k} w \cap Q_{k} \backslash Q_{k}} \sum_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in M_{k}^{Q} \backslash Q_{k}} \varphi(w \gamma \cdot n g) d n \\
=\int_{N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}} \sum_{\gamma \in N_{k}^{Q}} \varphi(w \gamma \cdot n g) d n=\int_{N_{\mathbb{A}}^{Q}} \varphi(w \cdot n g) d n
\end{gathered}
$$

Let $\varphi^{\prime}(g)$ be the latter integral. With $\varphi$ right $K_{v}$-invariant at all places, to understand $\varphi^{\prime}(g)$ it suffices to take

$$
g=m=\left(\begin{array}{cc}
m_{2} & 0 \\
0 & m_{1}
\end{array}\right)=w^{-1}\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) w
$$

by the Iwasawa decomposition.
Certainly $\varphi^{\prime}$ is left $N_{\mathbb{A}}^{Q}$-invariant and invariant under the center. It is left $M_{k}^{Q}$-invariant, since for $\gamma \in M_{k}^{Q}$

$$
\varphi^{\prime}(\gamma m)=\int_{N_{\mathbb{A}}^{Q}} \varphi(w \cdot n \gamma g) d n=\int_{N_{\mathbb{A}}^{Q}} \varphi(w \cdot \gamma n g) d n=\int_{N_{\mathbb{A}}^{Q}} \varphi\left(w \gamma w^{-1} \cdot w n g\right) d n=\int_{N_{\mathbb{A}}^{Q}} \varphi(w n g) d n=\varphi^{\prime}(m)
$$

by changing variables in the integral, and observing that the change-of-measure is 1 , by the product formula. Since the right translation action commutes with the integration along $N_{\mathbb{A}}^{Q}$ on the left, $m_{1} \times m_{2} \rightarrow \varphi^{\prime}(m)$ is a spherical Hecke eigenfunction on $G L\left(r_{1}\right) \times G L\left(r_{2}\right)$ with the same eigenvalues as $\varphi$.

To see the behavior of the $s$ parameter, it suffices to consider left translation by

$$
h=\left(\begin{array}{cc}
t \cdot 1_{r_{2}} & 0 \\
0 & 1_{r_{1}}
\end{array}\right)=w^{-1}\left(\begin{array}{cc}
1_{r_{1}} & 0 \\
0 & t \cdot 1_{r_{2}}
\end{array}\right) w
$$

with $t>0$ imbedded diagonally at archimedean places. Then
$\varphi^{\prime}(h m)=\int_{N_{\mathbb{A}}} \varphi(w \cdot n \cdot h m) d n=\left|\operatorname{det}\left(t \cdot 1_{r_{2}}\right)\right|^{r_{1}} \int_{N_{\mathbb{A}}^{Q}} \varphi(w \cdot h \cdot n m) d n=\left|\operatorname{det}\left(t \cdot 1_{r_{2}}\right)\right|^{r_{1}} \int_{N_{\mathbb{A}}} \varphi\left(w h w^{-1} \cdot w n m\right) d n$
by replacing $n$ by $h n h^{-1}$, picking up the indicated change-of-measure. The left equivariance of $\varphi$ under elements of the form $w h w^{-1}$ is

$$
\varphi\left(w h w^{-1} \cdot w n m\right)=\left|1 / \operatorname{det}\left(t \cdot 1_{r_{2}}\right)^{r_{1}}\right|^{s} \cdot \varphi(w n m)
$$

Thus,

$$
\varphi^{\prime}(h m)=\left|\operatorname{det}\left(t \cdot 1_{r_{2}}\right) / 1\right|^{1-s} \cdot \varphi^{\prime}(m)
$$

as claimed in the assertion of the theorem.
Finally, the multiplicity-one assumption says that $m_{1} \times m_{2} \rightarrow \varphi^{\prime}(m)$ must be a scalar multiple $c_{s, f}^{P}$ of $\varphi(m)$. Since $s \rightarrow E_{s, f}^{P}$ is a meromorphic smooth-function-valued function of $s$, composition with $c^{Q}$ gives a meromorphic smooth-function-valued function of $s$. Since it differs by the scalar $c_{s, f}^{Q}$ from $\varphi_{s, f}$, this scalar must be meromorphic in $s$.
[3.11.10] Remark: The meromorphic functions $c_{s, f}$ have Euler product expansions attached to $f_{1}$ and $f_{2}$, but we do not have immediate need of this fact. [32]

In parallel with the spherical Hecke algebra behavior of $E_{s, f}^{P}$ at finite places, keeping the assumption that $f_{1}, f_{2}$ have trivial central character and are right $K_{\mathbb{A}}$-invariant,
[3.11.11] Theorem: For $v$ archimedean, for $f=f_{1} \otimes f_{2}$ with $f_{1}, f_{2}$ eigenfunctions for the invariant Laplacians on the factors $G L_{r_{1}}\left(k_{v}\right)$ and $G L_{r_{2}}\left(k_{v}\right)$ of $M_{v}^{P}$, the function $\varphi_{s, f}^{P}$ is an eigenfunction for the invariant Laplacian on $G_{v}$, and, thus, $E_{s, f}^{P}$ is also an eigenfunction. In particular, letting $\lambda_{j}$ be the eigenvalue of $f_{j}$,

$$
\Omega \cdot E_{s, f}^{P}=\left(r_{1} r_{2}\left(r_{1}+r_{2}\right)\left(s^{2}-s\right)+\lambda_{1}+\lambda_{2}\right) \cdot E_{s, f}^{P}
$$

In particular, the eigenvalue is invariant under $s \longrightarrow 1-s$.
Proof: For simplicity, treat $G_{v} \approx G L_{r}(\mathbb{R})$. Accommodations for the complex case are illustrated in [4.6]. This is a purely local issue, and it suffices to consider arbitrary functions $f_{1} \otimes f_{2}$ on $G L_{r_{1}}(\mathbb{R}) \times G L_{r_{2}}(\mathbb{R})$ with trivial central character. That is, the possibility that $f_{1}, f_{2}$ are automorphic forms of any sort is irrelevant. Similarly, we have a purely locally defined function

$$
\varphi_{s}\left(\left(\begin{array}{cc}
a & * \\
0 & d
\end{array}\right) k\right)=\left|\frac{(\operatorname{det} a)^{r_{2}}}{(\operatorname{det} d)^{r_{1}}}\right|^{s} \cdot f_{1}(a) \cdot f_{2}(d) \quad\left(\text { for } a \in G L_{r_{1}}(\mathbb{R}), d \in G L_{r_{2}}(\mathbb{R}), k \in O(r, \mathbb{R})\right)
$$

As in [4.2] and [4.4], the invariant Laplacian on $G_{v} / K_{v}$ is Casimir $\Omega$ on $G_{v}$ descended to that quotient, and then to any further quotient. For any choice of basis $\left\{x_{i}\right\}$ of the Lie algebra $\mathfrak{g}$ of $G_{v}$, and and dual basis $\left\{x_{i}^{*}\right\}$ with respect to the pairing $\langle x, y\rangle=\operatorname{tr}(x y)$, Casimir is expressible as an element in the (center of the) universal enveloping algebra [4.2], [4.3] as $\Omega=\sum_{i} x_{i} x_{i}^{*}$. It is easy to exhibit a basis so that the summands separate into three pieces: Casimir $\Omega_{1}$ of $G L_{r_{1}}(\mathbb{R})$ acting only on $f_{1}$, Casimir $\Omega_{2}$ of $G L_{r_{2}}(\mathbb{R})$ acting only on $f_{2}$, and a leftover acting only on $\left|(\operatorname{det} a)^{r_{2}} /(\operatorname{det} d)^{r_{1}}\right|^{s}$.

Let $h_{i}$ be the diagonal matrix with 1 at the $i^{\text {th }}$ place 0 's otherwise. For $i<j$, let $x_{i j}$ be the matrix with a unique non-zero entry, a 1 , at the $i j^{t h}$ location, and for $i>j$ let $y_{i j}$ be the matrix with a unique non-zero entry, a 1 , at the $i j^{t h}$ location. The $\left\{x_{i j}\right\}$ and $\left\{y_{j i}\right\}$ are dual under $\langle$,$\rangle . Thus, the Casimir operators for$ $G L_{r_{1}}(\mathbb{R})$ and $G L_{r_{2}}(\mathbb{R})$ are

$$
\Omega_{1}=\sum_{i=1}^{r_{1}} h_{i}^{2}+\sum_{i<j \leq r_{1}}\left(x_{i j} y_{j i}+y_{j i} x_{i j}\right) \quad \Omega_{2}=\sum_{i=r_{1}+1}^{r} h_{i}^{2}+\sum_{r_{1}<j \leq r}\left(x_{i j} y_{j i}+y_{j i} x_{i j}\right)
$$

and $\Omega=\Omega_{1}+\Omega_{2}+\Omega^{\prime}$ with leftover

$$
\Omega^{\prime}=\sum_{1 \leq i \leq r_{1}, r_{1}<j \leq r}\left(x_{i j} y_{j i}+y_{j i} x_{i j}\right)
$$

[32] [Langlands 1971] considers consequences of the appearance of Euler products in constant terms of Eisenstein series series. A part of that program is completed in [Shahidi 1978] and [Shahidi 1985].

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

Use the fact that Casimir commutes with conjugation by $G_{v}$, so we can let the associated differential operators [4.1] act on the left on left $N_{v}^{P}$-invariant functions such as $\varphi_{s}$, so that the Lie algebra $\mathfrak{n}$ of $N^{P}$ annihilates such functions. There is a sign or order-of-operations issue: for a smooth function $\varphi$ on $G_{v}$, the effect of Casimir acting on the right is

$$
\Omega \varphi(g)=\sum_{i} \frac{\partial}{\partial t_{1}} \frac{\partial}{\partial t_{2}} \varphi\left(g \cdot e^{t_{1} x_{i}} \cdot e^{t_{2} x_{i}^{*}}\right)
$$

Invoking the invariance under conjugation by $G_{v}$, this is

$$
\Omega \varphi(g)=\sum_{i} \frac{\partial}{\partial t_{1}} \frac{\partial}{\partial t_{2}} \varphi\left(e^{t_{1} x_{i}} \cdot e^{t_{2} x_{i}^{*}} g\right)
$$

Thus, terms $x_{i j} y_{j i}$ with $1 \leq i \leq r_{1}$ and $r_{1}<j \leq r$ annihilate $\varphi_{s}$, because, after conjugating, $x_{i j}$ acts first and is in $\mathfrak{n}$. For $1 \leq i \leq r_{1}$ and $r_{1}<j \leq r$, we can move toward invocation of this annihilation property by noting that $x_{i j} y_{j i}-y_{i j} x^{i j}=h_{i}-h_{j}$, so

$$
x_{i j} y_{j i}+y_{j i} x_{i j}=2 x_{i j} y_{j i}+y_{j i} x_{i j}-x_{i j} y_{j i}=2 x_{i j} y_{j i}-\left[x_{i j}, y_{j i}\right]=2 x_{i j} y_{j i}-h_{i}+h_{j}
$$

which acts just by $-h_{i}+h_{j}$ on left $N_{v}^{P}$-invariant functions:

$$
\Omega^{\prime} \cdot \varphi_{s}=\sum_{1 \leq i \leq r_{1}, r_{1}<j \leq r}\left(x_{i j} y_{j i}+y_{j i} x_{i j}\right) \cdot \varphi_{s}=\sum_{1 \leq i \leq r_{1}, r_{1}<j \leq r}\left(-h_{i}+h_{j}\right) \cdot \varphi_{s}=-r_{2} \sum_{1 \leq i \leq r_{1}} h_{i} \cdot \varphi_{s}+r_{1} \sum_{r_{1}<j \leq r} h_{j} \cdot \varphi_{s}
$$

Thus, with $z_{1}=\sum_{1 \leq i \leq r_{1}} h_{i}$ and $z_{2}=\sum_{r_{1}<j \leq r} h_{j}$,

$$
\Omega \cdot \varphi_{s}=\left(\Omega_{1}-r_{2} z_{1}\right) \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s} f_{1}\left(m_{1}\right)+\left(\Omega_{2}+r_{1} z_{2}\right) \cdot\left|\operatorname{det} m_{2}\right|^{r_{1} s} f_{2}\left(m_{2}\right)
$$

Since $z_{1}$ is in the Lie algebra of the center of the $G L_{r_{1}}$ factor of $M^{P}$, and $f_{1}$ has trivial central character, $z_{1} \cdot f_{1}=0$, and

$$
\left(\Omega_{1}-r_{2} z_{1}\right) \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s} f_{1}\left(m_{1}\right)=\Omega_{1} \cdot\left(\left|\operatorname{det} m_{1}\right|^{r_{2} s} f_{1}\left(m_{1}\right)\right)-\left(r_{2} z_{1} \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s}\right) f_{1}\left(m_{1}\right)
$$

and similarly for $\Omega_{2}+r_{1} z_{2}$. The effect of $z_{1}$ on that power of determinant is straightforward:

$$
z_{1} \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s}=\left.\frac{\partial}{\partial t}\right|_{t=0}\left|\operatorname{det}\left(e^{t z_{1}} \cdot m_{1}\right)\right|^{r_{2} s}=\left.\frac{\partial}{\partial t}\right|_{t=0}\left(e^{t}\right)^{r_{1} r_{2} s} \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s}=r_{1} r_{2} s \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s}
$$

Similarly, $z_{2} \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1} s}=-r_{1} r_{2} s \cdot\left|\operatorname{det} m_{2}\right|^{-r_{1} s}$. Thus, the $-r_{2} z_{1}+r_{1} z_{2}$ terms combine to $-r_{1} r_{2}\left(r_{1}+\right.$ $\left.r_{2}\right) s \cdot \varphi_{s}$.

The effect of $\Omega_{1}$ on $f_{1}\left(m_{1}\right)$ adjusted by a power of determinant is only slightly more complicated, using Leibniz' rule. The terms $x_{i j} y_{j i}$ and $y_{j i} x_{i j}$ annihilate the determinant, and $h_{i} \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s}=r_{2} s \cdot\left|\operatorname{det} m_{1}\right|^{r_{2} s}$, so

$$
\begin{gathered}
\Omega_{1}\left(\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot f_{1}\left(m_{1}\right)\right)=\sum_{1 \leq i \leq r_{1}} h_{1}^{2} \cdot\left(\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot f_{1}\left(m_{1}\right)\right)+\left|\operatorname{det} m_{1}\right|^{r_{2} s} \sum_{1 \leq i<j \leq r_{1}}\left(x_{i j} y_{j i}+y_{j i} x_{i j}\right) \cdot f_{1}\left(m_{1}\right) \\
=\sum_{1 \leq i \leq r_{1}}\left(h_{i}^{2}\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot f_{1}\left(m_{1}\right)+2 h_{i}\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot h_{i} f_{1}\left(m_{1}\right)\right)+\left|\operatorname{det} m_{1}\right|^{r_{2} s} \Omega_{1} f_{1}\left(m_{1}\right) \\
=r_{1}\left(r_{2} s\right)^{2}\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot f_{1}\left(m_{1}\right)+2 r_{2} s\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot\left(\sum_{1 \leq i \leq r_{1}} h_{i}\right) \cdot f_{1}\left(m_{1}\right)+\left|\operatorname{det} m_{1}\right|^{r_{2} s} \Omega_{1} f_{1}\left(m_{1}\right) \\
=r_{1}\left(r_{2} s\right)^{2}\left|\operatorname{det} m_{1}\right|^{r_{2} s} \cdot f_{1}\left(m_{1}\right)+0+\left|\operatorname{det} m_{1}\right|^{r_{2} s} \Omega_{1} f_{1}\left(m_{1}\right)=\left|\operatorname{det} m_{1}\right|^{r_{2} s}\left(r_{1}\left(r_{2} s\right)^{2}+\Omega_{1}\right) \cdot f_{1}\left(m_{1}\right)
\end{gathered}
$$

since $\sum_{i} h_{i}$ annihilates $f_{1}$ due to the latter's trivial central character. A similar computation applies to $\Omega_{2}$ and $f_{2}$. Letting $\Omega_{1} f_{1}=\lambda_{1} \cdot f$ and $\Omega_{2} f_{2}=\lambda_{2} \cdot f_{2}$, these computations give

$$
\Omega \cdot \varphi_{s}=\left(r_{1} r_{2}\left(r_{1}+r_{2}\right)\left(s^{2}-s\right)+\lambda_{1}+\lambda_{2}\right) \cdot \varphi_{s}
$$

That is, $\varphi_{s}$ is an eigenfunction for Casimir on $G_{v}$, and by the invariance so is $E_{s, f}^{P}$.
[3.11.12] Remark: The argument can be recast as an application of induction in stages, in the archimedean case, analogous to the corresponding non-archimedean argument [6.9]. [33]

### 3.12 Continuation of minimal-parabolic Eisenstein series

We show that the meromorphic continuations of some simple types of minimal-parabolic Eisenstein series on $G L_{r}$ follow from the $G L_{2}$ case [2.B] via Bochner's Lemma [3.B]. That determines the $r$ ! functional equations corresponding to elements of the Weyl group $W$, the latter identified with permutation matrices in $G L_{r}$. We can also use this to compute the minimal-parabolic constant terms. To illustrate the points with minimal clutter, we consider just the simplest Eisenstein series

$$
E_{s}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s}^{o}(\gamma \cdot g) \quad\left(\text { where } \varphi_{s}^{o}(n m k)=\left|m_{1}\right|^{s_{1}}\left|m_{2}\right|^{s_{2}} \ldots\left|m_{r}\right|^{s_{r}} \text { with } s_{1}+\ldots+s_{r}=0\right)
$$

with $P=P^{\min }, n \in N_{\mathbb{A}}=N^{\min }, m \in M_{\mathbb{A}}^{\min }$, and $k \in K_{\mathbb{A}}$. Let $\xi(s)$ be the completed zeta function of the underlying number field. For $s \in \mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$, write $w \cdot s$ for the action of $w \in W$, that is, $\left(w m w^{-1}\right)^{s}=m^{w \cdot s}$. In the following, because the fixed point in $\mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$ of all the functional equations turns out to be the half-sum

$$
\rho=\left(\rho_{1}, \ldots, \rho_{r}\right)=\left(\frac{r-1}{2}, \frac{r-3}{2}, \frac{r-5}{2}, \ldots, \frac{3-r}{2}, \frac{1-r}{2}\right)
$$

of positive roots, we will express the functional equation in terms of

$$
E_{\rho+s}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{\rho+s}^{o}(\gamma \cdot g) \quad\left(\text { where } \varphi_{\rho+s}^{o}(n m k)=\left|m_{1}\right|^{\rho_{1}+s_{1}}\left|m_{2}\right|^{\rho_{2}+s_{2}} \ldots\left|m_{r}\right|^{\rho_{r}+s_{r}}\right)
$$

We prove the following theorem and the corollary together.
[3.12.1] Theorem: (Selberg, Langlands, et alia) Minimal-parabolic Eisenstein series $E_{s}$ have meromorphic continuations in $s \in \mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$, with functional equations

$$
E_{\rho+w \cdot s}=c_{w, s}^{-1} \cdot E_{\rho+s}
$$

[3.12.2] Corollary: The meromorphic continuation of $E_{\rho+s}$ is holomorphic for $s$ off the zero-sets of $\xi(\langle s, \beta\rangle+1), 1 \pm\langle s, \beta\rangle$, and $\langle s, \beta\rangle$, for positive roots $\beta$.
[3.12.3] Corollary: $c_{\tau, s}=\xi\langle s, \alpha\rangle / \xi(1+\langle s, \alpha\rangle)$ for reflections $\tau$, and the cocycle relation $c_{w^{\prime}, w \cdot s} \cdot c_{w, s}=c_{w w^{\prime}, s}$ holds for $w, w^{\prime} \in W$ and $s \in \mathfrak{a}^{*} \otimes_{R} \mathbb{C}$.

Proof: In brief, the idea is to view the minimal-parabolic Eisenstein series as an iterated object, variously as an Eisenstein series for all the next-to-minimal parabolics $Q^{i}=P^{1, \ldots, 1,2,1, \ldots, 1}$ with the 2 at the $i^{t h}$ place, formed from data including a suitably normalized $G L_{2}$ Eisenstein series $\widetilde{E}$ on the Levi component factor $G L_{2}$ of $Q^{i}$, rather than cuspidal data on that $G L_{2}$. Phragmén-Lindelöf gives boundedness of the analytic continuation of $\widetilde{E}$ in vertical strips, yielding convergence of the $Q^{i}$ Eisenstein series in a larger region

$$
\Omega_{i}=\left\{s \in \mathbb{C}^{n}: \operatorname{Re}\left(s_{j}\right)-\operatorname{Re}\left(s_{j+1}\right)>2 \text { for } j \neq i\right\}
$$

[33] The archimedean analogue of the Borel-Casselman-Matsumoto result [Borel 1976], [Casselman 1980], [Matsumoto 1977] is the sharper subrepresentation theorem [Casselman Miličić 1982], considerably improving the subquotient theorem of [Harish-Chandra 1954].

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

That is, in $\Omega_{i}$ there is no constraint on $\operatorname{Re}\left(s_{i}\right)-\operatorname{Re}\left(s_{i+1}\right)$. This applies to all the 2-by- 2 blocks along the diagonal, giving a meromorphic continuation of $E_{s}$ to $\bigcup_{i} \Omega_{i}$. Then Bochner's lemma [3.B] analytically continues the whole Eisenstein series to the convex hull of $\bigcup_{i} \Omega_{i}$, namely, $\mathbb{C}^{n}$.

Partial analytic continuation: Let $P=P^{\text {min }}$. For each fixed index $1 \leq i<r$, there is the next-to-minimal standard parabolic $Q=Q_{i}$ with standard Levi components and unipotent radicals given by $Q=N^{Q} \cdot M^{Q}$ with

$$
N^{Q}=\left(\begin{array}{cccccccc}
1 & * & * & * & * & * & * & * \\
& \ddots & * & * & * & * & * & * \\
& & 1 & * & * & * & * & * \\
& & & 1 & 0 & * & * & * \\
& & & 0 & 1 & * & * & * \\
& & & & & 1 & * & * \\
& & & & & & \ddots & * \\
& & & & & & 1
\end{array}\right) \quad M^{Q}=\left(\begin{array}{cccccccc}
* & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
& \ddots & 0 & 0 & 0 & 0 & 0 & 0 \\
& & * & 0 & 0 & 0 & 0 & 0 \\
& & & * & * & 0 & 0 & 0 \\
& & & & * & 0 & 0 & 0 \\
& & & & & & & 0 \\
& & & & & & & 0 \\
&
\end{array}\right)
$$

with the anomalous block at the $(i, i),(i, i+1),(i+1, i)$, and $(i+1, i+1)$ positions. The minimal-parabolic Eisenstein series can be written as an iterated sum

$$
E_{s}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s}^{o}(\gamma g)=\sum_{\gamma \in Q_{k} \backslash G_{k}}\left(\sum_{\delta \in P_{k} \backslash Q_{k}} \varphi_{s}^{o}(\delta \gamma g)\right)
$$

The quotient $P_{k} \backslash Q_{k}$ has representatives

$$
M_{k}^{P} \backslash M_{k}^{Q} \approx\left\{\delta=\left(\begin{array}{cccccccc}
1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\
& \ddots & 0 & 0 & 0 & 0 & 0 & 0 \\
& & 1 & 0 & 0 & 0 & 0 & 0 \\
& & & a & b & 0 & 0 & 0 \\
& & & c & d & 0 & 0 & 0 \\
& & & & & 1 & 0 & 0 \\
& & & & & & \ddots & 0 \\
& & & & & & & 1
\end{array}\right):\left(\begin{array}{cc}
a & b \\
c & d
\end{array}\right) \in P_{k}^{1,1} \backslash G L_{2}(k)\right\} \approx P_{k}^{1,1} \backslash G L_{2}(k)
$$

where $P^{1,1}$ is the standard upper-triangular parabolic in $G L_{2}$. Further,

$$
\begin{aligned}
& \quad\left(\begin{array}{cccccccc}
a_{1} & * & * & * & * & * & * & * \\
& \ddots & * & * & * & * & * & * \\
& & a_{i-1}^{o} & * & * & * & * & * \\
& & & a_{i} & * & * & * & * \\
& & & 0 & a_{i+1} & * & * & * \\
& & & & & a_{i+2} & * & * \\
& & & & & & \ddots & * \\
& & & & & & & a_{r}
\end{array}\right)=\left|a_{1}\right|^{s_{1}} \ldots\left|a_{r}\right|^{s_{r}} \\
& =\left|a_{1}\right|^{s_{1}} \ldots\left|a_{i-1}\right|^{s_{i-1}}\left|a_{i} / a_{i+1}\right|^{\frac{s_{i}-s_{i+1}}{2}}\left|a_{i} a_{i+1}\right|^{\frac{s_{i}+s_{i+1}}{2}}\left|a_{i+2}\right|^{s_{i+2}} \ldots\left|a_{r}\right|^{s_{r}}
\end{aligned}
$$

Thus, the inner sum in the expression for $E_{s}$ is

$$
\sum_{\delta \in P_{k} \backslash Q_{k}} \varphi_{s}^{o}\left(\delta \cdot\left(\begin{array}{cccccccc}
a_{1} & * & * & * & * & * & * & * \\
& \ddots & * & * & * & * & * & * \\
& & a_{i-1} & * & * & * & * & * \\
& & & a & b & * & * & * \\
& & & c & d & * & * & * \\
& & & & & a_{i+2} & * & * \\
& & & & & & \ddots & * \\
& & & & & & & a_{r}
\end{array}\right)\right.
$$

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$$
=\left|a_{1}\right|^{s_{1}} \ldots\left|a_{i-1}\right|^{s_{i-1}} \cdot E_{\frac{s_{i}-s_{i+1}}{2}}^{1,1}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left|\operatorname{det}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)\right|^{\frac{s_{i}+s_{i+1}}{2}} \cdot\left|a_{i+2}\right|^{s_{i+2}} \ldots\left|a_{r}\right|^{s_{r}}
$$

where $E^{1,1}$ is the usual $G L_{2}$ Eisenstein series with trivial central character. So let $g=n m k$ be an Iwasawa decomposition with $n \in N^{Q}, m \in M^{Q}$, and $k \in \prod_{v} K_{v}$ with $m$ in the form just displayed, and put

$$
\Phi_{s}^{i}(g)=\left|a_{1}\right|^{s_{1}} \ldots\left|a_{i=1}\right|^{s_{i-1}} \cdot E_{\frac{s_{i}-s_{i+1}}{2}}^{1,1}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left|\operatorname{det}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)\right|^{\frac{s_{i}+s_{i+1}}{2}} \cdot\left|a_{i+2}\right|^{s_{i+2}} \ldots\left|a_{r}\right|^{s_{r}}
$$

Then

$$
E_{s}(g)=\sum_{\gamma \in Q_{k} \backslash G_{k}} \Phi_{s}^{i}(\gamma g) \quad\left(\text { for } g \in G L_{r}\right)
$$

This expresses the $G L_{r}$ minimal-parabolic Eisenstein series as $Q$-Eisenstein series formed from the $P^{1,1}$ Eisenstein series on the $G L_{2}$ part of its Levi component.

The usual normalization of the $G L_{2}$ Eisenstein series to eliminate poles, for boundededness on vertical strips for $g$ in compacts in $G L_{2}(\mathbb{A})$, and to be invariant under $s \rightarrow 1-s$, is

$$
\widetilde{E}_{s}(g)=s(1-s) \cdot \xi(2 s) \cdot E_{s}^{1,1}(g) \quad(\text { for } s \in \mathbb{C})
$$

Thus, let

$$
\widetilde{\Phi}_{s}^{i}=\left(\frac{s_{i}-s_{i+1}}{2}\right)\left(1-\frac{s_{i}-s_{i+1}}{2}\right) \cdot \xi\left(s_{i}-s_{i+1}\right) \cdot \Phi_{s}^{i}
$$

An argument similar to [3.10.1] for convergence of the minimal-parabolic Eisenstein series $E_{s}$ and [3.11.1] for maximal-proper-parabolic Eisenstein series will prove the absolute convergence of

$$
\left(\frac{s_{i}-s_{i+1}}{2}\right)\left(1-\frac{s_{i}-s_{i+1}}{2}\right) \cdot \xi\left(s_{i}-s_{i+1}\right) \cdot E_{s}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \widetilde{\Phi}_{s}^{i}(\gamma g)
$$

for $\frac{\operatorname{Re}\left(s_{j}\right)-\operatorname{Re}\left(s_{j+1}\right)}{2}>1$ for $j \neq i$, with no condition on $s_{i}-s_{i+1}$, because we use the analytically-continued Eisenstein series on $G L_{2}$ rather than the expression of it as a series.

The convergence argument is as follows. For $g$ in a fixed compact and $s_{i}-s_{i+1}$ in a fixed vertical strip, $\Phi_{s}^{i}(g)$ is dominated by the function obtained by replacing $\widetilde{E}^{1,1}$ by a constant, namely, with $\sigma_{j}=\operatorname{Re}\left(s_{j}\right)$,

$$
\theta(g)=\left|a_{1}\right|^{\sigma_{1}} \ldots\left|a_{i-1}\right|^{\sigma_{i-1}} \cdot\left|\operatorname{det}\left(\begin{array}{cc}
a & b \\
c & d
\end{array}\right)\right|^{\frac{\sigma_{i}+\sigma_{i+1}}{2}} \cdot\left|a_{i+2}\right|^{\sigma_{i+2}} \ldots\left|a_{n}\right|^{\sigma_{n}}
$$

We prove the absolute convergence of the Eisenstein series $E(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \theta(\gamma g)$, which is degenerate in the same sense as the approximating Eisenstein series in the proof of [3.11.1]. As in the earlier convergence arguments, convergence is equivalent to convergence of an integrated form, namely

$$
\int_{Z_{\mathbb{A}} \backslash C} \sum_{\gamma \in P_{k} \backslash G_{k}} \theta(\gamma g) d g
$$

Shrinking $C$ sufficiently so that $\gamma \cdot C \cap C \neq \phi$ implies $\gamma=1$,

$$
\int_{Z_{\mathbb{A}} \backslash C} \sum_{\gamma \in P_{k} \backslash G_{k}} \theta(\gamma g) d g=\int_{Z_{\mathbb{A}} P_{k} \backslash G_{k} \cdot C} \theta(g) d g
$$

As in the earlier convergence arguments, letting $\eta_{j}$ be the norm of the determinant of the lower right $n-j$ minor, $G_{k} \cdot C$ is contained in

$$
Y=\left\{g \in G_{\mathbb{A}}: 1<_{C} \eta_{j}(g) \text { for } j=1, \ldots, n\right\}
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

To compare with $M^{Q}$ with $Q=Q^{i}$, drop the $(i+1)^{t h}$ condition: $G_{k} \cdot C$ is contained in

$$
Y^{\prime}=\left\{g \in G_{\mathbb{A}}: 1<_{C} \eta_{j}(g) \text { for } j \neq i+1\right\}
$$

Thus, convergence of the Eisenstein series is implied by convergence of

$$
\int_{Z_{\mathbb{A}} Q_{k} \backslash Y^{\prime}} \theta(g) d g
$$

As $Y^{\prime}$ is stable by right multiplication by the maximal compact subgroup $K_{v} \subset G_{v}$ at all places $v$, by an Iwasawa decomposition this integral is

$$
\int_{Z_{\mathbb{A}} Q_{k} \backslash\left(Y^{\prime} \cap Q_{\mathbb{A}}\right)} \theta(p) d p \quad \text { (left Haar measure on } Q \text { ) }
$$

Let $\alpha=\alpha_{i}$ be the $i^{t h}$ simple positive root, and $\rho$ the half-sum of positive roots. The left Haar measure on $Q_{\mathbb{A}}$ is $d(n m)=d n d m / m^{2 \rho-\alpha}$, where $d n$ is Haar measure on $N^{Q}$ and $d m$ is Haar measure on the Levi component $M^{Q}$. Since $\theta$ is left $N_{\mathbb{A}}^{Q}$-invariant and $N_{k}^{Q} \backslash N_{\mathbb{A}}^{Q}$ is compact, convergence of the latter integral is equivalent to convergence of

$$
\int_{Z_{\mathbb{A}} M_{k}^{Q} \backslash\left(Y^{\prime} \cap M_{\mathbb{A}}^{Q}\right)} \theta(m) \frac{d m}{m^{2 \rho-\alpha}}
$$

As in the earlier convergence argument, the compactness lemma [2.A] and right action of $M \cap \prod_{v} K_{v}$ reduce the convergence question to that of a simpler integral.

As in the proof of [3.10.1], parametrize a subgroup $H$ of $S L_{n}(\mathbb{A})$ by $r-1$ maps from $\left(G L_{1}(\mathbb{A})\right.$, namely,

$$
h_{j}: t \longrightarrow\left(\begin{array}{ccccccc}
1 & & & & & & \\
& \ddots & & & & & \\
& & 1 & & & & \\
\\
& & & t & & & \\
\\
& & & & t^{-1} & & \\
\\
& & & & & 1 & \\
\\
& & & & & & \ddots
\end{array}\right) \quad \quad \quad\left(\text { at } j^{t h} \text { and }(j+1)^{t h}\right. \text { positions) }
$$

with the $i^{\text {th }}$ replaced by the obvious map from $S L_{2}(\mathbb{A})$, namely,

$$
h_{i}^{\prime}:\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \longrightarrow\left(\begin{array}{ccccccc}
1 & & & & & & \\
& \ddots & & & & & \\
\\
& & 1 & & & & \\
\\
& & & a & b & & \\
\\
& & & c & d & & \\
\\
& & & & & 1 & \\
& & & & & & \ddots
\end{array}\right) \quad\left(\text { at } i^{t h} \text { and }(i+1)^{t h}\right. \text { positions) }
$$

Then

$$
Y^{\prime} \cap M_{\mathbb{A}}^{Q} \cap S L_{n}(\mathbb{A})=\left\{\prod_{j \neq i} h_{j}\left(t_{j}\right): t_{j} \in \mathbb{J} \text { and }\left|t_{j}^{-1}\right| \gg 1\right\} \times\left\{h_{i}^{\prime}(T): T \in S L_{2}(\mathbb{A}),|\operatorname{det} T| \gg 1\right\}
$$

Noting that $h_{j}(t)^{2 \rho}=|t|^{2}$, convergence is implied by convergence of

$$
\begin{cases}\int_{0}^{\ll 1} t^{\sigma_{j}-\sigma_{j+1}-2} \frac{d t}{t} & (\text { for } j \neq i) \\ \int_{S L_{2}(k) \backslash S L_{2}(\mathbb{A})} 1 d t & \text { (right invariant measure } d t)\end{cases}
$$

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The $G L_{1}$ integrals are absolutely convergent for $\sigma_{j}-\sigma_{j+1}-2>0$ for $j \neq i$. By reduction theory [2.2] and [3.3], for example, $S L_{2}(k) \backslash S L_{2}(\mathbb{A})$ has finite volume, so the $S L_{2}$ integral is convergent. This is the desired convergence conclusion: there is no constraint on $\sigma_{i}-\sigma_{i+1}$. Thus, the iterated expression for the Eisenstein series analytically continues as indicated.

Functional equations for reflections: In addition to a partial analytic continuation, the previous argument gives the functional equations for the reflections in $W$ attached to the simple roots, as follows. The main issue is making the functional equations understandable. In the iterated expression for $E_{s}$ above in terms of a $G L_{2}$ Eisenstein series, the functional equation $\widetilde{E}_{1-z}^{1,1}=\widetilde{E}_{z}^{1,1}$ of that $G L_{2}$ Eisenstein series gives

$$
\begin{aligned}
& \widetilde{\Phi}_{s}^{i}(g)=\left|a_{1}\right|^{s_{1}} \ldots\left|a_{i=1}\right|^{s_{i-1}} \cdot \widetilde{E}_{\frac{s_{i}-s_{i+1}}{2}}^{\sim_{2}}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left|\operatorname{det}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)\right|^{\frac{s_{i}+s_{i+1}}{2}} \cdot\left|a_{i+2}\right|^{s_{i+2}} \ldots\left|a_{r}\right|^{s_{r}} \\
& \quad=\left|a_{1}\right|^{s_{1}} \ldots\left|a_{i=1}\right|^{s_{i-1}} \cdot \widetilde{E}_{1-\frac{s_{i}-s_{i+1}}{1,1}}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \cdot\left|\operatorname{det}\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)\right|^{\frac{s_{i}+s_{i+1}}{2}} \cdot\left|a_{i+2}\right|^{s_{i+2}} \ldots\left|a_{r}\right|^{s_{r}}
\end{aligned}
$$

This is not presented immediately in terms of $s=\left(s_{1}, \ldots, s_{r}\right)$, but, instead, says

$$
\left(s_{1}, \ldots, s_{i-1}, \frac{s_{i}-s_{i+1}}{2}, \frac{s_{i}+s_{i+1}}{2}, s_{i+2}, \ldots, s_{r}\right) \longrightarrow\left(s_{1}, \ldots, s_{i-1}, 1-\frac{s_{i}-s_{i+1}}{2}, \frac{s_{i}+s_{i+1}}{2}, s_{i+2}, \ldots, s_{r}\right)
$$

We hope for clarification by identifying the simultaneous fixed point(s), if any, of all these, for $i=1, \ldots, r$, together with the condition $s_{1}+\ldots+s_{r}=0$ : the $i^{t h}$ transformation fixes all by the $i^{t h}$ and $(i+1)^{t h}$ coordinate, and in those two coordinates the fixed-point condition is

$$
\frac{s_{i}-s_{i+1}}{2}=1-\frac{s_{i}-s_{i+1}}{2} \quad \text { and } \quad \frac{s_{i}+s_{i+1}}{2}=\frac{s_{i}+s_{i+1}}{2}
$$

The second equation is a tautology, so the $i^{t h}$ fixed-point condition is simply $s_{i}-s_{i+1}=1$. These conditions for $i=1, \ldots, r-1$ and $s_{1}+\ldots+s_{r}=0$, give a unique fixed point,

$$
\text { fixed point }=\left(\frac{n-1}{2}, \frac{n-3}{2}, \frac{n-5}{2}, \ldots, \frac{3-n}{2}, \frac{1-n}{2}\right)
$$

This is the half-sum $\rho$ of positive roots

$$
\rho=\left(\rho_{1}, \ldots, \rho_{r}\right)=\frac{1}{2} \sum_{i<j}(0, \ldots, 0,1,0, \ldots, 0,-1,0, \ldots, 0) \quad\left(\text { at the } i^{t h} \text { and } j^{t h} \text { places }\right)
$$

Replacing $s$ by $\rho+s$ replaces $s_{i}-s_{i+1}$ by $\left(\rho_{i}+s_{i}\right)-\left(\rho_{i+1}+s_{i+1}\right)=s_{i}-s_{i+1}+1$, and the map $\frac{s_{i}-s_{i+1}}{2} \rightarrow 1-\frac{s_{i}-s_{i+1}}{2}$ becomes

$$
\frac{\left(\rho_{i}+s_{i}\right)-\left(\rho_{i+1}+s_{i+1}\right)}{2} \rightarrow 1-\frac{\left(\rho_{i}+s_{i}\right)-\left(\rho_{i+1} s_{i+1}\right)}{2}
$$

which simplifies to $s_{i} \longleftrightarrow s_{i+1}$. That is, in the $\rho+s$ coordinates, this functional equation is the interchange of $s_{i}$ and $s_{i+1}$ :

$$
\rho+\left(s_{1}, \ldots, s_{i-1}, s_{i}, s_{i+1}, s_{i+2}, \ldots, s_{r}\right) \longleftrightarrow \rho+\left(s_{1}, \ldots, s_{i-1}, s_{i+1}, s_{i}, s_{i+2}, \ldots, s_{r}\right)
$$

This is the same as the effect $m^{s} \longrightarrow\left(\tau_{i} m \tau_{i}^{-1}\right)^{s}$ with

$$
\tau_{i}=\left(\begin{array}{cccccccc}
1 & & & & & & & \\
& \ddots & & & & & & \\
& & 1 & & & & & \\
& & & 0 & 1 & & & \\
& & & 1 & 0 & & & \\
& & & & & 1 & & \\
& & & & & & \ddots & \\
& & & & & \\
&
\end{array}\right) \in W \quad\left(\text { at } i^{t h} \text { and }(i+1)^{t h}\right. \text { positions) }
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

This $\tau_{i} \in W$ is usually considered to be attached to the $i^{\text {th }}$ simple root $\alpha_{i}(m)=m_{i}-m_{i+1}$ on the Lie algebra, as it is characterized by interchanging $\pm \alpha_{i}$ and permuting the other positive roots $\beta_{j \ell}(m)=m_{j}-m_{\ell}$ : for $m$ in the Lie algebra of $M^{\text {min }}$, unless $j=i$ and $\ell=i+1$,

$$
\beta_{j \ell}\left(\tau_{i} m \tau_{i}^{-1}\right)==\left(\tau_{i} m \tau_{i}^{-1}\right)_{j}-\left(\tau_{i} m \tau_{i}^{-1}\right)_{\ell}=m_{j^{\prime}}-m_{\ell^{\prime}} \quad \quad\left(\text { for some } j^{\prime}<\ell^{\prime}\right)
$$

thus giving some other positive root evaluated on $m$. When $j=i$ and $\ell=i+1$, the effect is qualitatively different, reversing the sign, producing $-\alpha_{i}$.

To rewrite the above in more geometric terms, use the pairing $\langle x, y\rangle=\operatorname{tr}(x y)$ on the Lie algebra $\mathfrak{g}$ of $G L_{r}$ to identify $\mathfrak{a}$ and $\mathfrak{a}^{*}$ and give each of them a non-degenerate inner product. The Weyl group preserves this inner product, since

$$
\left\langle w x w^{-1}, w y w^{-1}\right\rangle=\operatorname{tr}\left(w x w^{-1} \cdot w y w^{-1}\right)=\operatorname{tr}\left(w x y w^{-1}\right)=\operatorname{tr}(x y)=\langle x, y\rangle
$$

by conjugation invariance of trace. The geometric characterization of the reflection $\tau=\tau_{\alpha}$ associated to a vector (here a simple positive root) $\alpha$ is that $\tau$ should fix the hyperplane orthogonal to $\alpha$, and should send $\alpha \rightarrow-\alpha$ : this is expressed by

$$
\tau x=x-2 \cdot \frac{\langle x, \alpha\rangle}{\langle\alpha, \alpha\rangle} \cdot \alpha \quad \text { (for } x \in \mathfrak{a}^{*} \approx \text { diagonal matrices) }
$$

Via this pairing, $\alpha_{i}$ is identified with

$$
\alpha_{i}=\left(\begin{array}{ccccccc}
0 & & & & & & \\
& \ddots & & & & & \\
\\
& & 0 & & & & \\
\\
& & & 1 & 0 & & \\
0 & -1 & & & \\
& & & & & 0 & \\
& & & & & \ddots & \\
& & & & & & \\
0
\end{array}\right) \quad\left(\text { at } i^{\text {th }} \text { and }(i+1)^{t h}\right. \text { positions) }
$$

because $\alpha_{i}(m)=\left\langle m, \alpha_{i}\right\rangle$. Similarly, $s_{i}-s_{i+1}=\left\langle s, \alpha_{i}\right\rangle$. It is immediate that $\tau \cdot \alpha=\tau \alpha \tau^{-1}=-\alpha$, as the reflection should. Since $\tau$ preserves $\langle$,$\rangle , it preserves the orthogonal complement to \alpha$, so truly is the associated reflection. Since $\tau$ flips the sign on $\alpha$ and permutes the other positive roots, we can compute, using $\tau^{-1}=\tau$,

$$
\langle 2 \rho, \alpha\rangle=-\langle 2 \rho, \tau \cdot \alpha\rangle=-\langle\tau \cdot 2 \rho, \tau \cdot \alpha\rangle
$$

and

$$
\tau \cdot 2 \rho=\tau \cdot \sum_{\beta>0} \beta=\tau \cdot\left(\sum_{\beta>0, \beta \neq \alpha} \beta\right)+\tau \cdot \alpha=\left(\sum_{\beta>0, \beta \neq \alpha} \beta\right)-\alpha=2 \rho-2 \alpha
$$

Thus, $\langle 2 \rho, \alpha\rangle=-\langle 2 \rho-2 \alpha, \alpha\rangle$, from which $\langle\rho, \alpha\rangle=\langle\alpha, \alpha\rangle / 2=1$. Thus, the $\alpha^{t h}$ functional equation, inherited from the $G L_{2}$ Eisenstein series, is

$$
\xi\langle\rho+\tau \cdot s, \alpha\rangle \cdot E_{\rho+\tau \cdot s}=\xi\langle\rho+s, \alpha\rangle \cdot E_{\rho+s} \quad\left(\text { reflection } \tau=\tau_{\alpha}\right)
$$

Using $\langle\rho, \alpha\rangle=1$, this is

$$
\left.\xi(1+\langle\tau \cdot s, \alpha\rangle) \cdot E_{\rho+\tau \cdot s}=\xi(1+\langle s, \alpha\rangle) \cdot E_{\rho+s} \quad \quad \text { reflection } \tau=\tau_{\alpha}\right)
$$

or, since $\langle\tau \cdot s, \alpha\rangle=\langle s, \tau \cdot \alpha\rangle=\langle s,-\alpha\rangle=-\langle s, \alpha\rangle$,

$$
E_{\rho+\tau \cdot s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi(1+\langle\tau \cdot s, \alpha\rangle)} \cdot E_{\rho+s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi(1-\langle s, \alpha\rangle)} \cdot E_{\rho+s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \cdot E_{\rho+s} \quad\left(\text { reflection } \tau=\tau_{\alpha}\right)
$$

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using the functional equation $\xi(1-z)=\xi(z)$.
Application of Bochner's Lemma: The $n-1$ partial analytic continuations can be organized to allow application of Bochner's Lemma. Above, for $\alpha=\alpha_{i}$ the $i^{t h}$ simple root, we showed that the function

$$
E_{\rho+s}^{\alpha}=\langle\rho+s, \alpha\rangle \cdot(1-\langle\rho+s, \alpha\rangle) \cdot \xi\langle\rho+s, \alpha\rangle \cdot E_{\rho+s}=(1+\langle s, \alpha\rangle) \cdot(1-\langle s, \alpha\rangle) \cdot \xi(1+\langle s, \alpha\rangle) \cdot E_{\rho+s}
$$

admits an analytic continuation in which $s_{i}-s_{i+1}=\langle s, \alpha\rangle$ is not constrained, and this normalized version of $E_{\rho+s}$ is invariant under $\rho+s \rightarrow \rho+\tau_{\alpha} \cdot s$ for the reflection $\tau_{\alpha}$. This might suggest adding normalization factors for all positive roots, to obtain an eventually $W$-invariant expression:

$$
E_{\rho+s} \cdot \prod_{\beta>0}(1+\langle s, \beta\rangle) \cdot(1-\langle s, \beta\rangle) \cdot \xi(1+\langle s, \beta\rangle)
$$

The intention is that $E_{\rho+s}^{\alpha}$ is invariant under the reflection $\tau_{\alpha}$ for each simple root $\alpha$, and the remaining factors should be permuted among themselves, since the other positive roots are permuted among themselves by $\tau_{\alpha}$.

A minor technical issue arises: to be sure to cancel the pole of $\xi(1+\langle s, \beta\rangle)$ at $1+\langle s, \beta\rangle=1$, in order to most easily justify application of Bochner's lemma, add additional polynomial factors, squared to avoid disturbing the sign in functional equations: let

$$
E_{\rho+s}^{\#}=E_{\rho+s} \cdot \prod_{\beta>0}(1+\langle s, \beta\rangle) \cdot(1-\langle s, \beta\rangle) \cdot\langle s, \beta\rangle^{2} \cdot \xi(1+\langle s, \beta\rangle)
$$

The exponential decay of the gamma factor in $\xi$ is more than sufficient to preserve boundedness in vertical strips for real part $s$ in compacts.
[3.12.4] Claim: $E_{\rho+s}^{\#}$ has an analytic continuation to a holomorphic function on $\mathbb{C}^{r}$, and is invariant under $s \rightarrow w \cdot s$ for all $w \in W$.
Proof: By the $G L_{2}$ discussion and the above adaptations, $E_{\rho+s}^{\#}$ has an analytic continuation to the tube domain $\Omega=\left\{z \in \mathbb{C}^{r}: \operatorname{Re}(z) \in \Omega_{o}\right\}$ over $\Omega_{o} \subset \mathbb{R}^{r}$ given by

$$
\Omega_{o}=\left\{\sigma \in \mathbb{R}^{r}:\langle\rho+\sigma, \alpha\rangle>1 \text { for all but possibly a single simple root } \alpha\right\}
$$

In $\Omega$, for $\operatorname{Re}(s)$ in compacts, $E_{\rho+s}^{\#}$ is bounded, so certainly has sufficiently modest growth for application of Bochner's Lemma, and $E_{\rho+s}^{\#}$ has an analytic continuation to the convex hull of $\Omega$, which is $\mathbb{C}^{r}$. It is invariant under all reflections attached to simple roots, and these generate $W$. This proves the claim. ///

Returning to the proof of the theorem, this last claim gives the meromorphic continuation of $E_{\rho+w}$, and the first corollary.

Given the meromorphic continuation of $E_{\rho+s}$, the functional equations

$$
\left.E_{\rho+\tau \cdot s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \cdot E_{\rho+s} \quad \text { (reflection } \tau=\tau_{\alpha}\right)
$$

of $E_{\rho+s}$ proven above for reflections $\tau=\tau_{\alpha}$ attached to simple roots $\alpha$ can be iterated. Taking constant terms gives, by the general form of the constant term [3.10.3],

$$
\sum_{w} c_{w, \tau \cdot s} m^{\rho+w \cdot \tau \cdot s}=c_{P} E_{\rho+\tau \cdot s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \cdot c_{P} E_{\rho+s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \sum_{w} c_{w, s} m^{\rho+w \cdot s}
$$

For generic $s \in \mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$, the coefficients of the various characters of $m \in M_{\mathbb{A}}$ must be equal, so by the identity principle are equal for all $s$. In particular, equating the coefficient of $m^{\rho+\tau \cdot s}$ gives

$$
1=c_{1, \tau \cdot s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \cdot c_{\tau, s}
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

That is, $c_{\tau, s}=\xi\langle s, \alpha\rangle / \xi(1+\langle s, \alpha\rangle)$, and

$$
E_{\rho+\tau \cdot s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \cdot E_{\rho+s}=\frac{1}{c_{\tau, s}} \cdot E_{\rho+s} \quad \quad \quad \text { (for reflection } \tau \text { ) }
$$

Equating the coefficients of $m^{\rho+w \cdot \tau \cdot s}$ gives

$$
c_{w, \tau \cdot s}=\frac{\xi(1+\langle s, \alpha\rangle)}{\xi\langle s, \alpha\rangle} \cdot c_{w \tau, s}=\frac{1}{c_{\tau, s}} \cdot c_{w \tau, s}
$$

from which $c_{w, \tau \cdot s} \cdot c_{\tau, s}=c_{w \tau, s}$. For two reflections $\sigma, \tau$,

$$
c_{w \sigma \tau, s}=c_{(w \sigma) \tau, s}=c_{w \sigma, \tau \cdot s} \cdot c_{\tau, s}=c_{w, \sigma \tau \cdot s} \cdot c_{\sigma, \tau \cdot s} \cdot c_{\tau, s}=c_{w, \sigma \tau \cdot s} \cdot c_{\sigma \tau, s}
$$

Induction gives $c_{w w^{\prime}, s}=c_{w, w^{\prime} \cdot s} \cdot c_{w^{\prime}, s}$. Then

$$
E_{\rho+\sigma \tau \cdot s}=E_{\rho+\sigma \cdot(\tau \cdot s)}=\frac{1}{c_{\sigma, \tau \cdot s}} \cdot E_{\rho+\tau \cdot s}=\frac{1}{c_{\sigma, \tau \cdot s}} \cdot \frac{1}{c_{\tau, s}} \cdot E_{\rho+s}=\frac{1}{c_{\sigma \tau, s}} \cdot E_{\rho+s}
$$

and a similar induction on the length of $w \in W$ gives the general functional equation. Qualitatively, the number of factors in both numerator and denominator of $c_{w}(s)$ is the length of $w$. This proves the theorem and corollary.
[3.12.5] Example: For $G=G L_{3}$ there are two simple positive roots,

$$
\langle x, \alpha\rangle=x_{1}-x_{2} \quad\langle x, \beta\rangle=x_{2}-x_{3} \quad\left(\text { for } x \in \mathfrak{a} \text { with diagonal entries } x_{i}\right)
$$

The other positive root is $\alpha+\beta$, so $\rho=\frac{1}{2}\left(\alpha+\beta+\left(\alpha_{\beta}\right)\right)=\alpha+\beta$. Let $\sigma, \tau$ be the reflections corresponding to $\alpha, \beta$, respectively. The whole Weyl group is $W=\{1, \sigma, \tau, \sigma \tau, \tau \sigma, \sigma \tau \sigma\}$ and $\sigma \tau \sigma=\tau \sigma \tau$. From the $G L_{2}$ computation,

$$
c_{\sigma, \rho+s}=\frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)} \quad c_{\tau, s}=\frac{\xi\langle s, \beta\rangle}{\xi(\langle s, \beta\rangle+1)}
$$

By the cocycle relation $c_{w r, s}=c_{w, r \cdot s} \cdot c_{r, s}$ for reflection $r \in W$ and $w \in W$, we have

$$
c_{\sigma \tau, s}=c_{\sigma, \tau \cdot s} \cdot c_{\tau, s}=\frac{\xi\langle\tau \cdot s, \alpha\rangle}{\xi(\langle\tau \cdot s, \alpha\rangle+1)} \cdot \frac{\xi\langle s, \beta\rangle}{\xi(\langle s, \beta\rangle+1)}
$$

Since $\langle\tau x, \alpha\rangle=\langle x, \tau \alpha\rangle=\langle x, \alpha+\beta\rangle$,

$$
c_{\sigma \tau, s}=\frac{\xi\langle s, \alpha+\beta\rangle}{\xi(\langle s, \alpha+\beta\rangle+1)} \cdot \frac{\xi\langle s, \beta\rangle}{\xi(\langle s, \beta\rangle+1)}
$$

Similarly,

$$
c_{\tau \sigma, s}=\frac{\xi\langle s, \alpha+\beta\rangle}{\xi(\langle s, \alpha+\beta\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)}
$$

Finally,

$$
c_{\tau \sigma \tau, s}=c_{\sigma \tau \sigma, s}=c_{\sigma \tau, \sigma \cdot s} \cdot c_{\sigma, s}=\frac{\xi\langle\sigma \cdot s, \alpha+\beta\rangle}{\xi(\langle\sigma \cdot s, \alpha+\beta\rangle+1)} \cdot \frac{\xi\langle\sigma \cdot s, \beta\rangle}{\xi(\langle\sigma \cdot s, \beta\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)}
$$

Using $\sigma \beta=\alpha+\beta$ and $\sigma(\alpha+\beta)=\beta$, this is

$$
c_{\tau \sigma \tau, s}=c_{\sigma \tau \sigma, s}=\frac{\xi\langle s, \beta\rangle}{\xi(\langle s, \beta\rangle+1)} \cdot \frac{\xi\langle s, \alpha+\beta\rangle}{\xi(\langle s, \alpha+\beta\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)}
$$

The latter example suggests that more can be said about $c_{w, s}$ :

## [3.12.6] Claim:

$$
c_{w, s}=\prod_{\beta>0: w \cdot \beta<0} \frac{\xi\langle s, \beta\rangle}{\xi(\langle s, \beta\rangle+1)}
$$

Proof: Induction on the length of $w$ in the generating reflections associated to simple roots. With $\tau=\tau_{\alpha}$ for simple root $\alpha$ the cocycle relation gives

$$
\begin{gathered}
c_{w \tau, s}=c_{w, \tau \cdot s} \cdot c_{\tau, s}=\prod_{\beta>0: w \cdot \beta<0} \frac{\xi\langle\tau \cdot s, \beta\rangle}{\xi(\langle\tau \cdot s, \beta\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)} \\
=\prod_{\beta>0: w \cdot \beta<0} \frac{\xi\langle s, \tau \cdot \beta\rangle}{\xi(\langle s, \tau \cdot \beta\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)}
\end{gathered}
$$

The effect of $\tau=\tau_{\alpha}$ on roots is to interchange $\pm \alpha$, permute the other positive roots, and permute the other negative roots. There are two cases.

First, if $w \cdot \alpha<0$, then $\alpha$ itself appears in the product, and $\left(w \tau_{\alpha}\right) \cdot \alpha=w(-\alpha)=-w \cdot \alpha>0$. So $\alpha$ will not appear in the corresponding product for $w \tau$. Using the functional equation $\xi(1-z)=\xi(z)$,

$$
\begin{aligned}
\frac{\xi\langle s, \tau \cdot \alpha\rangle}{\xi(\langle s, \tau \cdot \alpha\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)} & =\frac{\xi\langle s,-\alpha\rangle}{\xi(\langle s,-\alpha\rangle+1)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)}=\frac{\xi(1-(\langle s, \alpha\rangle+1)}{\xi(1-\langle s, \alpha\rangle)} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)} \\
& =\frac{\xi(\langle s, \alpha\rangle+1)}{\xi\langle s, \alpha\rangle} \cdot \frac{\xi\langle s, \alpha\rangle}{\xi(\langle s, \alpha\rangle+1)}=1
\end{aligned}
$$

Thus, the leftover factor from the product for $w$ cancels the new factor from the cocycle relation, and the desired relation holds for $w \tau_{\alpha}$ in the case that $w \cdot \alpha<0$.

Second, similarly but oppositely, suppose $w \cdot \alpha>0$. Then $\alpha$ does not appear in the product for $w$. But $\left(w \tau_{\alpha}\right) \alpha=w(-\alpha)<0$, so $\alpha$ should appear in the product for $w \tau_{\alpha}$. The extra term provides this, proving the relation in this case.

### 3.13 Continuation of cuspidal-data Eisenstein series

The functional equations of Eisenstein series attached to non-minimal parabolics $P^{d_{1}, \ldots, d_{\ell}}$ involve all the parabolics $P^{d_{1}^{\prime}, \ldots, d_{\ell}^{\prime}}$ with $d_{1}^{\prime}, \ldots, d_{\ell}^{\prime}$ a permutation of $d_{1}, \ldots, d_{\ell}$, called the associates of $P^{d_{1}, \ldots, d_{\ell}}$. This is so even with the simplifying assumption of cuspidal data, without which the situation is messier. Then, the expression of pseudo-Eisenstein series for such parabolics in terms of genuine Eisenstein series, even with the corresponding assumption of cuspidal data, involves all these. Thus, as in [3.11], we consider only maximal proper parabolics $P=P^{r_{1}, r_{2}}$, right $M_{\mathbb{A}}^{P} \cap K_{\mathbb{A}}$-invariant cuspidal data $f=f_{1} \otimes f_{2}$ on $M_{\mathbb{A}}^{P}$ with trivial central character. Assume $f_{1}, f_{2}$ are spherical Hecke eigenfunctions at all finite places, so by [3.11.6] the Eisenstein series $E_{s, f}^{P}$ is a spherical Hecke eigenfunction at all finite places. Similarly, at archimedean places $v$, we assume (at least) that $f_{1}$ and $f_{2}$ are eigenfunctions for the invariant Laplacians on the factors of the Levi component, so by [3.11.11] $E_{s, f}^{P}$ is an eigenfunction for the invariant Laplacian on $G_{v}$.

Assume strong multiplicity one for $f_{1} \otimes f_{2}$, as in [3.11.9], so that the constant terms of $E_{s, f}^{P}$ are as simple as possible. With $P=P^{r_{1}, r_{2}}$, let $Q=P^{r_{2}, r_{1}}$, so $Q=P$ for self-associate $P$ and otherwise is the unique other associate of $P$. Thus, the following is special, but perhaps more palatable than the general case. Write $f^{w}=\left(f_{1} \otimes f_{2}\right)^{w}=f_{2} \otimes f_{1}$.
[3.13.1] Theorem: (Langlands, Bernstein, Wong, et alia) With the constant-term conventions as in [3.11.9] such Eisenstein series $E_{s, f}^{P}$ have meromorphic continuations in $s$, with functional equation

$$
E_{1-s, f}^{P}=\left(c_{s, f^{w}}^{P}\right)^{-1} \cdot E_{s, f^{w}}^{Q} \quad \text { and } \quad c_{1-s, f}^{Q} \cdot c_{s, f^{w}}^{P}=1
$$

(Proofs in [11.10], [11.12].)

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

Although the proof of meromorphic continuation is postponed to [11.10], [11.12], if we grant meromorphic continuation then the form of the functional equation is determined by the constant term, using the theory of the constant term [8.3], as follows. As in [3.11.9], the constant terms of $E_{s, f}^{P}$ and $E_{1-s, f w}^{Q}$ are explicit. First, for $P$ self-associate,

$$
c_{P} E_{s, f^{w}}^{P}=\varphi_{s, f^{w}}^{P}+c_{s, f w} \varphi_{1-s, f}^{P} \quad \text { and } \quad c_{P} E_{1-s, f}^{P}=\varphi_{1-s, f}^{P}+c_{1-s, f} \varphi_{s, f}^{P}
$$

and all other constant terms are 0 . Thus,

$$
c_{P}\left(E_{1-s, f}^{P}-c_{s, f^{w}}^{-1} \cdot E_{s, f^{w}}^{P}\right)=\left(c_{1-s, f}-c_{s, f^{w}}^{-1}\right) \cdot \varphi_{s, f^{w}}^{P}
$$

and all other constant terms are 0 . The functions $f_{1}$ and $f_{2}$ are cuspforms in a strong sense, so by [8.2] are bounded. Thus,

$$
\left|\varphi_{s, f}^{P}(n m k)\right| \leq\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{\operatorname{Re}(s)}
$$

is bounded in standard Siegel sets for $\operatorname{Re}(s)$ sufficiently negative. By the theory of the constant term [8.3], $E_{1-s, f^{w}}^{P}-c_{s, f^{w}}^{-1} \cdot E_{s, f}^{P}$ is bounded in Siegel sets. Thus, this difference is in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. However, from [3.11.11], both $E_{1-s, f w}^{P}$ and $E_{s, f}^{P}$ have the same eigenvalues for invariant Laplacians at archimedean places, namely, $r_{1} r_{2}\left(r_{1}+r_{2}\right)\left(s^{2}-s\right)+\lambda_{1}+\lambda_{2}$ where $\lambda_{j}$ is the eigenvalue for $f_{j}$ on the corresponding archimedean factor $G L_{r_{j}}$ of $M^{P}$. Thus, the difference has that eigenvalue. There are many choices of $s$ with $\operatorname{Re}(s)<0$ which make this eigenvalue non-real, however, which is impossible for an $L^{2}$ eigenfunction other than 0 , as in the proof of [1.10.5]. Thus, this difference must be 0 , and have constant term 0 . This gives the functional equation, assuming the meromorphic continuation, in the self-associate case.

For the non-self-associate case, for $P$ and its other associate $Q$, both constant terms must be considered to invoke the theory of the constant term. Starting from

$$
c_{P} E_{1-s, f}^{P}=\varphi_{1-s, f}^{P} \quad c_{Q} E_{1-s, f}^{P}=c_{1-s, f}^{Q} \varphi_{s, f^{w}}^{Q} \quad c_{P} E_{s, f^{w}}^{Q}=c_{s, f^{w}}^{P} \varphi_{1-s, f}^{P} \quad c^{Q} E_{s, f^{w}}^{Q}=\varphi_{s, f^{w}}^{Q}
$$

we have

$$
\left\{\begin{array}{l}
c_{P}\left(E_{1-s, f}^{P}-\left(c_{s, f^{w}}^{P}\right)^{-1} \cdot E_{s, f^{w}}^{Q}\right)=0 \\
c_{Q}\left(E_{1-s, f}^{P}-\left(c_{s, f^{w}}^{P}\right)^{-1} \cdot E_{s, f^{w}}^{Q}\right)=\left(c_{1-s, f}^{Q}-\left(c_{s, f^{w}}^{P}\right)^{-1}\right) \cdot \varphi_{s, f^{w}}^{Q}
\end{array}\right.
$$

and all other constant terms are 0 . As usual, the cuspforms are bounded by [8.3], so, for $\operatorname{Re}(s)$ sufficiently negative, the constant term along $Q$ is square-integrable on Siegel sets. By the theory of the constant term [8.3], the difference $E_{1-s, f}^{P}-\left(c_{s, f}^{P}\right)^{-1} \cdot E_{s, f^{w}}^{Q}$ is square-integrable. However, again by [3.11.11], this difference is an eigenfunction for invariant Laplacians at archimedean places, with non-real eigenvalues for many choices of $s$. Thus, it is identically 0 .

### 3.14 Truncation and Maaß-Selberg relations

First, we make precise a notion of truncation of automorphic forms, relative to a choice of parabolic subgroup, especially maximal proper parabolics. For the self-associate maximal proper parabolic $P^{r, r}$ in $G L_{2 r}$, the computation of inner products of truncations of $P^{r, r}$ Eisenstein series with cuspidal data is parallel to the computation for $G L_{2}$. As in [1.11] and [2.10], corollaries give information about possible poles of Eisenstein series, and square-integrability of residues of Eisenstein series.

This bears upon the occurrence of non-trivial residual square-integrable automorphic forms coming from cuspforms on smaller groups, anticipating that such automorphic forms occur as residues of Eisenstein series. For example, there is no non-constant non-cuspidal discrete spectrum for $G L_{2}(\mathbb{Z})$ nor for $G L_{3}(\mathbb{Z})$, but only for $G L_{4}(\mathbb{Z})$ and larger. Namely, the Eisenstein series on $G L_{3}$ with $G L_{2}$ cuspidal data have no poles in the right half-plane, as follows immediately from the Maaß-Selberg relations below.

The simplest non-trivial examples of Maaß-Selberg relations and corollaries concern spherical Eisenstein series on $G L_{n}$ associated to cuspidal data on the Levi component of maximal (proper) parabolics $P=P^{r_{1}, r_{2}}$.

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For simplicity, we continue to consider only right $K_{\mathbb{A}}$-invariant Eisenstein series $E_{s, f_{1} \otimes f_{2}}^{P}$, where $f_{1}, f_{2}$ are cuspforms in the strong sense of being spherical Hecke eigenfunctions everywhere, with trivial central characters, allowing the simple outcomes of computations of constant terms as in [3.11.8] and [3.11.9]. When $r_{1} \neq r_{2}$, that is, when $P$ is not self-associate, let $Q=P^{r_{2}, r_{1}}$ be its other associate parabolic.

Let $\delta^{P}$ be the modular function of $P_{\mathbb{A}}$

$$
\delta\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)=\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|
$$

and extend this to a height function aligned with $P$, by making it right $K_{\mathbb{A}}$-invariant: $h^{P}(n m k)=\delta^{P}(n m)=$ $\delta^{P}(m)$ for $n \in N_{\mathbb{A}}^{P}, m \in M_{\mathbb{A}}^{P}$, and $k \in K_{\mathbb{A}}$. For fixed large real $T$, the $T$-tail of the $P$-constant term of an automorphic form $F$ is

$$
c_{P}^{T} F(g)=\left\{\begin{array}{cc}
c_{P} F(g) & \left(\text { for } h^{P}(g) \geq T\right) \\
0 & \left(\text { for } h^{P}(g)<T\right)
\end{array}\right.
$$

Similarly, the $T$-tail of the $Q$-constant term is

$$
c_{Q}^{T} F(g)=\left\{\begin{array}{cc}
c_{Q} F(g) & \left(\text { for } h^{Q}(g) \geq T\right) \\
0 & \left(\text { for } h^{Q}(g)<T\right)
\end{array}\right.
$$

Suitable truncations of these cuspidal-data Eisenstein series should be square integrable (potentially accomplished a number of ways), and their inner products calculable in explicit, straightforward terms. There should be no obstacle to meromorphic continuation of the tail in the truncation. These requirements are at odds with each other. Writing $\Psi^{P}(\varphi)=\Psi_{\varphi}^{P}$ for the pseudo-Eisenstein series attached to data $\varphi$, the truncation at height $T$ of the Eisenstein series is

$$
\wedge^{T} E_{s, f}^{P}= \begin{cases}E_{s, f}^{P}-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right) & \left(\text { for } n_{1}=n_{2}, \text { i.e., for } P \text { self-associate }\right) \\ E_{s, f}^{P}-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right)-\Psi^{Q}\left(c_{Q}^{T} E_{s, f}^{P}\right) & \left(\text { for } n_{1} \neq n_{2}, \text { i.e., for } P \text { not self-associate }\right)\end{cases}
$$

[3.14.1] Proposition: The truncated Eisenstein series $\wedge^{T} E_{s, f}^{P}$ is of rapid decay in Siegel sets.
Proof: The argument is simpler in the self-associate case, which we carry out first. From the computations in [3.11.8] and [3.11.9] of constant terms of such Eisenstein series, for self-associate maximal proper $P$ in $G L_{r}$ all such constant terms are 0 except that along $P$ itself. By the theory of the constant term, on standard Siegel sets $E_{s, f}^{P}-c_{P} E_{s, f}^{P}$ is of rapid decay. Thus, $E_{s, f}^{P}-c_{P}^{T} E_{s, f}^{P}$ is of rapid decay on standard Siegel sets, and then the automorphic form

$$
\wedge^{T} E_{s, f}^{P}=E_{s, f}^{P}-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right)
$$

is of rapid decay on all Siegel sets.
As in the discussion immediately prior to [3.10.2], for a root $\alpha$ of $G$, for $a \in M_{\mathbb{A}}^{\min }$, let $a^{\alpha}=e^{\alpha(\log |a|)}$. For $g \in G_{\mathbb{A}}$, in an Iwasawa decomposition let $g \in N^{\min } \cdot a_{g} \cdot K_{\mathbb{A}}$ with $a_{g} \in M_{\mathbb{A}}^{\min }$, so we can consider the functions $g \longrightarrow a_{g}^{\alpha}$. The ambiguity of $a$ by $M_{\mathbb{A}}^{\min } \cap K_{\mathbb{A}}$ does not affect the value of this function. In the non-self-associate case, let $\alpha, \beta$ be the simple positive roots corresponding to $P$ and $Q$, in the sense that $N^{P}$ contains the $\alpha$ root subgroup

$$
N^{\alpha}=\left\{n=e^{x}: x \in \mathfrak{n}, a x a^{-1}=a^{\alpha} \cdot x\right\} \quad\left(\mathfrak{n}=\text { Lie algebra of } N^{\min }\right)
$$

and $N^{Q}$ contains the $\beta$ root subgroup $N^{\beta}$. Because $f$ is a cuspform and $P$ is not self-associate, only a single Bruhat cell contributes to $c_{P} E_{s, f}^{P}$, and $c_{P} E_{s, f}^{P}=\varphi_{s, f}^{P}$, which is rapidly decreasing on standard Siegel sets as $a_{g}^{\gamma} \rightarrow+\infty$ for any simple positive root $\gamma \neq \alpha$, because $f$ is a cuspform in a strong sense. Similarly, only a single Bruhat cell (corresponding to the longest Weyl element) contributes to the constant term $c_{Q} E_{s, f}^{P}$, which similarly is rapidly decreasing on standard Siegel sets as $a_{g}^{\gamma} \rightarrow+\infty$ for any simple (positive) root $\gamma \neq \beta$. Thus, the truncation

$$
\wedge^{T} E_{\varphi}^{P}=E_{\varphi}^{P}-\Psi^{P}\left(c_{P}^{T} E_{\varphi}^{P}\right)-\Psi^{Q}\left(c_{Q}^{T} E_{\varphi}^{P}\right)
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

has decay properties as follows. If $a_{g}^{\gamma} \rightarrow+\infty$ for $\gamma$ other than $\alpha, \beta$, then all three terms on the right-hand side are of rapid decay in standard Siegel sets. If $\alpha \rightarrow+\infty$, then each of the two expressions $E_{\varphi}^{P}-\Psi^{P}\left(c_{P}^{T} E_{\varphi}^{P}\right)$ and $\Psi^{Q}\left(c_{Q}^{T} E_{\varphi}^{P}\right)$ is of rapid decay. If $\beta \rightarrow+\infty$, then each of the two expressions $E_{\varphi}^{P}-\Psi^{Q}\left(c_{Q}^{T} E_{\varphi}^{P}\right)$ and $\Psi^{P}\left(c_{P}^{T} E_{\varphi}^{P}\right)$ is of rapid decay. Thus, as the value of any simple positive root goes to $+\infty$ in a standard Siegel set, the truncation goes rapidly to zero.

Let $h=h_{1} \otimes h_{2}$ be a another cuspform on $M=M^{P}$. Let $\langle f, h\rangle_{M^{P}}$ be the inner product on the quotient $Z_{\mathbb{A}}^{M} M_{k} \backslash M_{\mathbb{A}}$. For brevity, write $f^{w}=\left(f_{1} \otimes f_{2}\right)^{w}=f_{2} \otimes f_{1}$.
[3.14.2] Theorem: (Maaß-Selberg relations) The hermitian inner product $\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{r, h}^{P}\right\rangle$ of truncations of two cuspidal-data Eisenstein series for maximal proper parabolic $P$ is given as follows. For $P$ self-associate,

$$
\begin{aligned}
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{r, h}^{P}\right\rangle= & \langle f, h\rangle_{M} \cdot \frac{T^{s+\bar{r}-1}}{s+\bar{r}-1}+\left\langle f, h^{w}\right\rangle_{M} \cdot \overline{c_{r, h}} \frac{T^{s+(1-\bar{r})-1}}{s+(1-\bar{r})-1} \\
& +\left\langle f^{w}, h\right\rangle_{M} \cdot c_{s, f} \frac{T^{(1-s)+\bar{r}-1}}{(1-s)+\bar{r}-1}+\left\langle f^{w}, h^{w}\right\rangle_{M} \cdot c_{s, f} \overline{c_{r, h}} \frac{T^{(1-s)+(1-\bar{r})-1}}{(1-s)+(1-\bar{r})-1}
\end{aligned}
$$

For $P$ not self-associate, that is, for $r_{1} \neq r_{2}$,

$$
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{r, h}^{P}\right\rangle=\langle f, h\rangle_{M^{P}} \cdot \frac{T^{s+\bar{r}-1}}{s+\bar{r}-1}+\left\langle f^{w}, h^{w}\right\rangle_{M^{Q}} \cdot c_{s, f}^{Q} \overline{c_{r, h}^{Q}} \frac{T^{(1-s)+(1-\bar{r})-1}}{(1-s)+(1-\bar{r})-1}
$$

[3.14.3] Remark: The expression for the not-self-associate case is that of the self-associate case with the middle two terms missing. In the non-self-associate case the inner products $\left\langle f^{w}, h\right\rangle$ and $\left\langle f, h^{w}\right\rangle$ would not make sense, because in that case $w M^{P} w^{-1} \neq M^{P}$, so the two functions live on different groups.
[3.14.4] Corollary: For maximal proper parabolics $P$ in $G L_{r}$, on the half-plane $\operatorname{Re}(s) \geq 1 / 2$ an Eisenstein series $E_{s, f}^{P}$ with cuspidal data $f$ has no poles if $P$ is not self-associate. If $P$ is self-associate, the only possible poles are on the real line, and only occur if $\left\langle f, f^{w}\right\rangle_{M} \neq 0$. In that case, any pole is simple, and the residue is square-integrable. In particular, taking $f=f_{o} \otimes f_{o}$

$$
\left\langle\operatorname{Res}_{s_{o}} E_{s, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle=\left\langle f_{o}, f_{o}\right\rangle_{M}^{2} \cdot \operatorname{Res}_{s_{o}} c_{s, f}
$$

Proof: (of theorem) The self-associate case admits a simpler argument, because in this case the truncated Eisenstein series $\wedge^{T} E_{s, f}^{P}$ is itself a pseudo-Eisenstein series

$$
\wedge^{T} E_{s, f}^{P}=\Psi^{P}\left(\varphi_{s, f}\right)-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right)=\Psi^{P}\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right)
$$

As in smaller cases [1.11], [2.10], the pseudo-Eisenstein series made from the tail of the constant term integrates to zero against the truncated Eisenstein series: fortunately, for cuspidal-data Eisenstein series this fact need not refer to subtle reduction theory, but only needs the $r_{1}=r_{2}$ instance of the following:
[3.14.5] Lemma: With $P=P^{r_{1}, r_{2}}, Q=P^{r_{2}, r_{2}}$, and $w=\left(\begin{array}{cc}0 & 1_{r_{2}} \\ 1_{r_{1}} & 0\end{array}\right)$,

$$
h^{Q}(w n \cdot m) \leq h^{P}(m)^{-1} \quad\left(\text { for all } n \in N_{\mathbb{A}}^{P}, m \in M_{\mathbb{A}}^{P}\right)
$$

Proof: This lemma observes a qualitative aspect of Iwasawa decompositions of special elements. First,

$$
h^{Q}(w n m)=h^{Q}\left(w m w^{-1} \cdot w m^{-1} n m\right)=\delta^{Q}\left(w m w^{-1}\right) \cdot h^{Q}\left(w m^{-1} n m\right)=\delta^{P}(m)^{-1} \cdot h^{Q}\left(w m^{-1} n m\right)
$$

since $w$ conjugates $M^{P}$ to $M^{Q}$, and inverts the modular function. Since $M^{P}$ normalizes $N^{Q}$, it suffices to prove $h^{Q}(w n) \leq 1$ for all $n \in N_{\mathbb{A}}^{P}$. Because $h^{Q}$ is right $K_{\mathbb{A}}$-invariant, this is equivalent to

$$
1 \geq h^{Q}\left(w n w^{-1}\right)=h^{Q}\left(\begin{array}{cc}
1_{r_{2}} & 0 \\
x & 1_{r_{1}}
\end{array}\right) \quad\left(\text { with } x r_{1} \text {-by- } r_{2}\right)
$$

In fact, we claim that the same inequality holds locally at every place $v$ with the local analogues $h_{v}^{Q}$ and $h_{v}^{P}$. At all places $v$, right multiplication of $w n w^{-1}$ by $K_{v}$ produces $q=\left(\begin{array}{cc}a & b \\ 0 & d\end{array}\right)$ with $a, d$ upper-triangular.

For finite $v$, we imagine achieving the Iwasawa decomposition in stages, first putting the bottom row into the correct shape, then the next-to-bottom, and so on. To begin, right multiply by $k_{1} \in K_{v}$ to put the $g c d$ of the bottom row of $x$ and of 1 into the far right entry of the bottom row of the whole, and to replace the bottom row of $x$ by 0 's. This can be done without disturbing the left $r_{1}$-by- $\left(r_{1}-1\right)$ part of the lower right block of $w n w^{-1}$. Next, without disturbing the adjusted bottom row, right multiply by $k_{2} \in K_{v}$ to put the $g c d$ of the next-to-bottom row of (the new) $x$ and 1 in the next-to-last entry of the next-to-bottom row of the whole, and to replace the next-to-bottom row of (the new) $x$ by 0 's. Continuing, every diagonal entry of $d$ will be a $g c d$ of some $v$-adic numbers and 1 , so not divisible by the local parameter $\varpi$. Thus, $|\operatorname{det} d|_{v} \geq 1$. At the same time, the entries of $a$ are among the entries of an element of $K_{v}$, so are all $v$-integral, and $|\operatorname{det} a|_{v} \leq 1$. Thus, $h_{v}^{Q}(w n)=h_{v}^{Q}\left(w n w^{-1}\right) \leq 1$.

Somewhat analogously, for archimedean $v$, the $i^{\text {th }}$ diagonal entry of $d$ is lengths of vectors with entries including the diagonal 1's in the lower-right block of $w n w^{-1}$. Thus, all the diagonal entries of $d$ will be at least 1 in size, and certainly $|\operatorname{det} d| \geq 1$. At the same time, the rows of $a$ are fragments of rows of a matrix in $K_{v}$, so have length at most 1 . The absolute value of the determinant of $a$ is the volume of the parallelopiped spanned by those rows, so is at most 1 .

Returning to the computation of the inner product in the self-associate case, the integral of $\Psi^{P}\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right)$ against $\Psi^{P}\left(c_{P}^{T} E_{r, h}^{P}\right)$ unwinds to the integral of $\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}$ against $\Psi^{P}\left(c_{P}^{T} E_{r, h}^{P}\right)$, which is then the integral of $\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}$ against $c_{P}\left(\Psi^{P}\left(c_{P}^{T} E_{r, h}^{P}\right)\right)$. By construction, $\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right)(m)$ is supported where $h^{P}(m) \geq T$, for $m \in M^{P}$. The proof of [3.11.3] and remarks in [3.11.8] apply as well to pseudo-Eisenstein series, so

$$
c_{P}\left(\Psi^{P}\left(c_{P}^{T} E_{r, h}^{P}\right)\right)(m)=c_{P}^{T} E_{r, h}^{P}+\int_{N_{\mathbb{A}}^{P}}\left(c_{P}^{T} E_{r, h}^{P}\right)(w n \cdot m) d n \quad\left(\text { for } m \in M^{P}\right)
$$

By definition of the truncation, the integrand is 0 unless $h^{P}(w n \cdot m) \geq T$. The lemma gives $h^{P}(w n \cdot m) \leq$ $h^{P}(m)^{-1}$. Thus, for $T>1$, there is no overlap of supports of $\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}$ and the second part of the constant term. That is,

$$
\begin{aligned}
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{r, h}^{P}\right\rangle & =\int \Psi^{P}\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right) \cdot \overline{\Psi^{P}\left(c_{P}^{T} E_{r, h}^{P}\right)}=\int\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right) \cdot \overline{c_{P}^{T} E_{r, h}^{P}} \\
& =\int \Psi^{P}\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right) \cdot \overline{E_{r, h}^{P}}=\left\langle\wedge^{T} E_{s, f}^{P}, E_{r, h}^{P}\right\rangle
\end{aligned}
$$

Unwind the truncated Eisenstein series:

$$
\begin{aligned}
\left\langle\wedge^{T} E_{s, f}^{P}, E_{r, h}^{P}\right\rangle= & \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi^{P}\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right) \overline{E_{r, h}^{P}}=\int_{Z^{+} \cdot P_{k} \backslash G_{\mathbb{A}}}\left(\varphi_{s, f}-c_{P}^{T} E_{s, f}^{P}\right) \overline{E_{r, h}^{P}} \\
& =\int_{Z^{+} \cdot P_{k} \backslash G_{\mathbb{A}}}\left\{\begin{array}{cc}
-c_{s, f} \varphi_{1-s, f w} & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{E_{r, h}^{P}}
\end{aligned}
$$

This is

$$
\begin{gathered}
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}}\left\{\begin{array}{cc}
-c_{s, f} \varphi_{1-s, f^{w}} & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{c_{P} E_{r, h}^{P}} \\
=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}}\left\{\begin{array}{cc}
-c_{s, f} \varphi_{1-s, f w} & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{\left(\varphi_{r, h}+c_{1-r, h^{w}} \varphi_{1-r, h^{w}}\right)}
\end{gathered}
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

Since the integrand is now left $N_{\mathbb{A}}$-invariant, $Z_{\mathbb{A}}$-invariant, and right $K_{\mathbb{A}}$-invariant, this integral may be computed as an integral over the Levi component $M^{P}$, using the Iwasawa decomposition, noting that the Haar integral on $G_{\mathbb{A}}$ in such coordinates is

$$
\int_{G_{\mathbb{A}}} f(g) d g=\int_{N_{\mathbb{A}}} \int_{M_{\mathbb{A}}} \int_{K_{\mathbb{A}}} f(n m k) d n \frac{d m}{\delta^{P}(m)} d k
$$

Then

$$
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{r, h}^{P}\right\rangle=\int_{Z_{\mathbb{A}} M_{k} \backslash M_{\mathbb{A}}}\left\{\begin{array}{cc}
-c_{s, f} \varphi_{1-s, f^{w}} & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{\left(\varphi_{r, h}+c_{1-r, h^{w}} \varphi_{1-r, h^{w}}\right)} \frac{d m}{\delta^{P}(m)}
$$

This gives the four terms of the theorem for the self-associate case. We evaluate one in detail, as follows.
Use $Z^{+} M_{k} \backslash M_{\mathbb{A}} \approx Z^{+} \backslash A_{P}^{+} \times M_{k} \backslash M^{1}$ and the $A_{P}^{+}$-invariance of $f$, parametrizing $Z^{+} \backslash A_{P}^{+}$by

$$
t \longrightarrow a_{t}=\left(\begin{array}{cc}
\delta(t)^{1 / r_{2}} \cdot 1_{r_{1}} & 0 \\
0 & 1_{r_{2}}
\end{array}\right) \quad(\text { for } t>0)
$$

with the diagonal imbedding of the ray $(0, \infty)$ into the archimedean part of the ideles, so that $\delta^{P}\left(a_{t}\right)=t$. Then, for example,

$$
\begin{gathered}
\int_{m \in Z_{\mathbb{A}} M_{k} \backslash M_{\mathbb{A}}: h^{P}(m)<T} \varphi_{s, f} \cdot \overline{\varphi_{1-r, h^{w}}} \frac{d m}{\delta^{P}(m)}=\int_{0<t<T, m \in Z_{\mathbb{A}} M_{k} \backslash M^{1}} t^{s} \cdot f\left(a_{t} m_{1}\right) \cdot \overline{t^{1-r} \cdot h^{w}\left(a_{t} m_{1}\right)} \frac{d t}{t} d m_{1} \\
=\int_{0}^{T} t^{s+(1-r)} \frac{d t}{t} \cdot \int_{Z_{\mathbb{A}} M_{k} \backslash M^{1}} f\left(m_{1}\right) \cdot \overline{h^{w}\left(m_{1}\right)} d m_{1}=\frac{T^{s+(1-r)}}{s+(1-r)} \cdot\left\langle f, h^{w}\right\rangle
\end{gathered}
$$

The other three integrals are evaluated in identical fashion.
In the non-self-associate case, invoking the Lemma in similar fashion,

$$
\begin{cases}\left\langle E_{s, f}^{P}-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right), \Psi^{P}\left(c_{P}^{T} E_{s, h}^{P}\right)\right\rangle & =0 \\ \left\langle E_{s, f}^{P}-\Psi^{Q}\left(c_{Q}^{T} E_{s, f}^{P}\right), \Psi^{Q}\left(c_{Q}^{T} E_{s, h}^{P}\right)\right\rangle & =0 \\ \left\langle\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right), \Psi^{Q}\left(c_{Q}^{T} E_{s, h}^{P}\right)\right\rangle & =0\end{cases}
$$

so the inner product of the truncated Eisenstein series is

$$
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{r, h}^{P}\right\rangle=\left\langle E_{s, f}^{P}-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right), E_{r, h}^{P}\right\rangle+\left\langle\Psi^{Q}\left(c_{Q}^{T} E_{s, f}^{P}\right), \Psi^{Q}\left(c_{Q}^{T} E_{r, h}^{P}\right)\right\rangle
$$

The pairings unwind. First,

$$
\left\langle E_{s, f}^{P}-\Psi^{P}\left(c_{P}^{T} E_{s, f}^{P}\right), E_{r, h}^{P}\right\rangle=\left\langle\Psi^{P}\left(c_{P} E_{s, f}^{P}-c_{P}^{T} E_{s, f}^{P}\right), E_{r, h}^{P}\right\rangle=\int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}}\left\{\begin{array}{cc}
0 & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f}^{P} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{E_{r, h}^{P}}
$$

Because of the left $N_{\mathbb{A}}$-invariance of the first part of the integral, this is

$$
\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}}\left\{\begin{array}{cc}
0 & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f}^{P} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{c_{P} E_{r, h}^{P}}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}}\left\{\begin{array}{cc}
0 & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f}^{P} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{\varphi_{r, h}^{P}}
$$

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Again, the integrand is left $N_{\mathbb{A}}$-invariant and right $K_{\mathbb{A}}$-invariant, so may be computed as an integral over the Levi component using the Iwasawa decomposition: it is

$$
\int_{Z_{\mathbb{A}}^{M} M_{k} \backslash M_{\mathbb{A}}}\left\{\begin{array}{cc}
0 & \left(\text { for } h^{P} \geq T\right) \\
\varphi_{s, f}^{P} & \left(\text { for } h^{P}<T\right)
\end{array}\right\} \cdot \overline{\varphi_{r, h}^{P}} \frac{d m}{\delta^{P}(m)}
$$

giving one term in the non-self-associate case. The other pairing unwinds similarly, and becomes

$$
\left\langle\Psi^{Q}\left(c_{Q}^{T} E_{s, f}^{P}\right), \Psi^{Q}\left(c_{Q}^{T} E_{r, h}^{P}\right)\right\rangle=\int_{Z_{\mathbb{A}} M_{k}^{Q} \backslash M_{\mathbb{A}}^{Q}}\left\{\begin{array}{cc}
c_{s, f}^{Q} \varphi_{1-s, f^{w}}^{Q} & \left(\text { for } h^{Q} \geq T\right) \\
0 & \left(\text { for } h^{Q}<T\right)
\end{array}\right\} \cdot \overline{c_{r, h}^{Q} \varphi_{1-r, h^{w}}^{Q}} \frac{d m}{\delta^{Q}(m)}
$$

giving the second term of the theorem for the not self-associate case.
Proof: (of corollary). From the theory of the constant term, the only possible poles of the Eisenstein series are at poles of the constant terms, which in this case means a pole of $c_{s, f}$. Invoke the Maaß-Selberg relation with $r=s$ and $h=f$. In the non-self-associate case this is

$$
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{s, f}^{P}\right\rangle=\langle f, f\rangle \frac{T^{2 \sigma-1}}{2 \sigma-1}+\left\langle f^{w}, h^{w}\right\rangle\left|c_{s}\right|^{2} \frac{T^{1-2 \sigma}}{1-2 \sigma} \quad \quad(\text { with } \sigma=\operatorname{Re}(s))
$$

The non-self-associate case is slightly unlike the simple case of $G L_{2}$, in that the inner product of truncated Eisenstein series is missing the two middle terms which for $G L_{2}$ made a pole possible. Specifically, in the non-self-associate case, let $s_{o}=\sigma_{o}+i t_{o}$ be an alleged pole $s_{o}$ of $c_{s}$ of order $\ell$ in that half-plane. Letting $s=\sigma_{o}+i t$ approach $s_{o}$ vertically the left-hand side of the relation is asymptotic to a positive multiple of $t^{-2 \ell}$, while on right-hand side only the second of the two terms blows up at all. In particular, that expression

$$
\left|c_{s}\right|^{2} \cdot\left\langle f^{w}, f^{w}\right\rangle \cdot \frac{T^{1-2 \sigma}}{1-2 \sigma}
$$

is asymptotic to a negative multiple of $t^{-2 \ell}$, since $\sigma=\operatorname{Re}(s)>\frac{1}{2}$. Thus, there is no pole in that half-plane.
Similarly, in the self-associate case, for there to be any pole in the right half-plane the two middle terms on the right-hand side of the relation must not vanish, or the same contradiction occurs, so $\left\langle f, f^{w}\right\rangle$ must be non-zero, and the alleged pole must be on the real axis, and must be simple: if any of these conditions fail, the middle terms cannot keep up with the negative value of the fourth term. Letting $f=f_{o} \otimes f_{o}$ with real-valued $f_{o}$, we have $f^{w}=f$ and $\langle f, f\rangle=\left\langle f_{o}, f_{o}\right\rangle \cdot\left\langle f_{o}, f_{o}\right\rangle$. Letting $s=\sigma+i t$,

$$
\left\langle\wedge^{T} E_{s, f}^{P}, \wedge^{T} E_{s, f}^{P}\right\rangle=\left\langle f_{o}, f_{o}\right\rangle^{2} \frac{T^{2 \sigma-1}}{2 \sigma-1}+\left\langle f_{o}, f_{o}\right\rangle^{2} \overline{c_{s, f}} \frac{T^{2 i t}}{2 i t}+\left\langle f_{o}, f_{o}\right\rangle^{2} c_{s, f} \frac{T^{-2 i t}}{-2 i t}+\left\langle f_{o}, f_{o}\right\rangle^{2}\left|c_{s, f}\right|^{2} \frac{T^{1-2 \sigma}}{1-2 \sigma}
$$

Multiplying through by $t^{2}=(i t)(-i t)$ and taking the limit as $t \rightarrow 0$ gives

$$
\left\langle\operatorname{Res}_{\sigma} \wedge^{T} E_{s, f}^{P}, \operatorname{Res}_{\sigma} \wedge^{T} E_{s, f}^{P}\right\rangle=\left\langle f_{o}, f_{o}\right\rangle^{2} \overline{\operatorname{Res}_{\sigma} c_{s, f}} \cdot \frac{1}{2}+\left\langle f_{o}, f_{o}\right\rangle^{2} \operatorname{Res}_{\sigma} c_{s, f} \cdot \frac{1}{2}+\left\langle f_{o}, f_{o}\right\rangle^{2}\left|\operatorname{Res}_{\sigma} c_{s, f}\right|^{2} \frac{T^{1-2 \sigma}}{1-2 \sigma}
$$

Letting $T \rightarrow+\infty$ causes the last term to go to zero, and yields the indicated finite limit in the self-associate case, since $c_{\bar{s}, f}=\overline{c_{s, f}}$ and the supposed pole is on the real axis.
[3.14.6] Claim: All residues of Eisenstein series $E_{s, f}^{P}$ are orthogonal to cuspforms.
Proof: There exists an orthogonal basis for cuspforms consisting of strong-sense cuspforms [3.7.3]. Thus, by the theory of the constant term [8.3], that basis consists of functions of rapid decay in Siegel sets. Eisenstein series are of moderate growth even when analytically continued, so the integral for $\left\langle E_{s, f}, F\right\rangle$ with strongsense cuspform $F$ is absolutely convergent, and is 0 . By properties of vector-valued integrals [14.1] and holomorphic vector-valued functions [15.2], taking residues commutes with the integral, so the integral of any residue against a cuspform is 0 , whether or not that residue is square-integrable.

### 3.15 Minimal-parabolic decomposition

The harmonic analysis required to express pseudo-Eisenstein series in terms of genuine Eisenstein series reduces to Fourier transform on Euclidean spaces. Here we treat the extreme case $P=P^{\text {min }}$, where no cuspidal data enters. We consider minimal-parabolic pseudo-Eisenstein series $\Psi_{\varphi}$ with test function data $\varphi$. All the $r$ ! functional equations [3.12.1] of the genuine Eisenstein series are needed to obtain the expression of pseudo-Eisenstein series as integrals of Eisenstein series.

Let $A^{+}$be the archimedean split component of $M=M^{\min }=M^{P}$, that is, the image of $r$ copies of $(0,+\infty)$ imbedded diagonally on the archimedean factors $M_{\infty}=\prod_{v \mid \infty} M_{v}$ of $M_{\mathbb{A}}$. With $M^{1}$ the subgroup of $M=M_{\mathbb{A}}^{\min }$ with diagonal entries $m_{i}$ all satisfying $\left|m_{i}\right|=1$, we have $M_{\mathbb{A}}^{\min }=A^{+} \cdot M^{1}$. Via the exponential $\mathbb{R} \rightarrow(0,+\infty)$, we have $A^{+} \approx \mathbb{R}^{r}$ and $Z^{+} \backslash A^{+} \approx \mathbb{R}^{r-1}$. Spectral decomposition along the Euclidean space $Z^{+} \backslash A^{+}$and the functional equations of the minimal-parabolic Eisenstein series $E_{s}=E_{s}^{P}$ yield the spectral decomposition of minimal-parabolic pseudo-Eisenstein series. Let $\langle$,$\rangle be the invariant pairing on the Lie$ algebra $\mathfrak{q}$ of $Z^{+} \backslash A^{+}$, as in [3.10.2], where it was shown that $E_{s}$ converges nicely in the cone

$$
\left\{s \in \mathfrak{q} \otimes_{\mathbb{R}} \mathbb{C}:\langle\alpha, \operatorname{Re}(s)-2 \rho\rangle>0, \text { for all simple positive roots } \alpha\right\}
$$

For simplicity, we only consider right $K_{\mathbb{A}}$-invariant $\Psi_{\varphi}$ with trivial central character formed from $\varphi \in$ $\mathcal{D}\left(Z^{+} \backslash A^{+}\right)$. Further, suppose that the pseudo-Eisenstein series $\Psi_{\varphi}$ is orthogonal to all residues of $E_{\rho+s}$ in the cone

$$
\left\{s \in \mathfrak{q} \otimes_{\mathbb{R}} \mathbb{C}:\langle\alpha, \operatorname{Re}(s)\rangle>0, \text { for all simple positive roots } \alpha\right\}
$$

[3.15.1] Theorem: $\Psi_{\varphi}$ is an integral of Eisenstein series:

$$
\Psi_{\varphi}=\frac{1}{r!(2 \pi i)^{r-1}} \int_{i \mathfrak{a}^{*}}\left\langle\Psi_{\varphi}, E_{\rho+s}\right\rangle \cdot E_{\rho+s} d s
$$

[3.15.2] Remark: From [3.8.1], the pseudo-Eisenstein series $\Psi_{\varphi}$ is compactly supported on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$, and $E_{\rho+s}$ is of moderate growth, so the integral

$$
\left\langle\Psi_{\varphi}, E_{\rho+s}\right\rangle=\int_{Z^{+} G_{k} \backslash G_{A}} \Psi_{\varphi} \cdot \overline{E_{\rho+s}}
$$

implied by $\left\langle\Psi_{\varphi}, E_{\rho+s}\right\rangle$, while not an inner product, converges well.
Proof: To decompose right $K_{\mathbb{A}}$-invariant pseudo-Eisenstein series as integrals of minimal-parabolic Eisenstein series, begin with Fourier transform on the Lie algebra $\mathfrak{q} \approx \mathbb{R}^{r-1}$ of $Z^{+} \backslash A^{+}$. Let $\langle\rangle:, \mathfrak{q}^{*} \times \mathfrak{q} \rightarrow \mathbb{R}$ be the $\mathbb{R}$-bilinear pairing of $\mathfrak{q}$ with its $\mathbb{R}$-linear dual $\mathfrak{q}^{*}$. For $f \in \mathcal{D}(\mathfrak{q})$, the Fourier transform and inversion are

$$
\widehat{f}(\xi)=\int_{\mathfrak{q}} e^{-i\langle x, \xi\rangle} f(x) d x \quad f(x)=\frac{1}{(2 \pi)^{r-1}} \int_{\mathfrak{q}^{*}} e^{i\langle x, \xi\rangle} \widehat{f}(\xi) d \xi
$$

Let $\exp : \mathfrak{q} \rightarrow Z^{+} \backslash A^{+}$be the Lie algebra exponential, and $\log : Z^{+} \backslash A^{+} \rightarrow \mathfrak{q}$ the inverse. Given $\varphi \in \mathcal{D}\left(Z^{+} \backslash A^{+}\right)$, let $f=\varphi \circ \exp$ be the corresponding function in $\mathcal{D}(\mathfrak{q})$. The Mellin transform $\mathcal{M} \varphi$ of $\varphi$ is the Fourier transform of $f$ :

$$
\mathcal{M} \varphi(i \xi)=\widehat{f}(\xi)
$$

Mellin inversion is Fourier inversion in these coordinates:

$$
\varphi(\exp x)=f(x)=\frac{1}{(2 \pi)^{r-1}} \int_{\mathfrak{q}^{*}} e^{i\langle\xi, x\rangle} \widehat{f}(\xi) d \xi=\frac{1}{(2 \pi)^{r-1}} \int_{\mathfrak{q}^{*}} e^{i\langle\xi, x\rangle} \mathcal{M} \varphi(i \xi) d \xi
$$

Extend the pairing $\langle$,$\rangle on \mathfrak{q}^{*} \times \mathfrak{q}$ to a $\mathbb{C}$-bilinear pairing on the complexification. Use the convention

$$
(\exp x)^{i \xi}=e^{i\langle\xi, x\rangle}=e^{\langle i \xi, x\rangle}
$$

With $a=\exp x \in Z^{+} \backslash A^{+}$, Mellin inversion is

$$
\varphi(a)=\frac{1}{(2 \pi)^{r-1}} \int_{\mathfrak{q}^{*}} a^{i \xi} \mathcal{M} \varphi(i \xi) d \xi=\frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{q}^{*}} a^{s} \mathcal{M} \varphi(s) d s \quad\left(\text { with } a \in Z^{+} \backslash A^{+} \text {and } s=i \xi\right)
$$

With this notation, the Mellin transform itself is

$$
\left.\mathcal{M} \varphi(s)=\int_{Z^{+} \backslash A^{+}} a^{-s} \varphi(a) d a \quad \text { (with } s \in i \mathfrak{q}^{*}\right)
$$

Since $\varphi$ is a test function, its Fourier-Mellin transform is entire on $\mathfrak{q}^{*} \otimes_{\mathbb{R}} \mathbb{C}$. (It is in the Paley-Wiener space.) Thus, for any $\sigma \in \mathfrak{q}^{*}$, Mellin inversion can be written

$$
\varphi(a)=\frac{1}{(2 \pi i)^{r-1}} \int_{\sigma+i \mathbf{q}^{*}} a^{s} \mathcal{M} \varphi(s) d s
$$

Via Iwasawa, identify $Z^{+} N_{\mathbb{A}} M^{1} \backslash G_{\mathbb{A}} / K_{\mathbb{A}} \approx A^{+}$, and let $g \rightarrow a(g)$ be the function that picks out the $A^{+}$ component. For $\sigma \in \mathfrak{q}^{*}$ suitable for convergence [3.10.1], the following rearrangement is legitimate:

$$
\begin{gathered}
\Psi_{\varphi}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(a(\gamma \circ g))=\sum_{\gamma \in P_{k} \backslash G_{k}} \frac{1}{(2 \pi i)^{r-1}} \int_{\sigma+i \mathbf{q}^{*}} \mathcal{M} \varphi(s) a(\gamma g)^{s} d s \\
=\frac{1}{(2 \pi i)^{r-1}} \int_{\sigma+i \boldsymbol{q}^{*}} \mathcal{M} \varphi(s)\left(\sum_{\gamma \in P_{k} \backslash G_{k}} a(\gamma g)^{s}\right) d s=\frac{1}{(2 \pi i)^{r-1}} \int_{\sigma+i \boldsymbol{q}^{*}} \mathcal{M} \varphi(s) \cdot E_{s}(g) d s
\end{gathered}
$$

Anticipating the invocation of the functional equations, using the rapid vertical decay of $\mathcal{M} \varphi(s)$, we can move the ( $r-1$ )-fold integration to $\rho+i q^{*}$. For simplicity, we assume $\Psi_{\varphi}$ is orthogonal to any (multi-) residues:

$$
\Psi_{\varphi}=\frac{1}{(2 \pi i)^{r-1}} \int_{i q^{*}} \mathcal{M} \varphi(\rho+s) \cdot E_{\rho+s} d s
$$

This does express the pseudo-Eisenstein series as a superposition of Eisenstein series. However, the coefficients $\mathcal{M} \varphi$ are not expressed in terms of $\Psi_{\varphi}$ itself. This is rectified using the functional equations of $E_{\rho+s}$, as follows.

Since $d n d m d a d k / a^{2 \rho}$ with $n \in N_{\mathbb{A}}, m \in M^{1}, a \in A^{+}, k \in K_{\mathbb{A}}$ is a Haar measure on $G_{\mathbb{A}}, d a d k / a^{2 \rho}$ is a right $G$-invariant measure on $N_{\mathbb{A}} M^{1} \backslash G_{\mathbb{A}}$, and $d a / a^{2 \rho}$ is the associated measure on $N_{A} M^{1} \backslash G_{\mathbb{A}} / K_{\mathbb{A}} \approx A^{+}$ and it descends to $Z^{+} \backslash A^{+}$. In the region of convergence, for $\varphi \in \mathcal{D}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$,

$$
\begin{aligned}
& \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f \cdot E_{\rho+s}=\int_{Z^{+} P_{k} \backslash G_{\mathbb{A}}} f \cdot a^{\rho+s}=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \int_{N_{k} \cap N_{\mathbb{A}}} f(n g) a(n g)^{\rho+s} d n d g \\
&=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} c_{P} f \cdot a^{\rho+s}=\int_{Z^{+} \backslash A^{+}} c_{P} f(a) \cdot a^{\rho+s} \frac{d a}{a^{2 \rho}}=\int_{Z^{+} \backslash A^{+}} c_{P} f(a) \cdot a^{-(\rho-s)} d a=\mathcal{M} c_{P} f(\rho-s)
\end{aligned}
$$

That is, with $f=\Psi_{\varphi}$,

$$
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot E_{\rho+s}=\mathcal{M} c_{P} \Psi_{\varphi}(\rho-s)
$$

On the other hand, a similar unwinding of the pseudo-Eisenstein series, and recollection of the constant term $c_{P} E_{\rho+s}$ from [3.10.3], gives

$$
\begin{gathered}
\int_{Z^{+} G_{k} \backslash G_{A}} \Psi_{\varphi} \cdot E_{\rho+s}=\int_{Z^{+} \backslash A^{+}} \varphi(a) \cdot c_{P} E_{\rho+s}(a) \frac{d a}{a^{2 \rho}}=\int_{Z^{+} \backslash A^{+}} \varphi(a) \cdot \sum_{w} c_{w, s} a^{\rho+w \cdot s} \frac{d a}{a^{2 \rho}} \\
=\sum_{w} c_{w, s} \int_{A^{+}} \varphi(a) a^{-(\rho-w \cdot s)} d a=\sum_{w} c_{w, s} \mathcal{M} \varphi(\rho-w \cdot s)
\end{gathered}
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

Combining these,

$$
\mathcal{M} c_{P} \Psi_{\varphi}(\rho-s)=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot E_{\rho+s}=\sum_{w} c_{w, s} \mathcal{M} \varphi(\rho-w \cdot s)
$$

Replacing $s$ by $-s$,

$$
\mathcal{M} c_{P} \Psi_{\varphi}(\rho+s)=\int_{Z+G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot E_{\rho-s}=\sum_{w} c_{w,-s} \mathcal{M} \varphi(\rho+w \cdot s)
$$

The Eisenstein series $E_{s}$ behaves reasonably under complex conjugation: $\overline{E_{s}}=E_{\bar{s}}$. This is visible in the region of convergence, and persists under analytic continuation, since $\overline{E_{\bar{s}}}=E_{s}$ is an equality of meromorphic functions. Thus, the previous equality becomes

$$
\mathcal{M} c_{P} \Psi_{\varphi}(\rho+s)=\int_{Z+G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot \overline{E_{\rho+s}}=\sum_{w} c_{w,-s} \mathcal{M} \varphi(\rho+w \cdot s)
$$

Behavior under complex conjugation is inherited by the constant term along $P$ :

$$
\sum_{w} \overline{c_{w, s}} \cdot \overline{a^{\rho+w \cdot s}}=c_{P} \overline{E_{\rho+s}}=c_{P} E_{\rho+\bar{s}}=\sum_{w} c_{w, \bar{s}} \cdot a^{\rho+w \cdot \bar{s}}
$$

Since $\overline{a^{\rho+w \cdot s}}=a^{\rho+w \cdot \bar{s}}$, this gives $\overline{c_{w, s}}=c_{w, \bar{s}}$. For $\rho+s$ on the unitary hyperplane $\rho+i \mathfrak{a}^{*}$, conveniently $\bar{s}=-s$, and $c_{w,-s}=\overline{c_{w, s}}$, so

$$
\mathcal{M} c_{P} \Psi_{\varphi}(\rho+s)=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\varphi} \cdot \overline{E_{\rho+s}}=\sum_{w} \overline{c_{w, s}} \mathcal{M} \varphi(\rho+w \cdot s)
$$

With these points in hand, average the relation

$$
\Psi_{\varphi}=\frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{q}^{*}} \mathcal{M} \varphi(\rho+s) \cdot E_{\rho+s} d s
$$

over $w \in W$ to convert it a $W$-symmetric expression, thereby to obtain an expression in terms of $c_{P} \Psi_{\varphi}$, using the functional equations:

$$
\Psi_{\varphi}=\frac{1}{|W|} \sum_{w} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{q}^{*}} \mathcal{M} \varphi(\rho+w \cdot s) \cdot E_{\rho+w \cdot s} d s=\frac{1}{|W|} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{a}^{*}}\left(\sum_{w} \frac{1}{c_{w, s}} \mathcal{M} \varphi(\rho+w \cdot s)\right) \cdot E_{\rho+s} d s
$$

Fortunately, from [3.12.6], $\left|c_{w, s}\right|=1$ for $s \in i \mathfrak{q}^{*}$, so this becomes

$$
\Psi_{\varphi}=\frac{1}{|W|} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{a}^{*}}\left(\sum_{w} \overline{c_{w, s}} \mathcal{M} \varphi(\rho+w \cdot s)\right) \cdot E_{\rho+s} d s=\frac{1}{|W|} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{a}^{*}}\left\langle\Psi_{\varphi}, E_{\rho+s}\right\rangle \cdot E_{\rho+s} d s
$$

The cardinality of the Weyl group $W$ is $r$-factorial.

### 3.16 Cuspidal-data decomposition

Now we treat the opposite extreme, the case of maximal proper parabolics $P=P^{r_{1}, r_{2}}$, where cuspforms on Levi components unavoidably enter. With $\delta:(0, \infty) \longrightarrow \mathbb{J}$ the diagonal imbedding at archimedean places, the split component of $P$ is

$$
A_{P}^{+}=\left\{\left(\begin{array}{cc}
\delta\left(t_{1}\right) \cdot 1_{r_{1}} & 0 \\
0 & \delta\left(t_{2}\right) \cdot 1_{r_{2}}
\end{array}\right): t_{1}>0, t_{2}>0\right\}
$$

Let

$$
M^{1}=\left\{m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) \in M_{\mathbb{A}}^{P}:\left|\operatorname{det} m_{1}\right|=1=\left|\operatorname{det} m_{2}\right|\right\}
$$

The family of pseudo-Eisenstein series $\Psi_{\varphi}$ with fixed cuspidal data $f=f_{1} \otimes f_{2}$ on $M^{1}$ with test-function data just on the quotient $Z+\backslash A_{P}^{+} \approx(0, \infty)$ of split components, as in [3.9], constitute the smallest natural vector spaces of functions expressible as integrals of genuine Eisenstein series. In contrast, the pseudo-Eisenstein series with test-function data on $Z_{\mathbb{A}} M_{k} \backslash M_{\mathbb{A}}$, as in [3.8], are smeared out across these smaller spaces of functions.

For simplicity, we only consider everywhere spherical automorphic forms with trivial central character, that is, right $K_{\mathbb{A}}$-invariant and left $Z_{\mathbb{A}}$-invariant functions. Thus, via Iwasawa decomposition, constant terms $c_{P} f$ are identifiable with functions on the quotient of the Levi component of $P$. allowing easier description of the cuspidal data, as follows. Let $f_{1}, f_{2}$ be cuspforms on $G L_{r_{1}}(\mathbb{A})$ and $G L_{r_{2}}(\mathbb{A})$, right invariant by the standard maximal compacts everywhere, themselves with trivial central characters. We require that $f_{1}$ and $f_{2}$ be eigenfunctions for all the spherical Hecke algebras, including the archimedean places. That is, $f_{1}$ and $f_{2}$ are cuspforms in a strong sense, beyond vanishing of constant terms. The theory of the constant term [8.3] shows that cuspforms in this strong sense are of rapid decay. Then $f=f_{1} \otimes f_{2}$ is a function on $G L_{r_{1}}(\mathbb{A}) \times G L_{r_{2}}(\mathbb{A}) \approx M_{\mathbb{A}}^{P}$. For a test function $\eta$ on the ray $(0, \infty)$, define

$$
\varphi(z n m k)=\varphi_{\eta, f}(z n m k)=\eta\left(\frac{\left|\operatorname{det} m_{1}\right|^{r_{2}}}{\left|\operatorname{det} m_{2}\right|^{r_{1}}}\right) \cdot f_{1}\left(m_{1}\right) \cdot f_{2}\left(m_{2}\right)
$$

with $m=\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right) \in M_{\mathbb{A}}^{P}, z \in Z^{+}, n \in N_{\mathbb{A}}, k \in K_{\mathbb{A}}$, with corresponding pseudo-Eisenstein series

$$
\Psi_{\varphi}^{P}(g)=\Psi_{\eta, f}^{P}=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi(\gamma \cdot g)
$$

Convergence follows from comparison to similarly-formed genuine Eisenstein series in their range of absolute convergence, in [3.11.2]. The decomposition of such pseudo-Eisenstein series in terms of the analogous genuine Eisenstein series reduces to Fourier inversion on $\mathbb{R}$ together with the functional equation (and analytic continuation) of the genuine Eisenstein series.

Without loss of generality, normalize so that $\int_{M_{k} \backslash M^{1}}|f|^{2}=1$, and $f$ is real-valued. In the self-associate case, we can assume that either $f_{1}=f_{2}$ or they are orthogonal.

## [3.16.1] Theorem:

$$
\Psi_{\eta, f}^{P}=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot E_{s, f}^{P} d s+\sum_{s_{o}}\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s, f}^{P}
$$

The residual part is non-zero only for self-associate $P$ and $f^{w}=f$, in which case there are at most finitelymany residues, all in $L^{2}$.
[3.16.2] Remark: The argument literally only proves the previous equality pointwise. In fact, a natural extension of the argument shows that the integral converges as a vector-valued integral, stemming from

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

corresponding convergence of Euclidean Fourier inversion, one instance of the latter proven in [14.3], and already exploited in [1.12] and [2.11-12] to prove Plancherel theorems for fragments of the spectrum. The current form of the issue is addressed in [3.17].
Proof: Euclidean Fourier-Mellin inversion as in [1.12] expresses $\eta \in \mathcal{D}(0, \infty)$ as

$$
\eta(y)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \mathcal{M} \eta(\sigma+i t) y^{\sigma+i t} d t \quad \text { (for any } \sigma \in \mathbb{R} \text { ) }
$$

Thus,

$$
\begin{gathered}
\eta\left(\frac{\left|\operatorname{det} m_{1}\right|^{r_{2}}}{\left|\operatorname{det} m_{2}\right|^{r_{1}}}\right) \cdot f_{1}\left(m_{1}\right) \cdot f_{2}\left(m_{2}\right)=\frac{1}{2 \pi} \int_{-\infty}^{\infty} \mathcal{M} \eta(\sigma+i t) f(m)\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{\sigma+i t} d t \\
=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \eta(s) f(m)\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{s} d s
\end{gathered}
$$

As usual, to see how a genuine Eisenstein series arises, let

$$
\varphi_{s, f}(z n m k)=\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{s} \cdot f(m)
$$

with $z \in Z_{\mathbb{A}}, n \in N_{\mathbb{A}}^{P}, m \in M_{\mathbb{A}}^{P}$, and $k \in K_{\mathbb{A}}$. Moving $\sigma$ to $\operatorname{Re}(s)>1$ for convergence of the sum, wind up to

$$
\Psi_{\eta, f}^{P}=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \eta(s) \sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{f, s}(\gamma g) d s=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \eta(s) \cdot E_{s, f}^{P}(g) d s
$$

with the genuine Eisenstein series

$$
E_{s, f}^{P}(g)=\sum_{\gamma \in P_{k} \backslash G_{k}} \varphi_{s, f}(\gamma g) \quad(\text { for } \operatorname{Re}(s)>1)
$$

Expression of $\Psi_{\eta, f}^{P}$ in terms of $\eta$ should be replaced by an intrinsic expression in terms of $\Psi_{\eta, f}^{P}$. The non-self-associate and self-associate cases are somewhat different from each other, due to the different behavior of the constant terms of Eisenstein series in those two cases. We treat the non-self-associate case first.

In the non-self-associate case, from [3.14.4], $E_{f, s}^{P}$ has no poles in $\operatorname{Re} s \geq \frac{1}{2}$, and has reasonable vertical behavior. Meanwhile, being essentially the Fourier transform of a compactly-supported smooth function, $\mathcal{M} \eta(s)$ is in the Paley-Wiener space, so is entire with rapid decay on vertical lines. Thus, we can shift the vertical integral to the line $\sigma=\frac{1}{2}$ without picking up any residues:

$$
\Psi_{\eta, f}^{P}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \eta(s) \cdot E_{s, f}^{P}(g) d s
$$

To obtain an intrinsic expression for $\mathcal{M} \eta(s)$ : unwinding the pseudo-Eisenstein series, using an Iwasawa decomposition, spherical-ness, and trivial central character, and the fact that $c_{P} E_{s, f}^{P}$ is just $\varphi_{s, f}$ in the non-self-associate case: with $\delta^{P}(m)=\left|\operatorname{det} m_{1}\right|^{r_{1}} /\left|\operatorname{det} m_{2}\right|^{r_{1}}$ the modular function of $P$,

$$
\begin{gathered}
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{f, \eta}^{P} \cdot \overline{E_{s, f}^{P}}=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} f(m) \eta\left(\frac{\left|\operatorname{det} m_{1}\right|^{r_{2}}}{\left|\operatorname{det} m_{2}\right|^{r_{1}}}\right) \cdot \overline{c_{P} E_{s, f}^{P}(m)} \frac{d m}{\delta^{P}(m)} \\
=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} f(m) \eta\left(\frac{\left|\operatorname{det} m_{1}\right|^{r_{2}}}{\left|\operatorname{det} m_{2}\right|^{r_{1}}}\right) \cdot \bar{f}(m)\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{1-s} \frac{d m}{\delta^{P}(m)}
\end{gathered}
$$

From

$$
Z^{+} M_{k} \backslash M_{\mathbb{A}} \approx Z^{+} \backslash A_{P}^{+} \times M_{k} \backslash M^{1}
$$

the pairing of pseudo-Eisenstein series against Eisenstein series becomes

$$
\begin{aligned}
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} & \Psi_{\eta, f}^{P} \cdot \overline{E_{s, f}^{P}}=\int_{Z^{+} \backslash A_{P}^{+} \times M_{k} \backslash M^{1}}|f(m)| \cdot \eta(\delta(a)) \cdot \delta(a)^{-s} d m d a \\
= & \int_{M_{k} \backslash M^{1}}|f|^{2} \cdot \int_{0}^{\infty} \eta(r) \cdot r^{s} \frac{1}{r} \frac{d r}{r}=\langle f, f\rangle \cdot \mathcal{M} \eta(s)
\end{aligned}
$$

yielding an intrinsic expression for $\mathcal{M} \eta$,

$$
\mathcal{M} \eta(s)=\frac{1}{\langle f, f\rangle} \int_{Z+G_{k} \backslash G_{A}} \Psi_{\eta, f}^{P} \cdot \overline{E_{s, f}^{P}}
$$

This computation incidentally demonstrates the absolute convergence of the integral. Unlike $G L_{2}$ and selfassociate cases, the previous computation of the integral of a pseudo-Eisenstein series against an Eisenstein series already gives an intrinsic expression for the coefficient $\mathcal{M} \eta(s)$ in the spectral decomposition, with no immediate reason to use functional equations to symmetrize the integral. With the normalization $\langle f, f\rangle=1$, the spectral decomposition is

$$
\Psi_{\eta, f}^{P}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \eta(s) \cdot E_{s, f}^{P} d s=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot E_{s, f}^{P} d s
$$

where, as usual, the pairing $\langle$,$\rangle cannot be the L^{2}$ pairing, because $E_{s, f}^{P}$ is not in $L^{2}$, but the implied integral converges absolutely, as the above unwinding argument demonstrates.

In the self-associate case, the subcase where $f$ and $f^{w}$ are orthogonal is similar to the non-self-associate case, as follows. First,

$$
\begin{aligned}
& \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{f, \eta}^{P} \cdot \overline{E_{s, f}^{P}}=\int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} f(m) \eta(\delta(m)) \cdot \overline{c_{P} E_{s, f}^{P}(m)} \frac{d m}{\delta^{P}(m)} \\
= & \int_{Z^{+} M_{k} \backslash M_{\mathbb{A}}} f(m) \eta(\delta(m)) \cdot \overline{\left(f(m) \cdot \delta(m)^{s}+c_{s, f}^{P} f^{w}(m) \delta(m)^{1-s}\right)} \frac{d m}{\delta^{P}(m)}
\end{aligned}
$$

Since $f$ and $f^{w}$ are orthogonal, the second summand in the constant term of the Eisenstein series integrates to 0 against the unwound pseudo-Eisenstein series. From [3.14.4], the Eisenstein series has no poles in $\operatorname{Re}(s) \geq \frac{1}{2}$, so we can move the contour to that line without picking up any residues. Again from

$$
Z^{+} M_{k} \backslash M_{\mathbb{A}} \approx Z^{+} \backslash A_{P}^{+} \times M_{k} \backslash M^{1}
$$

the pairing of pseudo-Eisenstein series against Eisenstein series with $\operatorname{Re}(s)=\frac{1}{2}$ becomes

$$
\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \Psi_{\eta, f}^{P} \cdot \overline{E_{s, f}^{P}}=\int_{M_{k} \backslash M_{\mathbb{A}}}|f|^{2} \cdot \int_{0}^{\infty} \eta(r) \cdot r^{1-s} \frac{1}{r} \frac{d r}{r}=\langle f, f\rangle \cdot \mathcal{M} \eta(s)
$$

yielding the same decomposition

$$
\Psi_{\eta, f}^{P}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \eta(s) \cdot E_{s, f}^{P} d s=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot E_{s, f}^{P} d s
$$

in the self-associate case with $f$ orthogonal to $f^{w}$.
In the subcase of the self-associate case where $f=f^{w}$, moving the contour from $\operatorname{Re}(s)>1$ to $\operatorname{Re}(s)=\frac{1}{2}$ may pick up finitely-many residues of $E_{s, f}^{P}$, which by [3.14.4] are in $L^{2}$. Thus,

$$
\Psi_{\eta, f}^{P}-(\text { residual part })=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \eta(s) \cdot E_{s, f}^{P} d s
$$

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

As suggested by both the $G L_{2}$ case [2.11] and the minimal-parabolic $G L_{n}$ case [3.15], average the original expression of $\Psi_{\eta, f}^{P}$ with its image under replacing $s$ by $1-s$, to symmetrize:

$$
\begin{gathered}
\Psi_{\eta, f}^{P}-(\text { residual part })=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left(\mathcal{M} \eta(s) \cdot E_{s, f}^{P}+\mathcal{M} \eta(1-s) \cdot E_{1-s, f}^{P}\right) d s \\
=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left(\mathcal{M} \eta(s)+\mathcal{M} \eta(1-s) \frac{1}{c_{s, f}}\right) \cdot E_{s, f}^{P} d s
\end{gathered}
$$

Because $f^{w}=f$, applying the functional equation twice gives $c_{s, f} \cdot c_{1-s, f}=1$, and $\left|c_{s, f}\right|=1$ on $\operatorname{Re}(s)=\frac{1}{2}$. Unwind and use Iwasawa:

$$
\begin{gathered}
\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle=\int_{Z^{+} N_{\mathbb{A}} M_{k} \backslash G_{\mathbb{A}}} \eta(\delta(m)) \cdot \overline{c_{P} E_{s, f}^{P}(m)} \frac{d m}{\delta^{P}(m)} \\
=\int_{Z^{+} \backslash A_{P}^{+} \times M_{k} \backslash M^{1}} \eta(\delta(a)) \cdot f(m) \cdot \overline{\left(f(m) \cdot \delta(a)^{s}+c_{s, f} f(m) \cdot \delta(a)^{1-s}\right)} \cdot \delta(a)^{-1} d m d a \\
=\int_{0}^{\infty} \eta(t) \cdot\left(t^{-s}+\overline{c_{s, f}} t^{s}\right) \frac{d t}{t}=\int_{0}^{\infty} \eta(t) \cdot\left(t^{-s}+\frac{1}{c_{s, f}} t^{s}\right) \frac{d t}{t}=\mathcal{M} \eta(s)+\mathcal{M} \eta(1-s) \frac{1}{c_{s, f}}
\end{gathered}
$$

so once again

$$
\Psi_{\eta, f}^{P}-(\text { residual part })=\frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot E_{s, f}^{P} d s
$$

To address the finitely-many residues of $E_{s, f}^{P}$ in the self-associate situation with $f^{w}=f$, recall that the poles $s_{o}$ of $E_{s, f}$ in $\operatorname{Re}(s)>\frac{1}{2}$ are poles of $c_{s, f}$ of the same order. Since $c_{s, f} \cdot c_{1-s, f}=1$, necessarily $c_{1-s, f}$ has a zero at $s_{o}$. Thus, from

$$
\mathcal{M} c_{P} \Psi_{\eta, f}^{P}(s)=\mathcal{M} \eta(s) f(m)+c_{1-s, f} \mathcal{M} \eta(1-s) f(m)
$$

at a pole $s_{o}$ of $E_{s}$

$$
\begin{aligned}
& \mathcal{M} c_{P} \Psi_{\eta, f}\left(s_{o}\right)=\left(\mathcal{M} \eta\left(s_{o}\right)+c_{1-s_{o}, f} \mathcal{M} \eta\left(1-s_{o}\right)\right) \cdot f(m) \\
& =\left(\mathcal{M} \eta\left(s_{o}\right)+0 \cdot \mathcal{M} \eta\left(1-s_{o}\right)\right) \cdot f(m)=\mathcal{M} \eta\left(s_{o}\right) f(m)
\end{aligned}
$$

That is, the value $\mathcal{M} c_{P} \Psi_{\eta, f}$ at $s_{o}$ is just the value of $\mathcal{M} \eta\left(s_{o}\right) f(m)$, so the coefficients appearing in the decomposition of $\Psi_{\eta, f}$ are intrinsic. Thus, the decomposition above has the form as in the statement of the theorem.

### 3.17 Plancherel for pseudo-Eisenstein series

For a fixed cuspform $f=f_{1} \otimes f_{2}$ on the Levi component $M^{P}$ of a maximal proper parabolic $P$, we show that the map from pseudo-Eisenstein series $\Psi_{\eta, f}^{P}$ attached to test functions $\eta$ on the ray $Z^{+} \backslash A_{P}^{+}$and that fixed cuspidal data $f$ to decomposition coefficients $\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle$ against genuine Eisenstein series with the same cuspidal data is an isometry to the image. This gives a Plancherel theorem for the fragment of $L^{2}$ given by (the closure of) the span of $\left\{\Psi_{\eta, f}^{P}\right.$ : test function $\left.\eta\right\}$. Similarly, we show that the map from minimalparabolic pseudo-Eisenstein series to their decomposition coefficients against minimal-parabolic Eisenstein series is an isometry to its image, giving Plancherel for this part of the spectrum. In both cases, retain the simplifying hypotheses of the last sections: trivial central character and right $K_{\mathbb{A}}$-invariance.

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Often, legitimization of natural procedures requires a linear map to commute with an integral. The best assurance of commuting is when the integral converges as a vector-valued integral, as in [14.1]. Thus, we will see that we want the spectral decomposition integrals to converge not merely pointwise, but as vector-valued integrals in spaces of functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ on which integration against cuspidal-data pseudo-Eisenstein series are continuous functionals. In earlier examples [1.12] and [2.11-12], the pseudo-Eisenstein series were compactly supported, so convergence of spectral integrals in $C^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, without growth or support constraints, was sufficient, and this was inherited from the analogous argument [14.3] for Fourier series on $\mathbb{R}$. The present discussion requires somewhat more, since cuspidal-data pseudo-Eisenstein series do not have compact support, but are of rapid decay for strong-sense cuspidal data, as in [3.9] and [3.16].

Fix a strong-sense cuspform $f=f_{1} \otimes f_{2}$ on $M^{P}$ with $P$ maximal proper. The Plancherel theorem for associated pseudo-Eisenstein series comes from
[3.17.1] Theorem: With $|f|^{2}=1$, letting $s_{o}$ run over poles of $E_{s, f}^{P}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$, with $\eta, \theta$ test functions on $\left.Z^{+} \backslash A_{P}^{+}\right) \approx(0, \infty)$,

$$
\left\langle\Psi_{\eta, f}^{P}, \Psi_{\theta, f}^{P}\right\rangle=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot \overline{\left\langle\Psi_{\theta, f}^{P}, E_{s, f}^{P}\right\rangle} d s+\sum_{s_{o}}\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle \cdot \overline{\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle}
$$

[3.17.2] Remark: There are no residues in $\operatorname{Re}(s) \geq \frac{1}{2}$ unless $P$ is self-associate and $\left\langle f, f^{w}\right\rangle \neq 0$, by [3.14.4]. Proof: Grant that the decomposition integral for $\Psi_{\eta, f}^{P}$ converges well enough to pass integration against $\Psi_{\theta, f}^{P}$ inside it: a vector-valued convergence certainly is sufficient, and more pedestrian arguments are also possible. Then

$$
\begin{aligned}
& \left\langle\Psi_{\eta, f}^{P}, \Psi_{\theta, f}^{P}\right\rangle=\left\langle\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot E_{f, s}^{P} d s+\sum_{s_{o}}\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle \cdot \operatorname{Res}_{s_{o}} E_{s, f}^{P}, \Psi_{\theta, f}^{P}\right\rangle \\
& \quad=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Psi_{\eta, f}^{P}, E_{s, f}^{P}\right\rangle \cdot \overline{\left\langle\Psi_{\theta, f}^{P}, E_{s, f}^{P}\right\rangle} d s+\sum_{s_{o}}\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle \cdot \overline{\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle}
\end{aligned}
$$

as claimed.
This spectral decomposition facilitates demonstration of the orthogonality of pseudo-Eisenstein series remarked upon in [3.9.1], beyond the non-self-associate case from [3.11.5]:
[3.17.3] Corollary: Pseudo-Eisenstein series $\Psi_{\eta, f}^{P}$ and $\Psi_{\theta, f^{\prime}}^{Q}$ for maximal proper parabolics $P, Q$, with test functions $\eta, \theta$ and cuspidal data $f, f^{\prime}$, are mostly mutually orthogonal: they are orthogonal if $P, Q$ are not associate, or if $P=Q$ but $\left\langle f, f^{\prime}\right\rangle=0$, or if $M^{P}=w M^{Q} w^{-1}$ but $\left\langle f^{w}, f^{\prime}\right\rangle=0$.
Proof: In all cases, the first part of the proof of the decomposition of pseudo-Eisenstein series in terms of genuine Eisenstein series yields a vector-valued integral

$$
\Psi_{\eta, f}^{P}=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \eta(s) \cdot E_{s, f}^{P} d s+\sum_{s_{o}} \mathcal{M} \eta\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}} E_{s, f}^{P}
$$

converging in the $C^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ topology, from the discussions just above. By [14.1], the inner product functional against $\Psi_{\theta, f^{\prime}}^{Q}$ can pass through the integral, and through the operation of residue as well [15.2]. Thus,

$$
\left\langle\Psi_{\eta, f}^{P}, \Psi_{\theta, f^{\prime}}^{Q}\right\rangle=\frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty} \mathcal{M} \varphi(s) \cdot\left\langle E_{s, f}^{P}, \Psi_{\theta, f^{\prime}}^{Q}\right\rangle d s+\sum_{s_{o}} \mathcal{M} \varphi\left(s_{o}\right) \cdot \operatorname{Res}_{s_{o}}\left\langle E_{s, f}^{P}, \Psi_{\theta, f^{\prime}}^{Q}\right\rangle
$$

The integrals of $\Psi_{\theta, f^{\prime}}^{Q}$ against genuine Eisenstein series unwind to integrals of $\varphi_{\theta, f^{\prime}}^{Q}$ against the $Q$-constant terms of the Eisenstein series, computed in [3.11.9]. If $Q$ is not associate to $P$, this is 0 . If $Q=P$, or if $P$ is not self-associate and $Q$ is the other associate to $P$, the integral includes an inner integral of $f$ against $f^{\prime}$ or $f^{w}$ against $f^{\prime}$, as in the proof of Maaß-Selberg relations [3.14.2], and these are 0 by assumption.

These formulas suggest the form of a Plancherel theorem for this fragment of $L^{2}$. For $P$ self-associate and $f^{w}=f$, let $s_{1}, \ldots, s_{n}$ be the poles of $E_{s, f}^{P}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$, and let

$$
\mathcal{F}_{f}^{P}: \Psi_{\eta, f}^{P} \longrightarrow\left(t \rightarrow\left\langle\Psi_{\eta, f}^{P}, E_{\frac{1}{2}+i t, f}^{P}\right\rangle\right) \oplus\left(\ldots,\left\langle\Psi_{\eta, f}^{P}, \operatorname{Res}_{s_{o}} E_{s, f}^{P}\right\rangle, \ldots\right) \in L^{2}\left(\frac{1}{2}+i \mathbb{R}\right) \oplus \mathbb{C}^{n}
$$

be the spectral coefficient map. It is necessary to identify the image of $\mathcal{F}_{f}^{P}$. Let

$$
V=\left\{\begin{array}{cc}
\left\{F \in L^{2}\left(\frac{1}{2}+i \mathbb{R}\right): F(1-s)=c_{s, f} F(s)\right\} \oplus \mathbb{C}^{n} & \quad\left(\text { for } P \text { self-associate and } f^{w}=f\right) \\
L^{2}\left(\frac{1}{2}+i \mathbb{R}\right) & \left(\text { for } P \text { not self-associate or }\left\langle f, f^{w}\right\rangle=0\right)
\end{array}\right.
$$

[3.17.4] Claim: The spectral coefficient $\operatorname{map} \mathcal{F}_{f}^{P}$ is an isometry to its image in $V$, and that image is dense in $V$.
Proof: The fact that it is an isometry to its image is [3.17.1]. The map $\eta \rightarrow \mathcal{M} \eta$ is essentially Fourier transform, so sends the dense subset of test functions inside the Schwartz space to the dense subset of PaleyWiener functions inside the Schwartz space. The Schwartz space is dense in $L^{2}$, so the functions $\mathcal{M} \eta$ are dense in $L^{2}\left(\frac{1}{2}+i \mathbb{R}\right)$.

For $P$ not self-associate, this is all we need, since there are no residues, and no folding-up of the spectral integral, since the functional equation of the Eisenstein series does not relate it to itself. The case of $P$ self-associate but $\left\langle f, f^{w}\right\rangle=0$ is similar.

Now consider $P$ self-associate and $f^{w}=f$. The residues of $E_{s, f}$ in $\operatorname{Re}(s) \geq \frac{1}{2}$ are orthogonal to cuspforms, by [3.14.6]. These residues have $Q$-constant terms 0 unless $Q$ is associate to $P$, by [3.11.3], since taking residues commutes with evaluation of constant terms, from generalities about vector-valued integrals [14.1] and holomorphic or meromorphic vector-valued functions [15.2]. With $Q$ associate to $P$, for cuspidal data $f^{\prime}$ with $\left\langle f^{\prime}, f\right\rangle=0=\left\langle f^{\prime}, f^{w}\right\rangle$, we have orthogonality by [3.17.3]. Thus, by a process of elimination, these residues must in the closure of the space of pseudo-cuspforms $\Psi_{\eta, f}^{P}$ with test-function data $\eta$.

The decomposition integral gets folded up via the functional equation of $E_{s, f}$ to obtain coefficients $\left\langle\Psi_{\eta, f}, E_{s, f}\right\rangle=\mathcal{M} \eta(s)+c_{s, f}^{-1} \mathcal{M} \eta(1-s)$. The map $F \rightarrow F(s)+c_{1-s, f} F(1-s)$ is a continuous map of $L^{2}$ to the subspace in $V$, and is continuous because $\left|c_{s, f}\right|$ is constant on $\operatorname{Re}(s)=\frac{1}{2}$. Since the residues are in the closure, the integral part of the spectral decomposition is in the closure. This proves that the spectral map has dense image.
[3.17.5] Remark: Continuing in this vein, the $L^{2}$ closure of the image of $\Psi_{\eta, f}^{P}$ for fixed $P$, for all cuspforms $f$, and for all test functions $\eta$, is the collection of functions orthogonal to cuspforms on $G$ and with all constant terms vanishing except $c_{P}$ and $c_{Q}$.

In the opposite case, with minimal parabolic $P=P^{\text {min }}$, in [3.10.1] the Eisenstein series $E_{s}^{P}$ was shown convergent for $s \in 2 \rho+C$, with cone

$$
C=\left\{s \in \mathfrak{q} \otimes_{\mathbb{R}} \mathbb{C}:\langle\alpha, \operatorname{Re}(s)\rangle>0, \text { for all simple positive roots } \alpha\right\}
$$

The argument for the spectral decomposition of pseudo-Eisenstein series for $P=P^{\text {min }}$ in [3.15] showed that only (multi-) residues of $E_{s}^{P}$ for in $s \in \rho+C$ are relevant to the spectral decomposition. Let $\mathfrak{q}$ be the Lie algebra of $Z^{+} \backslash A_{P}^{+}$, as in [3.15].
[3.17.6] Theorem: Let $\eta, \theta$ be test functions on $Z^{+} \backslash A_{P}^{+} \approx(0, \infty)^{r-1}$ such that $\Psi_{\varphi}^{P}$ is orthogonal to all residues of $E_{\rho+s}^{P}$ with $s \in C$. Then

$$
\left\langle\Psi_{\eta}^{P}, \Psi_{\theta}^{P}\right\rangle=\frac{1}{r!} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{q}^{*}}\left\langle\Psi_{\eta}, E_{\rho+s}\right\rangle \cdot \overline{\left\langle\Psi_{\theta}, E_{\rho+s}\right\rangle} d s
$$

Proof: In this example, since the pseudo-Eisenstein series have only test-function data, they are compactlysupported on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$, by [3.8.1]. Thus, integration against $\Psi_{\theta}^{P}$ in the following is justified by observing that the spectral expansion of $\Psi_{\eta}^{P}$ converges in the $C^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ topology. The latter follows from the corresponding assertion for Fourier transform on $\mathbb{R}$ and $\mathbb{R}^{n}$, proven in [14.3], simply augmenting the argument

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for the spectral decomposition [3.15] to retain that aspect, rather than mere pointwise convergence. More pedestrian arguments are possible, as usual. Granting that, compute directly:

$$
\left\langle\Psi_{\eta}^{P}, \Psi_{\theta}^{P}\right\rangle=\left\langle\frac{1}{r!} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{q}^{*}}\left\langle\Psi_{\eta}^{P}, E_{\rho+s}^{P}\right\rangle \cdot E_{\rho+s}^{P} d s, \quad \Psi_{\theta}^{P}\right\rangle=\frac{1}{r!} \frac{1}{(2 \pi i)^{r-1}} \int_{i \mathfrak{q}^{*}}\left\langle\Psi_{\varphi}^{P}, E_{\rho+s}^{P}\right\rangle \cdot \overline{\left\langle\Psi_{\theta}^{P}, E_{\rho+s}^{P}\right\rangle} d s
$$

as claimed.
Toward Plancherel: Any $F \in L^{2}\left(i \mathfrak{q}^{*}\right)$ satisfying $F(w \cdot s)=c_{w, s} \cdot F(s)$ for all $w \in W$ is in the closure of the image of the map $\Psi_{\eta}^{P} \rightarrow\left\langle\Psi_{\eta}^{P}, E_{\rho+s}\right\rangle$ ranging over all test functions $\eta$ on $i q^{*}$. Indeed, the map $\eta \rightarrow \mathcal{M} \eta$ is essentially Fourier transform, and as in other examples maps test functions to the Paley-Wiener space, dense in $L^{2}$. The averaging map

$$
F \longrightarrow \sum_{w \in W} \frac{1}{c_{w, s}} F(w \cdot s)
$$

surjects $L^{2}\left(i \mathfrak{q}^{*}\right)$ to its subspace where $F(w \cdot s)=c_{w, s} \cdot F(s)$ for all $w \in W$, since $\left|c_{w, s}\right|=1$ on $i \mathfrak{q}^{*}$, by [3.12.6]. However, we have not identified the (multi-) residues of $E_{s}$ that appear when moving contours, and cannot immediately distinguish the subspace of pseudo-Eisenstein series orthogonal to these residues.

Further, we would need to be able to argue that these multi-residues are entirely inside the closure of the images of the pseudo-Eisenstein series. For the latter, it seems necessary to invoke the complete spectral decomposition of $L^{2}(\Gamma \backslash G / K)$, that cuspforms and cuspidal data Eisenstein series attached to non-minimal parabolics, and their $L^{2}$ residues, as well as the minimal-parabolic pseudo-Eisenstein series, span $L^{2}(\Gamma \backslash G / K)$. Only then is the orthogonality of integrals of minimal-parabolic Eisenstein series to all the other spectral components clear.

Thus, while we did prove that the map from the space of pseudo-Eisenstein series to integrals of Eisenstein series is an isometry to its image, we did not quite identify that image.

### 3.18 Automorphic spectral expansions

We would like to express $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ as the closure of subspaces consisting of eigenfunctions for invariant differential operators and for spherical Hecke operators. Analogous decomposition of $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ needs more general integral operators.

First, we can decompose by central characters into pieces $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}, \omega\right)$, for general reasons [3.6].
Then the general pattern is that there are cuspforms, and we are left to sift through their orthogonal complement. That orthogonal complement is spanned by pseudo-Eisenstein series with cuspidal data attached to the various parabolics. The cuspidal-data pseudo-Eisenstein series themselves are not eigenfunctions for invariant differential operators or Hecke operators, but are essentially integrals of genuine Eisenstein series, and the latter are eigenfunctions [3.11.6] and [3.11.11]. The functional equations of genuine Eisenstein series show that parabolics $P, Q$ that are associate, in the sense that their Levi components are conjugate, produce the same functions on the group. Thus, a rough indexing of parts of $L^{2}$ is by associateclass of parabolics.

In general, residues of cuspidal-data Eisenstein series also enter the expression of pseudo-Eisenstein series. The relevant residues are square-integrable, and inherit eigenfunction properties from the genuine Eisenstein series. For $G L_{2}$, these residues are relatively uninteresting, as in [2.B]. For $G L_{3}$, the $P^{2,1}$ and $P^{1,2}$ parabolics' Eisenstein series have no relevant residues [3.14.4], so the only residues are those from $P^{1,1,1}=P^{\text {min }}$, which turn out to be constants. Granting the latter, fact, for example, with trivial central character, over groundfield $\mathbb{Q}$, since there are no unramified Hecke characters, functions $\Phi$ in $L^{2}\left(Z_{\mathbb{A}} G L_{3}(\mathbb{Q}) \backslash G L_{3}(\mathbb{A}) / G L_{3}(\mathbb{A})\right)$ have $L^{2}$ decompositions
$\Phi=\sum_{G L_{3} \operatorname{cfm} F}\langle\Phi, F\rangle \cdot F+\sum_{G L_{2} \operatorname{cfm} f} \frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Phi, E_{s, f}^{2,1}\right\rangle \cdot E_{s, f}^{2,1} d s+\frac{1}{3!\cdot 2 \pi i} \int_{i \mathbf{q}^{*}}\left\langle\Phi, E_{\rho+s}^{\min }\right\rangle \cdot E_{\rho+s}^{\min } d s+\frac{\langle\Phi, 1\rangle \cdot 1}{\langle 1,1\rangle}$
where the first sum is over an orthonormal basis of spherical cuspforms for $G L_{3}(\mathbb{Z})$ with trivial central character, and the second sum is over an orthonormal basis for spherical cuspforms for $G L_{2}(\mathbb{Z})$ with trivial

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

central character. The right-hand side is only promised to converge in an $L^{2}$ sense, and the explicit and implicit integrals involving Eisenstein series are merely isometric extensions of the corresponding literal integrals.

Non-trivial residual spectrum for $G L_{4}$ : This is the smallest $G L_{n}$ in which some Eisenstein series $E_{s, f}^{2,2}$ with real-valued $f=f^{w}$ have non-constant residues. ${ }^{[34]}$ The Maaß-Selberg relations do not exclude the possibility of a pole of such an Eisenstein series. A computation of the constant term of that Eisenstein series shows that it is a ratio of values of the Rankin-Selberg $L$-function attached to $f \times \bar{f}$, which definitely has a pole in $\operatorname{Re}(s)>\frac{1}{2}$, yielding a square-integrable residue. Granting this, for example, over $\mathbb{Q}$, spherical $\Phi$ in $L^{2}\left(Z_{\mathbb{A}} G L_{4}(\mathbb{Q}) \backslash G L_{4}(\mathbb{A})\right.$, that is, in $L^{2}\left(Z_{\mathbb{A}} G L_{4}(\mathbb{Q}) \backslash G L_{4}(\mathbb{A}) / K_{\mathbb{A}}\right)$, have $L^{2}$ decompositions described as follows. For a modicum of coherence, let $\Xi_{n}$ be a fixed orthonormal basis for spherical cuspforms with trivial central character for $G L_{n}(\mathbb{Z})$, consisting of spherical Hecke eigenfunctions, etc. We grant that there is a unique relevant residue $F_{f}$ of $E_{s, f \otimes \bar{f}}^{2,2}$ for cuspforms $f$ on $G L_{2}$. Then

$$
\begin{aligned}
& \Phi=\sum_{f \in \Xi_{4}}\langle\Phi, F\rangle \cdot F+\sum_{f \in \Xi_{3}} \frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Phi, E_{s, f}^{3,1}\right\rangle \cdot E_{s, f}^{3,1} d s+\sum_{f_{1}, f_{2} \in \Xi_{2}, f_{1} \neq \bar{f}_{2}} \frac{1}{2 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Phi, E_{s, f_{1} \otimes f_{2}}^{2,2}\right\rangle \cdot E_{s, f_{1} \otimes f_{2}}^{2,2} d s \\
&+\sum_{f \in \Xi_{2}} \frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Phi, E_{s, f \otimes \bar{f}}^{2,2}\right\rangle \cdot E_{s, f \otimes \bar{f}}^{2,2} d s+\sum_{f \in \Xi_{2}}\left\langle\Phi, F_{f}\right\rangle \cdot F_{f} \\
&+\sum_{f \in \Xi_{2}} \frac{1}{4 \pi i} \int_{\frac{1}{2}-i \infty}^{\frac{1}{2}+i \infty}\left\langle\Phi, E_{s, f}^{2,1,1}\right\rangle \cdot E_{s, f}^{2,1,1} d s+\frac{1}{4!\cdot 2 \pi i} \int_{i \mathfrak{q}^{*}}\left\langle\Phi, E_{\rho+s}^{1,1,1,1}\right\rangle \cdot E_{\rho+s}^{1,1,1,1} d s+\frac{\langle\Phi, 1\rangle \cdot 1}{\langle 1,1\rangle}
\end{aligned}
$$

Again, the explicit and implicit integrals involving Eisenstein series are in fact isometric extensions of the literal integrals.

## 3.A Appendix: Bochner's Lemma

Bochner's Lemma is a one-of-a-kind device for meromorphic continuation in two or more complex variables. [35] Let $\Omega_{o}$ be a non-empty, connected, open set in $\mathbb{R}^{n}$ with $n>1$. The tube domain $\Omega$ over $\Omega_{o}$ is $\Omega=\Omega_{o}+i \mathbb{R}^{n}$, that is, the collection of $z \in \mathbb{C}^{n}$ with real part in $\Omega_{o}$. Let $f$ be a holomorphic $\mathbb{C}$-valued function on $\Omega$, of not-too-awful vertical growth, in the sense that, for $x$ in fixed compact $C \subset \Omega_{o}$, there is $1 \leq N \in \mathbb{Z}$ such that

$$
|f(x+i y)|<_{C} e^{|y|^{N}} \quad\left(\text { with }\left|\left(y_{1}, \ldots, y_{n}\right)\right|^{2}=y_{1}^{2}+\ldots+y_{n}^{2}\right)
$$

[3.A.1] Claim: $f$ extends to a holomorphic function on the convex hull of $\Omega$.
Proof: First, let $x, \xi$ be two points in $\Omega_{o}$, such that the line segment connecting them lies entirely within $\Omega_{o}$. We will specify a rectangle inside $\Omega$ with $x, \xi$ the midpoints of opposite sides. Let $\gamma=\gamma_{x, \xi, R}$ parametrize

[^9]the rectangle with sides individually parametrized by
\[

\left\{$$
\begin{array}{lll}
\text { side through } x: & x+i t(x-\xi) & (\text { with }-R \leq t \leq R) \\
\text { top: } & (1-t)(x+i R(x-\xi))+t(\xi+i R(x-\xi)) & (\text { with } 0 \leq t \leq 1) \\
\text { side through } \xi: & \xi-i t(x-\xi) & (\text { with }-R \leq t \leq R) \\
\text { bottom: } & (1-t)(\xi-i R(\xi))+t(x-i R(x-\xi)) & (\text { with } 0 \leq t \leq 1)
\end{array}
$$\right.
\]

The expressions for the top and bottom simplify to

$$
\left\{\begin{array}{lll}
\text { top: } & (1-t) x+t \xi+i R(x-\xi) & (\text { with } 0 \leq t \leq 1) \\
\text { bottom: } & (1-t) \xi+t x-i R(x-\xi) & (\text { with } 0 \leq t \leq 1)
\end{array}\right.
$$

This rectangle lies inside $Z=x+\mathbb{C} \cdot(x-\xi) \approx \mathbb{C}$, and is contractible in $\Omega$. Let $j(\zeta)=x+\zeta \cdot(x-\xi)$. In $Z$, Cauchy's formula in one variable is

$$
f \circ j\left(\zeta_{o}\right)=\frac{1}{2 \pi i} \int_{\gamma} \frac{f \circ j(\zeta) d \zeta}{\zeta-\zeta_{o}}
$$

To legitimately push the top and bottom of the rectangle to infinity, use the growth assumption on $f$, and the modified integral expression

$$
f \circ j\left(\zeta_{o}\right)=e^{-\zeta_{o}^{2 N}} \frac{1}{2 \pi i} \int_{\gamma} \frac{e^{\zeta^{2 N}} \cdot f \circ j(\zeta) d \zeta}{\zeta-\zeta_{o}}
$$

Thus, taking the limit $R \rightarrow+\infty$,
$e^{\zeta_{o}^{2 N}} \cdot f\left(\zeta_{o} \cdot(x-\xi)\right)=\frac{1}{2 \pi} \int_{-\infty}^{+\infty} \frac{e^{(x+i t(x-\xi))^{2 N}} f(x+i t(x-\xi)) d t}{i t-\zeta_{o}}+\frac{1}{2 \pi} \int_{-\infty}^{+\infty} \frac{e^{(\xi-i t(x-\xi))^{2 N}} f(\xi-i t(x-\xi)) d t}{-1-i t-\zeta_{o}}$
The right-hand side makes sense for any $x, \xi \in \Omega_{o}$, whether or not the line segment connecting them lies in $\Omega_{o}$. Further, the right-hand side is holomorphic in $x, \xi \in \Omega$. Thus, the left-hand side is holomorphic, and gives the extension to the convex hull of $\Omega$.

## 3.B Appendix: Phragmén-Lindelöf theorem

This is from [Lindelöf 1908] and [Phragmén-Lindelöf 1908].
The maximum modulus principle can easily be misapplied on unbounded open sets. That is, while for an open set $U \subset \mathbb{C}$ with bounded closure $\bar{U}$, it does follow that the sup of a holomorphic function $f$ on $U$ extending continuously to $\bar{U}$ occurs on the boundary $\partial U$ of $U$, holomorphic functions on an unbounded set can be bounded by 1 on the edges but be violently unbounded in the interior.

The usual simple example is $f(z)=e^{e^{z}}$ :

$$
\left|e^{e^{x+i y}}\right|=e^{\operatorname{Re}\left(e^{x+i y}\right)}=e^{e^{x} \cdot \cos y}
$$

On one hand, for fixed $y=\operatorname{Im} z$ with $\cos y>0$, the function blows up as $x=\operatorname{Re} z \rightarrow+\infty$. On the other hand, for $\cos y=0$ the function is bounded. Thus, on the strip $-\frac{\pi}{2} \leq y \leq \frac{\pi}{2}$, the function $e^{e^{z}}$ is bounded on the edges but blows up as $x \rightarrow+\infty$.

This example suggests growth conditions under which a bound of 1 on the edges implies the same bound throughout the strip. In fact, the suggested bound is essentially sharp, in light of the example. For a half-strip, the theorem is

$$
\text { 3. } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

[3.B.1] Theorem: For $f$ a holomorphic function on the horizontal half-strip

$$
\left\{z:-\frac{\pi}{2} \leq y \leq \frac{\pi}{2} \text { and } 0 \leq x\right\}
$$

satisfying

$$
|f(z)| \ll e^{e^{C \cdot \operatorname{Re} z}} \quad(\text { for some constant } 0 \leq C<1)
$$

$|f(z)| \leq 1$ on the edges of the half-strip implies $|f(z)| \leq 1$ in the interior, as well.
Proof: Unsurprisingly, the proof is a reduction to the usual maximum modulus principle. Take any fixed $D$ in the range

$$
C<D<1
$$

The function

$$
F_{\varepsilon}(z)=f(z) / e^{\varepsilon e^{D \cdot z}} \quad(\text { for } \varepsilon>0)
$$

is bounded by 1 on the edges of the half-strip, and in the interior goes to 0 uniformly in $y$ as $x \rightarrow+\infty$, for fixed $\varepsilon>0$, exploiting the modification with $D$. Thus, on a rectangle

$$
R_{T}=\left\{z:-\frac{\pi}{2} \leq y \leq \frac{\pi}{2}, \text { and } 0 \leq x \leq T\right\}
$$

for sufficiently large $T>0$ depending upon $\varepsilon$, the function $F_{\varepsilon}$ is bounded by 1 on the edge. The usual maximum modulus principle implies that $F_{\varepsilon}$ is bounded by 1 throughout. That is, for each fixed $z_{o}$ in the half-strip,

$$
\left|f\left(z_{o}\right)\right| \leq e^{\varepsilon \cdot e^{D \operatorname{Re} z_{o}}} \quad(\text { for all } \varepsilon>0)
$$

Let $\varepsilon \rightarrow 0^{+}$, giving $\left|f\left(z_{o}\right)\right| \leq 1$.
[3.B.2] Remark: Analogous theorems on strips of other widths follow by using $e^{c \cdot e^{z}}$ with suitable constants c.

The theorem on a full strip follows by using $e^{\cosh z}$ in place of $e^{e^{z}}$, as follows.
[3.B.3] Theorem: For $f$ a holomorphic function on the horizontal strip

$$
\left\{z:-\frac{\pi}{2} \leq \operatorname{Im} z \leq \frac{\pi}{2}\right\}
$$

satisfying

$$
|f(z)| \ll e^{\cosh C \cdot \operatorname{Re} z} \quad(\text { for some constant } 0 \leq C<1)
$$

$|f(z)| \leq 1$ on the edges of the strip implies $|f(z)| \leq 1$ in the interior, as well.
Proof: Again, reduce to the maximum modulus principle. Fix $D$ in the range $C<D<1$. The function

$$
F_{\varepsilon}(z)=f(z) / e^{\varepsilon \cosh D z} \quad(\text { for } \varepsilon>0)
$$

is bounded by 1 on the edges of the strip, and in the interior goes to 0 uniformly in $y$ as $x \rightarrow \pm \infty$, for fixed $\varepsilon>0$. Thus, on a rectangle

$$
R_{T}=\left\{z:-\frac{\pi}{2} \leq y \leq \frac{\pi}{2}, \text { and }-T \leq x \leq T\right\} \quad(\text { for large } T>0, \text { depending upon } \varepsilon)
$$

the function $F_{\varepsilon}$ is bounded by 1 on the edge. The usual maximum modulus principle implies that $F_{\varepsilon}$ is bounded by 1 throughout. That is, for each fixed $z_{o}$ in the half-strip,

$$
\left|f\left(z_{o}\right)\right| \leq e^{\varepsilon \cosh D \operatorname{Re} z_{o}} \quad(\text { for all } \varepsilon>0)
$$

We can let $\varepsilon \rightarrow 0^{+}$, giving $\left|f\left(z_{o}\right)\right| \leq 1$.

## 4. Invariant differential operators

1. Derivatives of group actions: Lie algebras
2. Laplacians and Casimir operators
3. Details about universal algebras
4. Descending to $G / K$
5. Example computation: $S L_{2}(\mathbb{R})$ and $\mathfrak{H}$
6. Example computation: $S L_{2}(\mathbb{C})$
7. Example computation: $S p_{1,1}^{*}$
8. Example computation: $S L_{2}(\mathbb{H})$

Appendix A: brackets
Appendix B: existence and uniqueness
We want a method to determine natural Laplacian-like differential operators invariant under group actions in coordinate-free terms, and also exhibit the operators in convenient coordinate systems. That is, we do not want to specify operators in coordinate systems and check invariance, but, rather, know invariance a priori. The first example is the well-known operator

$$
\left.\Delta=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) \quad \text { (in coordinate(s) } z=x+i y \text { on } \mathfrak{H}\right)
$$

Although, once exhibited, this operator is certifiably invariant under the linear fractional action of $S L_{2}(\mathbb{R})$, it is oppressive and unenlightening to do this checking. Worse, it is misguided to think in terms of such verification. The relevant issue is the coordinate-independent origin of the operator, expressed subsequently in coordinates. No prior acquaintance with Lie groups or Lie algebras is assumed.

### 4.1 Derivatives of group actions: Lie algebras

Let $G$ be a subgroup of $G L_{n}(\mathbb{R})$ or $G L_{n}(\mathbb{C})$ or $G L_{n}(\mathbb{H})$ acting differentiably ${ }^{[36]}$ on the right on a smooth manifold ${ }^{[37]}$ thereby acting on functions $f$ on $M$ by

$$
(g \cdot f)(m)=f(m g)
$$

Our operational definition of the (real) Lie algebra $\mathfrak{g}$ of $G$ is

$$
\mathfrak{g}=\left\{\text { real } n \text {-by- } n \text { real matrices } x: e^{t x} \in G \text { for all real } t\right\}
$$

where the matrix exponential is

$$
\exp (x)=e^{x}=1+x+\frac{x^{2}}{2!}+\frac{x^{3}}{3!}+\ldots
$$

This definition makes clear that $\mathfrak{g}$ is closed under scalar multiplication, but not that it is closed under addition. When $x$ and $y$ are $n$-by- $n$ real or complex matrices with $x y=y x$, then $e^{x+y}=e^{x} \cdot e^{y}$, but this does not hold more generally, so we cannot easily conclude closed-ness under addition. We will prove closure under addition as a side effect of proof in [4.A] that such a Lie algebra is closed under Lie brackets

$$
x \times y \longrightarrow[x, y]=x y-y x \quad \text { (for } x, y \in \mathfrak{g}, \text { with } x \times y \rightarrow x y \text { matrix multiplication) }
$$

[^10]
## 4. Invariant differential operators

In any particular example, the vector space property is readily verified, as below. However, this binary operation $x \times y \rightarrow[x, y]$ is not similar to more elementary ring or algebra multiplications, as it is not associative, is anti-commutative, and $[x, x]=0$.

For each $x \in \mathfrak{g}$ we have a differentiation $X_{x}$ of functions $f$ on $M$ in the direction $x$, by

$$
\left(X_{x} f\right)(m)=\left.\frac{d}{d t}\right|_{t=0} f\left(m \cdot e^{t x}\right)
$$

This applies uniformly to any space $M$ on which $G$ acts (differentiably). [38] The differential operators $X_{x}$ for $x \in \mathfrak{g}$ do not typically commute with the action of $g \in G$, although the relation between the two is reasonable:

$$
\left(g \circ X_{x} \circ g^{-1}\right) f(h)=\left.\frac{d}{d t}\right|_{t=0} f\left(h \cdot g e^{t x} g^{-1}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(h \cdot e^{t \cdot g x g^{-1}}\right)=X_{g x g^{-1}} f(h)
$$

[4.1.1] Example: The condition $e^{t x} \in S L_{n}(\mathbb{R})$ for all real $t$ is that $\operatorname{det}\left(e^{t x}\right)=1$. To see what this requires of $x$, observe that for $n$-by- $n$ (real or complex) matrices $x$

$$
\operatorname{det}\left(e^{x}\right)=e^{\operatorname{tr} x} \quad(\text { where } \operatorname{tr} \text { is trace })
$$

To see why, both determinant and trace are invariant under conjugation $x \rightarrow g x g^{-1}$, so without loss of generality $x$ is upper-triangular. Then $e^{x}$ is upper-triangular, with diagonal entries $e^{x_{i i}}$, with diagonal entries $x_{i i}$ of $x$. Thus,

$$
\operatorname{det}\left(e^{x}\right)=e^{x_{11}} \cdots e^{x_{n n}}=e^{x_{11}+\ldots+x_{n n}}=e^{\operatorname{tr} x}
$$

Using this, the determinant-one condition is

$$
1=\operatorname{det}\left(e^{t x}\right)=e^{t \cdot \operatorname{tr} x}=1+t \cdot \operatorname{tr} x+\frac{(t \cdot \operatorname{tr} x)^{2}}{2!}+\ldots
$$

Taking the derivative with respect to $t$ and setting $t=0$ gives $0=\operatorname{tr} x$. Looking at the right-hand side of the expanded $1=\operatorname{det}\left(e^{t x}\right)$, this condition is also sufficient for $\operatorname{det}\left(e^{t x}\right)=1$. Thus,

$$
\text { Lie algebra } \mathfrak{s l}_{n}(\mathbb{R}) \text { of } S L_{n}(\mathbb{R})=\{n \text {-by- } n \text { real } x: \operatorname{tr} x=0\}
$$

[4.1.2] Example: Similarly,

$$
\text { Lie algebra } \mathfrak{s l}_{n}(\mathbb{C}) \text { of } S L_{n}(\mathbb{C})=\{n \text {-by- } n \text { complex } x: \operatorname{tr} x=0\}
$$

[4.1.3] Example: From $\operatorname{det}\left(e^{x}\right)=e^{\operatorname{tr} x}$, any matrix $e^{x}$ is invertible, so

$$
\text { Lie algebra } \mathfrak{g l}_{n}(\mathbb{R}) \text { of } G L_{n}(\mathbb{R})=\{\text { all real } n \text {-by- } n \text { matrices }\}
$$

[4.1.4] Example: For the simplest real orthogonal group $G=O(n, \mathbb{R})=\left\{g \in G L_{n}(\mathbb{R}): g^{\top} \cdot g=1_{n}\right\}$, using $\left(e^{t x}\right)^{\top}=e^{t x^{\top}}$,

$$
1=\left(e^{t x}\right)^{\top} \cdot e^{t x}=\left(1+t x^{\top}+\ldots\right) \cdot(1+t x+\ldots)=1+t\left(x+x^{\top}\right)+\ldots
$$

Thus, necessarily $x^{\top}+x=0$. On the other hand, when $x^{\top}+x=0$ we have $x^{\top}=-x$, so

$$
\left(e^{t x}\right)^{\top} \cdot e^{t x}=e^{-t x} \cdot e^{t x}=\left(e^{t x}\right)^{-1} \cdot e^{t x}=1
$$

[^11]
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This shows that the Lie algebra of $O(n, \mathbb{R})$ is skew-symmetric matrices.
[4.1.5] Example: Exponentiation of matrices with quaternion entries is similar:
Lie algebra $\mathfrak{g l}_{n}(\mathbb{H})$ of $G L_{n}(\mathbb{H})=\{$ all quaternion $n$-by- $n$ matrices $\}$
Slightly more subtly,

$$
\text { Lie algebra } \mathfrak{s l}_{n}(\mathbb{H}) \text { of } S L_{n}(\mathbb{H})=\left\{n \text {-by- } n \text { quaternionic } x: \sum_{i} \operatorname{tr} x_{i i}=0\right\} \quad \text { (quaternion trace) }
$$

[4.1.6] Example: For $G=S p_{1,1}^{*} \subset G L_{2}(\mathbb{H})$, let $S=\left(\begin{array}{cc}0 & 1 \\ 1 & 0\end{array}\right)$ and let $\sigma$ be quaternion-conjugate-transpose. The defining condition $S=\left(e^{t x}\right)^{\sigma} S\left(e^{t x}\right)$ is $S=e^{t x^{\sigma}} S e^{t x}$. Differentiating with respect to $t \in \mathbb{R}$ gives $0=x^{\sigma} e^{t x^{\sigma}} S e^{t x}+e^{t x^{\sigma}} S x e^{t x}$. Setting $t=0$ gives $x^{\sigma} S+S x=0$, which is $S x^{\sigma} S^{-1}=-x$. Multiplying out, this condition is

$$
\left(\begin{array}{ll}
-a & -b \\
-c & -d
\end{array}\right)=S\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right)^{\sigma} S^{-1}=S\left(\begin{array}{cc}
\bar{a} & \bar{c} \\
\bar{b} & \bar{d}
\end{array}\right) S^{-1}=\left(\begin{array}{cc}
\bar{d} & \bar{b} \\
\bar{c} & \bar{a}
\end{array}\right)
$$

The sufficiency of this necessary condition is seen by exponentiating, noting that exponentiation respects conjugation: first,

$$
\left(e^{t x}\right)^{\sigma} S\left(e^{t x}\right)=S \cdot S\left(e^{t x}\right)^{\sigma} S \cdot\left(e^{t x}\right) S \cdot S e^{t x^{\sigma}} S e^{t x}=S \cdot e^{t \cdot S x^{\sigma} S} e^{t x}=S \cdot e^{t \cdot(-x)} e^{t x}=S \cdot e^{-t x} e^{t x}=S
$$

Thus,

$$
\text { Lie algebra } \mathfrak{s p}_{1,1}^{*}(\mathbb{H}) \text { of } S p_{1,1}^{*}=\left\{2 \text {-by-2 quaternionic }\left(\begin{array}{cc}
a & b \\
c & -\bar{a}
\end{array}\right): b=-\bar{b}, c=-\bar{c}\right\}
$$

[4.1.7] Example: In the simple case that the space $M$ is $G$ itself, there is a second action of $G$ on itself in addition to right multiplication, namely left multiplication. The right differentiation by elements of $\mathfrak{g}$ does commute with the left multiplication by $G$, for the simple reason that

$$
F\left(h \cdot\left(g e^{t x}\right)\right)=F\left((h \cdot g) \cdot e^{t x}\right) \quad(\text { for } g, h \in G, x \in \mathfrak{g})
$$

That is, $\mathfrak{g}$ gives left $G$-invariant differential operators on $G$.
[4.1.8] Claim: The conjugation action of $G$ on $\mathfrak{g}$ stabilizes $\mathfrak{g}$, and $g \cdot X_{x} \cdot g^{-1}=X_{g x g^{-1}}$ for $g \in G$ and $x \in \mathfrak{g}$. Proof: For smooth $f$ on $M$,

$$
\begin{gathered}
\left(g \cdot X_{x} \cdot g^{-1} \cdot f\right)(m)=\left(g\left(X_{x}\left(g^{-1} f\right)\right)\right)(m)=\left(X_{x}\left(g^{-1} f\right)\right)(m g) \\
=\left.\frac{d}{d t}\right|_{t=0}\left(g^{-1} f\right)\left(m g e^{t x}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(m g e^{t x} g^{-1}\right)
\end{gathered}
$$

Again, conjugation and exponentiation interact well:

$$
g e^{t x} g^{-1}=g\left(1+t x+\frac{(t x)^{2}}{2!}+\ldots\right) g^{-1}=1+t g x g^{-1}+\frac{\left(t g x g^{-1}\right)^{2}}{2!}+\ldots=e^{\operatorname{tg} x g^{-1}}
$$

Thus,

$$
\left(g \cdot X_{x} \cdot g^{-1} \cdot f\right)(m)=\left.\frac{d}{d t}\right|_{t=0} f\left(m g e^{t x} g^{-1}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(m e^{t g x g^{-1}}\right)=\left(X_{g x g^{-1}} f\right)(m)
$$

as claimed, and $g x g^{-1} \in \mathfrak{g}$.
The commutant expression $[x, y]=x y-y x$, the Lie bracket, arises naturally:

## 4. Invariant differential operators

[4.1.9] Claim: Since $G$ is non-abelian in many cases of interest, typically $e^{x} \cdot e^{y} \neq e^{y} \cdot e^{x}$ for $x, y \in \mathfrak{g}$. Specifically,

$$
e^{t x} e^{t y} e^{-t x} e^{-t y}=1+t^{2}[x, y]+(\text { higher-order terms }) \quad(\text { where }[x, y]=x y-y x)
$$

Proof: This is a direct computation, easy if we drop cubic and higher-order terms.

$$
\begin{gathered}
e^{t x} e^{t y} e^{-t x} e^{-t y}=\left(1+t x+t^{2} x^{2} / 2\right)\left(1+t y+t^{2} y^{2} / 2\right)\left(1-t x+t^{2} x^{2} / 2\right)\left(1-t y+t^{2} y^{2} / 2\right) \\
=\left(1+t(x+y)+\frac{t^{2}}{2}\left(x^{2}+2 x y+y^{2}\right)\right)\left(1-t(x+y)+\frac{t^{2}}{2}\left(x^{2}+2 x y+y^{2}\right)\right) \\
=1+t^{2}\left(x^{2}+2 x y+y^{2}-(x+y)(x+y)\right)=\left(1+t^{2}(2 x y-x y-y x)\right)=1+t^{2}[x, y]
\end{gathered}
$$

as claimed.
[4.1.10] Claim: The conjugation/adjoint action of $G$ on $\mathfrak{g}$ respects brackets:

$$
\left[g x g^{-1}, g y g^{-1}\right]=g[x, y] g^{-1} \quad(\text { for } x, y \in \mathfrak{g} \text { and } g \in G)
$$

Proof: For Lie brackets expressed in terms of matrix operations, this is straightforward:

$$
\left[g x g^{-1}, g y g^{-1}\right]=g x g^{-1} g y g^{-1}-g y g^{-1} g x g^{-1}=g x y g^{-1}-g y x g^{-1}=g(x y-y x) g^{-1}=g[x, y] g^{-1}
$$

as claimed.
Composition of the derivatives $X_{x}$ operators mirrors the bracket in the Lie algebra:
[4.1.11] Theorem: The map $x \rightarrow X_{x}$ is a Lie algebra homomorphism, meaning it respects these commutants (brackets): $X_{x} \circ X_{y}-X_{y} \circ X_{x}=X_{[x, y]}$. (Proof: see [4.A].)

### 4.2 Laplacians and Casimir operators

As in the last theorem, commutants of differential operators coming from Lie algebras $\mathfrak{g}$ are again differential operators coming from the Lie algebra, namely ${ }^{[39]}$

$$
X_{x} \circ X_{y}-X_{y} \circ X_{x}=\left[X_{x}, X_{y}\right]=X_{[x, y]}=X_{x y-y x}
$$

However, the composition of differential operators has no analogue inside the Lie algebra. That is, typically,

$$
X_{x} \circ X_{y} \neq X_{\varepsilon} \quad(\text { for any } \varepsilon \in \mathfrak{g})
$$

We want an object associated to the Lie algebra that allows this composition. [40] That is, we want an associative algebra $U \mathfrak{g}$ universal in the sense that any linear map $\varphi: \mathfrak{g} \rightarrow B$ to an associative algebra $B$ respecting brackets

$$
\varphi([x, y])=\varphi(x) \varphi(y)-\varphi(y) \varphi(x) \quad(\text { for } x, y \in \mathfrak{g})
$$

[39] For matrix groups with Lie bracket described via matrix multiplication $[x, y]=x y-y x$, properties otherwise needing explicit declaration, such as the Jacobi identity
$[x,[y, z]]-[y,[x, z]]=[[x, y], z]$, can be verified directly by expanding the brackets. The content of the Jacobi identity is that the map ad $: \mathfrak{g} \rightarrow \operatorname{End}(\mathfrak{g})$ by $(\operatorname{ad} x)(y)=[x, y]$ is a Lie algebra homomorphism. That is, $[\operatorname{ad} x, \operatorname{ad} y]=\operatorname{ad}[x, y]$.
[40] For Lie algebras $\mathfrak{g}$ such as $\mathfrak{s o}(n), \mathfrak{s l}_{n}$, or $\mathfrak{g l}_{n}$ lying inside matrix rings, typically $X_{x} \circ X_{y} \neq X_{x y}$. That is, multiplication of matrices is definitely not multiplication in any sense that will match multiplication (composition) of differential operators.
should give a unique associative algebra homomorphism $\Phi: U \mathfrak{g} \longrightarrow B$. There must be a connection to the original $\varphi: \mathfrak{g} \rightarrow B$, so we require existence of a fixed map $i: \mathfrak{g} \rightarrow U \mathfrak{g}$ respecting brackets and commutativity of a diagram
where the labels tell the type of the maps.
Below, we see that $U \mathfrak{g}$ is a canonical quotient of the universal associative algebra $A V$ of a vector space $V$ over a field $k$, very often called the tensor algebra and denoted $\otimes^{\bullet} V$, although, unhelpfully, this name refers to details of a specific construction, rather than to the characterizing property of the algebra. The characterizing property of the universal associative algebra $A V$ is that there is a fixed linear $j: V \rightarrow A V$, and any linear map $V \rightarrow B$ to an (associative) algebra $B$ extends to a unique associative algebra map $A V \rightarrow B$. That is, there is a commutative diagram


Since the universal associative algebra $j: \mathfrak{g} \rightarrow A \mathfrak{g}$ is universal with respect to maps $\mathfrak{g} \rightarrow B$ that are merely linear, not necessarily preserving the Lie brackets, there is a (unique) natural (quotient) map $q: A \mathfrak{g} \rightarrow U \mathfrak{g}$.

The conjugation (Adjoint) action $x \rightarrow g x g^{-1}$ of $G$ on $\mathfrak{g}$ should extend to an action of $G$ on $U \mathfrak{g}$ (which we may still write as conjugation) compatible with the multiplication in $U \mathfrak{g}$. That is, we expect

$$
\left\{\begin{array}{rlr}
g(\alpha)=g \alpha g^{-1} & & (\text { for } \alpha \in \mathfrak{g} \text { and } g \in G) \\
g(\alpha \beta) & =g(\alpha) \cdot g(\beta) & (\text { for } \alpha, \beta \in U \mathfrak{g} \text { and } g \in G)
\end{array}\right.
$$

The action of $G$ on $\mathfrak{g}$ should extend to $A \mathfrak{g}$, too, and the quotient map $q: A \mathfrak{g} \rightarrow U \mathfrak{g}$ should respect that action. We also assume for the moment that we have a non-degenerate symmetric bilinear form $\langle$,$\rangle on \mathfrak{g}$, and that this form is $G$-invariant: $\left\langle g x g^{-1}, g y g^{-1}\right\rangle=\langle x, y\rangle$ for $x, y \in \mathfrak{g}$ and $g \in G$.

Granting these things, we can intrinsically describe the simplest non-trivial $G$-invariant element in $U \mathfrak{g}$, the Casimir element $\Omega$. In any action of $G$, the Casimir element gives rise to a $G$-invariant differential operator, the corresponding Casimir operator. In many situations the Casimir operator is the suitable notion of invariant Laplacian. Map $\zeta:$ End $_{\mathbb{C}}(\mathfrak{g}) \rightarrow U \mathfrak{g}$ by
where the first map is an instance of the inverse of the isomorphism $V \otimes V^{*} \rightarrow$ End $V$ for finite-dimensional vectorspaces $V$, the second map uses the inverse of the isomorphism $V \rightarrow V^{*}$ given by $v \rightarrow\langle-, v\rangle$ for a non-degenerate bilinear form $\langle$,$\rangle on V$. The action of $G$ respects all the maps. An obvious endomorphism of $\mathfrak{g}$ commuting with the action of $G$ on $\mathfrak{g}$ is the identity map $\mathrm{id}_{\mathfrak{g}}$. Thus,
[4.2.1] Claim: The Casimir element $\Omega=\zeta\left(\mathrm{id}_{\mathfrak{g}}\right)$ is a $G$-invariant element of $U \mathfrak{g}$.
Proof: Since $\zeta$ is $G$-equivariant by construction,

$$
g \zeta\left(\mathrm{id}_{\mathfrak{g}}\right) g^{-1}=\zeta\left(g \mathrm{id}_{\mathfrak{g}} g^{-1}\right)=\zeta\left(g g^{-1} \mathrm{id}_{\mathfrak{g}}\right)=\zeta\left(\mathrm{id}_{\mathfrak{g}}\right)
$$

since $\operatorname{id}_{g}$ commutes with any endomorphism of $\mathfrak{g}$. Thus, $\zeta\left(\mathrm{id}_{\mathfrak{g}}\right)$ is a $G$-invariant element of $U \mathfrak{g}$.

## 4. Invariant differential operators

The possible hazard is that $\zeta\left(\mathrm{id}_{\mathfrak{g}}\right)$ is accidentally 0 . This non-vanishing can be proven by demonstrating at least one associative algebra $B$ and $\mathfrak{g} \rightarrow B$ so that the induced image of Casimir is non-zero in $B$. [41]

The above prescription does implicitly tell how to express the Casimir element $\Omega=\zeta\left(\mathrm{id}_{\mathfrak{g}}\right)$ in various coordinates. Namely, for any basis $x_{1}, \ldots, x_{n}$ of $\mathfrak{g}$, let $\lambda_{1}, \ldots, \lambda_{n}$ be the corresponding dual dual basis of the dual $\mathfrak{g}^{*}: \lambda_{i}\left(x_{j}\right)$ is 0 or 1 as $i=j$ or not. Let $x_{1}^{*}, \ldots, x_{n}^{*}$ be the corresponding dual basis for $\mathfrak{g}$ in terms of $\langle$,$\rangle ,$ namely, $\left\langle x_{i}, x_{j}^{*}\right\rangle$ is 0 or 1 as $i=j$ or not. Then


The intrinsic description of the Casimir element as $\zeta\left(\mathrm{id}_{\mathfrak{g}}\right)$ shows that it does not depend upon the choice of basis $x_{1}, \ldots, x_{n}$. ${ }^{[42]}$

### 4.3 Details about universal algebras

We fill in details about $U \mathfrak{g}$ and $A \mathfrak{g}$, including constructions. Again, we want an associative algebra $U \mathfrak{g}$ such that any Lie algebra $\operatorname{map} \varphi: \mathfrak{g} \rightarrow B$ to an associative algebra $B$ with the property

$$
\varphi([x, y])=\varphi(x) \varphi(y)-\varphi(y) \varphi(x) \quad(\text { for } x, y \in \mathfrak{g})
$$

gives a unique associative algebra homomorphism $\Phi: U \mathfrak{g} \longrightarrow B$ fitting into a commutative diagram


Similarly, we want a universal associative algebra $A V$ of a vector space $V$ over a field $k$, with a specified linear $j: V \rightarrow A V$, such that any linear map $V \rightarrow B$ to an associative algebra $B$ extends to a unique associative algebra map $A V \rightarrow B$ fitting into a commutative diagram

Granting for a moment the existence of $A \mathfrak{g}$, construct $U \mathfrak{g}$ as the quotient of $A \mathfrak{g}$ by the two-sided ideal generated by all elements

$$
(j x \otimes j y-j y \otimes j x)-j[x, y] \quad(\text { where } x, y \in \mathfrak{g})
$$

[41] The non-vanishing is also a corollary of the Poincaré-Birkhoff-Witt theorem, but we need not invoke it.
[42] Some sources define the Casimir element as $\sum_{i} x_{i} x_{i}^{*}$ in the universal enveloping algebra, show by computation that it is $G$-invariant, and show by change-of-basis that the defined object is independent of the choice of basis. That element $\sum_{i} x_{i} x_{i}^{*}$ is of course the image in $U \mathfrak{g}$ of the tensor $\sum_{i} x_{i} \otimes x_{i}^{*}$ (discussed here) which is simply the image of $\mathrm{id}_{\mathfrak{g}}$ in coordinates.

The map $i: \mathfrak{g} \rightarrow U \mathfrak{g}$ is the obvious composite $q \circ j$. Given a Lie algebra map $\varphi: \mathfrak{g} \rightarrow B$ from $\mathfrak{g}$ to an associative algebra, we show that the induced map $\Phi: A \mathfrak{g} \longrightarrow B$ factors through $q: A \mathfrak{g} \longrightarrow U \mathfrak{g}$. Diagrammatically, we claim the existence of an arrow to fill in a commutative diagram


Indeed, the the Lie algebra homomorphism property $\varphi(x) \varphi(y)-\varphi(y) \varphi(x)-\varphi[x, y]=0$ and the commutativity imply that

$$
\Phi(j x \otimes j y-j y \otimes j x)-j[x, y])=0 \quad(\text { for all } x, y \in \mathfrak{g})
$$

That is, $\Phi$ vanishes on the kernel of the quotient map $q: A \mathfrak{g} \rightarrow U \mathfrak{g}$, so factors through this quotient map. This proves the existence of $U \mathfrak{g}$ in terms of $A \mathfrak{g}$.

The conjugation (Adjoint) action $x \rightarrow g x g^{-1}$ of $G$ on $\mathfrak{g}$ should extend to an action of $G$ on $U \mathfrak{g}$ (which we may still write as conjugation) compatible with the multiplication in $U \mathfrak{g}$. That is, we expect

$$
\left\{\begin{aligned}
g(\alpha)=g \alpha g^{-1} & (\text { for } \alpha \in \mathfrak{g} \text { and } g \in G) \\
g(\alpha \beta)=g(\alpha) \cdot g(\beta) & (\text { for } \alpha, \beta \in U \mathfrak{g} \text { and } g \in G)
\end{aligned}\right.
$$

The action of $G$ on $\mathfrak{g}$ should extend to $A \mathfrak{g}$, too, and the quotient map $q: A \mathfrak{g} \rightarrow U \mathfrak{g}$ should respect that action. Fulfillment of this requirement, or the observation that it is automatically fulfilled, is best understood from further details about $A \mathfrak{g}$, just below.
[4.3.1] Construction of universal associative algebras The tensor construction of $A \mathfrak{g}$ gives enough further information so that we can see that it inherits an action of $G$ from $\mathfrak{g}$, and that this action is inherited by $U \mathfrak{g}$. The construction of $A V$ in terms of tensors is

$$
A V=k \oplus V \oplus(V \otimes V) \oplus(V \otimes V \otimes V) \oplus \ldots
$$

with multiplication given by (the bilinear extension of) the obvious

$$
\left(v_{1} \otimes \ldots \otimes v_{m}\right) \cdot\left(w_{1} \otimes \ldots \otimes w_{n}\right)=v_{1} \otimes \ldots \otimes v_{m} \otimes w_{1} \otimes \ldots \otimes w_{n}
$$

The well-definedness of the multiplication follows from noting that there is a unique linear map $\bigotimes^{m} V \otimes \otimes^{n} V \longrightarrow \bigotimes^{m+n} V$ induced from the bilinear map

$$
\left(v_{1} \otimes \ldots \otimes v_{m}\right) \times\left(w_{1} \otimes \ldots \otimes w_{n}\right) \longrightarrow v_{1} \otimes \ldots \otimes v_{m} \otimes w_{1} \otimes \ldots \otimes w_{n}
$$

Distributivity of multiplication over addition follows from the fact that the multiplication maps are induced from bilinear maps. The map $V \rightarrow A V$ is to the summand $V \subset A V$, which shows that this map is injective. It is also true that $\mathfrak{g} \rightarrow U \mathfrak{g}$ is injective, but the latter fact is considerably less trivial to prove.

To verify that this constructed object has the requisite universal property, let $\varphi: V \rightarrow B$ be a linear map to an associative algebra. Then the linear map $\Phi_{n}: \bigotimes^{n} V \rightarrow B$ defined by

$$
\left.\Phi\left(v_{1} \otimes \ldots \otimes v_{n}\right)=\varphi\left(v_{1}\right) \ldots \varphi\left(v_{n}\right) \quad \text { (latter is multiplication in } B\right)
$$

is well-defined, being induced from the $n$-multilinear map

$$
\underbrace{V \times \ldots \times V}_{n} \longrightarrow B \quad \text { by } \quad v_{1} \times \ldots \times v_{n} \longrightarrow \varphi\left(v_{1}\right) \ldots \varphi\left(v_{n}\right)
$$

## 4. Invariant differential operators

Letting $k$ be the underlying field (probably either $\mathbb{C}$ or $\mathbb{R}$ ), there is also the map $\Phi_{0}: k \rightarrow B$ by $a \rightarrow 1_{B}$. The collection of maps $\Phi_{n}$ gives a linear map $\Phi: A V \rightarrow B$. It also obviously preserves multiplication. This proves that the tensor construction yields the universal associative algebra.
[4.3.2] $G$-action on $A \mathfrak{g}$ and $U \mathfrak{g}$ The notationally obvious $G$-action on $A \mathfrak{g}$ is

$$
g\left(x_{1} \otimes \ldots \otimes x_{m}\right) g^{-1}=g x_{1} g^{-1} \otimes \ldots \otimes g x_{m} g^{-1}
$$

This gives a well-defined linear map of each $\bigotimes^{n} \mathfrak{g}$ to itself, because it is the unique map induced by the multilinear map

$$
\underbrace{\mathfrak{g} \times \ldots \times \mathfrak{g}}_{n} \longrightarrow \bigotimes^{n} \mathfrak{g} \quad \text { by } \quad v_{1} \times \ldots \times v_{n} \longrightarrow g v_{1} g^{-1} \otimes \ldots \otimes g v_{n} g^{-1}
$$

The map is visibly compatible with multiplication. Since $\mathfrak{g}$ injects to $A \mathfrak{g}$, we can safely suppress the map $j$ in this discussion. The $G$-action stabilizes the kernel of the kernel of $q: A \mathfrak{g} \rightarrow U \mathfrak{g}$, since

$$
\begin{gathered}
g((x \otimes y-y \otimes x)-[x, y]) g^{-1}=g(x \otimes y) g^{-1}-g(y \otimes x) g^{-1}-g[x, y] g^{-1} \\
=g x g^{-1} \otimes g y g^{-1}-g y g^{-1} \otimes g x g^{-1}-\left[g x g^{-1}, g y g^{-1}\right]
\end{gathered}
$$

This gives a natural action of $G$ on $U \mathfrak{g}$, respecting the quotient $q: A \mathfrak{g} \rightarrow U \mathfrak{g}$, and, therefore, respecting the map $\mathfrak{g} \rightarrow U \mathfrak{g}$. The universal associative algebra $A \mathfrak{g}$ is sufficiently large that, roughly, it has no non-trivial relations. Thus, the notationally-obvious apparent definition of the $G$-action on $A \mathfrak{g}$ is well-defined. Then the $G$-action descends to $U \mathfrak{g}$.
[4.3.3] Killing's bilinear form The last necessary item is more special, and not possessed by all Lie algebras. We want a non-degenerate symmetric $\mathbb{R}$-bilinear map

$$
\langle,\rangle: \mathfrak{g} \times \mathfrak{g} \longrightarrow \mathbb{R}
$$

$G$-equivariant in the sense that

$$
\left\langle g x g^{-1}, g y g^{-1}\right\rangle=\langle x, y\rangle
$$

Happily, for $\mathfrak{s o}(n), \mathfrak{s l}_{n}(\mathbb{R})$, and $\mathfrak{g l}_{n}(\mathbb{R})$, the obvious trace form

$$
\langle x, y\rangle=\operatorname{tr}(x y)
$$

suffices. For $G$ described as a subgroup of $G L_{n}(\mathbb{C})$, take $\langle x, y\rangle=\operatorname{Re}(\operatorname{tr} x y)$. For $G$ described as a subgroup of $G L_{n}(\mathbb{H})$, take $\langle x, y\rangle=\sum_{i j}\left(x_{i j} y_{j i}+\overline{x_{i j} y_{j i}}\right)$. For notational simplicity, we write out the arguments only for $G \subset G L_{n}(\mathbb{R})$. The behavior under the action of $G$ is clear:

$$
\left\langle g x g^{-1}, g y g^{-1}\right\rangle=\operatorname{tr}\left(g x g^{-1} \cdot g y g^{-1}\right)=\operatorname{tr}\left(g x y g^{-1}\right)=\operatorname{tr}(x y)=\langle x, y\rangle
$$

The non-degeneracy and $G$-equivariance of $\langle$,$\rangle give a natural G$-equivariant isomorphism $\mathfrak{g} \rightarrow \mathfrak{g}^{*}$ by

$$
x \longrightarrow \lambda_{x} \quad \text { by } \quad \lambda_{x}(y)=\langle x, y\rangle \quad(\text { for } x, y \in \mathfrak{g})
$$

When $G$ acts on a vector space $V$ the action on the dual $V^{*}$ is by $(g \cdot \lambda)(v)=\lambda\left(g^{-1} \cdot v\right)$ for $v \in V$ and $\lambda \in V^{*}$. As usual, the inverse appears to preserve associativity. The equivariance of $\langle$,$\rangle gives$

$$
\lambda_{g \cdot x}(y)=\lambda_{g x g^{-1}}(y)=\left\langle g x g^{-1}, y\right\rangle=\left\langle x, g^{-1} y g\right\rangle=\lambda_{x}\left(g^{-1} y g\right)=\lambda_{x}\left(g^{-1} \cdot y\right)=\left(g \cdot \lambda_{x}\right)(y)
$$

proving that the map $x \rightarrow \lambda_{x}$ is a $G$-isomorphism.
Finally, recall that the natural isomorphism $V \otimes_{k} V^{*} \longrightarrow \operatorname{End}_{k} V$ for $V$ a finite-dimensional vector space over a field $k$ is given by the $k$-linear extension of the map $v \times \lambda \longrightarrow(w \rightarrow \lambda(w) \cdot v)$ for $v, w \in V$ and $\lambda \in V^{*}$. The fact that the map is an isomorphism follows by dimension counting, using the finite-dimensionality.

### 4.4 Descending to $G / K$

Now we see how the Casimir operator $\Omega$ on $G$ gives $G$-invariant Laplacian-like differential operators on quotients $G / K$. Let $\mathfrak{k} \subset \mathfrak{g}$ be the Lie algebra of a maximal compact $K \subset G$. Again, the action of $x \in \mathfrak{g}$ on the right on functions $f$ on $G$, by

$$
(x \cdot f)(g)=\left.\frac{d}{d t}\right|_{t=0} f\left(g e^{t x}\right)
$$

is left $G$-invariant for the straightforward reason that

$$
\left.f\left(h \cdot\left(g e^{t x}\right)\right)=f\left((h \cdot g) \cdot e^{t x}\right)\right) \quad(\text { for } g, h \in G, x \in \mathfrak{g})
$$

For a (closed) subgroup $K$ of $G$ let $q: G \rightarrow G / K$ be the quotient map. A function $f$ on $G / K$ gives the right $K$-invariant function $F=f \circ q$ on $G$. Given $x \in \mathfrak{g}$, the differentiation

$$
(x \cdot(f \circ q))(g)=\left.\frac{d}{d t}\right|_{t=0}(f \circ q)\left(g e^{t x}\right)
$$

makes sense. However, $x \cdot(f \circ q)$ is not usually right $K$-invariant. Indeed, the condition for right $K$-invariance is

$$
\left.\frac{d}{d t}\right|_{t=0} F\left(g e^{t x}\right)=(x \cdot F)(g)=(x \cdot F)(g k)=\left.\frac{d}{d t}\right|_{t=0} F\left(g k e^{t x}\right) \quad(k \in \mathfrak{k})
$$

Using the right $K$-invariance of $F=f \circ q$,

$$
F\left(g k e^{t x}\right)=F\left(g k e^{t x} k^{-1} k\right)=F\left(g e^{t \cdot k x k^{-1}}\right)
$$

Thus, unless $k x k^{-1}=x$ for all $k \in K$, it is unlikely that $x \cdot F$ is still right $K$-invariant. That is, the left $G$-invariant differential operators coming from $\mathfrak{g}$ usually do not descend to differential operators on $G / K$.

The differential operators in the $G$-invariant subalgebra

$$
\mathfrak{z}=\left\{\alpha \in U \mathfrak{g}: g \alpha g^{-1}\right\}
$$

do descend to $G / K$, exactly because of the commutation property, as follows. For any function $\varphi$ on $G$ let $(k \cdot \varphi)(g)=\varphi(g k)$. For $F$ right $K$-invariant on $G$, for $\alpha \in Z(\mathfrak{g})$ compute directly

$$
k \cdot(\alpha \cdot F)=\alpha \cdot(k \cdot F)=\alpha \cdot F
$$

showing the right $K$-invariance of $\alpha \cdot F$. Thus, $\alpha \cdot F$ gives a well-defined function on $G / K$.

### 4.5 Example computation: $S L_{2}(\mathbb{R})$ and $\mathfrak{H}$

Let $\mathfrak{g}=\mathfrak{s l}_{2}(\mathbb{R})$, the Lie algebra of $G=S L_{2}(\mathbb{R})$. A typical choice of basis for $\mathfrak{g}$ is ${ }^{\text {[43] }}$

$$
H=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad X=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right) \quad Y=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right)
$$

These have the easily verified relations

$$
[H, X]=H X-X H=2 X \quad[H, Y]=H Y-Y H=-2 Y \quad[X, Y]=X Y-Y X=H
$$

[^12]
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To see that the pairing $\langle x, y\rangle=\operatorname{tr}(x y)$ is non-degenerate, use the stability of $\mathfrak{g}$ under transpose $v \rightarrow v^{\top}$, and then

$$
\left\langle v, v^{\top}\right\rangle=\operatorname{tr}\left(v v^{\top}\right)=2 a^{2}+b^{2}+c^{2} \quad\left(\text { for } v=\left(\begin{array}{rr}
a & b \\
c & -a
\end{array}\right)\right)
$$

We easily compute that

$$
\langle H, H\rangle=2 \quad\langle H, X\rangle=0 \quad\langle H, Y\rangle=0 \quad\langle X, Y\rangle=1
$$

Thus, for the basis $H, X, Y$ we have dual basis $H^{*}=H / 2, X^{*}=Y$, and $Y^{*}=X$, and in these coordinates the Casimir operator is

$$
\Omega=H H^{*}+X X^{*}+Y Y^{*}=\frac{1}{2} H^{2}+X Y+Y X \quad \text { (now inside } U \mathfrak{g} \text { ) }
$$

Since $X Y-Y X=H^{[44]}$ the expression for $\Omega$ can be rewritten in various useful forms, such as

$$
\Omega=\frac{1}{2} H^{2}+X Y+Y X=\frac{1}{2} H^{2}+X Y-Y X+2 Y X=\frac{1}{2} H^{2}+H+2 Y X
$$

and, similarly,

$$
\Omega=\frac{1}{2} H^{2}+X Y+Y X=\frac{1}{2} H^{2}+X Y-(-Y X)=\frac{1}{2} H^{2}+2 X Y-(X Y-Y X)=\frac{1}{2} H^{2}+2 X Y-H
$$

To obtain a $G$-invariant differential operator on the upper half-plane $\mathfrak{H}$ from $\Omega$, use the $G$-space isomorphism $\mathfrak{H} \approx G / K$ where $K=S O_{2}(\mathbb{R})$ is the isotropy group of the point $i \in \mathfrak{H}$. Let $q: G \rightarrow G / K$ be the quotient map $q(g)=g K \longleftrightarrow g(i)$. A function $f$ on $\mathfrak{H}$ naturally yields the right $K$-invariant function $f \circ q$

$$
(f \circ q)(g)=f(g(i)) \quad(\text { for } g \in G)
$$

As above, for any $z \in \mathfrak{g}$ there is the corresponding left $G$-invariant differential operator on a function $F$ on $G$ by

$$
(z \cdot F)(g)=\left.\frac{d}{d t}\right|_{t=0} F\left(g e^{t z}\right)
$$

but these linear operators should not be expected to descend to operators on $G / K$. Nevertheless, $G$-invariant elements such as the Casimir operator $\Omega$ in $Z(\mathfrak{g})$ do descend.

The computation of $\Omega$ on $f \circ q$ can be simplified by using the right $K$-invariance of $f \circ q$, which implies that $f \circ q$ is annihilated by

$$
\mathfrak{s o}_{2}(\mathbb{R})=\text { Lie algebra of } S O_{2}(\mathbb{R})=\text { skew-symmetric 2-by-2 real matrices }=\left\{\left(\begin{array}{rr}
0 & t \\
-t & 0
\end{array}\right): t \in \mathbb{R}\right\}
$$

Thus, in terms of the basis $H, X, Y$ above, $X-Y$ annihilates $f \circ q$.
Among other possibilities, a point $z=x+i y \in \mathfrak{H}$ is the image

$$
x+i y=(n \cdot m)(i) \quad \text { where } \quad n_{x}=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right) \quad m_{y}=\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & \frac{1}{\sqrt{y}}
\end{array}\right)
$$

These are convenient group elements because they match the exponentiated Lie algebra elements:

$$
e^{t X}=n_{t} \quad e^{t H}=m_{e^{2 t}}
$$

In contrast, the exponentiated $Y$ has a more complicated action on $\mathfrak{H}$. This suggests invocation of the fact that $X-Y$ acts trivially on right $K$-invariant functions on $G$. Order matters: application of a differential operator typically disrupts right $K$-invariance. For right $K$-invariant $F$ on $G$,
[44] The identity $X Y-Y X=H$ holds both in the universal enveloping algebra and as matrices.

$$
\begin{gathered}
(\Omega F)\left(n_{x} m_{y}\right)=\left(\frac{1}{2} H^{2}+X Y+Y X\right) F\left(n_{x} m_{y}\right)=\left(\frac{1}{2} H^{2}+2 X Y-X Y+Y X\right) F\left(n_{x} m_{y}\right) \\
\quad=\left(\frac{1}{2} H^{2}+2 X^{2}+2 X(Y-X)-[X, Y]\right) F\left(n_{x} m_{y}\right)=\left(\frac{1}{2} H^{2}+2 X^{2}-H\right) F\left(n_{x} m_{y}\right)
\end{gathered}
$$

Compute the pieces separately. First, using the identity

$$
m_{y} n_{x}=\left(m_{y} n_{x} m_{y}^{-1}\right) m_{y}=\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & \frac{1}{\sqrt{y}}
\end{array}\right)\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & \frac{1}{\sqrt{y}}
\end{array}\right)^{-1} m_{y}=n_{y x} m_{y}
$$

we compute the effect of $X$

$$
(X \cdot F)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} F\left(n_{x} m_{y} n_{t}\right)=\left.\frac{d}{d t}\right|_{t=0} F\left(n_{x} n_{y t} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} F\left(n_{x+y t} m_{y}\right)=y \frac{\partial}{\partial x} F\left(n_{x} m_{y}\right)
$$

Thus,

$$
2 X^{2} \longrightarrow 2\left(y \frac{\partial}{\partial x}\right)^{2}=2 y^{2}\left(\frac{\partial}{\partial x}\right)^{2}
$$

The action of $H$ is

$$
(H \cdot F)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} F\left(n_{x} m_{y} m_{e^{2 t}}\right)=\left.\frac{d}{d t}\right|_{t=0} F\left(n_{x} m_{y e^{2 t}}\right)=2 y \frac{\partial}{\partial y} F\left(n_{x} m_{y}\right)
$$

Then

$$
\frac{1}{2} H^{2}-H=\frac{1}{2}\left(2 y \frac{\partial}{\partial y}\right)^{2}-\left(2 y \frac{\partial}{\partial y}\right)=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}+2 y \frac{\partial}{\partial y}-2 y \frac{\partial}{\partial y}=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}
$$

Altogether, on right $K$-invariant functions $F$,

$$
(\Omega F)\left(n_{x} m_{y}\right)=2 y^{2}\left(\left(\frac{\partial}{\partial x}\right)^{2}+\left(\frac{\partial}{\partial y}\right)^{2}\right) F\left(m_{x} n_{y}\right)
$$

That is, in the usual coordinates $z=x+i y$ on $\mathfrak{H}$, discarding the leading constant,

$$
\text { (image of) } \Omega=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)
$$

### 4.6 Example computation: $S L_{2}(\mathbb{C})$

Let $\mathfrak{g}=\mathfrak{s l}_{2}(\mathbb{C})$, the Lie algebra of $G=S L_{2}(\mathbb{C})$, with basis

$$
H=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad H^{\prime}=\left(\begin{array}{rr}
i & 0 \\
0 & -i
\end{array}\right) \quad X=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right) \quad X^{\prime}=\left(\begin{array}{ll}
0 & i \\
0 & 0
\end{array}\right) \quad Y=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right) \quad Y^{\prime}=\left(\begin{array}{cc}
0 & 0 \\
-i & 0
\end{array}\right)
$$

with $[X, Y]=H$ and $\left[X^{\prime}, Y^{\prime}\right]=H$, and so on. To see that the pairing $\langle x, y\rangle=\operatorname{Retr}(x y)$ is non-degenerate, use the stability of $\mathfrak{g}$ under conjugate-transpose $v \rightarrow v^{*}=\bar{v}^{\top}$, and then

$$
\left\langle v, v^{*}\right\rangle=\operatorname{Retr}\left(v v^{*}\right)=2|a|^{2}+|b|^{2}+|c|^{2} \quad\left(\text { for } v=\left(\begin{array}{rr}
a & b \\
c & -a
\end{array}\right)\right)
$$

We easily compute that

$$
\langle H, H\rangle=2 \quad\left\langle H^{\prime}, H^{\prime}\right\rangle=-2 \quad\langle X, Y\rangle=1 \quad\left\langle X^{\prime}, Y^{\prime}\right\rangle=1
$$

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and all other pairings give 0 . Thus, for the basis $H, X, Y, H^{\prime}, X^{\prime}, Y^{\prime}$ we have dual basis $H^{*}=H / 2, X^{*}=Y$, $Y^{*}=X, H^{\prime *}=-H^{\prime} / 2, X^{*}=Y^{\prime}, Y^{\prime *}=X^{\prime}$, In these coordinates the Casimir operator is
$\Omega=H H^{*}+X X^{*}+Y Y^{*}+H^{\prime} H^{* *}+X^{\prime} X^{* *}+Y^{\prime} Y^{\prime *}=\frac{1}{2} H^{2}+X Y+Y X-\frac{1}{2} H^{\prime 2}+X^{\prime} Y^{\prime}+Y^{\prime} X^{\prime}$
(in $U \mathfrak{g}$ )
Let $q: G \rightarrow G / K$ be the quotient map and use Iwasawa coordinates

$$
\left.n_{x} a_{y} k=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & \frac{1}{\sqrt{y}}
\end{array}\right) \cdot k \quad \text { (with } x \in \mathbb{C}, y>0, \text { and } k \in K=S U_{2}\right)
$$

For any $z \in \mathfrak{g}$ there is the corresponding left $G$-invariant differential operator on a function $F$ on $G$ by

$$
(z \cdot F)(g)=\left.\frac{d}{d t}\right|_{t=0} F\left(g e^{t z}\right)
$$

but these linear operators generally do not descend to operators on $G / K$, that is, are not right $K$-invariant. Nevertheless, $G$-invariant elements such Casimir do descend.

Differentiating the condition $1_{2}=\left(e^{t x}\right)^{*}\left(e^{t x}\right)$ with respect to $t$ gives

$$
0=\frac{d}{d t} 1_{2}=x^{*}\left(e^{t x}\right)^{*}\left(e^{t x}\right)+\left(e^{t x}\right)^{*}\left(e^{t x}\right) x
$$

and setting $t=0$ gives a necessary condition for $x$ to be in the Lie algebra $\mathfrak{k}$ of $K=S U_{2}$, namely, $0=x^{*}+x$. Conversely, for $x^{*}=-x$, exponentiation gives $1_{2}=\left(e^{t x}\right)^{*}\left(e^{t x}\right)$. The determinant-one condition gives $\operatorname{tr} x=0$. Thus, the Lie algebra $\mathfrak{k}=\mathfrak{s u}_{2}$ of $K=S U_{2}$ is

$$
\mathfrak{k}=\text { skew-hermitian 2-by-2 trace } 0 \text { complex matrices }=\left\{\left(\begin{array}{cc}
i t & \tau \\
-\bar{\tau} & -i t
\end{array}\right): t \in \mathbb{R}, \tau \in \mathbb{C}\right\}
$$

The computation of $\Omega$ on $G / K$ is simplified by using the right $K$-invariance: the right action of $\mathfrak{k}$ annihilates right $K$-invariant functions on $G$. In terms of the basis $H, X, Y, H^{\prime}, X^{\prime}, Y^{\prime}$ above, $H^{\prime}, X-Y$ and $X^{\prime}+Y^{\prime}$ are all in $\mathfrak{k}$, so annihilate right $K$-invariant functions. Order matters: application of a differential operator typically disrupts right $K$-invariance. First, rearrange $\Omega$ in anticipation of application to right $K$-invariant $f$ on $G$ :

$$
\begin{gathered}
\Omega=\frac{1}{2} H^{2}+X Y+Y X-\frac{1}{2} H^{\prime 2}+X^{\prime} Y^{\prime}+Y^{\prime} X^{\prime}=\frac{1}{2} H^{2}+2 X Y-H-\frac{1}{2} H^{\prime 2}+2 X^{\prime} Y^{\prime}-H \\
=\frac{1}{2} H^{2}-2 H+2 X^{2}+2 X(Y-X)-\frac{1}{2} H^{\prime 2}-2 X^{\prime 2}+2 X^{\prime}\left(Y^{\prime}+X^{\prime}\right)
\end{gathered}
$$

Since $H^{\prime}, X-Y$, and $X^{\prime}+Y^{\prime}$ annihilate right $K$-invariant functions $f$, this gives

$$
\Omega f=\left(\frac{1}{2} H^{2}-2 H+2 X^{2}-2 X^{\prime 2}\right) f
$$

Compute the pieces separately. Use coordinates $x=x_{1}+i x_{2} \in \mathbb{C}$. Using $m_{y} n_{x}=n_{y t} m_{y}$, the effects of $X$ and $X^{\prime}$ are

$$
\begin{gathered}
(X \cdot f)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} n_{t}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} n_{y t} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x+y t} m_{y}\right)=y \frac{\partial}{\partial x_{1}} f\left(n_{x} m_{y}\right) \\
\left(X^{\prime} \cdot f\right)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} n_{i t}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} n_{i y t} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x+i y t} m_{y}\right)=i y \frac{\partial}{\partial x_{2}} f\left(n_{x} m_{y}\right)
\end{gathered}
$$

Thus,

$$
2 X^{2} \longrightarrow 2\left(y \frac{\partial}{\partial x_{1}}\right)^{2}=2 y^{2}\left(\frac{\partial}{\partial x_{1}}\right)^{2} \text { and }-2 X^{\prime 2} \longrightarrow-2\left(i y \frac{\partial}{\partial x_{2}}\right)^{2}=2 y^{2}\left(\frac{\partial}{\partial x_{2}}\right)^{2}
$$

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As for $S L_{2}(\mathbb{R})$, the action of $H$ is

$$
(H \cdot f)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} m_{e^{2 t}}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y e^{2 t}}\right)=2 y \frac{\partial}{\partial y} f\left(n_{x} m_{y}\right)
$$

so

$$
\frac{1}{2} H^{2}-2 H=\frac{1}{2}\left(2 y \frac{\partial}{\partial y}\right)^{2}-2\left(2 y \frac{\partial}{\partial y}\right)=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}+2 y \frac{\partial}{\partial y}-4 y \frac{\partial}{\partial y}=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}-2 y \frac{\partial}{\partial y}
$$

Altogether, on right $K$-invariant functions $f$, discarding the irrelevant leading constant 2,

$$
\Omega \longrightarrow y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-y \frac{\partial}{\partial y}
$$

In contrast to the situation for $S L_{2}(\mathbb{R})$ and $\mathfrak{H}$, this is not merely a multiple of the Euclidean Laplacian.

### 4.7 Example computation: $S p_{1,1}^{*}$

Let $\mathfrak{g}=\mathfrak{s p}_{1,1}^{*}$, the Lie algebra of $G=S p_{1,1}^{*}$, with ten-element basis

$$
H=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad H_{\ell}=\left(\begin{array}{rr}
\ell & 0 \\
0 & -\ell
\end{array}\right) \quad X_{\ell}=\left(\begin{array}{ll}
0 & \ell \\
0 & 0
\end{array}\right) \quad Y_{\ell}=\left(\begin{array}{rr}
0 & 0 \\
-\ell & 0
\end{array}\right) \quad(\text { with } \ell=i, j, k)
$$

Note that

$$
\left[X_{i}, Y_{i}\right]=\left[X_{j}, Y_{j}\right]=\left[X_{k}, Y_{k}\right]=H
$$

With $\operatorname{tr}(a+b i+c j+d k)=a$, use the pairing

$$
\left\langle\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right),\left(\begin{array}{cc}
\alpha^{\prime} & \beta^{\prime} \\
\gamma^{\prime} & \delta^{\prime}
\end{array}\right)\right\rangle=\operatorname{tr}\left(\alpha \cdot \alpha^{\prime}+\beta \cdot \gamma^{\prime}+\gamma \cdot \beta^{\prime}+\delta \cdot \delta^{\prime}\right)
$$

The non-degeneracy follows from the stability of $\mathfrak{g}$ under conjugate-transpose $v \rightarrow v^{*}$, and then noting that $v \rightarrow\left\langle v, v^{*}\right\rangle$ is positive-definite. Compute that

$$
\langle H, H\rangle=2 \quad\left\langle H_{\ell}, H_{\ell}\right\rangle=-2 \quad\left\langle X_{\ell}, Y_{\ell}\right\rangle=1
$$

and all other pairings give 0 . Thus, we have dual basis $H^{*}=H / 2, H_{\ell}^{*}=-H_{\ell} / 2, X_{\ell}^{*}=Y_{\ell}, Y_{\ell}^{*}=X_{\ell}$. In these coordinates the Casimir operator is

$$
\Omega=H H^{*}+\sum_{\ell=i, j, k}\left(H_{\ell} H_{\ell}^{*}+X_{\ell} X_{\ell}^{*}+Y_{\ell} Y_{\ell}^{*}\right)=\frac{1}{2} H^{2}+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}+X_{\ell} Y_{\ell}+Y_{\ell} X_{\ell}\right)
$$

Use Iwasawa coordinates

$$
n_{x} a_{y} k=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & \frac{1}{\sqrt{y}}
\end{array}\right) \cdot k \quad\left(\text { with } x=i x_{1}+j x_{2}+k x_{3} \in \mathbb{H}, y>0, \text { and } k \in K\right)
$$

The Lie algebra $\mathfrak{k}$ of the compact subgroup $K \approx S p_{1}^{*} \times S p_{1}^{*}$ of $G$ can be identified by observing two copies of the Lie algebra of $S p_{1}^{*}$, as follows. By differentiating,

$$
0=\frac{d}{d t} 1=\frac{d}{d t} \overline{\left(e^{t} x\right)}\left(e^{t} x\right)=\frac{d}{d t}\left(e^{t} \bar{x}\right)\left(e^{t} x\right)=\bar{x} \cdot\left(e^{t} \bar{x}\right)\left(e^{t} x\right)+\left(e^{t} \bar{x}\right)\left(e^{t} x\right) x
$$

and at $t=0$ this is $\bar{x}+x=0$. We can observe two suitable copies of this inside $\mathfrak{s p}_{1,1}^{*}$, by

$$
\mathfrak{k}=\left\{\left(\begin{array}{cc}
\alpha & 0 \\
0 & -\alpha
\end{array}\right): \bar{\alpha}=-\alpha\right\} \oplus\left\{\left(\begin{array}{cc}
0 & \beta \\
-\beta & 0
\end{array}\right): \bar{\beta}=-\beta\right\}
$$

## 4. Invariant differential operators

The computation of $\Omega$ on $G / K$ is simplified by using the right $K$-invariance, which entails annihilation by $\mathfrak{k}$. In terms of the basis above, $H_{\ell}$ and $X_{\ell}+Y_{\ell}$, for $\ell=i, j, k$, are all in $\mathfrak{k}$. Rearrange $\Omega$ in anticipation of application to right $K$-invariant $f$ on $G$ :

$$
\begin{aligned}
& \Omega=\frac{1}{2} H^{2}+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}+X_{\ell} Y_{\ell}+Y_{\ell} X_{\ell}\right)=\frac{1}{2} H^{2}+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}+2 X_{\ell} Y_{\ell}-X_{\ell} Y_{\ell}+Y_{\ell} X_{\ell}\right) \\
= & \frac{1}{2} H^{2}+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}+2 X_{\ell} Y_{\ell}-\left[X_{\ell}, Y_{\ell}\right]\right)=\frac{1}{2} H^{2}-3 H+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}-2 X_{\ell}^{2}+X_{\ell}\left(X_{\ell}+Y_{\ell}\right)\right)
\end{aligned}
$$

Since the elements $H_{\ell}$ and $X_{\ell}+Y_{\ell}$ annihilate right $K$-invariant functions $f$, this gives

$$
\Omega f=\left(\frac{1}{2} H^{2}-3 H-\sum_{\ell=i, j, k} 2 X_{\ell}^{2}\right) f
$$

Compute the pieces separately, with coordinates $x=i x_{1}+j x_{2}+k x_{3} \in \mathbb{H}$. Using $m_{y} n_{x}=n_{y t} m_{y}$, the effects of $X_{i}, X_{j}, X_{k}$ are

$$
\begin{aligned}
& \left(X_{i} \cdot f\right)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} n_{i t}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} n_{i y t} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x+i y t} m_{y}\right)=y \frac{\partial}{\partial x_{1}} f\left(n_{x} m_{y}\right) \\
& \left(X_{j} \cdot f\right)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} n_{j t}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} n_{j y t} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x+j y t} m_{y}\right)=j y \frac{\partial}{\partial x_{2}} f\left(n_{x} m_{y}\right) \\
& \left(X_{k} \cdot f\right)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} n_{k t}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} n_{k y t} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x+k y t} m_{y}\right)=k y \frac{\partial}{\partial x_{3}} f\left(n_{x} m_{y}\right)
\end{aligned}
$$

Thus,

$$
-2 X_{i}^{2}-2 X_{j}^{2}-2 X_{k}^{2} \longrightarrow-2\left(i y \frac{\partial}{\partial x_{1}}\right)^{2}-2\left(j y \frac{\partial}{\partial x_{2}}\right)^{2}-2\left(k y \frac{\partial}{\partial x_{3}}\right)^{2}=2 y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial x_{3}^{2}}\right)
$$

As for $S L_{2}(\mathbb{R})$ and $S L_{2}(\mathbb{C})$, the action of $H$ is

$$
(H \cdot f)\left(n_{x} m_{y}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y} m_{e^{2 t}}\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(n_{x} m_{y e^{2 t}}\right)=2 y \frac{\partial}{\partial y} f\left(n_{x} m_{y}\right)
$$

so

$$
\frac{1}{2} H^{2}-3 H=\frac{1}{2}\left(2 y \frac{\partial}{\partial y}\right)^{2}-3\left(2 y \frac{\partial}{\partial y}\right)=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}+2 y \frac{\partial}{\partial y}-6 y \frac{\partial}{\partial y}=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}-4 y \frac{\partial}{\partial y}
$$

Altogether, on right $K$-invariant functions $f$, discarding the irrelevant leading constant 2,

$$
\Omega \longrightarrow y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial x_{3}^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-2 y \frac{\partial}{\partial y}
$$

Again, as with $S L_{2}(\mathbb{C})$, in contrast to the situation for $S L_{2}(\mathbb{R})$ and $\mathfrak{H}$, this is not merely a multiple of the Euclidean Laplacian.

### 4.8 Example computation: $S L_{2}(\mathbb{H})$

Let $\mathfrak{g}=\mathfrak{s l}_{2}(\mathbb{H})$, the Lie algebra of $G=S L_{2}(\mathbb{H})$. Letting $\ell$ run over $i, j, k$, take basis

$$
\begin{gathered}
H=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad H_{\ell}=\left(\begin{array}{rr}
\ell & 0 \\
0 & -\ell
\end{array}\right) \quad H_{\ell}^{\prime}=\left(\begin{array}{ll}
0 & 0 \\
0 & \ell
\end{array}\right) \\
X=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right) \quad X_{\ell}=\left(\begin{array}{ll}
0 & \ell \\
0 & 0
\end{array}\right) \quad Y=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right) \quad Y_{\ell}=\left(\begin{array}{rr}
0 & 0 \\
-\ell & 0
\end{array}\right)
\end{gathered}
$$

Note that

$$
[X, Y]=\left[X_{i}, Y_{i}\right]=\left[X_{j}, Y_{j}\right]=\left[X_{k}, Y_{k}\right]=H
$$

With $\operatorname{tr}(a+b i+c j+d k)=a$, as with $S p_{1,1}^{*}$, use the pairing

$$
\left\langle\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right),\left(\begin{array}{ll}
\alpha^{\prime} & \beta^{\prime} \\
\gamma^{\prime} & \delta^{\prime}
\end{array}\right)\right\rangle=\operatorname{tr}\left(\alpha \cdot \alpha^{\prime}+\beta \cdot \gamma^{\prime}+\gamma \cdot \beta^{\prime}+\delta \cdot \delta^{\prime}\right)
$$

Compute that

$$
\langle H, H\rangle=2 \quad\left\langle H_{\ell}, H_{\ell}\right\rangle=-2 \quad\left\langle H_{\ell}^{\prime}, H_{\ell}^{\prime}\right\rangle=-1 \quad\left\langle X_{\ell}, Y_{\ell}\right\rangle=1
$$

and all other pairings give 0 . Thus, we have dual basis $H^{*}=H / 2, H_{\ell}^{*}=-H_{\ell} / 2, H_{\ell}^{\prime *}=-H_{\ell}^{\prime}, X_{\ell}^{*}=Y_{\ell}$, $Y_{\ell}^{*}=X_{\ell}$. In these coordinates the Casimir operator is

$$
\begin{aligned}
& \Omega=H H^{*}+X X^{*}+Y Y^{*}+\sum_{\ell=i, j, k}\left(H_{\ell} H_{\ell}^{*}+H_{\ell}^{\prime} H_{\ell}^{\prime *}+X_{\ell} X_{\ell}^{*}+Y_{\ell} Y_{\ell}^{*}\right) \\
= & \frac{1}{2} H^{2}+X Y+Y X+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}-H_{\ell}^{\prime} H_{\ell}^{\prime}+X_{\ell} Y_{\ell}+Y_{\ell} X_{\ell}\right) \quad(\text { in } U \mathfrak{g})
\end{aligned}
$$

Use Iwasawa coordinates

$$
n_{x} a_{y} k=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & \frac{1}{\sqrt{y}}
\end{array}\right) \cdot k \quad\left(\text { with } x=x_{1}+i x_{2}+j x_{3}+k x_{4} \in \mathbb{H}, y>0, \text { and } k \in K\right)
$$

To determine the Lie algebra $\mathfrak{k}$ of the compact subgroup $K=S p_{2}^{*}$ of $G$, differentiate

$$
0=\frac{d}{d t} 1_{2}=\frac{d}{d t}\left(\left(e^{t x}\right)^{*}\left(e^{t x}\right)\right)=x^{*}\left(e^{t x}\right)^{*}\left(e^{t x}\right)+\left(e^{t x}\right)^{*}\left(e^{t x}\right) x
$$

and at $t=0$ obtain $x^{*}+x=0$. As usual, the converse follows by exponentiating. Thus,

$$
\mathfrak{k}=\mathfrak{s p}_{2}^{*}=\left\{\left(\begin{array}{rr}
a & b \\
-\bar{b} & d
\end{array}\right): \bar{a}=-a, \bar{d}=-d, a, b, d \in \mathbb{H}\right\}
$$

The computation of $\Omega$ on $G / K$ can be simplified by using the right $K$-invariance, which entails annihilation by $\mathfrak{k}$. The basis elements $H_{i}, H_{j}, H_{k}$ and $H_{i}^{\prime}, H_{j}^{\prime}, H_{k}^{\prime}$ are in $\mathfrak{k}$, as are $X-Y$ and $X_{i}+Y_{i}, X_{j}+Y_{j}$, and $X_{k}+Y_{k}$. Rearrange $\Omega$ anticipating application to right $K$-invariant $f$ on $G$ :

$$
\begin{gathered}
\Omega=\frac{1}{2} H^{2}+X Y+Y X+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}-H_{\ell}^{\prime} H_{\ell}^{\prime}+X_{\ell} Y_{\ell}+Y_{\ell} X_{\ell}\right) \\
=\frac{1}{2} H^{2}+2 X Y+Y X-X Y+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}-H_{\ell}^{\prime} H_{\ell}^{\prime}+2 X_{\ell} Y_{\ell}-X_{\ell} Y_{\ell}+Y_{\ell} X_{\ell}\right) \\
=\frac{1}{2} H^{2}+2 X Y-H+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}-H_{\ell}^{\prime} H_{\ell}^{\prime}+2 X_{\ell} Y_{\ell}-H\right) \\
=\frac{1}{2} H^{2}-4 H+2 X^{2}+2 X(-X+Y)+\sum_{\ell=i, j, k}\left(-\frac{1}{2} H_{\ell}^{2}-H_{\ell}^{\prime} H_{\ell}^{\prime}-2 X_{\ell}^{2}+2 X_{\ell}\left(X_{\ell}+Y_{\ell}\right)\right)
\end{gathered}
$$

Since the elements $H_{\ell}, H_{\ell}^{\prime}, X-Y$, and $X_{\ell}+Y_{\ell}$ annihilate right $K$-invariant functions $f$, this gives

$$
\Omega f=\left(\frac{1}{2} H^{2}-4 H-\sum_{\ell=i, j, k} 2 X_{\ell}^{2}\right) f
$$

## 4. Invariant differential operators

The individual terms are computed as in the previous three cases. For example,
$X^{2}-X_{i}^{2}-X_{j}^{2}-X_{k}^{2} \longrightarrow\left(y \frac{\partial}{\partial x_{1}}\right)^{2}-\left(i y \frac{\partial}{\partial x_{2}}\right)^{2}-\left(j y \frac{\partial}{\partial x_{3}}\right)^{2}-\left(k y \frac{\partial}{\partial x_{4}}\right)^{2}=y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial x_{3}^{2}}+\frac{\partial^{2}}{\partial x_{4}^{2}}\right)$
and

$$
\frac{1}{2} H^{2}-4 H=\frac{1}{2}\left(2 y \frac{\partial}{\partial y}\right)^{2}-4\left(2 y \frac{\partial}{\partial y}\right)=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}+2 y \frac{\partial}{\partial y}-8 y \frac{\partial}{\partial y}=2 y^{2}\left(\frac{\partial}{\partial y}\right)^{2}-6 y \frac{\partial}{\partial y}
$$

Altogether, on right $K$-invariant functions $f$, discarding the irrelevant leading constant 2,

$$
\Omega \longrightarrow y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\frac{\partial^{2}}{\partial x_{2}^{2}}+\frac{\partial^{2}}{\partial x_{3}^{2}}+\frac{\partial^{2}}{\partial x_{4}^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-3 y \frac{\partial}{\partial y}
$$

Again, as with $S L_{2}(\mathbb{C})$ and $S p_{1,1}^{*}$, in contrast to $S L_{2}(\mathbb{R})$, this is not merely a multiple of the Euclidean Laplacian.

## 4.A Appendix: brackets

Here we prove the basic result about intrinsic derivatives. Let $G$ act on itself by right translations, and on functions on $G$ by $(g \cdot f)(h)=f(h g)$, for $g, h \in G$. For $x \in \mathfrak{g}$, the corresponding differential operator $X_{x}$ on smooth functions $f$ on $G$ is $\left(X_{x} f\right)(h)=\left.\frac{d}{d t}\right|_{t=0} f\left(h \cdot e^{t x}\right)$.
[4.A.1] Theorem: $\left[X_{x}, X_{y}\right]=X_{x, y}$, that is, $X_{x} \circ X_{y}-X_{y} \circ, X_{x}=X_{[x, y]}$ for $x, y \in \mathfrak{g}$.
Proof: First, re-characterize the Lie algebra $\mathfrak{g}$ less formulaically. The tangent space $T_{m} M$ to a smooth manifold $M$ at a point $m \in M$ is intended to be the collection of first-order (homogeneous) differential operators, on functions near $m$, followed by evaluation of the resulting functions at the point $m$.

One way to make the description of the tangent space precise is as follows. Let $\mathcal{O}$ be the ring of germs ${ }^{[45]}$ of smooth functions at $m$. Let $e_{m}: f \rightarrow f(m)$ be the evaluation-at- $m$ map $\mathcal{O} \rightarrow \mathbb{R}$ on (germs of) functions in $\mathcal{O}$. Since evaluation is a ring homomorphism, (and $\mathbb{R}$ is a field) the kernel $\mathfrak{m}$ of $e_{m}$ is a maximal ideal in $\mathcal{O}$. A first-order homogeneous differential operator $D$ might be characterized by the Leibniz rule

$$
D(f \cdot F)=D f \cdot F+f \cdot D F
$$

Then $e_{m} \circ D$ vanishes on $\mathfrak{m}^{2}$, since

$$
\left(e_{m} \circ D\right)(f \cdot F)=f(m) \cdot D F(m)+D f(m) \cdot F(m)=0 \cdot D F(m)+D f(m) \cdot 0=0 \quad(\text { for } f, F \in \mathfrak{m})
$$

Thus, $D$ gives a linear functional on $\mathfrak{m}$ that factors through $\mathfrak{m} / \mathfrak{m}^{2}$. Define

$$
\text { tangent space to } M \text { at } m=T_{m} M=\left(\mathfrak{m} / \mathfrak{m}^{2}\right)^{*}=\operatorname{Hom}_{\mathbb{R}}\left(\mathfrak{m} / \mathfrak{m}^{2}, \mathbb{R}\right)
$$

To see that we have included exactly what we want, and nothing more, use the defining fact (for manifold) that $m$ has a neighborhood $U$ and a homeomorphism-to-image $\varphi: U \rightarrow \mathbb{R}^{n}$. [46] The precise definition of

[^13]$$
\varphi_{i} \circ \varphi_{j}^{-1}: \varphi_{j}\left(U_{i} \cap U_{j}\right) \longrightarrow U_{i} \cap U_{j} \longrightarrow \varphi_{i}\left(U_{i} \cap U_{j}\right)
$$
is a smooth map from the subset $\varphi_{j}\left(U_{i} \cap U_{j}\right)$ of $\mathbb{R}^{n}$ to the subset $\varphi_{i}\left(U_{i} \cap U_{j}\right)$.

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smoothness of a function $f$ near $m$ is that $f \circ \varphi^{-1}$ be smooth on some subset of $\varphi(U)$. [47] In brief, the nature of $\mathfrak{m} / \mathfrak{m}^{2}$ and $\left(\mathfrak{m} / \mathfrak{m}^{2}\right)^{*}$ can be immediately transported to an open subset of $\mathbb{R}^{n}$. From Maclaurin-Taylor expansions, the pairing

$$
v \times f \longrightarrow(\nabla f)(m) \cdot v \quad\left(\text { for } v \in \mathbb{R}^{n} \text { and } f \text { smooth at } m \in \mathbb{R}^{n}\right)
$$

induces an isomorphism $\mathbb{R}^{n} \rightarrow\left(\mathfrak{m} / \mathfrak{m}^{2}\right)^{*}$. Thus, $\left(\mathfrak{m} / \mathfrak{m}^{2}\right)^{*}$ is a good notion of tangent space.
[4.A.2] Claim: The Lie algebra $\mathfrak{g}$ of $G$ is naturally identifiable with the tangent space to $G$ at 1 , via

$$
x \times\left. f \longrightarrow \frac{d}{d t}\right|_{t=0} f\left(e^{t x}\right) \quad(\text { for } x \in \mathfrak{g} \text { and } f \text { smooth near } 1)
$$

Proof: The exponential map is a diffeomorphism of the Lie algebra $\mathfrak{g}$ to its image, and the image is a neighborhood of the identity in $G$. For linear Lie groups, the invertibility is immediate from existence of an explicit local inverse to the exponential near 1 , given by the usual logarithm.

Left translation action of $G$ on functions on $G$ is $\left(L_{g} f\right)(h)=f\left(g^{-1} h\right)$ with $g, h \in G$, with the inverse for associativity, as usual.
[4.A.3] Claim: The map $x \longrightarrow X_{x}$ gives an $\mathbb{R}$-linear isomorphism

$$
\mathfrak{g} \longrightarrow \text { left } G \text {-invariant vector fields on } G
$$

Proof: (of claim) On one hand, since the action of $x$ is on the right, it is not surprising that $X_{x}$ is invariant under the left action of $G$, namely

$$
\left(X_{x} \circ L_{g}\right) f(h)=X_{x} f\left(g^{-1} h\right)=\left.\frac{d}{d t}\right|_{t=0} f\left(g^{-1} h e^{t x}\right)=\left.L_{g} \frac{d}{d t}\right|_{t=0} f\left(h e^{t x}\right)=\left(L_{g} \circ X_{x}\right) f(h)
$$

On the other hand, for a left-invariant vector field $X$,

$$
(X f)(h)=\left(L_{h}^{-1} \circ X\right) f(1)=\left(X \circ L_{h}^{-1}\right) f(1)=X\left(L_{h}^{-1} f\right)(1)
$$

That is, $X$ is completely determined by what it does to functions at 1.
Let $\mathfrak{m}$ be the maximal ideal of functions vanishing at 1 , in the ring $\mathcal{O}$ of germs of smooth functions at 1 on $G$. The first-order nature of vector fields is captured by the Leibniz rule $X(f \cdot F)=f \cdot X F+X f \cdot F$. As above, the Leibniz rule implies that $e_{1} \circ X$ vanishes on $\mathfrak{m}^{2}$. Thus, we can identify $e_{1} \circ X$ with an element of

$$
\left(\mathfrak{m} / \mathfrak{m}^{2}\right)^{*}=\operatorname{Hom}_{\mathbb{R}}\left(\mathfrak{m} / \mathfrak{m}^{2}, \mathbb{R}\right)=\text { tangent space to } G \text { at } 1=\mathfrak{g}
$$

Thus, the map $x \rightarrow X_{x}$ is an isomorphism from $\mathfrak{g}$ to left invariant vector fields, proving the claim. ///
Now use the re-characterized $\mathfrak{g}$ to prove $\left[X_{x}, X_{y}\right]=X_{z}$ for some $z \in \mathfrak{g}$. Consider $\left[X_{x}, X_{y}\right]$ for $x, y \in \mathfrak{g}$. That this differential operator is left $G$-invariant is clear, since it is a difference of composites of such. It is less clear that it satisfies Leibniz' rule (and thus is first-order). But, indeed, for any two vector fields $X, Y$,

$$
\begin{gathered}
{[X, Y](f F)=X Y(f F)-Y X(F f)=X(Y f \cdot F+f \cdot Y F)-Y(X f \cdot F+f \cdot X F)} \\
=(X Y f \cdot F+Y f \cdot X F+X f \cdot Y F+f \cdot X Y F)-(Y X f \cdot F+X f \cdot Y F+Y f \cdot X F+f \cdot Y X F) \\
=[X, Y] f \cdot F+f \cdot[X, Y] F
\end{gathered}
$$

so $[X, Y]$ does satisfy the Leibniz rule. In particular, $\left[X_{x}, X_{y}\right]$ is again a left- $G$-invariant vector field, so is of the form $\left[X_{x}, X_{y}\right]=X_{z}$ for some $z \in \mathfrak{g}$.

In fact, the relation $\left[X_{x}, X_{y}\right]=X_{z}$ is the intrinsic definition of the Lie bracket on $\mathfrak{g}$, since we could define the element $z=[x, y]$ by the relation $\left[X_{x}, X_{y}\right]=X_{[x, y]}$. However, we are burdened by having the ad hoc
[47] The well-definedness of this definition depends on the maximality property of an atlas.

## 4. Invariant differential operators

but convenient definition $[x, y]=x y-y x$ in terms of matrix multiplication. However, our assumption that $G$ is a subgroup of some $G L_{n}(\mathbb{R})$ or $G L_{n}(\mathbb{C})$ allows us to use the explicit exponential and a local logarithm inverse to it, to determine the bracket $\left[X_{x}, X_{y}\right]$ somewhat more intrinsically, as follows.

Consider linear functions on $\mathfrak{g}$, locally transported to $G$ via locally inverting the exponential near $1 \in G$. Thus, for $\lambda \in \mathfrak{g}^{*}$, near $1 \in G$, define

$$
f\left(e^{x}\right)=\lambda(x)
$$

Then

$$
\left.\left[X_{x}, X_{y}\right] f(1)=\left.\left.\frac{d}{d t}\right|_{t=0} \frac{d}{d s}\right|_{s=0}\left(\lambda\left(\log \left(e^{s x} e^{t y}\right)\right)-\lambda\left(\log \left(e^{t y} e^{s x}\right)\right)\right)\right)
$$

Dropping $O\left(s^{2}\right)$ and $O\left(t^{2}\right)$ terms, this is

$$
\begin{gathered}
=\left.\left.\frac{d}{d t}\right|_{t=0} \frac{d}{d s}\right|_{s=0}(\lambda(\log (1+s x)(1+t y))-\lambda(\log (1+t y)(1+s x))) \\
=\left.\left.\frac{d}{d t}\right|_{t=0} \frac{d}{d s}\right|_{s=0} \lambda(\log (1+s x+t y+s t x y)-\log (1+t y+s x+s t y x)) \\
=\left.\left.\frac{d}{d t}\right|_{t=0} \frac{d}{d s}\right|_{s=0} \lambda\left(\left(s x+t y+s t x y-\frac{1}{2}(s x+t y)^{2}\right)-\left(t y+s x+s t y x-\frac{1}{2}(t y+s x)^{2}\right)\right) \\
=\left.\left.\frac{d}{d t}\right|_{t=0} \frac{d}{d s}\right|_{s=0} \lambda\left(\left(s t x y-\frac{1}{2} s t x y-\frac{1}{2} s t y x\right)-\left(s t y x-\frac{1}{2} s t x y-\frac{1}{2} s t y x\right)\right) \\
=\left.\left.\frac{d}{d t}\right|_{t=0} \frac{d}{d s}\right|_{s=0} s t \cdot \lambda(x y-y x)=\lambda(x y-y x)
\end{gathered}
$$

where the multiplication and commutator $x y-y x$ is in the ring of matrices. Thus, since $\mathfrak{g}^{*}$ separates points on $\mathfrak{g}$, we have the equality $\left[X_{x}, X_{y}\right]=X_{[x, y]}$. with the ad hoc definition of $[x, y]$.

## 4.B Appendix: existence and uniqueness

The characterization of the Lie algebra of a subgroup $G$ of $G L_{n}(F)$ with $F=\mathbb{R}, \mathbb{C}$, or $\mathbb{H}$, as $\mathfrak{g}=\{x$ : $e^{t x} \in G$, for all $\left.t \in \mathbb{R}\right\}$ produces a set $\mathfrak{g}$ closed under scalar multiplication, but not obviously closed under addition, although in the explicit examples above this latter property is clear.

One way to prove closure under addition is verification that the matrix exponentiation characterization of $\mathfrak{g}$ really does produce the tangent space to $G$ at 1 , since the tangent space is a vectorspace. For matrix $x$ such that $e^{t x} \in G$ for small real $t$, the map $t \rightarrow e^{t x}$ is a curve inside $G$, and in an extrinsic-geometry sense $x$ is a tangent vector to $G$ at 1 . The converse is not as elementary, namely, given a tangent vector $x$ to $G$ at 1 , show that $e^{t x} \in G$ for all $t \in \mathbb{R}$.

The curve $u(t)=e^{t x}$ certainly lies in $G L_{n}(F)$, and satisfies the differential equation $\frac{d}{d t} u=u \cdot x$ with initial condition $u(0)=1_{n}$. This differential equation can be viewed as a differential equation on $G$, and also as a differential equation on $G L_{n}(F)$, and $u(t)=e^{t x}$ a solution in $G L_{n}(F)$. To prove that in fact $e^{t x} \in G$ for small $t$ would follow from uniqueness of solutions to this differential equation on $G L_{n}(F)$, and from existence of a solution (for small $t$ ) on $G$. Then because $G$ is a group, $e^{t x} \in G$ for all $t \in \mathbb{R}$.
[4.B.1] Theorem: For smooth $F$ on an open subset $\Omega$ of $\mathbb{R}^{2}$, and for $\left(x_{o}, y_{o}\right) \in \Omega$, the equation $\frac{d f}{d x}=F(x, f(x))$ has a unique differentiable solution $f$ on a neighborhood of $\left(x_{o}, y_{o}\right)$ with $f\left(x_{o}\right)=y_{o}$. This solution is smooth.
Proof: Picard iteration converts the differential equation to an equivalent integral equation to prove existence. Uniqueness of fixed points of contractive mappings proves uniqueness. Assuming $f^{\prime}=d f / d x$ exists as a pointwise-valued function and $f$ is continuous, the relation $d f / d x=F(x, f(x))$ shows that $f^{\prime}$ is continuous. Thus, by the fundamental theorem of calculus,

$$
f(x)=f\left(x_{o}\right)+\int_{x_{o}}^{x} f^{\prime}(t) d t=y_{o}+\int_{x_{o}}^{x} F(t, f(t)) d t
$$

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That is, $f$ satisfies the integral equation

$$
f(x)=y_{o}+\int_{x_{o}}^{x} F(t, f(t)) d t
$$

Conversely, for continuous $f$ satisfying this integral equation, by the fundamental theorem of calculus $f$ is differentiable and

$$
f^{\prime}(x)=F(x, f(x)) \quad\left(\text { and } f \text { satisfies } f\left(x_{o}\right)=y_{o}\right)
$$

Thus, for continuous $f$, the integral equation is equivalent to the differential equation and initial value. Without loss of generality, $x_{o}=y_{o}=0$. Picard's iteration scheme is to take $f_{o}(x)=0$, and iterate:

$$
f_{n+1}(x)=\int_{0}^{x} F\left(t, f_{n}(t)\right) d t
$$

These are continuous functions. The claim is that, on a sufficiently small neighborhood of $x=0$, these $f_{n}$ approach a solution to the integral equation on that interval, proving existence.

We should check that, with $x$ restricted to a small-enough interval $|x| \leq \delta$, the pairs $\left(x, f_{n}(x)\right)$ stay inside $\Omega$. By smoothness of $F$, given a finite rectangle

$$
R=\{|x| \leq \delta,|y| \leq \eta\} \subset \Omega
$$

there is a constant $B$ such that $|F(x, y)| \leq B$ for all $(x, y) \in R$. Shrink $\delta$ so that $0<\delta<B^{-1}$. Assuming the pairs $\left(x, f_{n}(x)\right)$ are inside $R$,

$$
\left|f_{n+1}(x)\right| \leq \int_{0}^{x}\left|F\left(t, f_{n}(t)\right)\right| d t \leq \int_{0}^{x} B d t \leq \delta \cdot B
$$

Thus, further shrinking $\delta$ so that $\delta \cdot B \leq \eta$, the restriction $|x| \leq \delta$ produces functions $f_{n}$ with $\left(x, f_{n}(x)\right)$ staying inside $R \subset \Omega$ of $F$.

We show that, possibly further shrinking $\delta$, for $|x| \leq \delta$ the sequence of functions $f_{n}$ converges in sup norm to a solution of the integral equation. The natural estimate succeeds: first,

$$
\sup _{|x| \leq \delta}\left|f_{n+1}(x)-f_{n}(x)\right| \leq \int_{0}^{\delta}\left|F\left(t, f_{n}(t)\right)-F\left(t, f_{n-1}(t)\right)\right| d t
$$

Since $F$ is smooth, for a fixed compact-closure neighborhood $U$ of $(0,0)$ there is a constant $C$ such that

$$
\left|F(x, y)-F\left(x, y^{\prime}\right)\right| \leq C \cdot\left|y-y^{\prime}\right| \quad\left(\text { for }(x, y) \in U \text { and }\left(x, y^{\prime}\right) \in U\right)
$$

Thus,

$$
\sup _{|x| \leq \delta}\left|f_{n+1}(x)-f_{n}(x)\right| \leq|\delta| \cdot C \cdot \sup _{|x| \leq \delta}\left|f_{n}(x)-f_{n-1}(x)\right|
$$

Shrinking $\delta$ so that $\delta \cdot C \leq \frac{1}{2}$, for example, gives convergence in sup norm to a continuous function $f$. This further shrinking of $\delta$ occurs just once, not for each $n$.

To show that $f$ is a solution of the integral equation, given $\varepsilon>0$, take $N$ large enough so that the sup norm of $f_{m}-f_{n}$ is less than $\varepsilon$ for all $m, n \geq N$. Then the sup norm of $f-f_{N+1}$ is $\leq \varepsilon$, and

$$
\begin{gathered}
\left|f(x)-\int_{0}^{x} F(t, f(t)) d t\right| \leq\left|f(x)-f_{N+1}(x)\right|+\int_{0}^{x}\left|F\left(t, f_{N}(t)\right)-F(t, f(t))\right| d t \\
\leq \varepsilon+\int_{0}^{x} C \cdot\left|f_{N}(t)-f(t)\right| d t \leq \varepsilon+\delta \cdot C \cdot \varepsilon
\end{gathered}
$$

This holds for every $\varepsilon>0$, so we have equality, proving $f$ is a solution of the integral equation.

## 4. Invariant differential operators

Toward uniqueness, we showed above that the mapping

$$
T g(x)=\int_{0}^{x} F(t, g(t)) d t
$$

maps continuous functions $g$ on $|x| \leq \delta$ with bounds $|g(x)| \leq \eta$ to continuous functions with the same bound, for sufficiently small $\delta$ and $\eta$. Let $X$ be the set of such functions. With metric given by sup norm, $X$ is complete. Shrinking $\delta$ if necessary, as above we have the contractive mapping property

$$
\sup _{x}|T g(x)-T h(x)| \leq \frac{1}{2} \cdot \sup _{x}|g(x)-h(x)|
$$

Given two solutions $g, h$ to the integral equation,

$$
\sup _{x}|g(x)-h(x)| \leq \sup _{x}|T g(x)-T h(x)| \leq \frac{1}{2} \cdot \sup _{x}|g(x)-h(x)|
$$

proving that $g(x)=h(x)$, giving uniqueness of solution to the integral equation.
Smoothness follows from the differential equation: granting $f \in C^{k}$, the relation $d f / d x=F(x, f(x))$ exhibits the derivative as $C^{k}$. By induction, $f \in C^{\infty}$.

## 5. Integration on quotients

1. Surjectivity of averaging maps
2. Invariant measures and integrals on quotients $H \backslash G$

Appendix A: apocryphal lemma $X \approx G / G_{x}$
Appendix B: topology on quotients $H \backslash G$ or $G / H$

The simplest case of unwinding is for $f \in C_{c}^{o}(\mathbb{R})$ :

$$
\int_{\mathbb{R} / \mathbb{Z}}\left(\sum_{n \in \mathbb{Z}} f(x+n)\right) d x=\int_{\mathbb{R}} f(x) d x
$$

In fact, the integral on the quotient $\mathbb{R} / \mathbb{Z}$ is unequivocally characterized, ${ }^{[48]}$ by this relation, once we know that the averaged functions $\sum_{n} f(x+n)$ are at least dense in $C^{o}(\mathbb{R} / \mathbb{Z})$. As corollary, for $F \in C^{o}(\mathbb{R} / \mathbb{Z})$, since $F \cdot f \in C_{c}^{o}(\mathbb{R})$,

$$
\begin{gathered}
\int_{\mathbb{R} / \mathbb{Z}} F(x)\left(\sum_{n \in \mathbb{Z}} f(x+n)\right) d x=\int_{\mathbb{R} / \mathbb{Z}}\left(\sum_{n \in \mathbb{Z}} F(x) f(x+n)\right) d x \\
\quad=\int_{\mathbb{R} / \mathbb{Z}}\left(\sum_{n \in \mathbb{Z}} F(x+n) f(x+n)\right) d x=\int_{\mathbb{R}} F(x) f(x) d x
\end{gathered}
$$

We need analogous assertions with less elementary group actions and less transparent representatives for the quotients. For example, with $G, K, \Gamma$ as in our examples, integration on $\Gamma \backslash G$ is characterized by requiring, for all $f \in C_{c}^{o}(G)$,

$$
\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma} f(\gamma g)\right) d g=\int_{G} f(g) d g
$$

once we know that the averages $\sum_{\gamma \in \Gamma} f(\gamma z)$ are at least dense in $C_{c}^{o}(\Gamma \backslash G)$. In fact, such averaging maps are universally surjective on compactly-supported continuous functions, as demonstrated below.

An important variant ${ }^{[49]}$ uses $f \in C_{c}^{o}\left(\Gamma_{\infty} \backslash G\right)$ for a subgroup $\Gamma_{\infty}$ of $\Gamma$, by the surjectivity of averaging maps, take $\varphi \in C_{c}^{o}(G)$ such that

$$
\sum_{\beta \in \Gamma_{\infty}} \varphi \circ \beta=f
$$

so then

$$
\begin{gathered}
\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} f \circ \gamma\right)=\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}\left(\sum_{\beta \in \Gamma_{\infty}}(\varphi \circ \beta) \circ \gamma\right)\right)=\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma} \varphi \circ \gamma\right) \\
=\int_{G} \varphi=\int_{\Gamma_{\infty} \backslash G}\left(\sum_{\beta \in \Gamma_{\infty}} \varphi \circ \beta\right)=\int_{\Gamma_{\infty} \backslash G} f
\end{gathered}
$$

The corollary with $F \in C^{o}(\Gamma \backslash G)$ and $f \in C_{c}^{o}\left(\Gamma_{\infty} \backslash G\right)$ is

$$
\int_{\Gamma \backslash G} F \cdot\left(\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} f \circ \gamma\right)=\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}(F \cdot f) \circ \gamma\right)=\int_{\Gamma_{\infty} \backslash G} F \cdot f
$$

[48] The Riesz-Markov-Kakutani theorem asserts that every (continuous) functional on compactly-supported continuous functions on a reasonable topological space $X$ is $f \rightarrow \int_{X} f(x) d \mu(x)$ for some measure $\mu$. Relying on this, specification of a functional (integration) on $C_{c}^{o}(X)$ specifies a measure. In fact, we care more about the integral than about the measure.
[49] This variant of unwinding arose most prominently in the Rankin-Selberg method, where $\int_{\Gamma \backslash \mathfrak{H}}|f|^{2} \cdot E_{s}$ for cuspform $f$ and Eisenstein series $E_{s}$ is unwound using the definition of $E_{s}$ as wound up from $y^{s}$. This theme is pervasive in the theory of automorphic forms.

### 5.1 Surjectivity of averaging maps

By convention, a topological group is a locally compact, Hausdorff topological space $G$ with a continuous group operation $G \times G \rightarrow G$, and continuous inversion map $g \rightarrow g^{-1}$. To avoid pathologies with regard to measures on products, we require that topological groups have a countable basis.

Let $d g$ be a right $G$-invariant measure on $G$, meaning that for $f \in C_{c}^{o}(G)$

$$
\left.\int_{G} f(g h) d g=\int_{G} f(g) d\left(g h^{-1}\right)=\int_{G} f(g) d g \quad \quad \text { (for all } h \in G\right)
$$

Let $\delta_{G}: G \rightarrow(0,+\infty)$ be the modular function of $G$, gauging the discrepancy between left and right invariant measures, in the sense that meas $(g E)=\delta_{G}(g) \cdot$ meas $(E)$ for a measurable set $E \subset G$. It is immediate that $\delta_{G}$ is a group homomorphism $\delta_{G}: G \rightarrow(0,+\infty)$. It is continuous. And $\delta_{G}^{-1}(g) d g$ is a left invariant measure. A group with $\delta=1$ is unimodular.
[5.1.1] Claim: Every discrete, compact, or abelian group is unimodular.
Proof: For an abelian group, $d(h g)=d(g h)$. For a discrete group, one uses counting measure, which is both left and right invariant. For a compact groups, observe that the modular function is continuous, and has image in the multiplicative group of positive reals. The only compact subgroup is $\{1\}$.

Let $H$ be a closed subgroup of $G$, with right $H$-invariant measure $d h$. The quotient $H \backslash G$ has the quotient topology [5.B]. See also [5A].
[5.1.2] Claim: The averaging map $\alpha: C_{c}^{o}(G) \rightarrow C_{c}^{o}(H \backslash G)$ by

$$
(\alpha F)(g)=\int_{H} F(h g) d h \quad\left(\text { for } F \in C_{c}^{o}(G)\right)
$$

is surjective.
Proof: By design, the image consists of left $H$-invariant functions on $G$ :

$$
\left.\left.(\alpha F)(h g)=\int_{H} F\left(h^{\prime}(h g)\right) d h^{\prime}=\int_{H} F\left(\left(h^{\prime} h\right) g\right)\right) d h^{\prime}=\int_{H} F\left(h^{\prime} g\right)\right) d h^{\prime}
$$

by replacing $h^{\prime}$ by $h^{\prime} h^{-1}$. The surjectivity is much less trivial. Let $q: G \rightarrow H \backslash G$ be the quotient map. Let $U$ be a neighborhood of $1 \in G$ having compact closure $\bar{U}$. For each $g \in G, g U$ is a neighborhood of $g$. The images $q(g U)$ are open, by the characterization of the quotient topology. Given $f \in C_{c}^{o}(H \backslash G)$, the support $\operatorname{spt}(f)$ of $f$ is covered by the opens $q(g U)$, and admits a finite subcover $q\left(g_{1} U\right), \ldots, q\left(g_{n} U\right)$. The set

$$
C=\operatorname{spt}(f \circ q) \cap\left(g_{1} \bar{U} \cup \ldots \cup g_{n} \bar{U}\right) \subset G
$$

is compact, and $q(C)=\operatorname{spt}(f) \subset H \backslash G$. By Urysohn's lemma [9.E.2], let $\varphi \in C_{c}^{o}(G)$ be identically 1 on $C$, and non-negative real-valued everywhere. Let $F=\varphi \cdot(f \circ q)$. Since $\alpha \varphi$ is strictly positive on a neighborhood of the (compact) support of $F$, the quotient $F / \alpha \varphi$ is in $C_{c}^{o}(G)$. Since $f$ is already left $H$-invariant,

$$
\alpha F(g)=\int_{H} \varphi(h g) \cdot f(h g) d h=\int_{H} \varphi(h g) \cdot f(g) d h=\int_{H} \varphi(h g) d h \cdot f(g)=\alpha \varphi(g) \cdot f(g)
$$

Because $f$ and $\alpha \varphi$ are left $H$-invariant,
$\alpha\left(\frac{F}{\alpha \varphi}\right)(g)=\int_{H} \frac{\varphi(h g) \cdot f(h g)}{\alpha \varphi(h g)} d h=\int_{H} \frac{\varphi(h g) \cdot f(g)}{\alpha \varphi(g)} d h=\int_{H} \varphi(h g) d h \cdot \frac{f(g)}{\alpha \varphi(g)}=\alpha \varphi(g) \cdot \frac{f(g)}{\alpha \varphi(g)}=f$
giving the surjectivity.

### 5.2 Invariant measures and integrals on quotients $H \backslash G$

Temporarily, for clarity in the proofs of this section, we may let $\dot{g}$ denote an element of the quotient $H \backslash G$, where $H$ is a closed subgroup of $G$. Let $H$ have modular function $\delta_{H}$.
[5.2.1] Theorem: The quotient $H \backslash G$ has a right $G$-invariant measure if and only if $\left.\delta_{G}\right|_{H}=\delta_{H}$. In that case, the integral is unique up to scalars, and is characterized as follows. For given right Haar measure $d h$ on $H$ and for given right Haar measure $d g$ on $G$ there is a unique invariant measure $d \dot{g}$ on $H \backslash G$ such that for $f \in C_{c}^{o}(G)$

$$
\int_{H \backslash G}\left(\int_{H} f(h \dot{g}) d h\right) d \dot{g}=\int_{G} f(g) d g \quad\left(\text { for } f \in C_{c}^{o}(G)\right)
$$

Proof: First, prove the necessity of the condition on the modular functions. Suppose that there is a right $G$-invariant measure on $H \backslash G$. Let $\alpha$ be the averaging map $f \rightarrow \int_{H} f(h g) d h$. For $f \in C_{c}^{o}(G)$ the map

$$
f \longrightarrow \int_{H \backslash G} \alpha f(\dot{g}) d \dot{g}
$$

emphasizing the coordinate $\dot{g}$ on the quotient, is a right $G$-invariant functional (with the continuity property as above), so by uniqueness of right invariant measure on $G$ must be a constant multiple of the Haar integral

$$
f \longrightarrow \int_{G} f(g) d g
$$

The averaging map behaves in a straightforward manner under left translation $L_{h} f(g)=f\left(h^{-1} g\right)$ for $h \in H$ : for $f \in C_{c}^{o}(G)$ and for $h \in H$

$$
\alpha\left(L_{h} f\right)(g)=\int_{H} f\left(h^{-1} x g\right) d x=\delta_{H}(h) \int_{H} f(x g) d x
$$

by replacing $x$ by $h x$. Then

$$
\int_{G} f(g) d g=\int_{H \backslash G} \alpha(f)(g) d \dot{g}=\delta(h)^{-1} \int_{H \backslash G} \alpha\left(L_{h} f\right)(g) d \dot{g}=\delta(h)^{-1} \int_{G} f\left(h^{-1} g\right) d g
$$

by comparing the iterated integral to the single integral. Replacing $g$ by $h g$ in the integral gives

$$
\int_{G} f(g) d g=\delta(h)^{-1} \delta_{G}(h) \int_{G} f(g) d g
$$

Choosing $f$ such that the integral is not 0 implies the stated condition on the modular functions.
Proof of sufficiency starts from existence of Haar measures on $G$ and on $H$. For simplicity, first suppose that both groups are unimodular. As expected, attempt to define an integral on $C_{c}^{o}(H \backslash G)$ by

$$
\int_{H \backslash G} \alpha f(\dot{g}) d \dot{g}=\int_{G} f(g) d g
$$

invoking the fact that the averaging map $\alpha$ from $C_{c}^{o}(G)$ to $C_{c}^{o}(H \backslash G)$ is surjective. The potential problem is well-definedness. It suffices to prove that $\int_{G} f(g) d g=0$ for $\alpha f=0$. Indeed, for $\alpha f=0$, for all $F \in C_{c}^{o}(G)$, the integral of $F$ against $\alpha f$ is certainly 0 . Rearrange

$$
0=\int_{G} F(g) \alpha f(g) d g=\int_{G} \int_{H} F(g) f(h g) d h d g=\int_{H} \int_{G} F\left(h^{-1} g\right) f(g) d g d h
$$

by replacing $g$ by $h^{-1} g$. Replace $h$ by $h^{-1}$, so

$$
0=\int_{G} \alpha F(g) f(g) d g
$$

## 5. Integration on quotients

Surjectivity of $\alpha$ shows that $F$ can be chosen so that $\alpha F$ is identically 1 on the support of $f$. Then the integral of $f$ is 0 , as claimed, proving the well-definedness for unimodular $H$ and $G$.

For not-necessarily-unimodular $H$ and $G$, in the previous argument the left translation by $h^{-1}$ produces a factor of $\delta_{G}\left(h^{-1}\right)$. Then replacing $h$ by $h^{-1}$ converts right Haar measure to left Haar measure, so produces a factor of $\delta_{H}(h)^{-1}$, and the other factor becomes $\delta_{G}(h)$. If $\delta_{G}(h) \cdot \delta_{H}(h)^{-1}=1$, then the product of these two factors is 1 , and the same argument goes through, proving well-definedness.
[5.2.2] Corollary: Let $G$ be a group, with closed subgroups $\Theta \subset H \subset G$. Suppose that $\left.\delta_{G}\right|_{H}=\delta_{H}$ and $\left.\delta_{H}\right|_{\Theta}=\delta_{\Theta}$. Given any two of: right $H$-invariant measure on $\Theta \backslash H$, right $G$-invariant measure on $\Theta \backslash G$, or right $G$-invariant measure on $H \backslash G$, the other one of the three is uniquely determined so that, for all $f \in C_{c}^{o}(\Theta \backslash G)$,

$$
\int_{\Theta \backslash G} f(g) d g=\int_{H \backslash G}\left(\int_{\Theta \backslash H} f(h g) d h\right) d g
$$

Proof: With the surjectivity of averaging maps in hand, this proof is just an iterative application of the previous theorem. Namely, given $f \in C_{c}^{o}(\Theta \backslash G)$, let $F \in C_{c}^{o}(G)$ map to $f$ by the averaging-over- $\Theta$ map, so, by the theorem,

$$
\int_{\Theta \backslash G} f=\int_{\Theta \backslash G} \int_{\Theta} F(\theta g) d \theta d g=\int_{G} F
$$

Again by the theorem,

$$
\int_{G} F=\int_{H \backslash G} \int_{H} F(h g) d h d g
$$

At the same time, applying the theorem to $h \rightarrow F(h g)$,

$$
\int_{H} F(h g) d h=\int_{\Theta \backslash H} \int_{\Theta} F(\theta h g) d \theta d h
$$

The inner integral in the latter is $f(h g)$, giving the claim.

## 5.A Appendix: apocryphal lemma $X \approx G / G_{x}$

We prove that under mild hypotheses a topological space $X$ acted upon transitively by a topological group $G$ is homeomorphic to the quotient $G / G_{x}$, where $G_{x}$ is the isotropy group of a chosen point $x$ in $X$. Ignoring the topology, the bijection $G / G_{x} \approx X$ by $g \cdot G_{x} \leftrightarrow g x$ is easy to see. In contrast, the topological aspects are not trivial, but are very general.
[5.A.1] Proposition: Let $G$ be a locally compact, Hausdorff topological group and $X$ a locally compact Hausdorff topological space with a continuous transitive action of $G$ upon $X$. Suppose that $G$ has a countable basis. Fix any $x \in X$, and let $G_{x}$ be the isotropy group $G_{x}=\{g \in G: g x=x\}$. Then we have a homeomorphism $G / G_{x} \longrightarrow X$ given by the natural $g G_{x} \longrightarrow g x$.
Proof: A little systematic development of topological groups will allow a coherent argument.
[5.A.2] Claim: In a locally compact Hausdorff space $X$, given an open neighborhood $U$ of a point $x$, there is a neighborhood $V$ of $x$ with compact closure $\bar{V}$ and $\bar{V} \subset U$.
Proof: By local compactness, $x$ has a neighborhood $W$ with compact closure. Intersect $U$ with $W$ if necessary so that $U$ has compact closure $\bar{U}$. Note that the compactness of $\bar{U}$ implies that the boundary $\partial U$ of $U$ is compact. Using the Hausdorff-ness, for each $y \in \partial U$ let $W_{y}$ be an open neighborhood of $y$ and $V_{y}$ an open neighborhood of $x$ such that $W_{y} \cap V_{y}=\phi$. By compactness of $\partial U$, there is a finite list $y_{1}, \ldots, y_{n}$ of points on $\partial U$ such that the sets $U_{y_{i}}$ cover $\partial U$. Then $V=\bigcap_{i} V_{y_{i}}$ is open and contains $x$. Its closure is contained in $\bar{U}$ and in the complement of the open set $\bigcup_{i} W_{y_{i}}$, the latter containing $\partial U$. Thus, the closure $\bar{V}$ of $V$ is contained in $U$.
[5.A.3] Claim: The map $g G_{x} \rightarrow g x$ is a continuous bijection of $G / G_{x}$ to $X$.
Proof: First, $G \times X \rightarrow X$ by $g \times y \rightarrow g y$ is continuous by definition of the continuity of the action. Thus, with fixed $x \in X$, the restriction to $G \times\{x\} \rightarrow X$ is still continuous, so $G \rightarrow X$ by $g \rightarrow g x$ is continuous.

## Garrett: Modern Analysis of Automorphic Forms

The quotient topology on $G / G_{x}$ is the unique topology on the set (of cosets) $G / G_{x}$ such that any continuous $G \rightarrow Z$ constant on $G_{x}$ cosets factors through the quotient map $G \rightarrow G / G_{x}$. That is, we have a commutative diagram


Thus, the induced map $G / G_{x} \rightarrow X$ by $g G_{x} \rightarrow g x$ is continuous.
We need to show that $g G_{x} \rightarrow g x$ is open to prove that it is a homeomorphism.
[5.A.4] Claim: For a given point $g \in G$, every neighborhood of $g$ is of the form $g V$ for some neighborhood $V$ of 1.
Proof: First, again, $G \times G \rightarrow G$ by $g \times g \rightarrow g h$ is continuous, by assumption. Then, for fixed $g \in G$, the $\operatorname{map} h \rightarrow g h$ is continuous on $G$, by restriction. And this map has a continuous inverse $h \rightarrow g^{-1} h$. Thus, $h \rightarrow g h$ is a homeomorphism of $G$ to itself. In particular, since $1 \rightarrow g \cdot 1=g$, neighborhoods of 1 are carried to neighborhoods of $g$, as claimed.
[5.A.5] Claim: Given an open neighborhood $U$ of 1 in $G$, there is an open neighborhood $V$ of 1 such that $V^{2} \subset U$, where $V^{2}=\{g h: g, h \in V\}$.
Proof: From the continuity of multiplication $G \times G \rightarrow G$, given the neighborhood $U$ of 1 , the inverse image $W$ of $U$ under the multiplication $G \times G \rightarrow G$ is open. Since $G \times G$ has the product topology, $W$ contains an open of the form $V_{1} \times V_{2}$ for opens $V_{i}$ containing 1. With $V=V_{1} \cap V_{2}$, we have $V^{2} \subset V_{1} \cdot V_{2} \subset U$. ///

Similarly, but more simply, inversion $g \rightarrow g^{-1}$ is continuous and is its own (continuous) inverse, so the image $V^{-1}=\left\{g^{-1}: g \in V\right\}$ of an open $V$ is open. For example, given a neighborhood $V$ of 1 , replacing $V$ by $V \cap V^{-1}$ replaces $V$ by a smaller symmetric neighborhood: the new $V$ satisfies $V^{-1}=V$.

The following result is not strictly necessary, but sheds some light on the nature of topological groups. It has an analogue for topological vector spaces.
[5.A.6] Claim: The closure of $E \subset G$ is $\bigcap_{U} E \cdot U$, where $U$ runs over open neighborhoods of 1 .
Proof: A point $g \in G$ is in the closure of $E$ if and only if every neighborhood of $g$ meets $E$. That is, from just above, every set $g U$ meets $E$, for $U$ an open neighborhood of 1 . That is, $g \in E \cdot U^{-1}$ for every neighborhood $U$ of 1 . We have noted that inversion is a homeomorphism of $G$ to itself (and sends 1 to 1 ), so the map $U \rightarrow U^{-1}$ is a bijection of the collection of neighborhoods of 1 to itself. Thus, $g$ is in the closure of $E$ if and only if $g \in E \cdot U$ for every open neighborhood $U$ of 1 , as claimed.
[5.A.7] Corollary: Given a neighborhood $U$ of 1 in $G$, there is a neighborhood $V$ of 1 such that $\bar{V} \subset U$.
Proof: From the continuity of $G \times G \rightarrow G$, there is $V$ such that $V \cdot V \subset U$. From the previous claim, $\bar{V} \subset V \cdot V$, so $\bar{V} \subset V \cdot V \subset U$, as claimed.

We can improve the conclusion of the previous remark using the local compactness of $G$, as follows. Given a neighborhood $U$ of 1 in $G$, there is a neighborhood $V$ of 1 such that $\bar{V} \subset U$ and $\bar{V}$ is compact. Indeed, local compactness means exactly that there is a local basis at 1 consisting of opens with compact closures. Thus, given $V$ as in the previous remark, shrink $V$ if necessary to have the compact closure property, and still $\bar{V} \subset V \cdot V \subset U$, as claimed.
[5.A.8] Corollary: For an open subset $U$ of $G$, given $g \in U$, there is a compact neighborhood $V$ of $1 \in G$ such that $g V^{2} \subset U$.
Proof: The set $g^{-1} U$ is an open containing 1, so there is an open $W \ni 1$ such that $W^{2} \subset g^{-1} U$. Using the previous claim and remark, there is a compact neighborhood $V$ of 1 such that $V \subset W$. Then $V^{2} \subset W^{2} \subset g^{-1} U$, so $g V^{2} \subset U$ as desired.
[5.A.9] Claim: Given an open neighborhood $V$ of 1 , there is a countable list $g_{1}, g_{2}, \ldots$ of elements of $G$ such that $G=\bigcup_{i} g_{i} V$.
Proof: To see this, first let $U_{1}, U_{2}, \ldots$ be a countable basis. For $g \in G$, by definition of a basis,

$$
g V=\bigcup_{U_{i} \subset g V} U_{i}
$$

## 5. Integration on quotients

Thus, for each $g \in G$, there is an index $j(g)$ such that $g \in U_{j(g)} \subset g V$. Do note that there are only countably many such indices. For each index $i$ appearing as $j(g)$, let $g_{i}$ be an element of $G$ such that $j\left(g_{i}\right)=i$, that is,

$$
g_{i} \in U_{j\left(g_{i}\right)} \subset g_{i} \cdot V
$$

Then, for every $g \in G$ there is an index $i$ such that

$$
g \in U_{j(g)}=U_{j\left(g_{i}\right)} \subset g_{i} \cdot V
$$

This shows that the union of these $g_{i} \cdot V$ is all of $G$.
Now we can prove that $G / G_{x} \approx X$ :
Given an open set $U$ in $G$ and $g \in U$, let $V$ be a compact neighborhood of 1 such that $g V^{2} \subset U$. Let $g_{1}, g_{2}, \ldots$ be a countable set of points such that $G=\bigcup_{i} g_{i} V$. Let $W_{n}=g_{n} V x \subset X$. By the transitivity, $X=\bigcup_{i} W_{i}$.

We observed at the beginning of this discussion that $G \rightarrow X$ by $g \rightarrow g x$ is continuous, so $W_{n}$ is compact, being a continuous image of the compact set $g_{n} V$. So $W_{n}$ is closed since it is a compact subset of the Hausdorff space $X$. By the Baire category theorem [15.A] for locally compact Hausdorff spaces, some $W_{m}=g_{m} V x$ contains a non-empty open set $S$ of $X$. For $h \in V$ so that $g_{m} h x \in S$,

$$
g x=g\left(g_{m} h\right)^{-1}\left(g_{m} h\right) x \in g h^{-1} g_{m}^{-1} S
$$

Every group element $y \in G$ acts by homeomorphisms of $X$ to itself, since the continuous inverse is given by $y^{-1}$. Thus, the image $g h^{-1} g_{m}^{-1} S$ of the open set $S$ is open in $X$. Continuing,

$$
g h^{-1} g_{m}^{-1} S \subset g h^{-1} g_{m}^{-1} g_{m} V x \subset g h^{-1} V x \subset g V^{-1} \cdot V x \subset U x
$$

Therefore, $g x$ is an interior point of $U x$, for all $g \in U$.

## 5.B Appendix: topology on quotients $H \backslash G$ or $G / H$

As always, $G$ is a topological group, which requires that $G$ be locally compact and Hausdorff, with the group operation and inversion continuous.

Let $H$ be a closed subgroup of $G$. As sets, the left-quotient $H \backslash G$ is the set of cosets $H g$ with $g \in G$, and $G / H$ is the set of cosets $g H$ with $g \in G$. Let $q: G \rightarrow H \backslash G$ be the quotient map $q(g)=H g$. A subset $U \subset H \backslash G$ is open when the inverse image $q^{-1}(U)=\{g \in G: H g \subset U\}$ is open in $G$. Similarly for $G / H$.
[5.B.1] Claim: The quotient maps are open maps, meaning that they take open sets to open sets. Proof: For open $U \subset G$,

$$
q^{-1}(q(U))=\left\{g \in G: H g \subset \bigcup_{u \in U} H u\right\}=\{g \in G: g \in H \cdot U\}=H \cdot U=\bigcup_{h \in H} h \cdot U
$$

Since the group operation and inverse are continuous, for every $h \in H$ the map $g \rightarrow h \cdot g$ is a homeomorphism of $G$ to itself. Thus, every set $h \cdot U$ is open. An arbitrary union of opens is open.

## 6. Action of $G$ on function spaces on $G$

1. Action of $G$ on $L^{2}(\Gamma \backslash G)$
2. Action of $G$ on $C_{c}^{o}(\Gamma \backslash G)$
3. Test functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$
4. Action of $G_{\mathbb{A}}$ on $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$
5. Symmetry of invariant Laplacians
6. An instance of Schur's lemma
7. Duality of induced representations
8. An instance of Frobenius reciprocity
9. Induction in stages
10. Representations of compact $G / Z$
11. Gelfand-Kazhdan criterion

Appendix A: action of compact abelian groups

The function spaces here are more complicated versions of the very concrete examples of chapter 12 , where various spaces of functions on the real line were given metrics or topologies so that they would be complete or quasi-complete. In some of those concrete examples of spaces of functions on $\mathbb{R}$, the translation action of $\mathbb{R}$ on functions plays a role. Spaces of automorphic forms are less visualizable examples. Fortunately, most of the specifics of the concrete examples are irrelevant to proofs.

### 6.1 Action of $G$ on $L^{2}(\Gamma \backslash G)$

For this section, $G$ need merely be a topological group, unimodular in the sense that its right-invariant measure is left-invariant. Let $\Gamma$ be a discrete subgroup, and $K$ a compact subgroup. This includes the assumptions that $G$ is locally compact, Hausdorff, and countably-based. This applies to both classical situations and adelic, such as $G=S L_{n}(\mathbb{R})$ and $\Gamma=S L_{n}(\mathbb{Z})$, and also to $G=Z^{+} G L_{n}(\mathbb{A})$ and $\Gamma=G L_{n}(k)$ for number fields $k$.

Identify functions on $\Gamma \backslash G / K$ with right $K$-invariant functions on the overlying space $\Gamma \backslash G$ by composition with the quotient map. Unlike $\Gamma \backslash G / K$, the space $\Gamma \backslash G$ admits an action of $G$ by right translation. A right $G$-invariant measure $d g$ on $G$ (Haar measure) specifies a unique normalization for a right $G$-invariant measure on $\Gamma \backslash G$ by the unwinding characterization [5.2]

$$
\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma} f(\gamma g)\right) d g=\int_{G} f(g) d g \quad\left(\text { for } f \in C_{c}^{o}(G)\right)
$$

Uniqueness of the Haar measure on $G$ up to scalars, and uniqueness of a measure on $\Gamma \backslash G$ compatible with unwinding, are special cases of uniqueness of invariant distributions, as in [14.4]. In many examples, existence is not an issue, as it is established by reduction to simpler cases. We have an isometry:
[6.1.1] Claim: For $f \in L^{2}(\Gamma \backslash G)$ and $g \in G$, the right translate $(g \cdot f)(x)=f(x g)$ is still in $L^{2}(\Gamma \backslash G)$, and $|g \cdot f|_{L^{2}}=|f|_{L^{2}}$.
Proof: Directly computing,

$$
\int_{\Gamma \backslash G}|(g \cdot f)(x)|^{2} d x=\int_{\Gamma \backslash G}|f(x g)|^{2} d x=\int_{\Gamma \backslash G}|f(x)|^{2} d x
$$

by replacing $x$ by $x g^{-1}$, using the invariance of the measure.
The over-riding point is the continuity of the group action:
[6.1.2] Theorem: $G \times L^{2}(\Gamma \backslash G) \longrightarrow L^{2}(\Gamma \backslash G)$ by $g \times f \longrightarrow(x \rightarrow f(x g))$ for $x, g \in G$ is (jointly) continuous. That is, $L^{2}(\Gamma \backslash G)$ is a unitary representation space for $G$.
Proof: For the moment, write $|\cdot|$ for the $L^{2}$ norm. The crux of the matter is that $L^{2}$ functions can be approximated by continuous, compactly-supported functions, and the latter are uniformly continuous.

## 6. Action of $G$ on function spaces on $G$

Granting that approximation property for a moment, given $\varepsilon>0$, take $\varphi \in C_{c}^{o}(\Gamma \backslash G)$ such that $|f-\varphi|<\varepsilon$. Being compactly supported, $\varphi$ is uniformly continuous, and the topology of $\Gamma \backslash G$ descends from that of $G$, so for every $\varepsilon_{1}>0$ there is a neighborhood $U$ of 1 in $G$ such that $|\varphi(x u)-\varphi(x)|<\varepsilon_{1}$ for all $x \in \Gamma \backslash G$ and for all $u \in U$. For every $F \in L^{2}(\Gamma \backslash G)$ such that $|f-F|<\varepsilon$, for $h=g u \in g U$, the triangle inequality breaks things into atomic issues:

$$
\begin{gathered}
|g \cdot f-h \cdot F|=\left|f-g^{-1} h \cdot F\right|=|f-u \cdot F| \leq|f-\varphi|+|\varphi-u \cdot \varphi|+|u \cdot \varphi-u \cdot f|+|u \cdot f-u \cdot F| \\
=|f-\varphi|+|\varphi-u \cdot \varphi|+|\varphi-f|+|f-F|<3 \varepsilon+|u \cdot \varphi-\varphi|
\end{gathered}
$$

The support $S$ of $\varphi$ has finite measure $\mu$, so for $U$ a small enough neighborhood of 1 such that $\sup _{x}|\varphi(x u)-\varphi(x)|<\varepsilon / 2 \mu$ for $u \in U$,

$$
\left(\int_{\Gamma \backslash G}|\varphi(x u)-\varphi(x)|^{2} d x\right)^{\frac{1}{2}}<\varepsilon
$$

and $|g \cdot f-h \cdot F|<4 \varepsilon$, proving joint continuity.
Density of continuous compactly-supported functions in $L^{2}$ follows from general principles, but its importance justifies some attention to it. Now write $|\cdot|_{L^{2}}$ to distinguish this from the absolute value on real or complex numbers, and $|\cdot|_{L^{1}}$ for the $L^{1}$ norm. First, it suffices to approximate four pieces of $f$, the positive and negative parts of its real and imaginary parts, separately, by the triangle inequality. So without loss of generality, suppose $f$ is real-valued and non-negative. For given $\varepsilon>0$ there is a bound $M$ such that a truncated form of $f$, with maximum value replaced by $M$, satisfies

$$
\int_{\Gamma \backslash G}|f-\min (f, M)|^{2}<\varepsilon
$$

Indeed, the sum of integrals of $|f|^{2}$ over the sets $\{M \leq|f| \leq M+1\}$ converges, so the tails must go to zero. Replace $f$ by that truncation.

For $\mu$ a (positive, regular, Borel) measure on a locally compact Hausdorff space, a simple function is a finite real-linear combination $s=\sum_{i} c_{i} \cdot \chi_{E_{i}}$ of characteristic functions $\chi_{E_{i}}$ of $\mu$-measurable sets $E_{i}$. The integral of $s$ is $\int s=\sum_{i} c_{i} \cdot \mu\left(E_{i}\right)$. The integral of a non-negative real-valued (measurable) function $f$ is the sup of the integrals of simple functions $s$ such that $0 \leq s(x) \leq f$ for all $x$. Since $0 \leq f \leq M$, for such $s$

$$
\int|f-s|^{2}=\int|f-s| \cdot|f-s| \leq \int 2 M \cdot|f-s|=2 M \cdot\left(\int f-\int s\right)
$$

Thus, with $s$ such that $\int s$ is within $\varepsilon / 2 M$ of $\int f$, we have $|f-s|_{L^{2}}^{2} \leq 2 M \cdot \varepsilon / 2 M=\varepsilon$. By the triangle inequality, it suffices to approximate characteristic functions of measurable sets. Regularity of $\mu$ is that the measure of a set is the sup of the measures of compacts inside it, and is the inf of the opens containing it. Let $K \subset E \subset U$ with compact $K$ and open $U$ such that $\mu(U)-\mu(K)<\varepsilon$. Urysohn's lemma [9.E.2] yields a continuous function $\varphi$ with values in the range $[0,1]$ and 1 on $K$ and 0 outside $U$, and

$$
\int\left|\varphi-\chi_{E}\right|^{2}=\int_{K}\left|\varphi-\chi_{E}\right|^{2}+\int_{U-K}\left|\varphi-\chi_{E}\right|^{2} \leq 0+\int_{U-K} 1<\varepsilon
$$

Thus, continuous functions approximate simple functions, which are dense in $L^{2}$.
As in [13.13], the strong operator topology on the continuous linear endomorphisms $\operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$ is given by the collection of seminorms

$$
T \longrightarrow|T v|_{L^{2}} \quad\left(\text { for } v \in L^{2}(\Gamma \backslash G)\right)
$$

The strong operator topology is weaker than the (uniform) operator norm topology $\sup _{|v| \leq 1}|T v|$. The strong operator topology is not complete-metrizable but is quasi-complete as a special case of [13.12], and immediately locally convex since the topology is given by seminorms [13.11].
[6.1.3] Corollary: The map $G \rightarrow \operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$ by right translation is continuous when $\operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$ has the strong operator topology.
Proof: This is a paraphrase of the joint continuity of $G \times L^{2}(\Gamma \backslash G) \rightarrow L^{2}(\Gamma \backslash G)$.
[6.1.4] Remark: As mentioned in [13.13], $G \rightarrow \operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$ is not continuous when $\operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$ has the operator norm topology. To see this, for each neighborhood $N$ of $1 \in G$, we claim there is an $L^{2}$ function $f$ such that $|f-g \cdot f|_{L^{2}} \geq 1$, so certainly cannot be made arbitrarily small. Indeed, shrink $N$ if necessary so that it injects to the quotient $\Gamma \backslash G$, by the discreteness of $\Gamma$. Replace $N$ by $N \cap N^{-1}$, so that it is closed under inverses. Take $1 \neq g \in N$, and let $U \subset N$ be a small-enough neighborhood of 1 such that $g \notin U$, by Hausdorff-ness. By continuity of multiplication, there is an open $V \ni 1$ such that $V \cdot V \subset U$. Replace $V$ by $V \cap V^{-1}$ so that $V$ is stable under inverses. Then $g \notin V^{2}$ gives $g^{-1} \notin V^{2}$, and $g^{-1} V \cap V=\phi$. For an $L^{2}$ function $f$ of norm 1 and supported in $V$, the supports of $f$ and $g \cdot f$, namely, $V$ and $g^{-1} V$, are disjoint, so $|f-g \cdot f|_{L}^{2}=\sqrt{2}$. This quantity cannot be made arbitrarily small, so the representation is not continuous for the uniform operator norm topology.

The continuity property of the theorem gives a precise and useful sense to certain integral operators:
[6.1.5] Corollary: For $\varphi \in C_{c}^{o}(G)$, the integral operator $f \rightarrow \varphi \cdot f$ defined on $f \in L^{2}(\Gamma \backslash G)$ by a convergent vector-valued integral

$$
\varphi \cdot f=\int_{G} \varphi(g) g \cdot f d g
$$

is a continuous linear map $L^{2}(\Gamma \backslash G) \longrightarrow L^{2}(\Gamma \backslash G)$, with the natural property that for $F \in L^{2}(\Gamma \backslash G)$

$$
\langle\varphi \cdot f, F\rangle=\int_{G} \varphi(g)\langle g \cdot f, F\rangle d g
$$

In fact, letting $T_{g} f=g \cdot f$ and $T_{\varphi}$ for the expected operator, we have a vector-valued integral convergent in the strong operator topology:

$$
T_{\varphi}=\int_{G} \varphi(g) T_{g} d g
$$

with the property

$$
T_{\varphi} f=\int_{G} \varphi(g) T_{g} f d g
$$

Proof: This is a special case of properties of Gelfand-Pettis vector-valued integrals [14.1], and using the quasi-completeness of the strong operator topology [13.12], [13.13]. First, because $G \times L^{2}(\Gamma \backslash G) \rightarrow L^{2}(\Gamma \backslash G)$ is continuous, the function $g \rightarrow g \cdot f$ is a continuous $L^{2}(\Gamma \backslash G)$-valued function on $G$. Then $g \rightarrow \varphi(g) g \cdot f$ is a compactly-supported continuous $L^{2}(\Gamma \backslash G)$-valued functions, so by [14.8] the integral purporting to define $\varphi \cdot f$ converges in $L^{2}(\Gamma \backslash G)$, and enjoys the properties [14.1] of Gelfand-Pettis integrals. In particular, $f \rightarrow\langle f, F\rangle$ is a continuous linear functional on $L^{2}(\Gamma \backslash G)$, so by [14.1]

$$
\langle\varphi \cdot f, F\rangle=\left\langle\int_{G} \varphi(g) g \cdot f d g, F\right\rangle=\int_{G}\langle\varphi(g) g \cdot f, F\rangle d g=\int_{G} \varphi(g)\langle g \cdot f, F\rangle d g
$$

Similarly, the continuity of the action of $G$ on $L^{2}(\Gamma \backslash G)$ actually gives the stronger assertion that $g \rightarrow T_{g}$ is a continuous $\operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$-valued function on $G$. Multiplying by $\varphi$ makes the function compactlysupported, and again a Gelfand-Pettis integral exists, by [14.8]. The map $T \rightarrow T f$ is a continuous $L^{2}(\Gamma \backslash G)$ valued function on $\operatorname{End}_{\mathbb{C}}^{o}\left(L^{2}(\Gamma \backslash G)\right)$, so commutes with the integral, by [14.1].

### 6.2 Action of $G$ on $C_{c}^{o}(\Gamma \backslash G)$

The following instance of a general result is also a warm-up to the analogous result for test functions $C_{c}^{\infty}(\Gamma \backslash G)$. For this section, it still suffices that $G$ is a topological group and $\Gamma$ a discrete subgroup. Although any locally convex topological vector space topology can be given by a separating family of seminorms [13.11], this need not be the way the topology arises. The topology on $C_{c}^{o}(\Gamma \backslash G)$ most naturally arises from the expression of $C_{c}^{o}(\Gamma \backslash G)$ as an ascending union of subspaces

$$
C_{E}^{o}(\Gamma \backslash G)=\left\{f \in C_{c}^{o}(\Gamma \backslash G): \operatorname{spt} f \subset E\right\} \quad(\text { where } E \text { varies over compact subsets of } \Gamma \backslash G)
$$

Recall the usual
[6.2.1] Lemma: Each $C_{E}^{o}(\Gamma \backslash G)$ is a Banach space, with sup-norm

$$
\left.|f|_{C^{o}}=\sup _{x \in E}|f(x)| \quad \quad \text { for } f \in C_{E}^{o}(\Gamma \backslash G)\right)
$$

Proof: First, as proven in [13.1.1], the space $C^{o}(E)$ of all continuous functions on a compact subset $E$ of $\Gamma \backslash G$ is a Banach space. Then $C_{E}^{o}(\Gamma \backslash G)$ is a closed subspace defined by pointwise vanishing on the topological boundary of $E$ in $\Gamma \backslash G$ : each evaluation map $f \rightarrow f\left(x_{o}\right)$ is certainly continuous with the sup norm on $C^{o}(E)$, so the kernels are closed, and the intersection of all these closed subspaces for $x_{o}$ on the boundary of $E$ is a closed subspace of a Banach space, so is Banach.

Then $C_{c}^{o}(\Gamma \backslash G)$ is the ascending union, a colimit, as discussed in [13.8] and [13.9]:

$$
C_{c}^{o}(\Gamma \backslash G)=\bigcup_{E} C_{E}^{o}(\Gamma \backslash G)=\operatorname{colim}_{E} C_{E}^{o}(\Gamma \backslash G)
$$

There is a countable cofinal colimit over a countable collection of compact subsets $E_{1} \subset E_{2} \subset \ldots$ of $G$ whose union is $G$. We can take the $E_{i}$ to be closures of a nested family $U_{1} \subset U_{2} \subset \ldots$ of opens whose union is $G$. Cofinal colimits are isomorphic, for general reasons, so

$$
C_{c}^{o}(\Gamma \backslash G)=\bigcup_{E_{i}} C_{E_{i}}^{o}(\Gamma \backslash G)=\operatorname{colim}_{E_{i}} C_{E_{i}}^{o}(\Gamma \backslash G)
$$

Each inclusion $C_{E_{i}}^{o}(\Gamma \backslash G) \subset C_{E_{i+1}}^{o}(\Gamma \backslash G)$ is a homeomorphism to its image, and its image is closed, defined by vanishing at all $x \in E_{i+1}-E_{i}$. A countable colimit of such restricted inclusions is a strict colimit. [50] A strict colimit of Hilbert, Banach, or Fréchet spaces is an LF-space [13.8]. As in [13.12.4], LF-spaces are rarely complete in the strongest sense, but are quasi-complete [13.8.5], and this is sufficient for use.
[6.2.2] Claim: $G$ acts continuously on $C_{c}^{o}(\Gamma \backslash G)$ by right translation: this is the joint continuity of

$$
G \times C_{c}^{o}(\Gamma \backslash G) \longrightarrow C_{c}^{o}(\Gamma \backslash G) \quad \text { by } \quad g \times f \longrightarrow(x \rightarrow f(x g))
$$

Proof: Of course, right translation by $g \in G$ does not stabilize any single $C_{E}^{o}(\Gamma \backslash G)$, only the colimit. Let $\nu_{E}(f)$ be $\sup _{x \in E}|f(x)|$. On the other hand, we have the tautological

$$
\nu_{E}(g \cdot f)=\nu_{E g}(f) \quad \text { and } \quad g \cdot C_{E}^{o}(\Gamma \backslash G)=C_{E g^{-1}}^{o}(\Gamma \backslash G)
$$

Fix $f \in C_{E}^{o}(\Gamma \backslash G), \varepsilon>0$, and $g \in G$. By the uniform continuity of $f$ there is a small-enough neighborhood $U$ of $1 \in G$ such that $|f(x)-f(x u)|<\varepsilon$ for all $u \in U$ and for all $x \in \Gamma \backslash G$. Without loss of generality, $U$ has compact closure $V$, and then $E V^{-1} g^{-1}$ is compact. For all $h=g u \in g U$,

$$
h \cdot C_{E}^{o}(\Gamma \backslash G) \subset g U \cdot C_{E}^{o}(\Gamma \backslash G) \subset C_{E V^{-1} g^{-1}}^{o}(\Gamma \backslash G) \quad(h=g u \in g U)
$$

[50] Slightly older terminology is that a strict colimit is a strict inductive limit.

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For all $F \in C_{E}^{o}(\Gamma \backslash G)$ with $\nu_{E}(f-F)<\varepsilon$, with $h=g u \in g U$

$$
\begin{gathered}
\nu_{E V^{-1} g^{-1}}(g \cdot f-h \cdot F)=\nu_{E V^{-1}}\left(f-g^{-1} h \cdot F\right)=\nu_{E V^{-1}}(f-u \cdot F) \\
\leq \nu_{E V^{-1}}(f-u \cdot f)+\nu_{E V^{-1}}(u \cdot f-u \cdot F)=\nu_{E V^{-1}}(f-u \cdot f)+\nu_{E}(f-F)<\varepsilon+\varepsilon
\end{gathered}
$$

proving joint continuity at $g \times f$ of

$$
g U \times C_{E}^{o}(\Gamma \backslash G) \longrightarrow C_{E V^{-1} g^{-1}}^{o}(\Gamma \backslash G)
$$

The inclusion $C_{E V^{-1} g^{-1}}^{o}(\Gamma \backslash G) \rightarrow C_{c}^{o}(\Gamma \backslash G)$ is continuous, so $g U \times C_{E}^{o}(\Gamma \backslash G) \rightarrow C_{c}^{o}(\Gamma \backslash G)$ is continuous at $\{g\} \times\{f\}$. Since the colimit is stable under the action of $G$, now it makes sense to say that $G \times C_{E}^{o}(\Gamma \backslash G) \rightarrow C_{c}^{o}(\Gamma \backslash G)$ is continuous at $g \times f$. Since $g \in G$ and $f \in C_{E}^{o}(\Gamma \backslash G)$ were arbitrary, this shows that $G \times C_{E}^{o}(\Gamma \backslash G) \rightarrow C_{c}^{o}(\Gamma \backslash G)$ is continuous.

Maps from a colimit $X=\operatorname{colim}_{i} X_{i}$ to another object $Y$ are exactly compatible families of maps $X_{i} \rightarrow Y$ from the limitands to $Y$, as in [13.8], [13.9]. Using the countable cofinal family $E_{1} \subset \ldots \subset E_{i} \subset E_{i+1} \subset \ldots$ for notational convenience, the compatible family of jointly continuous maps

is the joint continuity of $G \times C_{c}^{o}(\Gamma \backslash G) \longrightarrow C_{c}^{o}(\Gamma \backslash G)$ for all $U$, as claimed.
[6.2.3] Corollary: The right translation action $G \times C_{c}^{\infty}(G) \rightarrow C_{c}^{\infty}(G)$ is jointly continuous.
Proof: Take $\Gamma=\{1\}$ in the previous.

### 6.3 Test functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$

As a preamble, we could consider $S L_{n}(\mathbb{Z}) \backslash S L_{n}(\mathbb{R})$, a smooth manifold due to the discreteness of $S L_{n}(\mathbb{Z})$, on which the notion of test function as compactly-supported smooth function has a general sense. The simple example $C_{c}^{\infty}(\mathbb{R})$ is in [13.9]. Edging toward generality, we could similarly consider $S L_{n}\left(\mathbb{Z}\left[\frac{1}{p}\right]\right) \backslash\left(S L_{n}(\mathbb{R}) \times S L_{n}\left(\mathbb{Q}_{p}\right)\right)$, where as usual $\mathbb{Z}\left[\frac{1}{p}\right]$ is $\mathbb{Z}$ with prime $p$ inverted. Again, $S L_{n}\left(\mathbb{Z}\left[\frac{1}{p}\right]\right)$ is demonstrably a discrete subgroup, basically because $\mathbb{Z}\left[\frac{1}{p}\right]$ is discrete in $\mathbb{R} \times \mathbb{Q}_{p}$. Certainly $S L_{n}(\mathbb{R}) \times S L_{n}\left(\mathbb{Q}_{p}\right)$ is not locally homeomorphic to any $\mathbb{R}^{N}$, but to $\left(\mathbb{R} \times \mathbb{Z}_{p}\right)^{N}$ for suitable $N$. Nevertheless, for a compact open subgroup $K_{p}$ of $S L_{n}\left(\mathbb{Q}_{p}\right)$, such as $K_{p}=S L_{n}\left(\mathbb{Z}_{p}\right)$ or any congruence subgroup, the quotient $S L_{n}\left(\mathbb{Q}_{p}\right) / K_{p}$ is discrete, so $S L_{n}\left(\mathbb{Z}\left[\frac{1}{p}\right]\right) \backslash S L_{n}(\mathbb{R}) \times S L_{n}\left(\mathbb{Q}_{p}\right) / K_{p}$ is a smooth manifold. But we object that $S L_{n}\left(\mathbb{Q}_{p}\right)$ no longer acts on this quotient, nor on functions on it.

To overcome this objection, and in anticipation of examination of the action of $G_{\mathbb{A}}$ on functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ in the next section, we can characterize differentiability slightly indirectly, taking advantage of the additional structure. At the same time, the appropriate notion of smoothness of functions $f$ on totallydisconnected groups such as $S L_{n}\left(\mathbb{Q}_{p}\right)$ is that there should exist an open subgroup $K^{\prime}$ such that $f$ is right $K^{\prime}$-invariant. We address these simultaneously. Let $G=G L_{n}(\mathbb{A})$ and $\Gamma=Z^{+} G L_{n}(k)$. ${ }^{[51]}$ Let $\mathfrak{g}$ be the Lie algebra [4.1] of $G_{\infty}=\prod_{v \mid \infty} G_{v}$, and $U \mathfrak{g}$ its universal enveloping algebra [4.3]. Each $\gamma \in \mathfrak{g}$ gives difference quotients for functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ :

$$
X_{\gamma} f(x)=\lim _{t \rightarrow 0} \frac{f\left(x e^{t \gamma}\right)-f(x)}{t} \quad\left(\text { for } x \in Z^{+} G_{k} \backslash G_{\mathbb{A}} \text { and } t \in \mathbb{R}\right)
$$

[^14]
## 6. Action of $G$ on function spaces on $G$

The limit may or may not exist, depending on $f$ and $x \in Z^{+} G_{k} \backslash G_{\mathbb{A}}$ and $\gamma \in \mathfrak{g}$. Given $f$ and $\gamma$, when the limit does exist for every $x$, it gives a compactly supported function on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. Say $f$ is $C^{1}$ if this limit exists for every $x \in Z^{+} G_{k} \backslash G_{\mathbb{A}}$ and $\gamma \in \mathfrak{g}$, and for each $\gamma \in \mathfrak{g}$ is a continuous function on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. Similarly, if $\ell$-fold limits exist and produce continuous functions, $f$ is $C^{\ell}$. The action of a monomial $\gamma_{1} \ldots \gamma_{\ell} \in U \mathfrak{g}$ is

$$
X_{\gamma_{1} \ldots \gamma_{\ell}} f=X_{\gamma_{1}} \circ \ldots \circ X_{\gamma_{\ell}} f
$$

when all the implied limits exist. Temporarily, say that a function $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ such that these limits exist for all elements of $U \mathfrak{g}$ is archimedean-smooth. Also temporarily, say $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is nonarchimedean smooth when $f$ is right $K^{\prime}$-invariant for some open subgroup $K^{\prime} \subset G_{\text {fin }}=\prod_{v<\infty} G_{v}$. [52] A function $f \in C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is smooth when it is both archimedean-smooth and non-archimedean smooth. Compactly supported smooth functions are test functions, denoted $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$.

The LF-space topology on $V=C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is described much as $C_{c}^{\infty}(\mathbb{R})$ in [13.9] and as the colimit $\mathbb{C}^{\infty}$ of finite-dimensional spaces $\mathbb{C}^{n}$ in [13.8]. For compact $E \subset Z^{+} G_{k} \backslash G_{\mathbb{A}}$ and compact-open subgroup $K \subset G_{\mathrm{fin}}$, let $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ be the collection of right $K$-invariant test functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ with support in $E$. Below, we see that each $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ is a Fréchet space, a (projective) limit of a countable collection of spaces $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ of right $K$-invariant $C^{\ell}$ functions with supports on $E$, suitably topologized. Then $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is the colimit of the spaces $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$. Details are as follows.

As always, the $C^{o}$ seminorm on $C_{E}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is the sup norm, and this is a Banach space. Since evaluation at points is a continuous linear functional, a requirement of right $K$-invariance for open subgroup $K \subset G_{\text {fin }}$ is a collection of closed conditions, so defines a closed subspace, giving a Banach space. The $C^{\ell}$ seminorm on $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ should be something like the sup of all derivatives of orders at most $\ell$, with derivatives specifically given by $\mathfrak{g}$ and $U \mathfrak{g}$. In contrast to $\mathbb{R}$, the action of the group here does not generally commute with the natural differential operators. To topologize $C_{E}^{1}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to behave well under the action of $G_{\infty}$ requires examination of the interaction of the right translation action with these right derivatives. Of course, $G_{\text {fin }}$ does commute with the action of $\mathfrak{g}$. For $G_{\infty}$, the interaction is by the conjugation action ${ }^{[53]}$ on $\mathfrak{g}$ :

$$
g \cdot e^{t \gamma} \cdot f=g e^{t \gamma} g^{-1} \cdot g \cdot f \quad\left(\text { for } g \in G_{\infty}\right)
$$

Conjugation interacts well with exponentiation:

$$
g e^{h} g^{-1}=g\left(\sum_{n \geq 0} \frac{h^{n}}{n!}\right) g^{-1}=\sum_{n \geq 0} g \frac{h^{n}}{n!} g^{-1}=\sum_{n \geq 0} \frac{\left(g h g^{-1}\right)^{n}}{n!}=e^{g h g^{-1}}
$$

so

$$
g \cdot e^{t \gamma} \cdot f=g e^{t g \gamma g^{-1}} \cdot g \cdot f \quad \text { and by differentiating } \quad g \cdot \gamma \cdot f=g \gamma g^{-1} \cdot g \cdot f
$$

The right translation action of $G_{\infty}$ on functions does not stabilize any individual differential operator coming from $\mathfrak{g}$, nor any finite subset, and does not stabilize the individual spaces $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, since it does not stabilize supports. Nevertheless, for any open subgroup $K \subset G_{\text {fin }}$, the condition of right $K$-invariance is a collection of closed conditions, so defines a Banach space $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$.

One approach to a suitable topology is as follows. For each bounded neighborhood bof 0 in $\mathfrak{g}$, and for each compact $E \subset Z^{+} G_{k} \backslash G_{\mathbb{A}}$, define a semi-norm

$$
\nu_{b, E}(f)=\sup _{\beta \in b} \sup _{x \in E}|(\beta \cdot f)(x)|
$$

The collection of these has the desirable stability property that

$$
\nu_{b, E}(f)=\nu_{g b g^{-1}, E g^{-1}}(g \cdot f)
$$

[52] For $f$ not necessarily compactly supported, the non-archimedean notion of smoothness would be local, allowing $K^{\prime}$ to vary. The compact support of $f$ implies uniform non-archimedean smoothness, so we may as well give the simpler definition.
[53] This is an instance of an Adjoint action of a Lie group on its Lie algebra.

Because $\mathfrak{g}$ is finite-dimensional, for every pair of bounded neighborhoods $b, b^{\prime}$ of $0 \in B$, there are constants $0<c<C<\infty$ such that $c \cdot b^{\prime} \subset b \subset C \cdot b^{\prime}$, so

$$
c \cdot \nu_{b^{\prime}, E}(f) \leq \nu_{b, E}(f) \leq C \cdot \nu_{b^{\prime}, E}(f) \quad(\text { for all } f)
$$

That is, the topologies are the same, for all bounded neighborhoods $b$ of 0 , for fixed $E$. That is, we can topologize each space $C_{E}^{1}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ by any one of the topologically equivalent seminorms $\nu_{b, E}$. As in the simplest Euclidean case in [13.1], for a fixed choice of $b \in B$ each $C_{E}^{1}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is complete with respect to the (semi-) norm $\nu_{b, E}$, so is a Banach space. However, here there is no canonical Banach space structure, only a canonical topology, given by any one of the topologically equivalent Banach-space structures. In this topology, pointwise evaluation is continuous, so for open subgroup $K \subset G_{\text {fin }}$ the requirement of right $K$-invariance is a collection of closed conditions, so $C_{E}^{1}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ is a closed subspace of any of these Banach spaces.

Similarly, to topologize $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, let $B$ be the collection of bounded neighborhoods of 0 in the graded piece $U \mathfrak{g}^{\leq \ell}$ of elements of degree $\leq \ell$ in $U \mathfrak{g}$. Each $b \in B$ and compact $E$ give a seminorm

$$
\nu_{b, E}(f)=\sup _{\alpha \in b} \sup _{x \in E}|\alpha f(x)|
$$

on $\ell$-times differentiable functions supported on $E$. Since $U \mathfrak{g} \leq \ell$ is finite-dimensional, these seminorms for varying bounded neighborhoods $b$ of $0 \in U \mathfrak{g} \leq k$ are all comparable, giving the same topology on $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. The collection of such seminorms is stabilized by the right action of $G$, by the extension of the conjugation (Adjoint) action, written as conjugation:

$$
\nu_{b, E}(f)=\nu_{g b g^{-1}, E g^{-1}}(g \cdot f)
$$

As for $C^{1}$, each space $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is complete, although there is no canonical Banach-space structure. Again, for open subgroup $K$ of $G_{\text {fin }}$, the $K$-fixed functions $C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ constitute a closed subspace, hence complete.

As in the simplest case [13.2], $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is a (projective) limit of topological vector spaces

$$
C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\bigcap_{\ell} C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\lim _{\ell} C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)
$$

This is equivalent to characterizing the topology on $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ by the seminorms $\nu_{b, E}$ with compact $b \subset U \mathfrak{g} \leq \ell$ for all $\ell$ and $E$. The completeness of the limitands implies completeness of the limit, for general reasons, as in [13.2] and other elementary examples in chapter 12. For each open subgroup $K \subset G_{\text {fin }}$, taking $K$-invariant subspaces commutes with the projective limit, for elementary reasons: the evaluation maps $f \rightarrow f\left(x_{o}\right)$ are continuous and commute with the restriction maps $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \longrightarrow C_{E^{\prime}}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ for $E \supset E^{\prime}$. Thus, we can unambiguously write

$$
C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}=\bigcap_{\ell} C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}=\lim _{\ell} C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}
$$

Then $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is a (strict) colimit

$$
C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)=\bigcup_{E, K} C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}=\operatorname{colim}_{E, K} C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}
$$

The strictness property resides first in the fact that there is the countable cofinal collection $E_{1}, E_{2}, \ldots$, and a countable local basis $K_{1}, K_{2}, \ldots$ for $G_{\text {fin }}$. For example, take $E_{i}$ to be closures of a nested family $U_{1} \subset U_{2} \subset \ldots$ of opens whose union is $Z^{+} G_{k} \backslash G_{A}$. Second, the strictness resides in the fact that the inclusion maps are isomorphisms to their images, which are closed subspaces. Thus, as in the more elementary examples [13.8] and [13.9], this colimit is an LF-space, and is quasi-complete [13.8.5].

The space of distributions $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{*}$ is the dual to $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, with the weak-dual (also called weak-*) topology, as in [13.14].

### 6.4 Action of $G_{\mathbb{A}}$ on $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$

[6.4.1] Theorem: $G_{\mathbb{A}}$ acts continuously on $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ by right translation. This is the joint continuity of

$$
G_{\mathbb{A}} \times C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \longrightarrow C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \quad \text { by } \quad g \times f \longrightarrow(x \rightarrow f(x g))
$$

Proof: As in the case of $C_{c}^{o}(\Gamma \backslash G)$ in the proof of [6.2.2], the collection of seminorms $\nu_{b, E}$ behaves reasonably under right translation by $G_{\infty}$ :

$$
\nu_{b, E}(f)=\nu_{g b g^{-1}, E g^{-1}}(g \cdot f) \quad\left(\text { compact neighborhood } b \text { of } 0 \in U \mathfrak{g}^{\leq k}\right)
$$

and the collection of spaces $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ behaves reasonably:

$$
g \cdot C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}=C_{E g^{-1}}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K} \quad\left(\text { for } g \in G_{\infty}, \text { fixed } K \subset G_{\text {fin }}\right)
$$

although $G_{\infty}$ does not stabilize any individual $C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Similarly, although right translation by $G_{\text {fin }}$ does not preserve right $K$-invariance for any individual open subgroup $K \subset G_{\text {fin }}$, for $g \in G_{\text {fin }}$, for fixed open subgroup $K \subset G_{\text {fin }}$, and for right $K$-invariant $f$, the translate $g \cdot f$ is $g K g^{-1}$-invariant: for $h \in K$,

$$
(g \cdot f)\left(x\left(g h g^{-1}\right)\right)=f\left(x g h g^{-1} g\right)=f(x g h)=f(x g)=(g \cdot f)(x)
$$

Still $g K g^{-1}$ is an open subgroup of $G_{\text {fin }}$, so

$$
g \cdot C_{E}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}=C_{E g^{-1}}^{\ell}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{g K g^{-1}} \quad\left(\text { for } g \in G_{\infty}, \text { fixed } K \subset G_{\text {fin }}\right)
$$

As in the proof of [6.2.2], for joint continuity, we further need comparisons for $g$ in a small open set containing a given $g_{o} \in G_{\mathbb{A}}$.

This uniformity is easiest to see for $G_{\text {fin }}$ :
[6.4.2] Claim: Given $g \in G_{\text {fin }}$, a compact neighborhood $C$ of $g$, and compact open subgroup $K$ of $G_{\text {fin }}$, $\bigcap_{h \in C} h K h^{-1}$ is still an open subgroup of $G_{\text {fin }}$.
Proof: (of claim) Since inversion and multiplication are continuous, $U=g^{-1} C$ is a compact neighborhood of 1 . We may as well enlarge it to $C \cdot K$. Since $K$ is compact, $C \cdot K$ is still compact. Thus, $C \cdot K$ consists of finitely-many cosets $c_{1} K, \ldots, c_{\ell} K$, and any $h \in C$ is $h=g c_{i} k$ for some $i$ and some $k \in K$, and

$$
h K h^{-1}=\left(g c_{i} k\right) \cdot K \cdot\left(g c_{i} k\right)^{-1}=g\left(c_{i} \cdot K \cdot c_{i}^{-1}\right) g^{-1}
$$

Thus, the indicated intersection is actually a finite intersection of open subgroups, so is open.
The uniformity at archimedean places is slightly more complicated, but is parallel to [6.2.2], with further details. Fix $f \in C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right), \varepsilon>0, g \in G_{\infty}$, and $1 \leq \ell \in \mathbb{Z}$. Given a compact neighborhood $b$ of $0 \in U \mathfrak{g} \leq \ell$, by the uniform continuity of $f$ and its derivatives $\beta f$ for $\beta \in b$, there is a small-enough neighborhood $U$ of $1 \in G_{\infty}$ such that $|\beta f(x)-\beta f(x u)|<\varepsilon$ for all $u \in U$, for all $x \in Z^{+} G_{k} \backslash G_{\mathbb{A}}$, and for all $\beta \in b$. Without loss of generality, $U$ has compact closure $V$. Certainly $E^{\prime}=E V^{-1} g^{-1}$ is compact. Being the continuous image of compact $V \times b$, the set $\bigcup_{v \in V} v b v^{-1}$ is itself a compact neighborhood of $0 \in U \mathfrak{g} \leq \ell$, so the seminorms $\nu_{g^{-1} b g, E^{\prime}}$ are uniformly comparable to $\nu_{b, E^{\prime}}$. For all $h=g u \in g U \subset G_{\infty}$,

$$
h \cdot C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K} \subset g U \cdot C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K} \subset C_{E V^{-1} g^{-1}}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K} \quad(h=g u \in g U)
$$

For all $F \in C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}$ with $\nu_{b, E V^{-1}}(f-F)<\varepsilon$, with $h=g u \in g U$, using the uniform comparability of these seminorms.

$$
\begin{gathered}
\nu_{b, E V^{-1} g^{-1}}(g \cdot f-h \cdot F)=\nu_{g^{-1} b g, E V^{-1}}\left(f-g^{-1} h \cdot F\right)=\nu_{g^{-1} b g, E V^{-1}}(f-u \cdot F) \ll \nu_{b, E V^{-1}}(f-u \cdot F) \\
\leq \nu_{b, E V^{-1}}(f-u \cdot f)+\nu_{b, E V^{-1}}(u \cdot f-u \cdot F)=\nu_{b, E V^{-1}}(f-u \cdot f)+\nu_{b, E}(f-F)<\varepsilon+\varepsilon
\end{gathered}
$$

proving archimedean joint continuity at $\{g\} \times\{f\}$ of

$$
g U \times C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K} \longrightarrow C_{E V^{-1} g^{-1}}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K}
$$

The inclusion $C_{E V^{-1} g^{-1}}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous, so $g U \times C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow$ $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous at $\{g\} \times\{f\}$. Since the colimit is stable under the action of $G$, now it makes sense to say that $G \times C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous at $\{g\} \times\{f\}$. Since $g \in G$ and $f \in C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ were arbitrary, this shows that $G \times C_{E}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is continuous.

Maps from a colimit $X=\lim _{i} X_{i}$ to another object $Y$ are exactly compatible families of maps $X_{i} \rightarrow Y$ from the limitands to $Y$, as in [13.8], [13.9], [13.10]. Thus, the compatible family of continuous maps

immediately gives the joint continuity of

$$
G \times C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \longrightarrow C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \quad(\text { for all } U)
$$

as claimed.
[6.4.3] Corollary: The right translation action $G \times C_{c}^{\infty}(G) \rightarrow C_{c}^{\infty}(G)$ is jointly continuous.
Proof: The proof did not use specific features of $Z^{+} G_{k}$ other than its discreteness in $Z^{+} \backslash G$, so we could as well take replace $Z^{+} G_{k}$ by $\{1\}$.
[6.4.4] Corollary: The contragredient or dual action of $G$ on distributions $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{*}$ defined by $(g \cdot u)(f)=u\left(g^{-1} \cdot f\right)$ for all $f \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ gives jointly continuous

$$
G \times C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{*} \longrightarrow C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{*}
$$

Proof: This is a special case of the continuity of $G \times V \rightarrow V$ giving that of $G \times V^{*} \rightarrow V^{*}$, where $V^{*}$ is the dual space of $V$ and is given the weak dual topology [13.11]. The group acts on the dual by $(g \cdot \lambda)(v)=\lambda\left(g^{-1} \cdot v\right)$. Given $g \in G, \lambda \in V^{*}, v \in V$, and $\varepsilon>0$, we want a neighborhood of $g U$ and a neighborhood $\lambda+N$ of $\lambda$ such that for $h=g u \in g U$ and $\mu=\lambda+\nu \in \lambda+N,|(h \mu-g \lambda)(v)|<\varepsilon$. This is

$$
\varepsilon>|(h \mu-g \lambda)(v)|=\left|\mu\left(h^{-1} v\right)-\lambda\left(g^{-1} v\right)\right|=\left|\nu\left(u^{-1} g^{-1} v\right)+\lambda\left(u^{-1} g^{-1} v-g^{-1} v\right)\right|
$$

By continuity of $G \times V \rightarrow V$, take $U$ small enough so that $u^{-1} g^{-1} v-g^{-1} v$ is in a small enough neighborhood $E$ of 0 in $V$ such that $|\lambda(E)|<\frac{\varepsilon}{2}$. Then take $N$ small enough so that $\left|\nu\left(U^{-1} g^{-1} v\right)\right|<\frac{\varepsilon}{2}$.

### 6.5 Symmetry of invariant Laplacians

Just as the density of $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ was used above to examine the representation of $G$ on $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, we need the density of test functions in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ to prove things about the invariant Laplacians on right $K_{\infty}$-invariant functions. As in [4.2], the invariant Laplacians $\Delta_{v}$ on $G_{v} / K_{v}$ with $v$ archimedean are the restrictions of the corresponding Casimir elements $\Omega_{v}$ to right $K_{v}$-invariant functions.
[6.5.1] Claim: In the four simplest examples from chapter $1, C_{c}^{\infty}(\Gamma \backslash G / K)$ is dense in $L^{2}(\Gamma \backslash G / K)$. Similarly, on adele groups $G_{\mathbb{A}}$ is in chapters 2 and $3, C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is dense in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$.
Proof: The argument for $\Gamma \backslash G / K$ is the same as that for $Z^{+} G_{k} \backslash G_{\mathbb{A}}$, simply dropping any reference to nonarchimedean phenomena. We give the adele-group argument, to include the non-archimedean aspects. By

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density of continuous, compactly-supported functions in either case, it suffices to $L^{2}$-approximate continuous, compactly-supported functions $f$ by smooth, compactly-supported functions. The standard device uses a smooth Dirac sequence ${ }^{[54]}\left\{\varphi_{i}\right\}$ on $G$, where smoothness at archimedean places means indefinitely continuously differentiable, and at non-archimedean places means locally constant. Put

$$
\left(\varphi_{i} \cdot f\right)(g)=\int_{G_{\mathbb{A}}} \varphi_{i}(h) f(g h) d h
$$

As in [14.5] and [14.6], via Gelfand-Pettis integrals, $\varphi_{i} \cdot f$ is in $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, and $\varphi_{i} \cdot f \rightarrow f$ in the $L^{2}$ topology, and, in fact, in the finer LF-topology on $C_{c}^{o}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. Changing variables in the integral,

$$
\left(\varphi_{i} \cdot f\right)(g)=\int_{G_{\mathbb{A}}} \varphi_{i}\left(g^{-1} h\right) f(h) d h
$$

That is, the function-valued function

$$
h \longrightarrow\left(g \rightarrow \varphi_{i}\left(g^{-1} h\right) f(h)\right)
$$

is a continuous, compactly-supported, $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$-valued function. Thus, it has a Gelfand-Pettis integral in $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. That is, each $\varphi_{i} \cdot f$ is a smooth, compactly-supported function. We saw that the sequence of these approaches $f$.
[6.5.2] Corollary: The Casimir operators $\Omega_{v}$ on the archimedean factors $G_{v}$ of $G_{\mathbb{A}}$ are symmetric on $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, that is,

$$
\left\langle\Omega_{v} f, F\right\rangle_{Z^{+} G_{k} \backslash G_{\mathbb{A}}}=\left\langle f, \Omega_{v} F\right\rangle_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \quad\left(\text { for } f, F \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

and negative semi-definite on right $K_{\infty}$-invariant functions $C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K_{\infty}}$ :

$$
\left\langle\Omega_{v} f, f\right\rangle_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \leq 0 \quad\left(\text { for } f \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)\right)
$$

Proof: As in [4.2], the Casimir element $\Omega_{v}$ on archimedean $G_{v}$, or a constant multiple of it, is described as follows. Let $\langle\alpha, \beta\rangle=\operatorname{Retr}(\alpha \beta)$ be the trace pairing on the Lie algebra $\mathfrak{g}$ of $G_{v}$, where trace is just matrix trace. As in [4.2], for any basis $\left\{x_{j}\right\}$ of $\mathfrak{g}$ and dual basis $\left\{x_{j}^{*}\right\}, \Omega_{v}=\sum_{j} x_{j} x_{j}^{*} \in U \mathfrak{g}$. In fact, the pairing is negative-definite on the Lie algebra $\mathfrak{k}$ of $K_{v}$, which is $O(n, \mathbb{R})$ or $U(n)$, and positive-definite on their complements $\mathfrak{s}$ consisting of symmetric real matrices and hermitian-symmetric complex matrices. So we can choose an orthogonal basis $\left\{x_{j}\right\}$ for $\mathfrak{s}$ with $x_{j}^{*}=x_{j}$, and an orthogonal basis $\left\{\theta_{j}\right\}$ for $\mathfrak{k}$ with $\theta_{j}^{*}=-\theta_{j}$. Thus,

$$
\Omega_{v}=\sum_{j} x_{j}^{2}-\sum_{j} \theta_{j}^{2}
$$

The action of $x \in \mathfrak{g}$ on $f \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ is

$$
(x f)(g)=\left.\frac{\partial}{\partial t}\right|_{t=0} f\left(g e^{t x}\right)
$$

To properly indicate the order of operations, the evaluation at $t=0$ should come after the derivative, so for clarity write

$$
(x f)(g)=\left.\right|_{t=0} \frac{\partial}{\partial t} f\left(g e^{t x}\right)
$$

[54] A smooth Dirac sequence or approximate identity $\left\{\varphi_{i}\right\}$ on a unimodular group $G$ is a sequence of smooth, compactly-supported real-valued functions $\varphi_{i}$ so that $\int_{G} \varphi_{i}(g) d g=1$ and $0 \leq \varphi \leq 1$ for all $i$, and for every neighborhood $N$ of the identity $e \in G$, there is $i_{o}$ such that for all $i \geq i_{o}$ the support of $\varphi_{i}$ is inside $N$.

For $f, F \in C_{c}^{\infty}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$,

$$
\langle x f, F\rangle=\left.\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}}\right|_{t=0} \frac{\partial}{\partial t} f\left(g e^{t x}\right) \cdot \bar{F}(g) d g
$$

It would be needlessly reckless to claim that the integrand is a compactly-supported, continuous $C_{c}^{\infty}(\mathbb{R})$ valued function

$$
g \longrightarrow\left(t \longrightarrow f\left(g e^{t x}\right) \cdot \bar{F}(g)\right)
$$

because the compact support in $t$ can easily fail. Instead, restrict $t$ to a small interval $[-\varepsilon, \varepsilon]$. The integrand is a compactly-supported, continuous $C^{\infty}[-\varepsilon, \varepsilon]$-valued function, where $C^{\infty}[-\varepsilon, \varepsilon]$ has its natural Fréchetspace structure, as in [13.8] and [13.9]. Evaluation at $t=0$ is a continuous linear functional on $C^{\infty}[-\varepsilon, \varepsilon]$, so by Gelfand-Pettis the evaluation commutes with the integral:

$$
\left.\int_{Z+G_{k} \backslash G_{\mathbb{A}}}\right|_{t=0} \frac{\partial}{\partial t} f\left(g e^{t x}\right) \cdot \bar{F}(g) d g=\left.\right|_{t=0} \int_{Z+G_{k} \backslash G_{\mathbb{A}}} \frac{\partial}{\partial t} f\left(g e^{t x}\right) \cdot \bar{F}(g) d g
$$

Similarly, $\partial / \partial t$ is a continuous map of $C^{\infty}[-\varepsilon, \varepsilon]$ to itself, so commutes with the integral:

$$
\left.\right|_{t=0} \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \frac{\partial}{\partial t} f\left(g e^{t x}\right) \cdot \bar{F}(g) d g=\left.\right|_{t=0} \frac{\partial}{\partial t} \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f\left(g e^{t x}\right) \cdot \bar{F}(g) d g
$$

This legitimizes the change of variables replacing $g$ by $g e^{-t \gamma}$ :

$$
\left.\right|_{t=0} \frac{\partial}{\partial t} \int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f\left(g e^{t x}\right) \cdot \bar{F}(g) d g=\left.\right|_{t=0} \frac{\partial}{\partial t} \int_{Z+G_{k} \backslash G_{\mathbb{A}}} f(g) \cdot \bar{F}\left(g e^{-t \gamma}\right) d g
$$

The differentiation and evaluation can be moved back inside the integral for the same reasons. Thus, $\langle x f, F\rangle=-\langle f, x F\rangle$, and

$$
\left\langle x^{2} f, F\right\rangle=-\langle x f, x F\rangle=\left\langle f, x^{2} F\right\rangle
$$

Thus,

$$
\left\langle\left(\sum_{j} x_{j}^{2}-\sum_{j} \theta_{j}^{2}\right) f, F\right\rangle=\left\langle f,\left(\sum_{j} x_{j}^{2}-\sum_{j} \theta_{j}^{2}\right) F\right\rangle
$$

This is the symmetry. On $f, F$ right $K_{v}$-invariant, $\theta_{j} \in \mathfrak{k}$ acts by 0 , so

$$
\left\langle\left(\sum_{j} x_{j}^{2}-\sum_{j} \theta_{j}^{2}\right) f, f\right\rangle=\left\langle\left(\sum_{j} x_{j}^{2}\right) f, f\right\rangle=\sum_{j}\left\langle x_{j}^{2} f, f\right\rangle=-\sum_{j}\left\langle x_{j} f, x_{j} f\right\rangle \leq 0
$$

giving the non-positiveness.

### 6.6 An instance of Schur's lemma

The general idea of Schur's lemma is that endomorphisms of an irreducible G-representation space commuting with the $G$-action must be scalar.

The group $G=G L_{n}\left(\mathbb{Q}_{p}\right)$ is abstracted to consider $G$ totally disconnected, in the sense that for every $x \neq y$ in $X$ there are open sets $x \in U, y \in V$ so that $U \cap V=\phi$ and $U \cup V=X$. Since we only consider Hausdorff topological groups, the sets $U, V$ in the definition are not only open but closed. The first assertions admit explicit arguments for $G L_{n}\left(\mathbb{Q}_{p}\right)$, but, in fact, use only general features of totally disconnected groups.
[6.6.1] Claim: At every point $x$ of a locally compact totally disconnected space $X$ there is a local basis consisting of compact open sets.

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Proof: By local compactness, take an open set $V$ containing $x$ so that the closure $\bar{V}$ is compact. The boundary $\partial V=\bar{V} \cap(X-V) \subset \bar{V}$ is closed, so is compact. For $y \in \partial V$, there are open (and closed) sets $x \in U_{y}$ and $y \in V_{y}$ so that $U_{y} \cap V_{y}=\emptyset$ and $U_{y} \cup V_{y}=X$. Take a finite subcover $V_{y_{1}}, \ldots, V_{y_{n}}$ of $\partial V$. The set

$$
V-\left(\bigcup_{i} \bar{V}_{y_{i}}\right)=\bar{V}-\left(\bigcup_{i} V_{y_{i}}\right)
$$

contains $x$, and is both open and closed. Being a closed subset of the compact set $\bar{V}$ in a Hausdorff space, it is compact.
[6.6.2] Claim: A locally compact totally disconnected topological group $G$ has a local basis at $1=1_{G}$ consisting of compact open subgroups. Further, for a fixed compact open subgroup $K_{1}$, there are normal compact open subgroups of $K_{1}$ such that $K_{1} \supset K_{2} \supset K_{3} \supset \ldots$ and these $K_{j}$ are a local basis at 1 .
Proof: Let $V$ be a compact open subset of $G$ containing 1, by the previous. The set

$$
K=\left\{x \in G: x V \subset V \text { and } x^{-1} V \subset V\right\}
$$

is a subgroup of $G$, and

$$
K=\left(\bigcap_{v \in V} V v^{-1}\right) \cap\left(\bigcap_{v \in V} V v^{-1}\right)^{-1}
$$

so $K$ is the continuous image of a compact set, so is compact. What remains to be shown is that $K$ is open.
To the latter end, show that the compact-open topology on $G$ constructed from the original topology on $G$ is the original topology on $G$. That is, show that, for compact $C$ in $G$ and for open $V$ in $G$, the set

$$
U=U_{C, V}=\{x \in G: x C \subset V\}
$$

is open in $G$. Take $U$ non-empty, and $x \in U$. For all points $x y \in x C$ for $y \in C$, there is a small-enough open neighborhood $U_{y}$ of 1 so that the open neighborhood $x U_{y} y$ of $x y$ is contained in $V$. By continuity of the multiplication in $G$, there is an open neighborhood $W_{y}$ of 1 so that $W_{y} W_{y} \subset U_{y}$. The sets $x W_{y} y$ cover $x C$. Let $x W_{y_{1}} y_{1}, \ldots, x W_{y_{n}} y_{n}$ be a finite subcover. Put $W=\bigcap_{i} W_{y_{i}}$. Then $x W$ is a neighborhood of $x$ and

$$
x W \cdot C \subset x W \cdot \bigcup_{i} W_{y_{i}} y_{i}
$$

and $x W W_{y_{i}} y_{i} \subset x W_{y_{i}} W_{y_{i}} y_{i} \subset x U_{y_{i}} y_{i}$. Thus, $U$ is open.
Let $H_{1} \supset H_{2} \supset \ldots$ be a local basis of compact open subgroups. Put $K_{1}=H_{1}$. Of course, $K_{1}$ is a union of cosets $K_{1}=\bigcup_{x \in K_{1}} x H_{2}$. By compactness, there is a finite subcover $K_{1}=x_{1} H_{2} \cup \ldots \cup x_{n} H_{2}$. Thus, $H_{2}$ is of finite index in $K_{1}$. Then

$$
K_{2}=\bigcap_{x \in K_{1}} x H_{2} x^{-1}=\bigcap_{x \in K_{1} / H_{2}} x H_{2} x^{-1}=x_{1} H_{2} x_{1}^{-1} \cap \ldots \cap x_{n} H_{2} x_{n}^{-1}
$$

is a finite intersection of compact open subgroups, so is compact and open, and is normal in $K_{1}$. Replace $H_{3}$ by $H_{3} \cap K_{2}$ and continue inductively.

A representation $G \times V \rightarrow V$ of $G$ on a complex vector space $V$ is smooth when, for all $v \in V$, the isotropy group $G_{v}=\{g \in G: g \cdot v=v\}$ is open. Because of the total-disconnectedness, this condition is equivalent to an expression of $V$ as an ascending union, or strict colimit [13.8], [13.9],

$$
V=\bigcup_{K} V^{K}=\operatorname{colim}_{K} V^{K} \quad\left(\text { with } V^{K}=K \text {-fixed vectors in } V\right)
$$

where $K$ runs over compact open subgroups of $G$. All our topological groups are countably based, so the colimit has a countable cofinal subsystem. Since every $G$-homomorphism $V \rightarrow W$ preserves $K$-fixed-ness, smoothness is automatically preserved by $G$-homomorphisms. The representation $V$ is admissible when each
$V^{K}$ is finite-dimensional. In that case, $V$ is a strict colimit of finite-dimensional topological vector spaces. [55] Finite-dimensional topological vectors spaces have unique topologies [13.4], and every linear map from a finite-dimensional space is continuous. Thus, admissible $V$ has a topology so that every linear $T: V \rightarrow W$ is continuous. In particular, the continuous dual $V^{*}$ of admissible $V$ consists of all linear functions on $V$. Similarly, for a topological vector space $V$ that is a strict colimit of finite-dimensional spaces, every group action $G \times V \rightarrow V$ is continuous. [56]

Generally, given an arbitrary (continuous) representation $G \times V \rightarrow V$ on a vectorspace $V$, the subspace $V^{\infty}$ of smooth vectors is

$$
V^{\infty}=\left\{v \in V: G_{v} \text { is open }\right\}
$$

This subspace is $G$-stable, so the restriction of $G \times V \rightarrow V$ to $G \times V^{\infty} \rightarrow V^{\infty}$ is a $G$-subrepresentation of $V$. For a smooth representation $G \times V \rightarrow V$, the (smooth) dual or (smooth) contragredient $V^{\vee}$ is the representation of $G$ on the smooth vectors $\left(V^{*}\right)^{\infty}$ in the continuous linear dual $V^{*}$, where, as always, $G$ acts on $V^{*}$ by $(g \cdot \lambda)(v)=\lambda\left(g^{-1} \cdot v\right)$. A smooth representation $G \times V \rightarrow V$ is (topologically) irreducible when it contains no proper closed subspace $W$ stable under $G$. It is (algebraically) irreducible when it contains no proper (not necessarily closed) subspace $W$ stable under $G$.
[6.6.3] Claim: Every subspace of a strict colimit of finite-dimensional spaces is closed.
Proof: Let $W$ be a subspace of $\bigcup_{i} V_{i}$ where $V_{i} \subset V_{i+1}$ and every $V_{i}$ is finite-dimensional. Given $v \in V$ and $v \notin W$, we can make a linear functional $\lambda_{v}$ on $V$ that vanishes on $W$ and is non-zero on $v$, as follows. Let $j$ be large-enough so that $v \in V_{j}$. Let $\lambda_{j}: V_{j} \rightarrow \mathbb{C}$ be linear with $\lambda_{j}\left(W \cap V_{j}\right)=0$ and $\lambda_{j} v=1$. Extend $\lambda_{j}$ to a linear functional $\lambda_{j+1}$ on $V_{j+1}$ by choosing a complementary subspace $X_{j+1}$ to $V_{j}$ inside $V_{j+1}$, and making $\lambda_{j+1} X_{j+1}=0$, while agreeing with $\lambda_{j}$ on the copy of $V_{j}$ inside $V_{j+1}$. Continue by induction. This defines $\lambda_{v}$ on the strict colimit, and continuity is automatic. Thus, $W$ is the intersection of the kernels of all such $\lambda_{v}$ over all $v \notin W$, so is closed.

The (full) Hecke algebra on $G$ is the space $\mathcal{H}=\mathcal{H}(G)=C_{c}^{\infty}(G)$ of smooth, compactly-supported complexvalued functions on $G$. Here smooth means locally constant, in the sense that, given $\eta \in \mathcal{H}$ and $g \in G$, there is a neighborhood $U$ of $g$ such that $\eta\left(g^{\prime}\right)=\eta(g)$ for all $g^{\prime} \in U$. The Hecke algebra acts on $V$ for any representation $G \times V \rightarrow V$ as usual by

$$
\eta \cdot v=\int_{G} \eta(g) g \cdot v d g \quad \text { (with right-invariant measure) }
$$

and convolution $\eta * \psi$ on $\mathcal{H}$ is characterized by

$$
(\eta * \psi) v=\eta \cdot(\psi \cdot v) \quad(\text { for } v \in V \text { and } \eta, \psi \in \mathcal{H})
$$

That is,

$$
\begin{aligned}
(\eta * \psi) v=\eta \cdot(\psi \cdot v)= & \eta \cdot \int_{G} \psi(x) x \cdot v d x=\int_{G} \eta(y) y \cdot\left(\int_{G} \psi(x) x \cdot v d x\right) d y \\
& =\int_{G} \int_{G} \eta(y) \psi(x) y \cdot x \cdot v d x d y
\end{aligned}
$$

using properties of Gelfand-Pettis integrals of vector-valued functions [14.1]. Changing the order of integration and replacing $y$ by $y x^{-1}$, then changing the order back, and again using Gelfand-Pettis, this is

$$
\int_{G} \int_{G} \eta\left(y x^{-1}\right) \psi(x) y \cdot v d y d x=\int_{G}\left(\int_{G} \eta\left(y x^{-1}\right) \psi(x) d x\right) y \cdot v d y
$$

[^15]
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Thus, we discover an expression for the convolution, without assuming unimodularity of $G$,

$$
(\eta * \psi)(y)=\int_{G} \eta\left(y x^{-1}\right) \psi(x) d x
$$

Unless $G$ is discrete, which is not of interest here, the algebra $\mathcal{H}$ with convolution has no unit element 1. Fortunately, there are sufficiently many idempotents: for compact open subgroup $K$ of $G$, let

$$
\left.e_{K}=\frac{\operatorname{ch}_{K}}{\operatorname{meas}(K)} \in \mathcal{H} \quad \text { (with characteristic function } \operatorname{ch}_{K} \text { of } K\right)
$$

These are idempotents: $e_{K} * e_{K}=e_{K}$. Indeed, for any $K^{\prime} \subset K$, we have $e_{K}^{\prime} * e_{K}=e_{K} * e_{K^{\prime}}=e_{K}$. There are sufficiently many of these idempotents in $\mathcal{H}$ in the sense that, given $\eta \in \mathcal{H}$ there is $e_{K}$ such that $e_{K} * \eta=\eta=\eta * e_{K}$. Smoothness of a $G$-representation space $V$ is that for all $v \in V$ there is a small-enough $K$ so that $e_{K} \cdot v=v$.

A complex vectorspace $V$ which is a module over the $\operatorname{ring} \mathcal{H}$ is smooth when for every $v \in V$ there is a small-enough compact open subgroup $K$ so that $e_{K} v=v$.
[6.6.4] Theorem: The category of smooth $\mathcal{H}$-modules is the same as the category of smooth $G$ representations.
Proof: We have already seen how to get smooth $\mathcal{H}$-modules from smooth $G$-representations. We need to recover the action of $G$ from the $\mathcal{H}$-module structure. Slightly generalizing previous notation, for a compact open subset $X$ of $G$, let

$$
e_{X}=\frac{\operatorname{ch}_{X}}{\operatorname{meas}(X)}
$$

For $v$ in a smooth $\mathcal{H}$-module $V$ with $K$ a small-enough compact open subgroup so that $e_{K} v=v$, and for $g \in G$, try go define the action of $g$ by $g \cdot v=e_{g K} v$. To see that this is well-defined, check that $K$ may be made smaller without altering $e_{g K} \cdot v$. Since $e_{K} \cdot v=v$, by associativity it suffices to show that $e_{g K^{\prime}} * e_{K}=e_{g K}$ for $K^{\prime} \subset K$ :

$$
\operatorname{meas}\left(g K^{\prime}\right) \cdot \operatorname{meas}(K) \cdot\left(e_{g K^{\prime}} * e_{K}\right)(x)=\int_{G} \operatorname{ch}_{g K^{\prime}}\left(x y^{-1}\right) \operatorname{ch}_{K}(y) d y=\int_{K} \operatorname{ch}_{g K^{\prime}}\left(x y^{-1}\right) d y
$$

The integrand is non-zero exactly for $y \in K$ and $x y^{-1} \in g K^{\prime}$, that is, for $y \in K$ and $y \in K^{\prime} g^{-1} x$. Since $K^{\prime} \subset K$, the set $K^{\prime} g^{-1} x \cap K$ is non-empty exactly for $g^{-1} x \in K$, which is $x \in g K$, in which case the intersection has measure meas $\left(K^{\prime}\right)$. That is, with modular function $\delta$,

$$
e_{g K^{\prime}} * e_{K}(x)=\frac{\operatorname{meas}\left(K^{\prime}\right)}{\operatorname{meas}\left(g K^{\prime}\right) \cdot \operatorname{meas}(K)} \cdot \operatorname{ch}_{g K}(x)=\frac{\operatorname{ch}_{g K}(x)}{\delta(g) \cdot \operatorname{meas}(K)}=\frac{\operatorname{ch}_{g K}(x)}{\operatorname{meas}(g K)}
$$

as claimed. Second, check that this gives a group homomorphism: given compact open $K$ and given $h \in G$, take $K^{\prime}$ small-enough so that $e_{K^{\prime}} * e_{h K}=e_{h K}$. This condition is

$$
\begin{gathered}
\operatorname{ch}_{h K}(x)=\operatorname{meas}\left(K^{\prime}\right) \cdot \int_{G} \operatorname{ch}_{K^{\prime}}\left(x y^{-1}\right) \operatorname{ch}_{h K}(y) d y=\operatorname{meas}\left(K^{\prime}\right) \cdot \int_{G} \operatorname{ch}_{K^{\prime}}\left(y^{-1}\right) \operatorname{ch}_{h K}(y x) d y \\
=\operatorname{meas}\left(K^{\prime}\right) \cdot \int_{K^{\prime}} \operatorname{ch}_{h K}(y x) d y
\end{gathered}
$$

It suffices that $\operatorname{ch}_{h K}$ be left $K^{\prime}$-invariant, since then the integrand is $\operatorname{ch}_{h K}(x)$ for all $y \in K^{\prime}$. This left invariance is that $k^{\prime} x \in h K$ for $x \in h K$ and $k^{\prime} \in K^{\prime}$. That is, $k^{\prime} h K \subset h K$, or $k^{\prime} h \in h K$, or $k^{\prime} \in h K h^{-1}$. Thus, $K^{\prime}=h K h^{-1}$ suffices. We must show that $e_{g K^{\prime}} * e_{h K}=e_{g h K}$. Noting that

$$
\operatorname{meas}\left(K^{\prime}\right)=\operatorname{meas}\left(h K h^{-1}\right)=\operatorname{meas}(h K)=\delta(h) \cdot \operatorname{meas}(K)
$$

we have

$$
\begin{aligned}
& \operatorname{meas}\left(g K^{\prime}\right) \cdot \operatorname{meas}(h K) \cdot\left(e_{g K^{\prime}} * e_{h K}\right)(x)=\int_{G} \operatorname{ch}_{g K^{\prime}}\left(x y^{-1}\right) \operatorname{ch}_{h K}(y) d y \\
& \begin{array}{r}
=\delta(h) \int_{G} \operatorname{ch}_{g K^{\prime}}\left(x y^{-1} h^{-1}\right) \operatorname{ch}_{K}(y) d y=\delta(h) \int_{K} \operatorname{ch}_{g h K h^{-1}}\left(x y^{-1} h^{-1}\right) d y \\
=\delta(h) \int_{K} \operatorname{ch}_{g h K}(x y) d y
\end{array}
\end{aligned}
$$

The integrand is non-zero exactly when $x y \in g h K$, which is $y \in K \cap x^{-1} g h K$. This is non-zero exactly for $x^{-1} g h K \subset K$, which is equivalent to $x \in g h K$. Then $y$ is integrated over $K$, giving $\delta(h) \operatorname{meas}(K) \cdot \mathrm{ch}_{g h K}(x)$. Thus,

$$
e_{g K^{\prime}} * e_{h K}=\frac{\delta(h) \operatorname{meas}(K) \cdot \mathrm{ch}_{g h K}}{\operatorname{meas}\left(g K^{\prime}\right) \cdot \operatorname{meas}(h K)}=\frac{\delta(h) \operatorname{meas}(K) \cdot \mathrm{ch}_{g h K}}{\delta(g) \delta(h) \operatorname{meas}(K) \cdot \delta(h) \operatorname{meas}(K)}=\frac{\mathrm{ch}_{g h K}}{\operatorname{meas}(g h K)}
$$

This is the homomorphism property. Visibly, $G$-homomorphisms and $\mathcal{H}$-module homomorphisms are interchanged under this bijection.

In the following, recall from above that for $G \times V \rightarrow V$ with $V$ a strict colimit of finite-dimensional complex vectorspaces, for example admissible, algebraic irreducibility and topological irreducibility are equivalent. For larger, not necessarily admissible, representations, algebraic irreducibility is usually strictly stronger.
[6.6.5] Theorem: (Schur's Lemma) Let $G \times V \rightarrow V$ be an algebraically irreducible smooth representation of $G$. Let $T$ be a $\mathbb{C}$-linear endomorphism of $V$ commuting with all maps $v \rightarrow g \cdot v$ with $g \in G$. Then $T$ is a scalar, that is, multiplication by an element of $\mathbb{C}$.
Proof: (Jacquet) Since $G$ has a countable basis, $\mathcal{H}$ has countable dimension as $\mathbb{C}$-vectorspace. Algebraic irreducibility implies $\mathcal{H} \cdot v=V$ for $v \neq 0$ in $V$, so $V$ is of countable $\mathbb{C}$-dimension. An $\mathcal{H}$-endomorphism $T$ is completely determined by $T v$ for one $v \neq 0$, since $T(\eta v)=\eta T(v)$ for $\eta \in \mathcal{H}$. Thus, the ring $D$ of $\mathcal{H}$-endomorphisms of $V$ has countable $\mathbb{C}$-dimension. As $V$ is algebraically irreducible, for all $T \in D$ both the kernel and image of $T$ are $\mathcal{H}$-submodules, so can be only 0 or $V$. Thus, $D$ is a division ring with $\mathbb{C}$ in its center.

Since $\mathbb{C}$ is algebraically closed, non-scalar $T \in D$ are transcendental over $\mathbb{C}$. Therefore, for $T \in D$ not a scalar the elements $S_{\lambda}=(T-\lambda)^{-1} \in D$ with $\lambda$ varying over $\mathbb{C}$ are linearly independent over $\mathbb{C}$, by uniqueness of partial fraction expansions in $\mathbb{C}(T)$. As $\mathbb{C}$ is uncountable, this would yield an uncountable set of linearly-independent elements of $D$, contradiction.

### 6.7 Duality of induced representations

The group $G$ is still totally disconnected, and representations $G \times V \rightarrow V$ are smooth. For a closed subgroup $H$ of $G$, there is the forgetful functor $\operatorname{Res}_{H}^{G}$ from smooth $G$-representations to smooth $H$-representations, by forgetting all but the action of $H$. We eventually want a right adjoint ${ }^{[57]} \operatorname{Ind}_{H}^{G}$ to this forgetful functor from representations of $H$ to representations of $G$, in the sense that there should be natural bijections

$$
\operatorname{Hom}_{G}\left(V, \operatorname{Ind}_{H}^{G} W\right) \approx \operatorname{Hom}_{H}\left(\operatorname{Res}_{H}^{G} V, W\right) \quad(\text { for representation } W \text { of } H)
$$

for all smooth $G$-representation $V$ and smooth $H$-representation $W$. This adjunction is Frobenius reciprocity, proven in the next section, and is viewed there as approximately analogous to the Cartan-Eilenberg adjunction

$$
\operatorname{Hom}_{\mathbb{Z}}\left(A, \operatorname{Hom}_{\mathbb{Z}}(B, C)\right) \approx \operatorname{Hom}_{\mathbb{Z}}\left(A \otimes_{\mathbb{Z}} B, C\right) \quad(\text { for } \mathbb{Z} \text {-modules } A, B, C)
$$

with $\varphi_{\Phi} \leftarrow \Phi$ by $\varphi_{\Phi}(a)(b)=\Phi(a \otimes b)$, and $\varphi \rightarrow \Phi_{\varphi}$ by $\Phi_{\varphi}(a \otimes b)=\varphi(a)(b)$. The relatively simple argument for this will also be recalled in the next section.
[57] It is not obvious, but there is no left adjoint to $\operatorname{Res}_{H}^{G}$ in this situation.

## 6. Action of $G$ on function spaces on $G$

The present section constructs induced representations $\operatorname{Ind}_{H}^{G} W$ of $G$ made from representations $W$ of $H$, and compactly-induced representations c- $\operatorname{Ind}_{H}^{G} W$ as spaces of $W$-valued functions on $G$, and proves a duality

$$
\left(\mathrm{c}-\operatorname{Ind}_{H}^{G} W, \mathbb{C}\right)^{\vee}=\operatorname{Hom}_{G}\left(\mathrm{c}-\operatorname{Ind}_{H}^{G} W, \mathbb{C}\right) \approx \operatorname{Ind}_{H}^{G}\left(\rho \cdot \operatorname{Hom}_{G}\left(W^{\vee}, \mathbb{C}\right)\right)
$$

with $\rho=\delta_{H} / \delta_{G}$ for modular functions on $H$ and $G$, (proof below) whose relevance to Frobenius reciprocity is suggested by the simpler Cartan-Eilenberg adjunction.

For a smooth representation $W$ of a closed subgroup $H$ of $G$, let $C_{c}^{\infty}(H \backslash G, W)$ be the space of $W$-valued functions $f$ on $G$ that are compactly-supported left-modulo $H$ (in the sense that the images in $H \backslash G$ of their supports are compact), locally constant, and so that

$$
f(h g)=h \cdot f(g) \quad(\text { for } h \in H \text { and } g \in G)
$$

The compact-induced representation c-Ind ${ }_{H}^{G} W$ has representation space $C_{c}^{\infty}(H \backslash G, W)$ with the right translation action of $G$

$$
(g \cdot f)\left(g^{\prime}\right)=f\left(g^{\prime} g\right) \quad\left(\text { for } g, g^{\prime} \in G\right)
$$

The induced representation $\operatorname{Ind}_{H}^{G} W$ has representation space consisting of uniformly locally constant $W$ valued functions $f$ on $G$ satisfying

$$
f(h g)=h \cdot f(g) \quad(\text { for } h \in H \text { and } g \in G)
$$

with the right translation action

$$
(g \cdot f)\left(g^{\prime}\right)=f\left(g^{\prime} g\right) \quad\left(\text { for } g, g^{\prime} \in G\right)
$$

The uniform locally-constant condition on $f$ in this space of functions is that there is a compact open subgroup $\Theta$ so that

$$
f(g \theta)=f(g) \quad(\text { for all } g \in G \text { and for all } \theta \in \Theta)
$$

As a variant of [5.2.1] about invariant measures on quotients, for the present discussion we need a slightly different unwinding:
[6.7.1] Lemma: Let $\delta_{H}, \delta_{G}$ be the modular functions of $H, G$, for a closed subgroup $H$ of $G$, and $\rho=\delta_{H} / \delta_{G}$. There is a non-trivial right $G$-invariant functional $u$ on $c$ - $\operatorname{Ind}_{H}^{G} \rho$, unique up to scalar multiples.
Proof: Let
be the appropriate averaging map, as in [5.1], [5.2], using right Haar measure on $H$. For totally disconnected groups and locally constant, compactly supported functions, the surjectivity of $\alpha$ allows an even simpler argument than the case treated in [5.1].

Attempt to define a right $G$-invariant $\mathbb{C}$-valued $u$ on $c-\operatorname{Ind}_{H}^{G} \rho$ by $u(\alpha f)=\int_{G} f$ with right Haar measure on $G$. The telling issue is well-definedness: replacing $g \rightarrow f(g)$ by $g \rightarrow f(h g)$ gives, on one hand

$$
\int_{G} f(h g) d g=\delta_{H}\left(h^{-1}\right) \cdot \int_{G} f(g) d g=u(\alpha f)
$$

On the other hand,

$$
\alpha\left(h^{\prime} \rightarrow f\left(h h^{\prime}\right)\right)=\int_{H} \rho^{-1}\left(h^{\prime}\right) f\left(h h^{\prime} g\right) d h^{\prime}=\delta_{G}\left(h^{-1}\right) \rho^{-1}\left(h^{-1}\right) \int_{H} \rho^{-1}\left(h^{\prime}\right) f\left(h^{\prime} g\right) d h^{\prime}
$$

Thus, $\rho$ is the only possible choice. As in the proof of [5.2.1], it suffices to show that $\int_{G} f(g) d g=0$ for $\alpha f=0$. Indeed, for $\alpha f=0$, for all $F \in C_{c}^{\infty}(G)$, the integral of $F$ against $\alpha f$ is certainly 0 . Rearrange

$$
\begin{aligned}
0=\int_{G} F(g) \alpha f(g) d g= & \int_{G} \int_{H} F(g) \rho(h)^{-1} f(h g) d h d g=\int_{H} \int_{G} F(g) \rho(h)^{-1} f(h g) d g d h \\
& =\int_{H} \int_{G} F\left(h^{-1} g\right) f(g) \delta_{G}\left(h^{-1}\right) d g d h
\end{aligned}
$$

by replacing $g$ by $h^{-1} g$. Replacing $h$ by $h^{-1}$ replaces right Haar measure by $\delta_{H}(h) d h$, so

$$
0=\int_{G}\left(\int_{H} \rho(h)^{-1} \cdot F(h g) d h\right) f(g) d g=\int_{G} \alpha F(g) f(g) d g
$$

Surjectivity of $\alpha$ shows that $F$ can be chosen so that $\alpha F$ is identically 1 on the support of $f$. Then the integral of $f$ is 0 , as claimed, proving the well-definedness.
[6.7.2] Claim: Let $W$ be a smooth representation of a closed subgroup $H$ of $G$. Let $\rho=\delta_{H} / \delta_{G}$. The smooth dual of the compactly-induced $c-\operatorname{Ind}_{H}^{G} W$ is the induced representation $\operatorname{Ind}_{H}^{G}\left(\rho \cdot W^{\vee}\right)$, by the map described as follows. With $\langle$,$\rangle the duality pairing on W \times W^{\vee}$, for $F \in \operatorname{Ind}_{H}^{G}\left(\rho \cdot W^{\vee}\right)$, and $u: \operatorname{c-Ind}_{H}^{G} \rho$ as in the previous lemma, define

$$
\lambda_{F}(f)=u(g \longrightarrow\langle f(g), F(g)\rangle) \quad\left(\text { for all } f \in \mathrm{c}-\operatorname{Ind}_{H}^{G} W\right)
$$

Proof: First, claim that the $\mathbb{C}$-valued function $\varphi(g)=\langle f(g), F(g)\rangle$ is in c-Ind ${ }_{H}^{G} \rho$. Since $F$ is uniformly locally constant and $f$ is locally constant and compactly-supported left modulo $H, \varphi$ is locally constant and compactly-supported left modulo $H$. Further, from the definition of the action on the dual $W^{\vee}$,

$$
\varphi(h g)=\langle f(h g), F(h g)\rangle=\langle h \cdot f(g), \rho(h) \cdot h \cdot F(g)\rangle=\rho(h) \cdot\langle f(g), F(g)\rangle=\rho(h) \cdot \varphi(g)
$$

Therefore, these functionals constitute a $\mathbb{C}$-linear subspace of the smooth dual of c - $\operatorname{Ind}_{H}^{G} W$. That $F \rightarrow \lambda_{F}$ is a $G$-homomorphism is also apparent.

To see that $F \rightarrow \lambda_{F}$ is injective, take $F \in\left(\operatorname{Ind}_{H}^{G} \rho \cdot W^{\vee}\right)^{K}$ for some compact open $K$, take $x \in G$ so that $0 \neq F(x) \in W^{\vee}$. Let $w \in W$ with $\langle w, F(x)\rangle=1$. Without loss of generality, $w \in W^{K}$ : by properties [14.1] of Gelfand-Pettis integrals,

$$
\left\langle\int_{K} k \cdot w d k, F(x)\right\rangle=\int_{K}\langle k \cdot w, F(x)\rangle d k=\int_{K}\left\langle w, k^{-1} \cdot F(x)\right\rangle d k=\operatorname{meas}(K) \cdot\langle w, F(x)\rangle
$$

Define $f \in \mathrm{c}-\operatorname{Ind}_{H}^{G} W$ by

$$
f(g)=\int_{H} \operatorname{ch}_{x K}(h g) \cdot h \cdot w d h
$$

Again using [14.1] to move the integration outside the pairing $\langle$,$\rangle ,$

$$
\langle f(g), F(g)\rangle=\int_{H} \operatorname{ch}_{x K}(h g) \cdot\langle h \cdot w, F(g)\rangle d h=\int_{H} \rho(h)^{-1} \cdot\left(\operatorname{ch}_{x K}(h g) \cdot\langle h \cdot w, F(h g)\rangle\right) d h
$$

This expresses $g \rightarrow\langle f(g), F(g)\rangle \in \mathrm{c}-\operatorname{Ind}_{H}^{G} \rho$ as an image of the averaging map of [6.7.1], so

$$
\begin{aligned}
\lambda_{F}(f) & =u(g \rightarrow\langle f(g), F(g)\rangle)=\int_{G} \operatorname{ch}_{x K}(g) \cdot\langle w, F(g)\rangle d g=\int_{x K}\langle w, F(g)\rangle d g \\
& =\left\langle w, \int_{x K} F(g) d g\right\rangle=\operatorname{meas}(x K) \cdot\langle w, F(x)\rangle=\operatorname{meas}(x K) \cdot 1 \neq 0
\end{aligned}
$$

Prove surjectivity by proving surjectivity to each $\left(\left(c-\operatorname{Ind}_{H}^{G} W\right)^{*}\right)^{K}$, for compact open subgroups $K$. The quotient $(H \backslash G) / K$ is discrete since $K$ is open, and Hausdorff since $K$ is closed. Fix a set of representatives

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$x_{i}$ for $H \backslash G / K$, and let $f_{i}$ be the characteristic function of $x_{i} K \subset G$. For $\lambda \in\left(\left(c-\operatorname{Ind}_{H}^{G} W\right)^{K}\right)^{*}$, define a smooth functional $\lambda_{i} \in W^{\vee}$ by

$$
\lambda_{i}(w)=\lambda\left(\int_{H} f_{i}(h) \cdot h \cdot w d h\right) \quad(\text { for } w \in W)
$$

Define $F$ piecewise by $F\left(h x_{i} \theta\right)=\rho(h) \lambda_{i}$, so $\lambda_{F}(f)=u(g \rightarrow\langle f(g), F(g)\rangle)$ is the original $\lambda$.

### 6.8 An instance of Frobenius reciprocity

The following simplest instance of the fundamental adjunction is a precursor to the adjunction that is the assertion of Frobenius reciprocity proven just below:
[6.8.1] Claim: (Cartan-Eilenberg adjunction) For $\mathbb{Z}$-modules $A, B, C$,

$$
\operatorname{Hom}_{\mathbb{Z}}\left(A, \operatorname{Hom}_{\mathbb{Z}}(B, C)\right) \approx \operatorname{Hom}_{\mathbb{Z}}\left(A \otimes_{\mathbb{Z}} B, C\right)
$$

with $\varphi_{\Phi} \leftarrow \Phi$ by $\varphi_{\Phi}(a)(b)=\Phi(a \otimes b)$, and $\varphi \rightarrow \Phi_{\varphi}$ by $\Phi_{\varphi}(a \otimes b)=\varphi(a)(b)$.
Proof: Given $\Phi \in \operatorname{Hom}_{\mathbb{Z}}(A \otimes B, C)$, certainly $\varphi_{\Phi}(a)(b)=\Phi(a \otimes b)$ is immediately well-defined. Oppositely, the universal property of tensor products produces a unique linear map $A \otimes B \rightarrow C$ for each bilinear $A \times B \rightarrow C$. Applied to $a \times b \rightarrow \varphi(a)(b)$ produces a well-defined $\Phi_{\varphi} \in \operatorname{Hom}(A \otimes B, C)$ by $\Phi_{\varphi}(a \otimes b)=\varphi(a)(b)$. ///

The duality of compact-induced and induced proven in the previous section gives

$$
\left(\mathrm{c}-\operatorname{Ind}_{H}^{G} W\right)^{\vee}=\operatorname{Hom}_{G}\left(\mathrm{c}-\operatorname{-ind}_{H}^{G} W, \mathbb{C}\right) \approx \operatorname{Ind}_{H}^{G}\left(\rho \cdot W^{\vee}\right)
$$

However, to apply such ideas in the present context, there would be several technical complications: tensor products of smooth bi-modules over non-commutative rings without units, even with sufficiently idempotents as $\mathcal{H}$, are not as simply-behaved as over $\mathbb{Z}$. These complications cannot be avoided in the following section, but we can prove Frobenius reciprocity more directly:
[6.8.2] Theorem: (Frobenius Reciprocity) There is a natural $\mathbb{C}$-vectorspace isomorphism

$$
\operatorname{Hom}_{G}\left(V, \operatorname{Ind}_{H}^{G} W\right) \longrightarrow \operatorname{Hom}_{H}\left(\operatorname{Res}_{H}^{G} V, W\right) \quad\left(\text { by } \Phi \rightarrow \varphi_{\Phi} \text { where } \varphi_{\Phi}(v)=\Phi(v)(1)\right)
$$

The inverse is $\Phi_{\varphi} \leftarrow \varphi$ where $\Phi_{\varphi}(v)(g)=\varphi(g \cdot v)$.
Proof: Once the formula for the inverse is conceived, the several things to be checked are fairly straightforward. The $H$-homomorphism property of $\varphi_{\Phi}$ follows from

$$
h \cdot \varphi_{\Phi}(v)=h \cdot \Phi(v)(1)=\Phi(v)(1 \cdot h)=\Phi(h \cdot v)(1)=\varphi_{\Phi}(h \cdot v)
$$

The $G$-homomorphism property of $\Phi_{\varphi}$ follows from

$$
\left(g \cdot \Phi_{\varphi}(v)\right)(x)=\Phi_{\varphi}(v)(x g)=\varphi(x g \cdot v)=\varphi(x \cdot(g \cdot v))=\Phi_{\varphi}(g \cdot v)(x) \quad(\text { for } g, x \in G)
$$

That the two maps are mutual inverses is easy in one direction:

$$
\varphi_{\Phi_{\varphi}}(v)=\Phi_{\varphi}(v)(1)=\varphi(1 \cdot v)=\varphi(v)
$$

and in the other direction

$$
\Phi_{\varphi_{\Phi}}(v)(x)=\varphi_{\Phi}(x \cdot v)=\Phi(x \cdot v)(1)=x \cdot \Phi(v)(1)=\Phi(v)(1 \cdot x)=\Phi(v)(x)
$$

This proves Frobenius reciprocity in this situation.

### 6.9 Induction in stages

The description just below of compactly-induced representations c- $\operatorname{Ind}_{H}^{G} W$ as tensor products $\mathcal{H}_{G} \otimes_{\mathcal{H}_{H}} W$, suggests the possibility of inducing in stages, from the associativity of tensor products.

First, we describe a purely algebraic context which encourages optimism about the more complicated situation at hand. Consider not-necessarily commutative rings with a copy of $\mathbb{C}$ in their centers, and $1_{\mathbb{C}}$ is the multiplicative unit $1_{R}$ of $R$. All $R$-modules will be unital, in the sense that $1_{R} \cdot v=v$ for all $v \in V$. The tensor product $M \otimes_{R} N$ of a right $R$-module $M$ and left $R$-module $N$ is the quotient of $M \otimes_{\mathbb{C}} N$ by the submodule generated by all expressions

$$
m \cdot r \otimes n-m \otimes r \cdot n
$$

In general, this tensor product is no longer an $R$-module. However, when $M$ is both a right $R$-module and a left $S$-module for another ring $S$, that is, is an $S, R$-bimodule, there does remain the left multiplication by $S$ on the tensor product over $R$, namely, $s \cdot(m \otimes n)=(s \cdot m) \otimes n$. In particular, when $M=S$, for a left $R$-module $N$ the tensor product $S \otimes_{R} N$ is a left $S$-module. This is one notion of induced module. However, for present purposes, we only need a comparison:
[6.9.1] Claim: Let $R, S, T$ be $\mathbb{C}$-algebras, with $S$ a right $R$-algebra, and $T$ a right $S$-algebra. Make $T$ a right $R$-module by $t \cdot r=t \cdot\left(1_{S} \cdot r\right)$. Let $M$ be a left $R$-module. Then

$$
T \otimes_{S}\left(S \otimes_{R} M\right) \approx T \otimes_{R} M
$$

by

$$
t \otimes(s \otimes m) \longrightarrow t \cdot s \otimes m \quad \text { and } \quad t \otimes\left(1_{S} \otimes m\right) \longleftarrow t \otimes m
$$

Proof: The two maps are well-defined as maps

$$
\varphi: T \otimes_{\mathbb{C}} S \otimes_{\mathbb{C}} M \longrightarrow T \otimes_{S} M \quad \psi: T \otimes_{\mathbb{C}} M \longrightarrow T \otimes_{S}\left(S \otimes_{R} M\right)
$$

Thus, it suffices to show that these maps factor through the corresponding quotients. For $t \in T, s, s^{\prime} \in S$, $r \in R$, and $m \in M$,

$$
\varphi\left(t \otimes\left(s s^{\prime} \otimes m\right)-t \cdot s \otimes\left(s^{\prime} \otimes m\right)\right)=t \cdot\left(s s^{\prime}\right) \otimes m-(t \cdot s) \cdot s^{\prime} \otimes m=0
$$

by associativity of the right action of $S$ on $T$. Similarly,

$$
\varphi(t \otimes(s \otimes r \cdot m)-t \otimes(s \cdot r \otimes m))=(t \cdot s) \cdot r \otimes m-t \cdot(s \cdot r) \otimes m=0
$$

by the associativity

$$
(t \cdot s) \cdot r=(t \cdot s) \cdot\left(1_{S} \cdot r\right)=t \cdot\left(s \cdot\left(1_{S} \cdot r\right)\right)=t \cdot\left(\left(s \cdot 1_{S}\right) \cdot r\right)=t \cdot(s \cdot r)
$$

In the other direction,

$$
\psi(t \otimes r \cdot m-t \cdot r \otimes m)=t \otimes 1_{S} \otimes r \cdot m-t \cdot\left(1_{S} \cdot r\right) \otimes m=\left(t \otimes 1_{S}\right) \cdot r \otimes m-t \cdot\left(1_{S} \cdot r\right) \otimes m=0
$$

again by the associativity

$$
\left(t \cdot 1_{S}\right) \cdot r=\left(t \cdot 1_{S}\right) \cdot\left(1_{S} \cdot r\right)=t \cdot\left(1_{S} \cdot\left(1_{S} \cdot r\right)\right)=t \cdot\left(1_{S} \cdot r\right)=t \cdot\left(1_{S} \cdot r\right)
$$

Thus, the two maps are mutual inverses.
The intention is to use the same idea in application to Hecke algebras $\mathcal{H}_{B}, \mathcal{H}_{P}, \mathcal{H}_{G}$ for closed subgroups $B \subset P \subset G$, with suitable right $\mathcal{H}_{P}$ structure on $\mathcal{H}_{G}$, and right $\mathcal{H}_{B}$ structure on $\mathcal{H}_{P}$. These are rings without units, but with sufficiently many idempotents.

## 6. Action of $G$ on function spaces on $G$

We first describe relevant right $\mathcal{H}_{H}$-module structure on $\mathcal{H}_{G}$, for closed subgroup $H$ of $G$. Let $\delta_{H}$ be the modular function on $H$. Give $\mathcal{H}_{G}=C_{c}^{\infty}(G)$ a right $\mathcal{H}_{H}$-module structure by

$$
(f \cdot \eta)(g)=\int_{H} f(h g) \eta(h) \delta_{H}(h) d h \quad\left(\text { with } \eta \in \mathcal{H}_{H} \text { and } f \in \mathcal{H}_{G}\right)
$$

The insertion of the modular function is a normalization choice which becomes sensible in hindsight, just below. The action of $h \in H$ on the argument of $f$ is literally left multiplication in $G$, but the $H$-structure or $\mathcal{H}_{H}$-module structure should be notated on the right, to have associativity

$$
\left(f \cdot \eta_{1}\right) \cdot \eta_{2}=f \cdot\left(\eta_{1} * \eta_{2}\right) \quad\left(\text { with } \eta_{1}, \eta_{2} \in \mathcal{H}_{H} \text { and } f \in \mathcal{H}_{G}\right)
$$

To check this associativity:

$$
\begin{gathered}
\left(\left(f \cdot \eta_{1}\right) \cdot \eta_{2}\right)(g)=\int_{H} \int_{H} f(x y g) \eta_{1}(x) \delta(x) \eta_{2}(y) \delta_{H}(y) d x d y \\
=\int_{H} \int_{H} f(x g) \eta_{1}\left(x y^{-1}\right) \delta\left(x y^{-1}\right) \eta_{2}(y) \delta_{H}(y) d x d y \\
=\int_{H} f(x g)\left(\int_{H} \eta_{1}\left(x y^{-1}\right) \eta_{2}(y) d y\right) \delta_{H}(x) d x=\left(f \cdot\left(\eta_{1} * \eta_{2}\right)\right)(g)
\end{gathered}
$$

A smooth representation $H \times W \rightarrow W$ is a smooth $\mathcal{H}_{H}$-module as in [6.6.4]. The tensor product $\mathcal{H}_{G} \otimes_{\mathcal{H}_{H}} W$ is $\mathcal{H}_{G} \otimes_{\mathbb{C}} W$ modulo all relations

$$
(f \cdot \eta) \otimes w=f \otimes(\eta \cdot w)
$$

and has left $\mathcal{H}_{G}$-module structure

$$
\zeta \cdot(f \otimes w)=(\zeta * f) \otimes w \quad\left(\text { convolution in } \mathcal{H}_{G}\right)
$$

[6.9.2] Theorem: We have an $\mathcal{H}_{G}$-isomorphism $\mathcal{H}_{G} \otimes_{\mathcal{H}_{H}} W \approx{\mathrm{c}-\operatorname{Ind}_{H}^{G} W \text { by }}$

$$
\beta(f \otimes w)(g)=\int_{H} f(h g) \cdot h^{-1} w d h \quad \quad\left(\text { for } f \in \mathcal{H}_{G} \text { and } w \in W\right)
$$

Proof: To see that $\beta$ commutes with the $\mathcal{H}_{G}$ action, that is, that $\beta(\zeta * f \otimes w)=\zeta * \beta(f \otimes w)$ for $\zeta, f \in \mathcal{H}_{G}$, is just a change of order of integration:

$$
\begin{aligned}
& \beta(\zeta * f \otimes w)(x)=\beta\left(\left(y \rightarrow \int_{G} \zeta\left(y z^{-1}\right) f(z) d z\right) \otimes w\right)(x)=\int_{H}\left(\int_{G} \zeta\left(h x z^{-1}\right) f(z) d z\right) \cdot h^{-1} w d h \\
& =\int_{G} \int_{H} \zeta\left(h x z^{-1}\right) \cdot h^{-1} w d h f(z) d z=\int_{G} \beta\left(y \rightarrow \zeta\left(y z^{-1}\right) \otimes w\right)(x) f(z) d z=(\zeta * \beta(f \otimes w))(x)
\end{aligned}
$$

as claimed. A change of variables in the integral shows that the image $\beta(f \otimes w)$ lies in the indicated compact-induced representation space:

$$
\begin{gathered}
\beta(f \otimes w)(h g)=\int_{H} f\left(h^{\prime} h g\right) \cdot\left(h^{\prime}\right)^{-1} w d h^{\prime}=\int_{H} f\left(h^{\prime} g\right) \cdot\left(h^{\prime} h^{-1}\right)^{-1} w d h^{\prime}=\int_{H} f\left(h^{\prime} g\right) \cdot h \cdot\left(h^{\prime}\right)^{-1} w d h^{\prime} \\
=h \cdot \int_{H} f\left(h^{\prime} g\right) \cdot\left(h^{\prime}\right)^{-1} w d h^{\prime}=h \cdot \beta(f \otimes w)(g)
\end{gathered}
$$

To show that the map factors through the tensor product over $\mathcal{H}_{H}$ we must show that

$$
\beta(f \cdot \eta \otimes v)=\beta(f \otimes \eta \cdot v) \quad\left(\text { for } \eta \in \mathcal{H}_{H}\right)
$$

Here the normalization by insertion of the modular function will play a role. Take $\eta \in \mathcal{H}_{H}, w \in W$, and $f \in \mathcal{H}_{G}$ :

$$
\begin{aligned}
& \beta(f \cdot \eta \otimes w)(g)=\beta\left(x \rightarrow \int_{H} f(h x) \eta(h) \delta_{H}(h) d h \otimes w\right)(g)=\int_{H}\left(\int_{H} f(h y g) \eta(h) \delta_{H}(h) d h\right) y^{-1} w d y \\
& \quad=\int_{H} \eta(h)\left(\int_{H} f(h y g) y^{-1} w d y\right) \delta_{H}(h) d h=\int_{H} \eta(h)\left(\int_{H} f(y g) \delta_{H}\left(h^{-1}\right) y^{-1} h w d y\right) \delta_{H}(h) d h
\end{aligned}
$$

by replacing $y$ by $h^{-1} y$. The $\delta_{H}\left(h^{-1}\right)$ and $\delta_{H}(h)$ cancel, giving

$$
\begin{gathered}
\int_{H} \eta(h)\left(\int_{H} f(y g) y^{-1} h w d y\right) d h=\int_{H} f(y g) y^{-1} \cdot\left(\int_{H} \eta(h) h w d h\right) d y \\
=\int_{H} f(y g) y^{-1} \cdot(\eta \cdot w) d y=\beta(f \otimes \eta \cdot w)(g)
\end{gathered}
$$

Thus, $\beta$ factors through a map $\gamma: \mathcal{H}_{G} \otimes_{\mathcal{H}_{H}} W \rightarrow \mathrm{c}-\operatorname{Ind}_{H}^{G} W$.
To make an inverse map, make an inverse on right $K$-fixed elements for each compact open subgroup $K$ of $G$. Given $K$, fix representatives $\left\{x_{i}\right\}$ for $H \backslash G / K$, and let $\mathrm{ch}_{i}$ be the characteristic function of $x_{i} K$. Let $q: \mathcal{H}_{G} \otimes_{\mathbb{C}} W \rightarrow \mathcal{H}_{G} \otimes_{\mathcal{H}_{H}} W$ be the quotient map. For $F \in \mathrm{c}-\operatorname{Ind}_{H}^{G} W$, put

$$
\Phi(F)=q\left(\sum_{i} \operatorname{ch}_{i} \otimes F\left(x_{i}\right)\right) \in \mathcal{H}_{G} \otimes_{\mathcal{H}_{H}} W
$$

This is the inverse to $\gamma$.
[6.9.3] Corollary: Let $B \subset P \subset G$ be closed subgroups of a totally disconnected group $G$, and $W$ a smooth representation of $B$. Then inducing $W$ from $B$ to $G$ produces the same outcome as inducing from $B$ to $P$ and then from $P$ to $G$ :

$$
{\mathrm{c}-\operatorname{Ind}_{B}^{G} W} \mathrm{cc}_{\mathrm{c}-\operatorname{Ind}_{P}^{G}}\left(\mathrm{c}-\operatorname{Ind}_{B}^{P} W\right)
$$

Proof: Grant for a moment that for a class of rings without units containing Hecke algebras $\mathcal{H}_{G}$, $\mathcal{H}_{P}$, and $\mathcal{H}_{B}$, and and containing smooth representations $V$ of them,

$$
T \otimes_{S}\left(S \otimes_{R} V\right) \approx T \otimes_{R} V
$$

The theorem gives

$$
{\mathrm{c}-\operatorname{Ind}_{B}^{G} W}^{2} \mathcal{H}_{G} \otimes_{\mathcal{H}_{B}} V \approx \mathcal{H}_{G} \otimes_{\mathcal{H}_{P}}\left(\mathcal{H}_{P} \otimes_{\mathcal{H}_{B}} V\right) \approx \mathrm{c}-\operatorname{Ind}_{P}^{G}\left(\mathrm{c}-\operatorname{-ind}_{B}^{P} W\right)
$$

The collapsing of tensor products does hold for idempotented rings $R, S, T$ where $T$ is a smooth right $S$ module, and $S$ is a smooth right $R$-module. A ring $R$ is idempotented when, for every finite $X \subset R$, there is an idempotent element $e \in R$ so that $e x=x=x e$ for all $x \in X$. This property holds for these Hecke algebras, using idempotents $e_{K}$ as in the proof of [6.6.4]. Smoothness of a module $V$ over an idempotent ring $R$ is that for every finite $X \subset V$ there is an idempotent $e \in R$ such that $e x=x=x e$ for all $x \in X$.

Then $T$ has a smooth right $R$-module structure $t \cdot r=t \cdot\left(e^{\prime} \cdot t\right)$ where $e^{\prime} \in S$ is an idempotent in $S$ fixing $t$, invoking smoothness of $T$ over $S$. Unlike the case of rings with units, we must check that this is well-defined. Let $e_{1}, e_{2}$ be idempotents in $S$ both fixing $t$, let $e^{\prime}$ be an idempotent in $S$ such that $e^{\prime} e_{1}=e_{1} e^{\prime}=e_{1}$ and $e^{\prime} e_{2}=e_{2} e^{\prime}=e_{2}$. Then

$$
t \cdot\left(e_{1} \cdot r\right)=t \cdot\left(e_{1} e^{\prime} \cdot r\right)=\left(t \cdot e_{1}\right) \cdot\left(e^{\prime} \cdot r\right)=t \cdot\left(e^{\prime} \cdot r\right)=\left(t \cdot e_{2}\right) \cdot\left(e^{\prime} \cdot r\right)=t \cdot\left(e_{2} e^{\prime} \cdot r\right)=t \cdot\left(e_{2} \cdot r\right)
$$

which gives the well-definedness. For an idempotent $e \in R$ such that $e^{\prime} \cdot e=e^{\prime}$, we have smoothness $t \cdot e=t \cdot\left(e^{\prime} \cdot e\right)=t \cdot e^{\prime}=t$. The isomorphisms are

$$
t \otimes s \otimes v \longrightarrow t \cdot s \otimes v \quad \text { and } \quad t \otimes e^{\prime} \otimes \longleftarrow t \otimes v
$$

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for any idempotent $e \in S$ such that $t \cdot e=t$. Certification of the well-definedness of the second map is similar to the previous argument: let $t \cdot e_{1}=t=t \cdot e_{2}$ for two idempotents $e_{1}, e_{2} \in S$, and let $e$ be another idempotent in $S$ such that $e_{1} e=e_{1}$ and $e_{2} e=e_{2}$. It suffices to compute in $T \otimes_{S} S$ :

$$
t \otimes e_{1}=t \otimes e_{1} e=t \cdot e_{1} \otimes e=t \otimes e=t \cdot e_{2} \otimes e=t \otimes e_{2} e=t \otimes e_{2}
$$

This gives the well-definedness.

### 6.10 Representations of compact $G / Z$

We still consider totally disconnected $G$ and smooth representations. The general case of representations of compact groups on topological vector spaces is treated in [9.C].

Let $Z$ be a closed subgroup of $G$ inside the center of $G$, and suppose that $G / Z$ is compact. Consider representations $V$ with central character $\omega: Z \rightarrow \mathbb{C}^{\times}$, i.e., so that $z \cdot v=\omega(z) v$ for all $v \in V$ and $z \in Z$. The simple situation that $G$ is compact and $Z=\{1\}$ is already useful.
[6.10.1] Proposition: Every finitely-generated smooth representation $V$ of $G$ with central character $\omega$ is finite-dimensional.
Proof: Take a compact open subgroup $K$ small enough so that a (finite) set $X$ of generators for $V$ lies inside $V^{K}$. Let $Y$ be a choice of a set of representatives for $G / Z K$; since $G / Z$ is compact, $Y$ is finite. The set of all vectors $g \cdot v$ with $v \in X$ and $g \in G$ is contained in the span of the finite set of vectors $y \cdot x$ for $y \in Y$ and $x \in X$.
[6.10.2] Corollary: Every irreducible smooth representation of $G$ having a central character for $Z$ is finitedimensional
[6.10.3] Proposition: Let $f: M \rightarrow N$ be a surjective $G$-homomorphism of two smooth $G$-representation spaces, both with central character $\omega$. Suppose there is a small-enough compact open subgroup $K$ of $G$ so that $M^{K}=M$ and $N^{K}=N$, as for $M, N$ finitely-generated. There is a $G$-homomorphism $\varphi: N \rightarrow M$ so that $f \circ \varphi$ is the identity $\operatorname{map}^{\operatorname{id}}{ }_{N}$ on $N$.
Proof: Let $n$ be the cardinality of $G / Z K$. Let $\psi: N \rightarrow M$ be any $k$-vectorspace map so that $f \circ \psi=\operatorname{id}_{N}$ : take any $k$-vectorspace $N_{1}$ in $M$ complementary to the kernel of $f$, and let $\psi$ be the inverse of the restriction of $f$ to $N_{1}$. Define

$$
\varphi v=\frac{1}{n} \sum_{h \in G / Z K} h^{-1} \psi h v
$$

The hypotheses assure that this $\varphi$ is independent of the choice of representatives for $G / Z K$, and it is immediate (by changing variables in the sum) that this averaged version of $\psi$ is a $G$-homomorphism providing a one-sided inverse to $f$.
[6.10.4] Corollary: Let $f: M \rightarrow N$ be an injective $G$-homomorphism of two $G$-representation spaces, both with central character $\omega$ (for $Z$ ). Suppose that there is a compact open subgroup $K$ of $G$ so that $M^{K}=M$ and $N^{K}=N$ (as is the case if $M, N$ are finitely-generated). There is a $G$-homomorphism $\varphi: M \rightarrow N$ so that $\varphi \circ f$ is the identity map $\operatorname{id}_{M}$ on $M$. In particular, every $G$-submodule of $N$ has a complementary submodule.
Proof: Let $Q=N / f M$ be the quotient, and $q: N \rightarrow Q$ the quotient map. The previous proposition yields $\psi: Q \rightarrow N$ so that $q \circ \psi=\operatorname{id}_{Q}$. Since $N=f M \oplus \psi Q$ and $f M \approx M, N / \psi Q$ is naturally isomorphic to $M$, and the composition

$$
N \longrightarrow N / \psi \approx M
$$

is the desired $\varphi$.
[6.10.5] Corollary: (Complete Reducibility) Every smooth representation of $G$ with central character $\omega$ (for $Z)$ is a direct sum of irreducible smooth representations, each with central character $\omega$ for $Z$.
Proof: This will follow from the previous and from Zorn's Lemma.

First, show that a finite-dimensional smooth representation $M$ contains a non-zero irreducible. Since $M$ is finite-dimensional it is finitely-generated, so has an irreducible quotient $q: M \rightarrow Q$. By the above discussion, there is a $G$-subspace $M^{\prime}$ of $M$ so that as $G$-spaces $M \approx M^{\prime} \oplus Q$. Thus, $M$ contains the irreducible $Q$.

Let $M=\oplus_{\alpha} M_{\alpha}$ be a maximal direct sum of (necessarily finite-dimensional) irreducibles inside $N$, and suppose that $M \neq N$. Take $x \in N$ not lying in $M$, and let $X$ be the $G$-subspace of $N$ generated by $x$. Then $X$ is finitely-generated, so is finite-dimensional, and has a non-zero irreducible quotient $Q$. From above, there is a copy $Q^{\prime}$ of $Q$ inside $X$ and $X=Q^{\prime} \oplus X^{\prime}$ for some $X^{\prime}$. By the maximality of $M, Q$ must be inside $M$ already. Apply the same argument to $X^{\prime}$, so by induction on dimension conclude that $X$ was 0 .

### 6.11 Gelfand-Kazhdan criterion

Subgroups $H$ of groups $G$ with the property that restrictions of most or all irreducibles $V$ of $G$ to $H$ are multiplicity-free, that is, so that $\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{H}\left(W, \operatorname{Res}_{H}^{G} V\right) \leq 1$ for most or all irreducibles $W$ of $H$, are called (strongly) Eulerian, or $G, H$ is a (strong) Gelfand pair.

It is important to know that induced representations are multiplicity-free, meaning contain at most one copy of a given irreducible representation, whenever this is the case, to produce Euler factorization of global integrals, for example. The idea of the proof, useful already in the representation theory of finite groups, is that if no irreducible occurs twice inside a representation, then the endomorphism algebra should be commutative, and vice-versa. Unfortunately, this principle is not quite valid more generally, for infinitedimensional representations of non-finite groups. After sufficient adaptations are made, we have the GelfandKazhdan criterion below. As in the Gelfand criterion for commutativity of the spherical Hecke algebra [2.4.5] (echoed in [3.5.4]), the criterion for multiplicity-free-ness depends upon identifying an involutive antiautomorphism to interchange the order of factors, but which nevertheless acts as the identity on suitable subalgebras.

Let $G$ be a unimodular, totally disconnected group. The space $\mathcal{D}=\mathcal{D}(G)$ of test functions on $G$ is the space of a compactly-supported, locally constant complex-valued functions on $G$. As a colimit (that is, direct limit) of finite-dimensional complex vector spaces, this space has a uniquely determined topology. The space $\mathcal{D}^{*}=\mathcal{D}(G)^{*}$ of distributions on $G$ is the space of continuous complex-linear functionals on $\mathcal{D}$.

Let $H$ be a closed, unimodular subgroup of $G$. For a one-dimensional complex representation $\psi$ of $H$,

$$
\operatorname{Ind}_{H}^{G} \psi=\left\{\begin{array}{l}
\mathbb{C} \text {-valued functions } f \text { on } G, \text { uniformly locally constant, so that } f(h g)=\psi(h) f(g) \\
\text { for all } g \in G \text { and } h \in H
\end{array}\right\}
$$

The case of trivial $\psi$ is already interesting. Let $\mathbb{C}$ denote the trivial representation of $G$ or of $H$. Say $(G, H)$ is a Gelfand pair, or equivalently, $H$ is an Euler subgroup of $G$ if, for all irreducible admissible representations $\pi$ of $G$,

$$
\operatorname{dim} \operatorname{Hom}_{G}\left(\pi, \operatorname{Ind}_{H}^{G} \mathbb{C}\right) \times \operatorname{dim} \operatorname{Hom}_{G}\left(\pi^{\vee}, \operatorname{Ind}_{H}^{G} \mathbb{C}\right) \leq 1
$$

where $\pi^{\vee}$ is the contragredient of $\pi$. By Frobenius Reciprocity, this is equivalent to

$$
\operatorname{dim} \operatorname{Hom}_{H}\left(\operatorname{Res}_{H}^{G} \pi, \mathbb{C}\right) \times \operatorname{dim} \operatorname{Hom}_{H}\left(\operatorname{Res}_{H}^{G} \pi^{\vee}, \mathbb{C}\right) \leq 1
$$

An involutive anti-automorphism $\sigma$ on a group $G$ is a bijection $G \rightarrow G$ so that $(g h)^{\sigma}=h^{\sigma} g^{\sigma}$. The action of $\sigma$ on functions is by $f^{\sigma}(g)=f\left(g^{\sigma}\right)$, and on distributions $u$ by $u^{\sigma}(f)=u\left(f^{\sigma}\right)$ for $f \in \mathcal{D}$.
[6.11.1] Theorem: Let $\psi$ be a one-dimensional representation of a closed unimodular subgroup $H$ of $G$. Suppose $H$ is stabilized by an involutive anti-automorphism $\sigma$ of $G$. Put $\psi^{\sigma}(h)=\psi\left(h^{\sigma}\right)$. Suppose that $u^{\sigma}=u$ for all distributions $u$ on $G$ possessing equivariance

$$
u\left(L_{h} \eta\right)=\psi(h) \cdot u(\eta) \quad \text { and } \quad u\left(R_{h} \eta\right)=\psi^{\sigma}(h)^{-1} \cdot u(\eta) \quad(\text { for } \eta \in \mathcal{D}(G))
$$

Then

$$
\operatorname{dim} \operatorname{Hom}_{G}\left(\pi, \operatorname{Ind}_{H}^{G} \psi\right) \times \operatorname{dim} \operatorname{Hom}_{H}\left(\operatorname{Res}_{H}^{G} \pi^{\vee}, \psi^{\sigma}\right) \leq 1
$$

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Proof: By Frobenius reciprocity $\operatorname{Hom}_{G}\left(\pi, \operatorname{Ind}_{H}^{G} \psi\right) \approx \operatorname{Hom}_{H}\left(\operatorname{Res}_{H}^{G} \pi, \psi\right)$, existence of non-trivial $G$-maps $\pi \rightarrow \operatorname{Ind}_{H}^{G} \psi$ and $\pi^{\vee} \rightarrow \operatorname{Ind}_{H}^{G} \psi^{\sigma}$ is equivalent to existence of non-zero $H$-homomorphisms $s: \pi \longrightarrow \psi$ and $t: \pi^{\vee} \longrightarrow \psi^{\sigma}$. We obtain $G$-homomorphisms $S: \mathcal{D} \longrightarrow \pi^{\vee}$ and $T: \mathcal{D} \longrightarrow\left(\pi^{\vee}\right)^{\vee}$ by integrating:

$$
(S \eta)(v)=\int_{G} \eta(x) s(x \cdot v) d x \quad(T \eta)(\lambda)=\int_{G} \eta(x) t(x \cdot \lambda) d x \quad(\text { for } \eta \in \mathcal{D}(G), v \in \pi, \lambda \in \check{\pi})
$$

The assumed admissibility of $\pi$ implies that $\pi$ is reflexive, that is, that $\left(\pi^{\vee}\right)^{\vee} \approx \pi$. By direct computation, right translation $R_{g}$ by $g \in G$, and left translation $L_{h}$ by $h \in H$ interact with $S$ and $T$ by

$$
S\left(R_{g} \eta\right)=g \cdot(S \eta) \quad T\left(R_{g} \eta\right)=g \cdot(T \eta) \quad S\left(L_{h} \eta\right)=\psi(h) \cdot S \eta \quad T\left(L_{H} \eta\right)=\psi^{\sigma}(h) \cdot T \eta
$$

The first assertion, for example, is verified as follows: for $v \in \pi$,

$$
S\left(R_{g} \eta\right)(v)=\int_{G} \eta(x g) s(x \cdot v) d x=\int_{G} \eta(x) s\left(x g^{-1} \cdot v\right) d x=S \eta\left(g^{-1} v\right)
$$

by replacing $x$ by $x g^{-1}$. The last expression is simply the contragredient action of $g$, that is, on $\pi^{\vee}$. The left $H$-invariance follows by

$$
S\left(L_{h} \eta\right)(v)=\int_{G} \eta\left(h^{-1} x\right) s(x \cdot v) d x=\int_{G} \eta(x) s(h x \cdot v) d x=\int_{G} \eta(x) \psi(h) s(x \cdot v) d x=S(\eta)(v)
$$

by replacing $x$ by $h x$ and then invoking the $H$-equivariance of $s$. The corresponding assertions for $T$ are proven similarly. That is, both $S$ and $T$ are left $H$-equivariant as indicated, and are right $G$-equivariant, giving $G$-homomorphisms from $\mathcal{D}(G)$ (with right regular representation) to $\pi$ and $\pi^{\vee}$.

Let $\langle\rangle:, \pi \times \pi^{\vee} \rightarrow \mathbb{C}$ be the canonical complex-bilinear pairing $\langle v, \lambda\rangle=\lambda(v)$. Let the induced complexlinear map on the tensor product be $\beta: \pi \otimes_{\mathbb{C}} \pi^{\vee} \longrightarrow \mathbb{C}$. Define

$$
B=\beta \circ(T \otimes S): \mathcal{D}(G) \otimes \mathcal{D}(G) \longrightarrow \pi \otimes \pi^{\vee} \longrightarrow \mathbb{C}
$$

Note the reversal of $S$ and $T$. The functional $B$ is in the space of distributions $\mathcal{D}(G \times G)^{*}$, is left $\left(H, \psi^{\sigma}\right) \times(H, \psi)$-equivariant and right $G^{\Delta}$-invariant, where $G^{\Delta}$ is the diagonal copy of $G$ in $G \times G$. Reversal of $\psi$ and $\psi^{\sigma}$ due to the reversal of $S$ and $T$.
[6.11.2] Lemma: With $B, t, S$ as above, for $\eta, \varphi$ in $\mathcal{D}(G), B(\eta \otimes \varphi)=t(S(\varphi * \eta))$.
Proof: Apart from an issue of interchange of integration and linear operators, this is a direct computation:

$$
B(\eta \otimes \varphi)=T \eta(S \varphi)=\int_{G} \eta(x) t(x \cdot S \varphi) d x=\int_{G} \eta(x) t\left(S\left(R_{x}^{-1} \cdot \varphi\right)\right) d x
$$

by the $G$-equivariance of $S$. Moving the integral inside $t \circ S$, this becomes

$$
(t \circ S)\left(\int_{G} \eta(x) R_{x}^{-1} \cdot \varphi d x\right)=(t \circ S)(\varphi * \eta)
$$

Exchange of integration and application of the operator $t \circ S$ is easily justified, since the indicated integral is actually a finite sum. More generally, $\mathcal{D}(G)$ is an an LF-space [13.10], so is quasi-complete [13.12], and Gelfand-Pettis integrals of compactly-supported continuous $\mathcal{D}(G)$-valued function $f$ exist, and, as in [14.1], for any continuous linear operator $L$ on $\mathcal{D}(G)$,

$$
L\left(\int_{G} f(x) d x\right)=\int_{G} L(f(x)) d x
$$

The desired exchange is a special case of this.
[6.11.3] Corollary: The distribution $u$ on $G$ defined by $u(\eta)=t(S(\eta))$ is left $H$-equivariant by $\psi$ and right $H$-equivariant by $\left(\psi^{\sigma}\right)^{-1}$, meaning that

$$
u\left(L_{h} \eta\right)=\psi(h) \cdot u(\eta) \quad \text { and } \quad u\left(R_{h} \eta\right)=\psi^{\sigma}(h)^{-1} \cdot u(\eta)
$$

Proof: Given $\eta \in \mathcal{D}(G)$ and given $h \in H$, we claim that there is $\varphi \in \mathcal{D}(G)$ so that $R_{h} \eta * \varphi=R_{h} \eta$. For example, for $K$ a small-enough compact open subgroup of $G$ so that $R_{h} \eta$ is left $K$-invariant, take $\varphi$ to be meas $(K)^{-1}$ on $K$ and 0 off $K$. Then

$$
u\left(R_{h} \eta\right)=(t \circ S)\left(R_{h} \eta\right)=(t \circ S)\left(\left(R_{h} \eta\right) * \varphi\right)=(t \circ S)\left(\eta * L_{h}^{-1} \varphi\right)=B\left(L_{h}^{-1} \varphi \otimes \eta\right)
$$

by the way that convolutions and translations interact. By the left $H$-equivariance of $B$ by $\psi^{\sigma}$ in its first argument,

$$
B\left(L_{h}^{-1} \varphi \otimes \eta\right)=\psi^{\sigma}(h)^{-1} \cdot B(\varphi \otimes \eta)
$$

Going back by the same procedure, $u\left(R_{h} \eta\right)=\psi^{\sigma}(h)^{-1} \cdot u(\eta)$. Even more simply, for $\varphi$ so that $\left(L_{h} \eta\right) * \varphi=L_{h} \eta$, we compute that

$$
u\left(L_{h} \eta\right)=B\left(\varphi \otimes L_{h} \eta\right)=\psi(h) \cdot B(\varphi \otimes \eta)=\psi(h) \cdot u(\eta)
$$

giving the equivariance.
As usual,
[6.11.4] Lemma: For $\eta, \varphi$ in $\mathcal{D}(G)$,

$$
(\eta * \varphi)^{\sigma}=\varphi^{\sigma} * \eta^{\sigma}
$$

Proof: For $g \in G$,

$$
\begin{aligned}
(\eta * \varphi)^{\sigma}(g) & =(\eta * \varphi)\left(g^{\sigma}\right)=\int_{G} \eta\left(g^{\sigma} x^{-1}\right) \varphi(x) d x=\int_{G} \eta\left(x^{-1}\right) \varphi\left(x g^{\sigma}\right) d x \\
& =\int_{G} \eta(x) \varphi\left(x^{-1} g^{\sigma}\right) d x=\int_{G} \eta\left(x^{\sigma}\right) \varphi\left(\left(g x^{-1}\right)^{\sigma}\right) d x
\end{aligned}
$$

replacing $x$ successively by $x g, x^{-1}$, and $x^{\sigma}$. This is

$$
\int_{G} \eta^{\sigma}(x) \varphi^{\sigma}\left(g x^{-1}\right) d x=\left(\varphi^{\sigma} * \eta^{\sigma}\right)(g)
$$

as claimed.
[6.11.5] Corollary: $B(\eta \otimes \varphi)=B\left(\varphi^{\sigma} \otimes \eta^{\sigma}\right)$.
Proof: $u\left(\eta^{\sigma}\right)=u(\eta)$ for all $\eta \in \mathcal{D}(G)$, so

$$
B(\eta \otimes \varphi)=u(\varphi * \eta)=u\left((\varphi * \eta)^{\sigma}\right)=u\left(\eta^{\sigma} * \varphi^{\sigma}\right)=B\left(\varphi^{\sigma} \otimes \eta^{\sigma}\right)
$$

[6.11.6] Corollary: For $\eta$ in $\mathcal{D}(G), T \eta=0$ implies $S\left(\eta^{\sigma}\right)=0$, and similarly $S \eta=0$ implies $T\left(\eta^{\sigma}\right)=0$.
Proof: For, $T \eta=0$, for all $\varphi$ in $\mathcal{D}(G)$

$$
0=\langle T \eta, S \varphi\rangle=B(\eta \otimes \varphi)=B\left(\varphi^{\sigma} \otimes \eta^{\sigma}\right)=\left\langle T \varphi^{\sigma}, S \eta^{\sigma}\right\rangle
$$

by the previous corollary. That is, $S \eta^{\sigma}$ gives the trivial linear functional on $\pi$, so must be 0 in $\pi^{\vee}$. The other assertion is similarly proven.

That is, $\operatorname{ker} T$ determines $\operatorname{ker} S$ and vice-versa.
Since $\pi$ is irreducible, by Schur's lemma the kernel of $S: \mathcal{D}(G) \rightarrow \pi$ determines $S$ uniquely up to a constant, and the same assertion holds for $T$. We can recover $s: \pi \rightarrow \mathbb{C}$ unambiguously from $S$. Given

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$v \in \pi$, let $\eta$ be meas $(K)^{-1}$ times the characteristic function of $K$, where $K$ is any sufficiently small compact open subgroup of $G$. Then

$$
(S \eta)(v)=\int_{G} \eta(x) s(x \cdot v) d x=s(v)
$$

That is, from $\operatorname{ker} S$ we recover $S$ uniquely up to a constant, and then recover $s$ uniquely up to a constant. The analogous assertion holds for $\operatorname{ker} T, T$, and $t$.

Then $t$ certainly determines $T$, which determines $\operatorname{ker} T$. From above, $\operatorname{ker} T$ determines $\operatorname{ker} S$, which (by the previous paragraph) determines $s$ up to a constant. We could have fixed $t$ and let $s$ be arbitrary, which would show that if the space of $t$ 's were positive-dimensional then the space of $s$ 's would be at most onedimensional. The symmetrical argument reversing the role of $s$ and $t$ goes through in the same manner, wherein we use the assumed admissibility of $\pi$ and, thus, of $\pi^{\vee}$. This proves the theorem.

## 6.A Appendix: action of compact abelian groups

Let $A$ be a compact, abelian topological group, and countably-based. We grant that $A$ has a translationinvariant measure, that is, a Haar measure. Let $|\cdot|$ be the corresponding norm on $L^{2}(A)$. The analogue of Fourier series expansion is
[6.A.1] Theorem: $L^{2}(A)$ is the completion of the direct $\operatorname{sum} \bigoplus_{\chi} \mathbb{C} \cdot \chi$ as $\chi$ ranges over continuous homomorphisms $\chi: A \rightarrow \mathbb{C}^{\times}$.
Proof: Certainly every $\chi$ is in $L^{2}(A)$. On the other hand, $a \in A$ acts on $F \in L^{2}(A)$ by translation $(a \cdot F)(b)=F(a+b)$, and every continuous $f$ on $A$ acts correspondingly on $F \in L^{2}(A)$ by an integral operator

$$
(f \cdot F)(b)=\int_{A} f(a) F(a+b) d a
$$

Replacing $a$ by $a-b$, this is

$$
(f \cdot F)(b)=\int_{A} f(a-b) F(a) d a
$$

expressing the map $F \rightarrow f \cdot F$ as a Hilbert-Schmidt operator with integral kernel $K(a, b)=f(a-b)$. HilbertSchmidt operators are compact [9.A.4]. For $f$ real-valued and even, in the sense that $f(-a)=f(a)$, the corresponding integral kernel is symmetric and real-valued, so gives a self-adjoint operator. The spectral theorem [9.A.6] for compact self-adjoint operators gives an eigenspace decomposition of $L^{2}(A)$ with respect to the operator given by $f$, and all eigenspaces are finite-dimensional except possibly the 0 -eigenspace. As usual, composition of two such operators is by the action of their convolution, as illustrated already in [2.4] for non-abelian groups: for $f, g$ real-valued in $C^{o}(A)$ and $F \in L^{2}(A)$,

$$
f \cdot(g \cdot F)=\int_{A} f(a) a \cdot\left(\int_{A} g(b) b \cdot F d b\right) d a=\int_{A} f(a)\left(\int_{A} g(b) a \cdot(b \cdot F) d b\right) d a
$$

The integrand $b \rightarrow g(b) b \cdot F$ of the inner integral is a continuous, compactly-supported, $L^{2}(A)$-valued function of $b \in A$, by the continuity of the action of $A$ on $L^{2}(A)$ and the continuity of $g$ on the compact $A$. Thus, the action of $a$ passes inside the vector-valued integral for general reasons [14.1]. By Fubini, this is

$$
\int_{A} \int_{A} f(a) g(b)(a+b) \cdot F d a d b=\int_{A} \int_{A} f(a-b) g(b) a \cdot F d a d b
$$

by replacing $a$ by $a-b$. This is

$$
\int_{A}\left(\int_{A} f(a-b) g(b) d b\right) a \cdot F d a=\int_{A}(f * g)(a) a \cdot F d a=(f * g) \cdot F
$$

Since the group is abelian, the convolution product is abelian:

$$
(f * g)(a)=\int_{A} f(a-b) g(b) d b=\int_{A} f(-b) g(a+b) d b=\int_{A} f(b) g(a-b) d b
$$

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by replacing $b$ by $b+a$ and then replacing $b$ by $-b$. Since the group is abelian, the Haar measure is invariant under $b \rightarrow-b$ as well. Thus, these self-adjoint compact operators commute with each other. Commuting operators preserve each others' eigenspaces: for $v$ in the $\lambda$-eigenspace for $T$, and for $S T=T S$,

$$
T(S v)=(T S) v=(S T) v=S(T v)=S(\lambda v)=\lambda \cdot S V
$$

Thus, $L^{2}(A)$ decomposes into simultaneous eigenspaces for all these operators. For that matter, the action $a \times F \longrightarrow a \cdot F$ stabilizes eigenspaces:

$$
\begin{aligned}
a \cdot(f \cdot F) & =a \cdot \int_{A} f(b) b \cdot F d b=\int_{A} a \cdot f(b) b \cdot F d b=\int_{A} f(b)(a+b) \cdot F d b \\
& =\int_{A} f(b)(b+a) \cdot F d b=\int_{A} f(b) b \cdot(a \cdot F) d b=f \cdot(a \cdot F)
\end{aligned}
$$

Again, the action of $a \in A$ passes inside the integral by Gelfand-Pettis theory [14.1].
The non-degeneracy result [14.1.5], that for every $0 \neq v \in V$ there is $\varphi \in C_{c}^{o}(G)$ such that $\varphi \cdot v \neq 0$, implies that the simultaneous 0-eigenspace for all the integral operators is trivial.

Since the integral operators are self-adjoint, distinct eigenspaces are mutually orthogonal: given two distinct eigenspaces, let real-valued $f \in C^{o}(A)$ be such that the two eigenvalues $\lambda, \mu \in \mathbb{R}$ are different. For $v, w$ in the respective eigenspaces, using the self-adjointness,

$$
\lambda \cdot\langle v, w\rangle=\langle f \cdot v, w\rangle=\langle v, f \cdot w\rangle=\langle v, \mu \cdot w\rangle=\mu \cdot\langle v, w\rangle
$$

We claim that each of these finite-dimensional spaces $V_{o}$ decomposes into simultaneous eigenspaces for $A$ itself. Again, as above, each is stabilized by the (unitary) action of the abelian group $A$. We do a descending induction on dimension. If $V_{o}$ is a simultaneous eigenspace for $A$, we are done. Otherwise, there is $a_{1} \in A$ such that $V_{o}$ decomposes properly into $a_{1}$-eigenspaces, noting that $a_{1}$ acts unitarily. Let $V_{1}$ be a proper $a_{1}$ eigenspace inside $V_{o}$. If $V_{1}$ is a simultaneous eigenspace for all $A$, we are done. Otherwise, take $a_{2} \in A$ such that $V_{1}$ decomposes properly into $a_{2}$-eigenspaces. Continuing, by the finite-dimensionality of $V_{o}$, the process must stop, producing a simultaneous eigenspace for $A$. That is, $L^{2}(A)$ has an orthogonal decomposition into simultaneous eigenspaces for $A$.

Let $V^{\lambda}$ be such a simultaneous eigenspace, where $a \cdot v=\lambda_{a} \cdot v$ for $\lambda_{a} \in \mathbb{C}$. The collection of eigenvalues $\lambda_{a}$ is a group homomorphism $A \rightarrow C^{\times}$:

$$
\lambda_{a} \cdot \lambda_{b} \cdot v=\lambda_{b} \cdot \lambda_{a} \cdot v=\lambda_{b} \cdot a \cdot v=a \cdot \lambda_{b} \cdot v=a \cdot b \cdot v=(a+b) \cdot v=\lambda_{a+b} \cdot v
$$

We claim that the $\lambda$-eigenspace $V^{\lambda}$ is simply all scalar multiples of $\lambda$ itself: for $F \in V^{\lambda}$,

$$
F(a)=F(0+a)=(a \cdot F)(0)=\left(\lambda_{a} \cdot F\right)(0)=\lambda_{a} \cdot F(0)
$$

That is, $F=F(0) \cdot \lambda$. Last, we show that such $\lambda$ is continuous: from the continuity of $A \times L^{2}(A) \rightarrow L^{2}(A)$, the restriction $A \times V^{\lambda} \rightarrow V^{\lambda}$ is continuous, with one-dimensional $V^{\lambda}=\mathbb{C} \cdot \lambda$. Thus, $a \rightarrow(a \cdot \lambda)=\lambda_{a} \cdot \lambda$ is continuous, so $a \rightarrow \lambda_{a}$ must be continuous. Thus, $L^{2}(A)$ is the completion of the orthogonal direct sum of $\mathbb{C} \cdot \chi$ with $\chi$ running over all continuous group homomorphisms $A \rightarrow \mathbb{C}^{\times}$.

Let compact, abelian $A$ act on a Hilbert space $V$ by unitary operators. For a continuous group homomorphism $\chi: A \rightarrow \mathbb{C}^{\times}$, the $\chi$-isotype $V^{\chi}$ in $V$ is

$$
V^{\chi}=\{v \in V: a \cdot v=\chi(a) \cdot v, \text { for all } a \in A\}
$$

[6.A.2] Claim: $V$ decomposes as the completion of the direct sum

$$
V=(\text { completion of }) \bigoplus_{\chi} V^{\chi} \quad\left(\chi \text { ranging over continuous homomorphisms } A \rightarrow \mathbb{C}^{\times}\right)
$$

## 6. Action of $G$ on function spaces on $G$

Proof: First, we claim that the orthogonal projection $P_{\chi}: V \rightarrow V^{\chi}$ is given by the integral operator

$$
P_{\chi} v=\int_{A} \bar{\chi}(a) a \cdot v d a
$$

with compact $A$ having total measure normalized to 1 . Indeed, the image is in $V^{\chi}$ : using properties of vector-valued integrals [14.1] to allow the action to pass inside the integral,

$$
\begin{gathered}
a \cdot P_{\chi} v=a \cdot \int_{A} \bar{\chi}(b) b \cdot v d b=\int_{A} \bar{\chi}(b) a \cdot(b \cdot v) d b=\int_{A} \bar{\chi}(b)(a+b) \cdot v d b=\int_{A} \bar{\chi}(b-a) b \cdot v d b \\
=\int_{A} \chi(a) \bar{\chi}(b) b \cdot v d b=; \chi(a) \cdot \int_{A} \bar{\chi}(b) b \cdot v d b=\chi(a) \cdot P_{\chi} v
\end{gathered}
$$

Next, $P_{\chi} \circ P_{\chi}=P_{\chi}$ : using properties of vector-valued integrals [14.1], and using the normalization that the total measure of $A$ is 1 ,

$$
\left.\left.\begin{array}{rl}
P_{\chi}\left(P_{\chi} v\right)= & \int_{A} \int_{A} \bar{\chi}(a) \bar{\chi}(b) a \cdot b \cdot v d a d b
\end{array}=\int_{A} \int_{A} \bar{\chi}(a) \bar{\chi}(b)(a+b) \cdot v d a d b=\int_{A} \int_{A} \bar{\chi}(a-b) \bar{\chi}(b) a \cdot v d a d b\right)=\int_{A} \int_{A} \bar{\chi}(a) \chi(b) \bar{\chi}(b) a \cdot v d a d b=\int_{A} \bar{\chi}(a) a \cdot v d a d b=\int_{A} \bar{\chi}(a) a \cdot v d a \cdot \int_{A} 1 d b\right]=\int_{A} \bar{\chi}(a) a \cdot v d a=P_{\chi} v .
$$

Last, we show that the orthogonal complement of all the images $P_{\chi} V$ is just $\{0\}$. Let $\left\{\eta_{i}\right\}$ be an approximate identity in $C^{o}(A)$, as in [6.5] and [14.1], meaning that $0 \leq \eta_{i}(a) \leq 1$ for all $i$ and $a \in A$, that $\int_{A} \eta_{i}=1$ for all $i$, and the supports shrink to $\left\{1_{A}\right\}$. By [14.1.4], $\eta_{i} \cdot v \rightarrow v$ in $V$ for each fixed $v \in V$. Each $\eta_{i}$ has a Fourier expansion $\eta_{i}=\sum_{\chi} \widehat{\eta}_{i}(\chi) \cdot \chi$ converging in $L^{2}(A)$, summing over continuous homomorphisms $\chi: A \rightarrow \mathbb{C}^{\times}$. For a finite set $X$ of $\chi \mathrm{s}$, let $\eta_{i}^{X}$ be the corresponding finite partial sum $\sum_{\chi \in X} \widehat{\eta}_{i}(\chi) \cdot \chi$ of the Fourier expansion of $\eta_{i}$. These are finite linear combinations of continuous functions, so are certainly continuous.

We claim that $\eta_{i}^{X} \cdot v \longrightarrow v$ in $V$. Since the partial sums $\eta_{i}^{X}$ approach $\eta_{i}$ in $L^{2}(A)$, it suffices to show that for $\eta \in C^{o}(A)$ with $|\eta|_{L^{2}(A)}$ small, $|\eta \cdot v|_{V}$ is small, for fixed $v \in V$. Indeed, invoking properties of vector-valued integrals [14.1] to exchange inner products and integrals,

$$
|\eta \cdot v|_{V}^{2}=\left\langle\int_{A} \eta(a) a \cdot v d a, \int_{A} \eta(b) b \cdot v d b\right\rangle=\int_{A} \int_{A} \eta(a) \cdot \bar{\eta}(b) \cdot\langle a \cdot v, b \cdot v\rangle d a d b
$$

By Cauchy-Schwarz-Bunyakowsky and the unitariness, $|\langle a \cdot v, b \cdot v\rangle| \leq|a v| \cdot|b v|=|v| \cdot|v|$, so

$$
|\eta \cdot v|_{V}^{2} \leq \int_{A} \int_{A}|\eta(a)| \cdot|\eta(b)| d a d b=\int_{A}|\eta(a)| d a \cdot \int_{A}|\eta(b)| d b \leq|\eta| \cdot|\eta|
$$

again by Cauchy-Schwarz-Bunyakowsky, since $|\eta(a)|=\eta(a) \cdot \mu(a)$, where $|\mu(a)|=1$, and the total measure of $A$ is 1 . This shows that $\eta_{i}^{X} \cdot v \rightarrow v$.

The action $\chi \cdot v$ is that of the projector $P_{\bar{\chi}}$ :

$$
\chi \cdot v=\int_{A} \chi(a) a \cdot v d a=\int_{A} \overline{\bar{\chi}}(a) a \cdot v d a=P_{\bar{\chi}} v
$$

Thus,

$$
\eta_{i}^{X} \cdot v=\sum_{\chi \in X} \widehat{\eta}_{i}(\chi) \cdot \chi \cdot v=\sum_{\chi \in X} \widehat{\eta}_{i}(\chi) \cdot P_{\bar{\chi}} v
$$

Since $\eta_{i} \cdot v \rightarrow v$ and $\eta_{i}^{X} \cdot v \rightarrow \eta_{i} \cdot v$, not all $P_{\chi} v$ can be 0 for non-zero $v$.

## 7. Discrete decomposition of cuspforms

1. The four simplest examples
2. $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$
3. $Z^{+} G L_{r}(k) \backslash G L_{r}(\mathbb{A})$

Appendix A: dualities
Appendix B: compact quotients $\Gamma \backslash G$

The general result here is that test functions $\varphi \in C_{c}^{\infty}(G)$ act as compact operators on spaces of squareintegrable cuspforms. Similar arguments and outcomes hold for both the four simplest archimedean examples from chapter 1 , adelic $Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})$ from chapter 2 , and archimedean or adelic versions of $G L_{n}$ from chapter 3 , with appropriate senses of test functions depending on context.

The most general argument refers to non-commutative rings of compact operators, closed under adjoints, and the irreducible representations of such rings are generally infinite-dimensional. In a larger context, this is the truth of the matter, but without knowing more about the irreducibles it is not easy to recover more tangible information about the behavior of the Laplacian or spherical Hecke operators.

Thus, we also give more special arguments using commutative rings of compact operators closed under adjoints, so that a more tangible notion of simultaneous eigenspace takes the place of (infinite-dimensional) irreducible representation. This gives the decomposition of spaces of right $K$-invariant functions by Laplacians and spherical Hecke algebras. The result for spherical Hecke algebras obtained in this fashion is nearly optimal, but there is still some imprecision in corollaries on eigenfunctions for Laplacians. The method of chapter 10 , directly considering the spectral behavior of pseudo-Laplacians and pseudo-cuspforms, gives better results of that sort.

Beyond the spectral decomposition theorems [7.1.1], [7.2.1], [7.3.1] and their immediate corollaries, we also conclude that there is an orthonormal basis for cuspforms consisting of smooth functions of rapid decay in Siegel sets: [7.1.20], [7.2.19], [7.2.20], [7.3.19].

For perspective, appendix [7.B] gives the much easier argument for discreteness of decomposition of $L^{2}(\Gamma \backslash G)$ for compact quotients $\Gamma \backslash G$, for a unimodular topological group $G$ and discrete subgroup $\Gamma$, although we do not give explicit examples of such $G, \Gamma$. In fact, this is a variant on the discrete decomposition [9.C.2] of $L^{2}(K)$ for compact topological groups $K$.

### 7.1 The four simplest examples

For this section, let $G, \Gamma, K$ be as in any of the four examples $S L_{2}(\mathbb{R}), S L_{2}(\mathbb{C}), S p_{1,1}^{*}, S L_{2}(\mathbb{H})$ from chapter 1. As earlier, test functions $C_{c}^{\infty}(G)$ are compactly supported, smooth functions, and act on functions $f$ on $\Gamma \backslash G$ by

$$
(\eta \cdot f)(g)=\int_{G} \eta(h) f(g h) d h \quad\left(\text { for } \eta \in C_{c}^{\infty}(G)\right)
$$

As earlier, functions on $\Gamma \backslash G / K$ are identified with right $K$-invariant functions on $\Gamma \backslash G$, and the action of the spherical convolution algebra of left-and-right $K$-invariant test functions $C_{c}^{\infty}(K \backslash G / K)$ stabilizes the subspace of such functions. As we recall below, the spherical convolution algebra $C_{c}^{\infty}(K \backslash G / K)$ is commutative. The main results of this section are
[7.1.1] Theorem: The spherical convolution algebra $C_{c}^{\infty}(K \backslash G / K)$ of left and right $K$-invariant test functions on $G$ acts on square-integrable right $K$-invariant cuspforms $L_{o}^{2}(\Gamma \backslash G / K)$ by compact operators, the collection of such operators is closed under adjoints, and is non-degenerate in the sense that for every $f \in L_{o}^{2}(\Gamma \backslash G / K)$ there is $\varphi \in C_{c}^{\infty}(K \backslash G / K)$ such that $\varphi \cdot f \neq 0$. (Proof in the sequel.)
[7.1.2] Corollary: The space $L_{o}^{2}(\Gamma \backslash G / K)$ of right $K$-invariant square-integrable cuspforms decomposes into simultaneous eigenspaces for operators in the commutative convolution algebra $C_{c}^{\infty}(K \backslash G / K)$, each finite-dimensional. The simultaneous eigenfunctions are smooth. (Proof below.)
[7.1.3] Corollary: There is an orthonormal Hilbert-space basis for the space of $K$-invariant square-integrable cuspforms consisting of simultaneous eigenfunctions for the invariant Laplacian $\Delta$. (Proof below.)

## 7. Discrete decomposition of cuspforms

[7.1.4] Remark: The argument here does not immediately prove that the eigenspaces for $\Delta$ are of finite multiplicity, since it only indirectly refers to $\Delta$. The unbounded-operator argument in chapter 9 gives a stronger result about $\Delta$-eigenspaces.

Let $P, M, N, A^{+}$be as in chapter 1 , and let $\eta$ be a height function

$$
\eta\left(n a_{t} k\right)=t^{r_{o}} \quad\left(n \in N, k \in K, \text { with } a_{t}=\left(\begin{array}{cc}
\sqrt{t} & 0 \\
0 & 1 / \sqrt{t}
\end{array}\right) \text { for } t>0\right)
$$

with $r_{o}=1,2,3,4$ in the respective cases. Thus, $\eta(m)=\delta(m)$, the modular function for $P$. In an Iwasawa decomposition $G=N \cdot M \cdot K$ for $x \in G$, write $x=n_{x} \cdot m_{x} \cdot k_{x}$ with $n_{x} \in N, m_{x} \in M$, and $k_{x} \in K$. This notation $n_{x}$ is in conflict with the use of that notation in the earlier discussion of the smaller examples, but those former uses will not be needed here. For $t>0$, compact $C \subset N$, the corresponding Siegel set in $G$ is

$$
\mathfrak{S}=\mathfrak{S}_{C, \tau}=C \cdot\left\{a_{t} \in A^{+}: t \geq \tau\right\} \cdot K=C \cdot\{m \in M: \eta(m) \geq \tau\} \cdot K
$$

[7.1.5] Claim: A point in a Siegel set is well approximated by its $M$-component in an Iwasawa decomposition, in the sense that, for $x \in \mathfrak{S}_{C, \tau}$ with $\tau>0$ and $C$ compact in $N$, there is another compact subset $C^{\prime}$ of $N$ such that $x \in m_{x} \cdot C^{\prime} \cdot K$.
Proof: $x \in \mathfrak{S}_{C, \tau}$ gives

$$
x=n_{x} \cdot m_{x} \cdot k_{x} \in C \cdot m_{x} \cdot K=m_{x} \cdot m_{x}^{-1} C m_{x} \cdot K
$$

The lower bound $\eta(m) \geq \tau$ gives a compact set $C^{\prime}$ in $N$ depending only upon $\tau$ and $C$ such that

$$
m^{-1} C m \subset C^{\prime} \quad(\text { for } m \in M \text { with } \eta(m) \geq \tau)
$$

In particular, $m_{x}^{-1} C m_{x} \subset C^{\prime}$. Thus, $x \in m_{x} \cdot C^{\prime} \cdot K$.
For strictly upper-triangular square matrices $x$ with entries in any field of characteristic zero, the series for the matrix exponential $e^{x}=\exp (x)=\sum_{\ell \geq 0} x^{\ell} / \ell$ ! is finite. Let $\mathfrak{n}$ be the Lie algebra of $N$ :

$$
\mathfrak{n}=\left\{\left(\begin{array}{ll}
0 & v \\
0 & 0
\end{array}\right)\right\}
$$

where $v$ is in $\mathbb{R}, \mathbb{C},\{w \in \mathbb{H}: w+\bar{w}=0\}$, or $\mathbb{H}$ in the four cases. Because $N$ is abelian, the exponential map is a diffeomorphism $\exp : \mathfrak{n} \longrightarrow N$, and there is a discrete additive subgroup $\Lambda$ in $\mathfrak{n}$ such that $\exp (\Lambda)=\Gamma \cap N$ :

$$
\exp \left(\begin{array}{ll}
0 & v \\
0 & 0
\end{array}\right)=\left(\begin{array}{ll}
1 & v \\
0 & 1
\end{array}\right)
$$

namely, take $v$ in $\mathbb{Z}, \mathbb{Z}[i], \mathbb{Z} i+\mathbb{Z} j+\mathbb{Z} k$, or the Hurwitz quaternion integers $\mathfrak{o}$, respectively. For test function $\varphi \in C_{c}^{\infty}(G)$, wind up the integral for $f \rightarrow \varphi \cdot f$ along $\exp (\Lambda)=N \cap \Gamma$ : for $y \in G$,

$$
\begin{gathered}
(\varphi \cdot f)(y)=\int_{G} \varphi(x) f(y x) d x=\int_{G} \varphi\left(y^{-1} x\right) f(x) d x=\int_{\exp (\Lambda) \backslash G}\left(\sum_{\gamma \in \exp (\Lambda)} \varphi\left(y^{-1} \gamma x\right)\right) f(x) d x \\
=\int_{\exp (\Lambda) \backslash G}\left(\sum_{\nu \in \Lambda} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)\right) f(x) d x
\end{gathered}
$$

The kernel function for the $(N \cap \Gamma)$-wound-up operator is the latter inner sum:

$$
K_{\varphi}(x, y)=\sum_{\nu \in \Lambda} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)
$$

[7.1.6] Claim: For a fixed Siegel set $\mathfrak{S}$ and for fixed compact $E \subset C_{c}^{\infty}(G)$, there is compact $C_{M} \subset M$ such that if there exist $n \in N, x, y \in \mathfrak{S}$, and $\varphi \in E$ with $\varphi\left(y^{-1} \cdot n \cdot x\right) \neq 0$, then $m_{x} \in m_{y} \cdot C_{M}$. That is, $K_{\varphi}(x, y)=0$ for all $x, y \in \mathfrak{S}$ and all $\varphi \in E$ unless $m_{x} \in m_{y} \cdot C_{M}$.

Proof: From the previous claim, there is compact $C^{\prime} \subset N$ such that $m_{y}^{-1} y \in C^{\prime} \cdot K$. A compact set of test functions has a common compact support $C_{G}$, because a compact set is bounded in the topological vector space sense, and a bounded subset of an LF-space such as $C_{c}^{\infty}(G)$ lies in some Fréchet limitand, by [13.8.5]. Non-vanishing of $\varphi\left(y^{-1} n x\right)$ implies $y^{-1} n x \in C_{G}$, so

$$
n x \in y \cdot C_{G} \subset m_{y} \cdot C^{\prime} \cdot K \cdot C_{G} \subset m_{y} \cdot C_{G}^{\prime} \quad \text { (with } C_{G}^{\prime}=C^{\prime} K C_{G}=\text { compact) }
$$

That is,

$$
C_{G}^{\prime} \ni m_{y}^{-1} \cdot n x=m_{y}^{-1} \cdot n n_{x} \cdot m_{x} \cdot k_{x}=\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x} \cdot k_{x}
$$

That is,

$$
\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x} \in C_{G}^{\prime} \cdot K=\text { compact }
$$

Since $M$ normalizes $N$, the element $m_{y}^{-1} n n_{x} m_{y}$ is in $N$. Since $N \cap M=\{1\}$, the multiplication map $N \times M \rightarrow N M$ is a homeomorphism. Thus, for the product $\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x}$ to lie in a compact set in $G$ requires that its $N$-component lies in a compact set in $N$ and its $M$-component lies in a compact set in $M$. Thus, there is compact $C_{M} \subset M$ such that $m_{y}^{-1} m_{x} \in C_{M}$, as claimed.
[7.1.7] Corollary: For fixed Siegel set $\mathfrak{S}$ and fixed compact $E \subset C_{c}^{\infty}(G)$, there is a compact $C_{M} \subset M$ such that, if $\varphi\left(y^{-1} n x\right) \neq 0$ for some $x, y \in \mathfrak{S}$, some $n \in N$ and some $\varphi \in E$, then $m_{y}^{-1} x \in C_{M}$.
Proof: By [7.1.5], there is a compact $C_{G}$ in $G$ such that $x \in m_{x} \cdot C_{G}$. By [7.1.6], there is a compact $C_{M}$ in $M$ such that $m_{x} \in m_{y} C_{M}$. Thus, $x \in m_{x} C_{G} \subset\left(m_{y} C_{M}\right) C_{G}^{\prime}$, rearranging to give the claim.
[7.1.8] Corollary: With $\omega_{y}=y^{-1} m_{y}$ and $\omega_{x, y}=m_{y}^{-1} x$, the functions $\nu \longrightarrow \varphi_{x, y}(\nu)=\varphi\left(\omega_{y} \cdot \exp (\nu) \cdot \omega_{x, y}\right)$ for $x, y$ in fixed Siegel set, $\varphi\left(y^{-1} n x\right) \neq 0, \varphi \in E$, constitute a compact subset of the Schwartz space $\mathscr{S}(\mathfrak{n})$.
Proof: The left and right translation actions of $G$ on test functions are continuous $G \times G \times C_{c}^{\infty}(G) \rightarrow C_{c}^{\infty}(G)$ by [6.4]. With fixed Siegel set $\mathfrak{S}$, by [7.1.5] and [7.1.7], $\left\{\omega_{y}: y \in \mathfrak{S}\right\}$ and $\left\{\omega_{x, y}: x \in \mathfrak{S}, y \in \mathfrak{S}\right\}$ are compact. This gives compactness of the image of

$$
\left\{\omega_{y}: y \in \mathfrak{S}\right\} \times\left\{\omega_{x, y}: x \in \mathfrak{S}, y \in \mathfrak{S}\right\} \times E \longrightarrow C_{c}^{\infty}(G) \quad \quad\left(\text { with } \varphi\left(y^{-1} n x\right) \neq 0\right)
$$

Since $N$ is closed in $G$, the restriction map $C_{c}^{\infty}(G) \rightarrow C_{c}^{\infty}(N) \approx C_{c}^{\infty}(\mathfrak{n})$ is continuous. A continuous image of a compact set is compact. Comparing the topologies from [13.7] and [13.9], as in [13.9.3], the inclusion $C_{c}^{\infty}(\mathfrak{n}) \subset \mathscr{S}(\mathfrak{n})$ is continuous, giving compactness of the image.

Poisson summation for the lattice $\Lambda \subset \mathfrak{n}$ gives

$$
K_{\varphi}(x, y)=\sum_{\nu \in \Lambda} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)=\sum_{\psi \in \Lambda^{*}} \int_{\mathfrak{n}} \bar{\psi}(\nu) \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu
$$

where $\Lambda^{*}$ is the collection of $\mathbb{C}^{\times}$-valued characters on $\mathfrak{n}$ that are trivial on $\Lambda$. Rearrange slightly to

$$
K_{\varphi}(x, y)=\sum_{\psi \in \Lambda^{*}} \int_{\mathfrak{n}} \bar{\psi}(\nu) \cdot \varphi\left(y^{-1} m_{y} \cdot \exp \left(m_{y}^{-1} \nu m_{y}\right) \cdot m_{y}^{-1} x\right) d \nu
$$

Replacing $\nu$ by $m_{y} \nu m_{y}^{-1}$ and letting $\varphi_{x, y}(\nu)=\varphi\left(y^{-1} m_{y} \cdot \nu \cdot m_{y}^{-1} x\right)$,

$$
K_{\varphi}(x, y)=\delta\left(m_{y}\right) \sum_{\psi \in \Lambda^{*}} \widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right)
$$

where $\widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right)$ is the Fourier transform $\widehat{\varphi}_{x, y}$ of $\varphi_{x, y}$ evaluated at $\psi$, and $\psi^{m_{y}}(\nu)=\psi\left(m_{y} \nu m_{y}^{-1}\right)$.
[7.1.9] Theorem: Fix compact $E \subset C_{c}^{\infty}(G)$ and Siegel set $\mathfrak{S}$. For any given $q>0$ there is a uniform bound

$$
|(\varphi \cdot f)(y)|<_{q} \eta(y)^{-q} \cdot|f|_{L^{2}(\Gamma \backslash G)} \quad\left(\text { for all } y \in \mathfrak{S}, \text { for all } \varphi \in E, \text { for all } L^{2} \text { cuspforms } f\right. \text { ) }
$$

## 7. Discrete decomposition of cuspforms

Proof: For convenient discussion of the Schwartz seminorms on $\mathfrak{n}$, give the real vector space $\mathfrak{n}$ a positivedefinite inner product $\langle$,$\rangle invariant under conjugation action of M \cap K$, allowing identification of $\mathfrak{n}$ with its dual $\mathfrak{n}^{*}$ when desired, for simplicity. Let $|\cdot|$ be the associated norm on either. The compactness [7.1.8] and continuity of Fourier transform on $\mathscr{S}(\mathfrak{n})$ give a uniform estimate on Fourier transforms: for fixed Siegel set $\mathfrak{S}$, for given $q>0,\left|\widehat{\varphi}_{x, y}(\psi)\right|<_{r}(1+|\psi|)^{-q}$ for all $x, y \in \mathfrak{S}$. Then

$$
\delta\left(m_{y}\right)\left|\widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right)\right|<_{q} \delta\left(m_{y}\right) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-q} \quad(\text { for all } x, y \in \mathfrak{S})
$$

Next, toward [7.1.9] we need
[7.1.10] Claim: For fixed Siegel set $\mathfrak{S}$ and $r \gg 1$, the kernel $K_{\varphi}(x, y)$ with its $0^{\text {th }}$ Fourier component removed satisfies

$$
\left|K_{\varphi}(x, y)-\widehat{\varphi}_{x, y}\left(\psi_{0}\right)\right|<_{q}|\eta(y)|^{-q} \quad(\text { for } x, y \in \mathfrak{S})
$$

Proof: First, we claim that, given $\Lambda$ and given a Siegel set $\mathfrak{S}=\mathfrak{S}_{C, \tau}$, there is an implied constant such that

$$
\eta(y) \ll 1+\left|\psi^{m_{y}}\right|^{r_{o}} \quad\left(\text { for all } y \in \mathfrak{S}, \text { for all } 0 \neq \psi \in \Lambda^{*}\right)
$$

Again, $\eta(y)=t^{r_{o}}$ for $y=n a_{t} k$ with $t>0$, and $\left|\psi^{m_{y}}\right|=t \cdot|\psi|$. Since the norms of non-zero elements of $\Lambda^{*}$ have a positive inf,

$$
\left|\psi^{m_{y}}\right|^{r_{o}}=\left(\eta(y)^{1 / r_{o}} \cdot|\psi|\right)^{r_{o}} \geq \eta(y) \cdot \inf _{0 \neq \psi \in \Lambda^{*}}|\psi|^{r_{o}}
$$

Since $\widehat{\varphi}_{x, y}$ is a Schwartz function, $\left|\widehat{\varphi}_{x, y}(\psi)\right|<_{\ell}(1+|\psi|)^{-\ell}$ for every $\ell>0$. By the comparison of $\eta(y)$ to $\left|\psi^{m_{y}}\right|$, for $0 \neq \psi \in \Lambda^{*}$,

$$
\begin{aligned}
\delta\left(m_{y}\right) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-\ell} & =\eta(y) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-\ell} \ll \eta(y) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-(q+1) \cdot r_{o}-\left(\ell-(q+1) \cdot r_{o}\right)} \\
& \ll \eta(y) \cdot \eta(y)^{-(q+1)} \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-\left(\ell-(q+1) \cdot r_{o}\right)}
\end{aligned}
$$

For $\ell$ sufficiently large depending on $q$, the latter sum over $0 \neq \psi \in \Lambda^{*}$ converges, giving the claim.
[7.1.11] Claim: Cuspforms $f$ ignore the $\psi_{0}^{t h}$ Fourier component of $K_{\varphi}(x, y)$ :

$$
(\varphi \cdot f)(y)=\int_{\exp (\Lambda) \backslash G}\left(K_{\varphi}(x, y)-\widehat{\varphi}_{x, y}\left(\psi_{0}\right)\right) \cdot f(x) d x
$$

Proof: For trivial character $\psi_{0}$ on $\mathfrak{n}$, the corresponding function $\widehat{\varphi}_{x, y}\left(\psi_{0}\right)$ is left $N$-invariant in $x$ : let $n \in N$, and replace $x$ by $n x$ in the original integral defining $\widehat{\varphi}_{x, y}\left(\psi_{0}\right)$, with $n=\exp \left(\nu^{\prime}\right)$, obtaining

$$
\int_{\mathfrak{n}} \bar{\psi}(\nu) \cdot \varphi\left(y^{-1} \exp (\nu) \cdot n x\right) d \nu=\int_{\mathfrak{n}} \bar{\psi}(\nu) \cdot \varphi\left(y^{-1} \exp \left(\nu+\nu^{\prime}\right) \cdot x\right) d \nu
$$

using the abelian-ness of $N$. Replacing $\nu$ by $\nu-\nu^{\prime}$ in the integral gives the left $N$-invariance in $x$ :

$$
\int_{\mathfrak{n}} \bar{\psi}(\nu) \cdot \varphi\left(y^{-1} \exp (\nu) \cdot n x\right) d \nu=\psi\left(\nu^{\prime}\right) \cdot \widehat{\varphi}_{x, y}(\psi)=\widehat{\varphi}_{x, y}(\psi)
$$

Therefore, in

$$
(\varphi \cdot f)(y)=\delta\left(m_{y}\right) \sum_{\psi \in \Lambda^{*}} \int_{(N \cap \Gamma) \backslash G} \widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right) \cdot f(x) d x
$$

the integral for trivial character $\psi_{0}$ is

$$
\begin{gathered}
\int_{(N \cap \Gamma) \backslash G} \widehat{\varphi}_{x, y}\left(\psi_{0}^{m_{y}}\right) \cdot f(x) d x=\int_{N \backslash G} \int_{(N \cap \Gamma) \backslash N} \widehat{\varphi}_{n x, y}\left(\psi_{0}^{m_{y}}\right) \cdot f(n x) d n d x \\
=\int_{N \backslash G} \widehat{\varphi}_{x, y}\left(\psi_{0}^{m_{y}}\right) \cdot\left(\int_{(N \cap \Gamma) \backslash N} f(n x) d n\right) d x=\int_{N \backslash G} \widehat{\varphi}_{x, y}\left(\psi_{0}^{m_{y}}\right) \cdot 0 d x=0 \quad \quad(\text { cuspform } f)
\end{gathered}
$$

proving the claim.
The proof of [7.1.9] is almost complete. From above, for $y$ in a fixed Siegel set $\mathfrak{S}$ and for fixed test function $\varphi$, there is a compact $C_{M} \subset A^{+}$such that, for $\varphi\left(y^{-1} n x\right)$ to be nonzero, the Iwasawa component $m_{x}$ of $x$ must lie in $m_{y} \cdot C_{M}$. Thus,

$$
\left\{x \in \mathfrak{S}: \varphi\left(y^{-1} n x\right) \neq 0 \text { for some } n \in N\right\} \subset m_{y} C_{M} \cdot K
$$

Combining this with the estimate just obtained, for cuspform $f$,

$$
|(\varphi \cdot f)(y)|<_{q}|\eta(y)|^{-q} \cdot \int_{\Gamma \backslash \Gamma\left(C \cdot m_{y} C_{M} \cdot K\right)}|f(x)| d x
$$

By Cauchy-Schwarz-Bunyakowsky,

$$
\left(\int_{\Gamma \backslash \Gamma\left(C \cdot m_{y} C_{M} \cdot K\right)}|f(x)| d x\right)^{2} \leq \int_{\Gamma \backslash \Gamma\left(C \cdot m_{y} C_{M} \cdot K\right)} 1 d x \cdot \int_{\Gamma \backslash \Gamma\left(C \cdot m_{y} C_{M} \cdot K\right)}|f(x)|^{2} d x \ll|f|_{L^{2}}^{2}
$$

This gives the desired decay, proving theorem [7.1.9].
We are getting closer to the compactness of the operators $f \rightarrow \varphi \cdot f$ on cuspforms $f$. Recall that a collection $E$ of continuous functions on $G$ or $\Gamma \backslash G$ is (uniformly) equicontinuous when, given $\varepsilon>0$, there is a neighborhood $U$ of 1 in $G$ such that

$$
|f(x)-f(y)|<\varepsilon \quad \text { (for all } f \in E, \text { for all } x \in G, \text { for all } y \in x \cdot U)
$$

We have the expected
[7.1.12] Lemma: For $X \in \mathfrak{g}$, the left-derivative map

$$
C_{c}^{\infty}(G) \longrightarrow C_{c}^{\infty}(G) \quad \text { by } \quad \varphi \longrightarrow X^{\text {left }} \cdot \varphi=\left(\left.g \rightarrow \frac{d}{d t}\right|_{t=0} \varphi\left(e^{-t X} g\right)\right)
$$

is continuous.
Proof: $C_{c}^{\infty}(G)$ is an LF-space, a (strict) colimit of Fréchet spaces, the limit being taken over spaces $\mathcal{D}_{\Omega}$ of smooth functions on $G$ supported on compact $\Omega \subset G$. The topology on each $\mathcal{D}_{\Omega}$ is given by seminorms taking sups of derivatives of all orders, but of course the notion of derivative is ambiguous, since there are at least two different choices of global vectorfields, left derivatives by $\mathfrak{g}$, and right derivatives by $\mathfrak{g}$.

But a reasonable general assertion is true, for fairly elementary reasons. On a smooth manifold, let $X^{i}$ be a tuple of (smooth) vector fields on an open $U$ containing a given compact set $B$ such that, for every $x$ in $U$, the values $X_{x}^{i}$ at $x$ are a basis for the tangent space at $x$. Another such tuple $Y^{j}$ can be expressed (smoothly, pointwise) as linear combinations of the $X^{i}$, and vice-versa. Every entry of the matrix of coefficients is bounded, the determinant of the matrix of coefficients does not vanish on the compact $B$, so is bounded and bounded away from 0 , the coefficients are smooth functions, and the inverse is smooth on $B$. Thus, the two sets of seminorms, corresponding to left or right first derivatives,

$$
\sup _{x \in B} \sup _{i}\left(X_{i} \varphi\right)(x) \quad \text { or } \quad \sup _{x \in B} \sup _{j}\left(Y_{j} \varphi\right)(x)
$$

are topologically equivalent. With this ambiguity removed, a (left or right) derivative is a continuous map $C^{k}(B) \rightarrow C^{k-1}(B)$, and so gives a continuous map on the limit: $C^{\infty}(B) \rightarrow C^{\infty}(B)$. The subspace $\mathcal{D}_{B}$ is a closed subspace of $C^{\infty}(B)$ described by closed conditions: all derivatives vanishing on the boundary. Thus, a derivative map is still continuous $\mathcal{D}_{B} \rightarrow \mathcal{D}_{B}$. This induces a continuous map on the colimit.
[7.1.13] Corollary: For a compact set $E$ of test functions on $G$, for a compact $C_{\mathfrak{g}}$ in $\mathfrak{g}$, and for $f$ ranging over cuspforms in the unit ball in $L^{2}(\Gamma \backslash G)$, there is a uniform implied constant such that

$$
|(X \cdot(\varphi \cdot f))(g)| \ll 1 \quad\left(\text { for all } g \in G, \text { for all } \varphi \in E, \text { for all } X \in C_{\mathfrak{g}}\right)
$$

Proof: The differentiation of $\varphi \cdot f$ can be rewritten as a differentiation of $\varphi$, followed by action of the resulting function on $f$ :

$$
\begin{gathered}
(X \cdot \varphi \cdot f)(x)=\left.\frac{d}{d t}\right|_{t=0} \int_{G} \varphi(y) f\left(x e^{t X} y\right) d y=\left.\frac{d}{d t}\right|_{t=0} \int_{G} \varphi\left(e^{-t X} y\right) f(x y) d y \\
\left.=\int_{G}\left(\left.\frac{d}{d t}\right|_{t=0} \varphi\left(e^{-t X} y\right)\right) f(x y) d y \quad \text { (replacing } y \text { by } e^{-t X} y\right)
\end{gathered}
$$

justifying interchange of differentiation and integration by the continuity of differentiation on test functions $\varphi$, and Gelfand-Pettis integral properties [14.1]. That is, $X \cdot \varphi \cdot f=\left(X^{\text {left }} \varphi\right) \cdot f$ with $X^{\text {left }}$ the left action. Since $\mathfrak{g}$ is a finite-dimensional real vector-space and the action is linear in $X$, this gives the continuity in $X \in \mathfrak{g}$. Thus, the collection of test functions $X^{\text {left }} \varphi$ with $X \in C_{\mathfrak{g}}$ and $\varphi \in E$ is again compact in $C_{c}^{\infty}(\mathcal{G})$. Thus, by the bound of theorem [7.1.9],

$$
\left.|(X \cdot \varphi \cdot f)(y)|=\mid\left(X^{\text {left }} \varphi\right) \cdot f\right)\left.(y)\left|<_{r} \eta(y)^{-r} \cdot\right| f\right|_{L^{2}(\Gamma \backslash G)} \quad\left(\text { for all } y \in \mathfrak{S}, X \in C_{\mathfrak{g}}, \varphi \in E\right)
$$

For large-enough Siegel set to cover the quotient, and any $r>0$, this gives

$$
\sup |X \cdot \varphi \cdot f(y)|<_{r}|f|_{L^{2}}
$$

proving uniform boundedness for $|f|_{L^{2}} \leq 1$.
The smoothing property of $f \rightarrow \varphi \cdot f$ as in [14.5] assures that each $\varphi \cdot f$ is in $C^{\infty}(G)$. A uniform bound on derivatives implies uniform continuity:
[7.1.14] Lemma: Let $F$ be a smooth function on $G$, with a uniform pointwise bound on all $X \cdot F$ with $X$ in a compact neighborhood $C_{\mathfrak{g}}$ of 0 in $\mathfrak{g}$, namely,

$$
\left.|(X \cdot F)(x)| \leq B \quad \text { (for all } x \in G, \text { all } X \in C_{\mathfrak{g}}\right)
$$

Then $F$ is uniformly continuous: for every $\varepsilon>0$ there is a neighborhood $U$ of 1 in $G$ such that $|F(x)-F(y)|<\varepsilon$ for all $x \in G$ and $y \in x U$.
Proof: Let $V$ be a small enough open containing 1 such that $V$ is contained in $\exp C_{\mathfrak{g}}$, and that the exponential map on $\exp ^{-1} V$ is injective to $V$. Let $y=x \cdot e^{s X}$ for $X \in C_{\mathfrak{g}}$ and $0 \leq s \leq 1$. By hypothesis, the function $h(t)=F\left(x \cdot e^{s X}\right)$ has

$$
h^{\prime}(s)=\left.\frac{d}{d t}\right|_{t=0} h(s+t)=F\left(x \cdot e^{s X} \cdot e^{t X}\right)
$$

bounded by $B$. From the mean value theorem, $\left|F\left(x \cdot e^{t X}\right)-F(x)\right| \leq t \cdot B$. Thus, for all $|t|<B \cdot \varepsilon$ we have the desired inequality.
[7.1.15] Corollary: For a compact set $E$ of test functions on $G$, and for $f$ ranging over cuspforms in the unit ball in $L^{2}(\Gamma \backslash G)$, the family of images $\varphi \cdot f$ is equicontinuous on $G$.

We are almost done with the proof that $f \rightarrow \varphi \cdot f$ is compact on cuspforms.
[7.1.16] Lemma: Let $E$ be a equicontinuous, uniformly bounded, set of functions on $\Gamma \backslash G$. Then $E$ has compact closure in $L^{2}(\Gamma \backslash G)$.
[7.1.17] Remark: Superficially, this lemma is reminiscent of the Arzela-Ascoli theorem, which asserts that an equicontinuous and uniformly bounded subset of $C^{o}(K)$ (with sup norm) for a compact topological space $K$ has compact closure in $C^{O}(K)$. Indeed is common to end a sketch of the discrete decomposition of cuspforms by an allusion to Arzela-Ascoli. However, we need less, fortunately, since adaptation of the literal Arzela-Ascoli result to the present circumstance seems awkward.
Proof: The proof is a maneuver to invoke the fact that a totally bounded subset of a complete metric space has compact closure. If the quotient $\Gamma \backslash G$ were compact then we could simply invoke Arzela-Ascoli, but this is perhaps exactly the difficulty. Without loss of generality, all the functions in $E$ are bounded (in absolute value) by 1 , and the total measure of $\Gamma \backslash G$ is 1 . Take $\varepsilon>0$. Using the equicontinuity, let $U$ be a

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small enough neighborhood (with compact closure) of 1 in $G$ such that for any $x \in G$ and $y \in x U$ we have $|F(x)-F(y)|<\varepsilon$ for all $F \in E$.

Let $\left\{x_{i}\right\}$ be a countable set dense in $G$. Let $U_{i}=x_{i} U$. Let $q: G \rightarrow \Gamma \backslash G$ be the quotient map. Let $V_{1}$ be the image of $U_{1}$ in $\Gamma \backslash G$, and recursively take

$$
V_{n+1}=\left\{x \in \Gamma \backslash G: x \in q U_{i} \text { but } x \notin q\left(U_{1} \cup \ldots \cup U_{n}\right)\right\}
$$

Since the $V_{i}$ are disjoint and their union is $\Gamma \backslash G$, which has finite measure, $\sum_{i}$ meas $\left(V_{i}\right)<+\infty$. In particular, the measures meas $V_{i}$ go to 0 as $i \rightarrow \infty$. Take $m$ large enough such that $\sum_{i>m}$ meas $\left(V_{i}\right)<\varepsilon$. Let $X$ be a finite set of complex numbers such that any complex number of absolute value at most 1 is within $\varepsilon / 2$ of an element of $X$. For each $m$-tuple $\xi=\left(\xi_{1}, \ldots, \xi_{m}\right)$ of elements of $X$, define a function $F_{i, \xi}$ on $\Gamma \backslash G$ by

$$
F_{i, \xi}(x)= \begin{cases}\xi_{i} & \left(\text { for } x \in V_{i}, i \leq m\right) \\ 0 & \left(\text { for } x \in V_{i}, i>m\right)\end{cases}
$$

Given $F \in E$, for each $i \leq m$ choose $\xi_{i}$ such that $\left|F\left(x_{i}\right)-\xi_{i}\right|<\varepsilon$. By the choice of $U$,

$$
\left.\left|F(x)-\xi_{i}\right| \leq\left|F(x)-F\left(x_{i}\right)\right|+\left|F\left(x_{i}\right)-\xi_{i}\right|<2 \varepsilon \quad \text { (for } x \in V_{i}\right)
$$

Then

$$
\int_{\Gamma \backslash G}\left|F-F_{\xi}\right|^{2}<\int_{V_{1} \cup \ldots \cup V_{m}}(2 \varepsilon)^{2}+\int_{V_{m+1} \cup \ldots} 1 \leq 4 \varepsilon^{2} \cdot \operatorname{meas} \Gamma \backslash G+\operatorname{meas}\left(V_{m+1} \cup \ldots\right) \leq 4 \varepsilon^{2}+\varepsilon
$$

Tweaking the estimates as desired, for given $\varepsilon>0$ we can cover $E$ by a finite number of balls of radius $\varepsilon>0$ in $L^{2}(\Gamma \backslash G)$. Total boundedness implies compact closure.
[7.1.18] Corollary: For $\varphi \in C_{c}^{\infty}(G)$, the operator $f \rightarrow \varphi \cdot f$ is a compact operator $L_{o}^{2}(\Gamma \backslash G) \longrightarrow L_{o}^{2}(\Gamma \backslash G)$.
Proof: The asymptotics of the kernels prove pointwise boundedness of the image of the unit ball $B$ of $L_{o}^{2}(\Gamma \backslash G)$. Consideration of derivatives proves equicontinuity of the image of $B$. The faux-Arzela-Ascoli compactness lemma above proves compactness of the closure of $\varphi \cdot B$. Being integrated versions of right translations, these operators stabilize the subspace of cuspforms, as the latter is defined by a left integral condition. Thus, $\varphi$ maps the unit ball to a set with compact closure, so is a compact operator.

In these examples, the space of right $K$-invariant cuspforms $L_{o}^{2}(\Gamma \backslash G / K)=L_{o}^{2}(\Gamma \backslash G)^{K}$ is of main interest, and the left-and-right $K$-invariant test functions act there:
[7.1.19] Corollary: For $\varphi \in C_{c}^{\infty}(K \backslash G / K)$, the operator $f \rightarrow \varphi \cdot f$ is a compact operator $L_{o}^{2}(\Gamma \backslash G / K) \longrightarrow L_{o}^{2}(\Gamma \backslash G / K)$.
Proof: A restriction of a compact operator is still compact. It suffices to show that the $K$-fixed subspace is closed in $L_{o}^{2}(\Gamma \backslash G)$, since the spherical Hecke algebra $C_{c}^{\infty}(K \backslash G / K)$ stabilizes it, by direct computation. From [6.1], the right action of $G$ on $L^{2}(\Gamma \backslash G)$ is unitary, so continuous. Thus, the condition of right $K$-invariance is a closed condition.

Adjoints of the operators $f \rightarrow \varphi \cdot f$ are easily computed: letting $\langle$,$\rangle be the inner product on L^{2}(\Gamma \backslash G)$,

$$
\begin{aligned}
\langle\varphi \cdot f, F\rangle & =\int_{\Gamma \backslash G} \int_{G} \varphi(x) f(y x) \bar{F}(y) d x d y=\int_{\Gamma \backslash G} \int_{G} \varphi(x) f(y) \bar{F}\left(y x^{-1}\right) d x d y \\
& =\int_{\Gamma \backslash G} \int_{G} f(y) \varphi\left(x^{-1}\right) \bar{F}(y x) d x d y=\int_{\Gamma \backslash G} f(y) \overline{\varphi^{\vee} \cdot \bar{F}}(y) d y
\end{aligned}
$$

where $\varphi^{\vee}(x)=\bar{\varphi}\left(x^{-1}\right)$ as suggested by the computation. The space $C_{c}^{\infty}(K \backslash G / K)$ is stable under the operation $\varphi \rightarrow \varphi^{\vee}$.

The non-degeneracy is essentially [14.1.5], but we need a right $K$-averaged version. Let $\varphi_{i}$ be an approximate identity, so that $\varphi_{i} \cdot f \rightarrow f$, by [14.1.4]. For $f$ right-invariant, the $K$-averaged versions

$$
\left.\alpha\left(\varphi_{i} \cdot f\right)(x)=\int_{K}\left(\varphi_{i} \cdot f\right)(x k) d k \quad \quad \text { (giving } K \text { total measure } 1\right)
$$

## 7. Discrete decomposition of cuspforms

of $\varphi_{i} \cdot f$ must approach $f$, since $K$-averaging is an orthogonal projector to the space of $K$-invariant functions:

$$
\begin{gathered}
\langle\alpha f, \alpha F\rangle=\int_{K} \int_{K} \int_{G} f(x h) F(x k) d x d k d h=\int_{K} \int_{K} \int_{G} f(x) F\left(x h^{-1} k\right) d x d k d h \\
=\int_{K} \int_{K} \int_{G} f(x) F(x k) d x d k d h=\langle f, \alpha F\rangle
\end{gathered}
$$

Thus, for $f \neq 0$, for all sufficiently large $i$ we have $\varphi_{i} \cdot f \neq 0$. These averaged versions are obtained by averaging $\varphi_{i}$ :

$$
\begin{aligned}
& \alpha\left(\varphi_{i} \cdot f\right)(x)=\int_{K}\left(\varphi_{i} \cdot f\right)(x k) d k=\int_{K} \int_{G} \varphi_{i}(y) \cdot f(x k y) d y d k \\
= & \int_{K} \int_{G} \varphi_{i}\left(k^{-1} y\right) \cdot f(x y) d y d k=\int_{G}\left(\int_{K} \varphi_{i}\left(k^{-1} y\right) d k\right) \cdot f(x y) d y
\end{aligned}
$$

Again using the right $K$-invariance of $f$,

$$
\begin{aligned}
& (\varphi \cdot f)(x)=\int_{G} \varphi(y) f(x y) d y=\int_{G} \varphi(y) f(x y) d y \cdot \int_{K} 1 d k \\
& =\int_{K} \int_{G} \varphi(y) f(x y k) d y d k=\int_{G}\left(\int_{K} \varphi\left(y k^{-1}\right) d k\right) f(x y) d y
\end{aligned}
$$

Thus, for $K$-invariant $f$, there is left and right $K$-invariant $\varphi$ such that $\varphi \cdot f \neq 0$. This is the non-degeneracy of the action of $C_{c}^{\infty}(K \backslash G / K)$ on $L_{o}^{2}(\Gamma \backslash G / K)$. This proves theorem [7.1.1].

Proof: (of corollary [7.1.2]) Just as in [2.4.5], the Gelfand commutativity criterion for the convolution algebra $C_{c}^{\infty}(K \backslash G / K)$ is that there should be an involutive anti-automorphism $\sigma$ on $G$, that is, $g \rightarrow g^{\sigma}$ such that $(g h)^{\sigma}=h^{\sigma} g^{\sigma}$ and $\left(g^{\sigma}\right)^{\sigma}=g$ for all $g, h \in G$, and stabilizing double cosets for $K$, that is, $(K g K)^{\sigma}=K g K$ for all $g \in G$. Then the convolution algebra $C_{c}^{\infty}(K \backslash G / K)$ is commutative. Again, transpose-conjugation $g^{\sigma}=g^{*}$ is such an anti-automorphism, because the Cartan decomposition $G=K A^{+} K$ from [1.2] shows that left and right $K$-invariant test functions are determined by their values on $A^{+}$, and $A^{+}$is pointwise fixed under transpose-conjugation.

By this commutativity and by the theorem [7.1.1], we have a commutative ring of compact operators, closed under adjoints, acting on a Hilbert space $V=L_{o}^{2}(\Gamma \backslash G / K)$ non-degenerately. The closed-ness under adjoints assures that any operator $T$ in that commutative algebra is a complex linear combination of self-adjoint operators in that algebra:

$$
T=\frac{T+T^{*}}{2}+i \cdot \frac{T-T^{*}}{2 i}
$$

For a non-zero self-adjoint operator $T$ on $V$, by the spectral theorem for self-adjoint compact operators [9.A], $V$ decomposes into finite-dimensional eigenspaces $V_{\lambda}$ for $T$ with non-zero, real eigenvalues, and an orthogonal complement $V^{\prime}$ :

$$
V=\left(\text { completion of } \bigoplus_{\lambda \neq 0} V_{\lambda}\right) \oplus V^{\prime}
$$

The eigenspaces $V_{\lambda}$ for non-zero eigenvalues $\lambda$ of a non-zero operator self-adjoint compact $T$ in that algebra are stabilized by every operator commuting with it, by the usual argument:

$$
T(S v)=(T S) v=(S T) v=S(T v)=S(\lambda v)=\lambda \cdot S v \quad\left(\text { for } v \in V_{\lambda}\right)
$$

Then $V_{\lambda}$ decomposes further into $S$-eigenspaces. By finite-dimensionality, we can do a downward induction to decompose $V_{\lambda}$ into simultaneous eigenspaces for all of $C_{c}^{\infty}(K \backslash G / K)$. The eigenvalue map from operators to eigenvalues on a given simultaneous eigenspace must be a ring homomorphism.

This motivates one formulation of a proof of the corollary. Let $X$ be the collection of not-identically-zero commutative-ring homomorphisms $\chi$ of $C_{c}^{\infty}(K \backslash G / K)$ to $\mathbb{C}$ such that

$$
V_{\chi}=\left\{v \in V: S v=\chi(S) \cdot v \text { for all } S \in C_{c}^{\infty}(K \backslash G / K)\right\}
$$

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is not $\{0\}$. Each $V_{\chi}$ is finite-dimensional, by the spectral theorem, since $\chi$ is not identically 0 . The orthogonal complement of the sum of all the $V_{\chi}$ is stable under the action of $C_{c}^{\infty}(K \backslash G / K)$, and every $\varphi \in C_{c}^{\infty}(K \backslash G / K)$ acts by 0 there. But this contradicts the non-degeneracy of the action, from the theorem. Thus, indeed, $V$ is the completion of the orthogonal direct sum of the $V_{\chi}$ with $\chi$ not identically 0 .

For non-zero $f$ in $V_{\chi}$, invoking the non-degeneracy, let $\varphi$ be a test function such that $\chi(\varphi) \neq 0$, and put $\eta=\varphi / \chi(\phi)$. Then $\eta \cdot f=f$. By smoothing, as in [14.5] for example, this entails that $f$ is smooth. This finishes the proof of [7.1.2].
Proof: (of corollary [7.1.3]) As in [4.2], the invariant Laplacian $\Delta$ on $\Gamma \backslash G / K$ or $G / K$ is the Casimir operator $\Omega$ restricted to right $K$-invariant functions. Since $\Omega$ commutes with both right and left translation action of $G$, it commutes with the integrated action of $\varphi \in C_{c}^{\infty}(G)$ :

$$
\Omega(\varphi \cdot f)(x)=\Omega \int_{G} \varphi(y) f(x y) d y=\int_{G} \varphi(y) \Omega_{x} f(x y) d y=\int_{G} \varphi(y)\left(\Omega_{x} f\right)(x y) d y=(\varphi \cdot \Omega f)(x)
$$

using the Gelfand-Pettis characterization [14.1] and the fact that $y \rightarrow(x \rightarrow \varphi(y) f(x y))$ is a continuous, compactly-supported, $C^{\infty}(\Gamma \backslash G)$-valued function and $\Omega$ is a continuous map of $C^{\infty}(\Gamma \backslash G)$ to itself. Thus, for right $K$-invariant smooth $f$ on $\Gamma \backslash G$ and $\varphi \in C_{c}^{\infty}(K \backslash G / K)$,

$$
\Delta(\varphi \cdot f)=\Omega(\varphi \cdot f)=\varphi \cdot(\Omega f)=\varphi \cdot \Delta f
$$

Thus, $\Delta$ stabilizes each of the finite-dimensional simultaneous eigenspaces $V_{\chi}$, each of which consists of smooth functions. The restriction of $\Delta$ to $V_{\chi}$ still has the symmetry proven in [6.5], so $V_{\chi}$ has an orthonormal basis of $\Delta$-eigenfunctions.

Further implications of the following corollary will be apparent in chapter 8:
[7.1.20] Corollary: The space of $L^{2}$ cuspforms has an orthonormal basis of cuspforms $f$ such that there is a test function $\varphi \in C_{c}^{\infty}(K \backslash G / K)$ such that $\varphi \cdot f=f$. Such a cuspform $f$ is smooth and of rapid decay in the sense that, on a standard Siegel set $\mathfrak{S}$, for every $q>0$,

$$
|f(g)|<_{q} \quad \eta(g)^{-q} \quad(\text { for all } g \in \mathfrak{S})
$$

Proof: By the commutativity of $\mathcal{H}=C_{c}^{\infty}(K \backslash G / K)$ from above, its irreducible representations consist of simultaneous eigenfunctions, and are one-dimensional. On each such $\mathcal{H}$ necessarily acts by an algebra homomorphism $\lambda: \mathcal{H} \rightarrow \mathbb{C}$. By the spectral consequences of the compactness above, the $\lambda^{t h}$ simultaneous eigenspace $V_{\lambda}$ has finite dimension.

That is, for $f \in V_{\chi}$ and for test function $\varphi, \varphi \cdot f=\lambda(\varphi) \cdot f$. From the non-degeneracy result in [7.1.1], there is $\varphi$ such that $\lambda(\varphi) \neq 0$. Replace $\varphi$ by $\varphi / \lambda(\varphi)$ for the result. Then [7.1.9] applies to $f=\varphi \cdot f$ to prove rapid decay. Smoothness follows as in [14.5] and [14.6].

$$
7.2 Z^{+} G L_{2}(k) \backslash G L_{2}(\mathbb{A})
$$

The next example shows how to adapt the argument of the previous section to adele groups $G=G L_{2}$ over number fields $k$. This incorporates spherical Hecke operators at good primes, and shows how to decouple bad primes. Use notation as in chapter 2 . Test functions $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ on $G_{\mathbb{A}}$ are compactly supported and smooth, where smoothness at archimedean places means indefinite continuous differentiability, and at finite places means local constancy. Test functions $\varphi$ act on functions $f$ on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ as usual by

$$
(\varphi \cdot f)(x)=\int_{G_{\mathbb{A}}} \varphi(y) f(x y) d y \quad\left(\text { for } x \in G_{\mathbb{A}}\right)
$$

The compactness of suitable operators on cuspforms is a global property, and the kernel function is a global object. Thus, we cannot expect purely local arguments to suffice. In particular, the (purely local) Hecke operators of chapter 2 are not quite adequate.

Let $K_{\infty}=\prod_{v \mid \infty} K_{v}$. As in chapter 2, for simplicity, we will eventually restrict attention to right $K_{\infty^{-}}$ invariant functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ rather than track $K_{\infty}$-types. Commensurately, we will eventually restrict attention to left and right $K_{\infty}$-invariant test functions on $G_{\mathbb{A}}$.

The main results of this section can be specialized to situations involving commutative Hecke algebras, which admit simultaneous eigenfunctions. The non-commutative Hecke algebras entering more general assertions mostly do not admit simultaneous eigenfunctions, and need a more complicated notion of irreducible representation, as follows.
[7.2.1] Theorem: $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ acts on square-integrable cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ by compact operators. The collection of such operators is closed under adjoints, and is non-degenerate in the sense that for every $f \in L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ there is $\varphi \in C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ such that $\varphi \cdot f \neq 0$. (Proof below.)
[7.2.2] Corollary: The space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ decomposes discretely with finite multiplicities into irreducibles for $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$. (Proof below.)

The assertion deserves clarification. Thinking of $A=C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, let $A$ be a (not necessarily commutative) associative algebra over $\mathbb{C}$, not necessarily having a unit. In the present context, a representation of $A$ is a Hilbert space $V$ on which $A$ acts by continuous linear operators, with the expected associativity

$$
\varphi \cdot(\psi \cdot v)=(\varphi * \psi) \cdot v \quad(\text { for } \varphi, \psi \in A, v \in V)
$$

where $*$ is the multiplication in $A$ (convolution in $C_{c}^{o}(G)$ ). The space $V$ is (topologically) irreducible (with respect to $A$ ) when it has no proper closed subspace stable under the action of $A$. An $A$-homomorphism $T: V \rightarrow W$ of $A$-representation spaces $V, W$ is a continuous linear map $T$ commuting with $A$, in the sense that $T(a v)=a T(v)$ for $a \in A$ and $v \in V$. The multiplicity of an $A$-irreducible $V$ in a larger $A$ representation (Hilbert) space $H$ is

$$
\text { multiplicity of } V \text { in } H=\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{A}(V, H)
$$

A form of Schur's lemma [9.D.12] shows that $\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{A}(V, V)=1$ for irreducibles $V$, and its corollary [9.D.14] shows that this removes potential ambiguity or ill-definedness in the definition of multiplicity. For Hilbert spaces, also

$$
\text { multiplicity of } V \text { in } H=\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{A}(H, V)
$$

because closed subspaces of Hilbert spaces admit orthogonal complements.
Yet, lacking further information about the irreducible representations of these non-commutative convolution algebras, or of the groups $G_{v}$ and $G_{\mathbb{A}}$, variant results for commutative Hecke algebras may be more immediately informative. Fix $K^{\prime}=K_{\infty} \prod_{v<\infty} K_{v}^{\prime}$, a compact subgroup of $G_{\mathbb{A}}$, with $K_{v}^{\prime}$ equal to $K_{v}=G L_{2}\left(\mathfrak{o}_{v}\right)$ for almost all $v$. The finite primes $v$ for which $K_{v}^{\prime}$ is $K_{v}$ are good primes, while the finite $v$ for which $K_{v}^{\prime}$ is strictly smaller than $K_{v}$ are bad primes. Let $S$ be the set of bad finite primes, of course depending on $K^{\prime}$. With notation differing from chapter 2, a suitable spherical Hecke algebra, depending on $K^{\prime}$, that does not attempt to do anything with bad primes is the collection $\mathcal{H}$ of left and right $K^{\prime}$-invariant test functions $\varphi$ on $G_{\mathbb{A}}$ which vanish at $g \in G_{\mathbb{A}}$ unless the $v^{t h}$ component $g_{v}$ is in $K_{v}^{\prime}$ for every $v \in S$.

Gelfand's criterion [2.4.5] and the $p$-adic and archimedean Cartan decompositions show that this Hecke algebra $\mathcal{H}$ is commutative, as in [2.4] and in the proof of [7.1.2].
[7.2.3] Corollary: The space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K^{\prime}}$ of right $K^{\prime}$-invariant cuspforms has an orthonormal basis of simultaneous eigenfunctions for the spherical Hecke algebra $\mathcal{H}$ attached to $K^{\prime}$, with each eigenspace finite-dimensional. The simultaneous eigenfunctions are smooth. (Proof below.)
[7.2.4] Corollary: The space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ of right $K_{\mathbb{A}}$-invariant square-integrable cuspforms decomposes into simultaneous eigenspaces for operators in the maximal spherical Hecke algebra $C_{c}^{\infty}\left(K_{\mathbb{A}} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$, with finite multiplicities. The simultaneous eigenfunctions are smooth. (Proof below.)
[7.2.5] Corollary: There is an orthonormal Hilbert-space basis for the space of $K_{\mathbb{A}}$-invariant squareintegrable cuspforms consisting of simultaneous eigenfunctions for the invariant Laplacians on the archimedean factors $G_{v}$. (Proof below.)

For strictly upper-triangular square matrices $x$ with entries in any field of characteristic zero, the series for the matrix exponential $e^{x}=\exp (x)=\sum_{\ell \geq 0} x^{\ell} / \ell!$ is finite. Thus, such matrices give an entirely algebraic notion of Lie algebra $\mathfrak{n}$ of $N$. Here,

$$
\mathfrak{n}_{\mathbb{A}}=\left\{\left(\begin{array}{cc}
0 & u \\
0 & 0
\end{array}\right): u \in \mathbb{A}\right\} \quad \mathfrak{n}_{v}=\left\{\left(\begin{array}{cc}
0 & u \\
0 & 0
\end{array}\right): u \in k_{v}\right\}
$$

Since $N$ is abelian, the exponential map is an isomorphism. Locally at archimedean places, the exponential map is a diffeomorphism $\exp : \mathfrak{n}_{v} \longrightarrow N_{v}$. Locally at non-archimedean places, it is a homeomorphism and preserves local-constant-ness. The discrete subgroup $\Lambda=n_{k} \subset \mathfrak{n}_{\mathbb{A}}$ is mapped isomorphically to $N_{k}$.

For $x \in G_{\mathbb{A}}$ write $x=n_{x} \cdot m_{x} \cdot k_{x}$ with $n_{x} \in N_{\mathbb{A}}, m_{x} \in M_{\mathbb{A}}$, and $k_{x} \in K_{\mathbb{A}}$. We can further decompose $m_{x}=m_{x}^{1} \cdot a_{x}$ with $m_{x} \in M^{1}$ and $a_{x}$ in the archimedean split component

$$
\left.A^{+}=\left\{\left(\begin{array}{cc}
t^{1 / r} & 0 \\
0 & 1
\end{array}\right): t>0\right\} \quad \text { (on the diagonal in } M_{\infty}=\prod_{v \mid \infty} M_{v}\right)
$$

Let $\eta$ be the height function as in chapter 2: in Iwasawa decomposition,

$$
\eta\left(n\left(\begin{array}{cc}
\alpha & 0 \\
0 & \beta
\end{array}\right) k\right)=\left|\frac{\alpha}{\beta}\right| \quad \quad \text { (idele norm, with } n \in N_{\mathbb{A}}, \alpha \in \mathbb{J}, \text { and } k \in K_{\mathbb{A}} \text { ) }
$$

Let

$$
A_{\tau}=\left\{a \in A^{+}: \eta(a) \geq \tau\right\}
$$

For this section, a slightly refined notion of Siegel set is convenient: for compact $C_{N} \subset N_{\mathbb{A}}$, compact $C_{M} \subset M^{1}$, and $\tau>0$ the corresponding Siegel set is

$$
\mathfrak{S}=\mathfrak{S}_{C_{N}, C_{M}, \tau}=Z^{+} \cdot C_{N} \cdot C_{M} \cdot A_{\tau} \cdot K_{\mathbb{A}}
$$

From [2.A], $M_{k} \subset M^{1}$, and $M_{k} \backslash M^{1}$ is compact. Thus, sufficiently large $C_{M}$ surjects to $M_{k} \backslash M^{1}$, and reduction theory [2.2] showed that for sufficiently small $\tau>0$ the Siegel set $\mathfrak{S}$ surjects to $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.
[7.2.6] Claim: A point in a Siegel set is well approximated by its $M$-component in an Iwasawa decomposition, in the sense that, for $x \in \mathfrak{S}_{C_{N}, C_{M}, \tau}$ with $\tau>0$, compact $C_{N} \subset N_{A}$, and compact $C_{M} \subset M^{1}$, there is another compact subset $C^{\prime}$ of $N_{\mathbb{A}}$ such that $x \in m_{x} \cdot C^{\prime} \cdot C_{M} \cdot K_{\mathbb{A}}$.
Proof: $x \in \mathfrak{S}$ gives

$$
x=n_{x} \cdot m_{x} \cdot k_{x} \in N_{\mathbb{A}} \cdot m_{x} \cdot K_{\mathbb{A}}=m_{x} \cdot m_{x}^{-1} N_{\mathbb{A}} m_{x} \cdot K
$$

The lower bound on $A_{\tau}$ gives a compact set $C^{\prime}$ in $N_{\mathbb{A}}$ depending only upon $\tau, C_{N}$, and $C_{M}$ such that

$$
m^{-1} C_{N} m \subset C^{\prime} \quad\left(\text { for } m=m_{1} a \text { with } m_{1} \in C_{M} \text { and } a \in A_{\tau}\right)
$$

In particular, $m_{x}^{-1} C_{N} m_{x} \subset C^{\prime}$. Thus, $x \in m_{x} \cdot C^{\prime} \cdot K_{\mathbb{A}}$ as claimed.

## 7. Discrete decomposition of cuspforms

Given test function $\varphi$, wind $u p$ the operator $f \rightarrow \varphi \cdot f$ along $N_{k}=\exp \left(\mathfrak{n}_{k}\right)$ : for $y \in G$,

$$
\begin{aligned}
(\varphi \cdot f)(y)=\int_{G_{\mathbb{A}}} \varphi(x) f(y x) d x & =\int_{G_{\mathbb{A}}} \varphi\left(y^{-1} x\right) f(x) d x=\int_{\exp (\Lambda) \backslash G_{\mathbb{A}}}\left(\sum_{\nu \in \mathbf{n}_{k}} \varphi\left(y^{-1} \exp (\nu) x\right)\right) f(x) d x \\
& =\int_{N_{k} \backslash G}\left(\sum_{\nu \in \Lambda} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)\right) f(x) d x
\end{aligned}
$$

The kernel function for the $(N \cap \Gamma)$-wound-up operator is the latter inner sum:

$$
K_{\varphi}(x, y)=\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)
$$

[7.2.7] Claim: For $x, y$ both in a Siegel set $\mathfrak{S}$ the $N_{k}$-wound-up kernel vanishes unless the $M$-components of the two are close. That is, for fixed compact $E \subset C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ and for $x, y \in \mathfrak{S}$, there is compact $C_{M} \subset M_{\mathbb{A}}$ such that if there exist $n \in N_{\mathbb{A}}$ and $\varphi \in E$ with $\varphi\left(y^{-1} \cdot n \cdot x\right) \neq 0$, then $m_{x} \in m_{y} \cdot C_{M}$.
Proof: From the previous claim, there is compact $C^{\prime} \subset N$ such that $m_{y}^{-1} y \in C^{\prime} \cdot K_{\mathbb{A}}$. A compact set of test functions has a common compact support $C_{G}$, because a compact set is bounded in the topological vector space sense, and a bounded subset of an LF-space such as $C_{c}^{\infty}(G)$ lies in some Fréchet limitand, by [13.8.5]. Non-vanishing of $\varphi\left(y^{-1} n x\right)$ implies $y^{-1} n x \in C_{G}$, so

$$
n x \in y \cdot C_{G} \subset m_{y} \cdot C^{\prime} \cdot K_{\mathbb{A}} \cdot C_{G} \subset m_{y} \cdot C_{G}^{\prime} \quad\left(\text { with } C_{G}^{\prime}=C^{\prime} K_{\mathbb{A}} C_{G}=\text { compact }\right)
$$

That is,

$$
C_{G}^{\prime} \ni m_{y}^{-1} \cdot n x=m_{y}^{-1} \cdot n n_{x} \cdot m_{x} \cdot k_{x}=\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x} \cdot k_{x}
$$

That is,

$$
\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x} \in C_{G}^{\prime} \cdot K_{\mathbb{A}}=\text { compact }
$$

Since $M_{\mathbb{A}}$ normalizes $N_{\mathbb{A}}$, the element $m_{y}^{-1} n n_{x} m_{y}$ is in $N_{\mathbb{A}}$. Since $N_{\mathbb{A}} \cap M_{\mathbb{A}}=\{1\}$, the multiplication map $N_{\mathbb{A}} \times M_{\mathbb{A}} \rightarrow N_{\mathbb{A}} M_{\mathbb{A}}$ is a homeomorphism. Thus, for the product $\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x}$ to lie in a compact set in $G_{\mathbb{A}}$ requires that $n$ lies in a compact set in $N_{\mathbb{A}}$ and $m$ lies in a compact set in $M_{\mathbb{A}}$. Thus, there is compact $C_{M} \subset M_{\mathbb{A}}$ such that $m_{y}^{-1} m_{x} \in C_{M}$, as claimed.
[7.2.8] Corollary: For $x, y$ in a fixed Siegel set, and for fixed compact $E \subset C_{c}^{\infty}\left(G_{\mathrm{A}}\right)$, there is a compact $C_{M} \subset M_{\mathbb{A}}$ such that, if $\varphi\left(y^{-1} n x\right) \neq 0$ for some $n \in N_{\mathbb{A}}$ and some $\varphi \in E$, then $m_{y}^{-1} x$ lies in $C_{M}$.
Proof: By the first of the two claims, there is a compact $C_{G}$ in $G_{\mathbb{A}}$ such that $x \in m_{x} \cdot C_{G}$. By the second, there is a compact $C_{M}$ in $M_{\mathbb{A}}$ such that $m_{x} \in m_{y} C_{M}$. Thus, $x \in m_{x} C_{G} \subset\left(m_{y} C_{M}\right) C_{G}^{\prime}$, rearranging to give the claim.

The notion of Schwartz function on an archimedean vectorspace such as $\mathfrak{n}_{\infty}=\prod_{v \mid \infty} \mathfrak{n}_{v}$ is as in [13.7]. On non-archimedean vectorspaces $\mathfrak{n}_{v}$, Schwartz functions are the same as test functions, namely, locally constant, compactly supported. Similarly, on the finite-adeles part of an adelic vectorspace, Schwartz functions are simply test functions, that is, locally constant, compactly supported. Then Schwartz functions on adelic vector spaces $\mathfrak{n}_{\mathbb{A}}$ are finite sums $\sum_{i} f_{\infty, i} \otimes f_{o, i}$ where the functions $f_{\infty, i}$ are Schwartz functions on the archimedean part, and the functions $f_{o, i}$ are Schwartz/test functions on the non-archimedean part. Topologies on such spaces are as in [6.2], [6.3], and as simpler examples in [13.7], [13.8], and [13.9].
[7.2.9] Claim: With $\omega_{y}=y^{-1} m_{y}$ and $\omega_{x, y}=m_{y}^{-1} x$, the functions $\nu \longrightarrow \varphi_{x, y}(\nu)=\varphi\left(\omega_{y} \cdot \exp (\nu) \cdot \omega_{x, y}\right)$ for $x, y$ in fixed Siegel set, $\varphi\left(y^{-1} n x\right) \neq 0, \varphi \in E$, constitute a compact subset of the Schwartz space $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$.
Proof: The left and right translation actions $G_{\mathbb{A}} \times G_{\mathbb{A}} \times C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ are continuous, by [6.4]. With fixed Siegel set $\mathfrak{S}$, by [7.2.6], [7.2.7], and [7.1.7], $\left\{\omega_{y}: y \in \mathfrak{S}\right\}$ and $\left\{\omega_{x, y}: x \in \mathfrak{S}, y \in \mathfrak{S}\right\}$ are compact. This gives compactness of the image of

$$
\left\{\omega_{y}: y\right\} \times\left\{\omega_{x, y}: x, y\right\} \times E \longrightarrow C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \quad\left(x, y \text { in fixed Siegel set, } \varphi\left(y^{-1} n x\right) \neq 0\right)
$$

## Garrett: Modern Analysis of Automorphic Forms

Since $N_{\mathbb{A}}$ is closed in $G_{\mathbb{A}}$, the restriction $\operatorname{map} C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(N_{\mathbb{A}}\right) \approx C_{c}^{\infty}\left(\mathfrak{n}_{\mathbb{A}}\right)$ is continuous. From [13.9.3], the inclusion $C_{c}^{\infty}\left(\mathfrak{n}_{\mathbb{A}}\right) \subset \mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$ of test functions to Schwartz functions is continuous, giving compactness of the image.

Poisson summation for the lattice $\mathfrak{n}_{k} \subset \mathfrak{n}_{\mathbb{A}}$ gives

$$
\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)=\sum_{\psi \in \mathfrak{n}_{k}^{*}} \int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}(\nu) \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu
$$

with suitably normalized measure, where $\mathfrak{n}_{k}^{*}$ is the collection of $\mathbb{C}^{\times}$-valued characters on $\mathfrak{n}_{\mathbb{A}} \approx \mathbb{A}$ trivial on the lattice $\mathfrak{n}_{k} \approx k$. By appendix [7.A], given a non-trivial character $\psi_{1}$ on $\mathfrak{n}_{\mathbb{A}} / \mathfrak{n}_{k} \approx \mathbb{A} / k$, every other character is of the form $\psi_{\xi}(\nu)=\psi_{1}(\xi \cdot \nu)$ with $\xi \in k$. Thus, choice of that non-trivial character identifies $\mathbb{A}^{*} \approx \mathfrak{n}_{\mathbb{A}}^{*}$ with $\mathbb{A} \approx \mathfrak{n}_{\mathbb{A}}$ and the dual lattice of $k \approx \mathfrak{n}_{k}$ with $k$ itself. Thus,

$$
\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)=\sum_{\xi \in k} \int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{1}(\xi \cdot \nu) \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu
$$

We have

$$
\int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{1}(\xi \cdot \nu) \cdot \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu=\int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{1}(\xi \cdot \nu) \cdot \varphi\left(y^{-1} m_{y} \cdot \exp \left(m_{y}^{-1} \nu m_{y}\right) \cdot m_{y}^{-1} x\right) d \nu
$$

Replacing $\nu$ by $m_{y} \nu m_{y}^{-1}$ and letting $\varphi_{x, y}(\nu)=\varphi\left(y^{-1} m_{y} \cdot \nu \cdot m_{y}^{-1} x\right)$,

$$
K_{\varphi}(x, y)=\delta\left(m_{y}\right) \sum_{\xi \in k} \widehat{\varphi}_{x, y}\left(\psi_{\xi}^{m_{y}}\right)
$$

where $\widehat{\varphi}_{x, y}\left(\psi_{\xi}^{m_{y}}\right)$ is the Fourier transform $\widehat{\varphi}_{x, y}$ of $\varphi_{x, y}$ evaluated at $\psi_{\xi}$, and $\psi_{\xi}^{m_{y}}(\nu)=\psi_{\xi}\left(m_{y} \nu m_{y}^{-1}\right)$.
[7.2.10] Theorem: Fix compact $E \subset C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ and Siegel set $\mathfrak{S}$. Given $q>0$ there is a uniform bound

$$
|(\varphi \cdot f)(y)|<_{q} \eta(y)^{-q} \cdot|f|_{L^{2}(\Gamma \backslash G)} \quad \text { (for all } y \in \mathfrak{S}, \text { for all } \varphi \in E \text {, for all } L^{2} \text { cuspforms } f \text { ) }
$$

Proof: With the self-duality identifications as above, Fourier transform is a continuous automorphism of $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$ to itself, by the archimedean case [13.15] and the $p$-adic case [13.17]. Thus, $\left\{\widehat{\varphi}_{x, y}: \varphi \in E, x, y \in \mathfrak{S}\right\}$ is a compact subset of $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$. The adelic Schwartz space $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$ is an LF-space, a strict colimit of Fréchet spaces, characterized as a countable ascending union of Fréchet subspaces described by restricting support at finite primes and by requiring uniform local constancy at finite primes. That is, for $U$ (large) compact open subgroup of the finite-adele part $\mathfrak{n}_{\mathbb{A}_{\text {fin }}}$ of $\mathfrak{n}_{\mathbb{A}}$, and for $H$ a (small) compact open subgroup of $\mathfrak{n}_{\mathbb{A}_{\text {fin }}}$, let $\mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$ be the space of $H$-invariant Schwartz functions supported on $\mathfrak{n}_{\mathbb{A}_{\infty}} \times U$. Then

$$
\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)=\bigcup_{H, U} \mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}=\operatorname{colim}_{H, U} \mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H} \quad \text { (as } H \text { shrinks and } U \text { grows) }
$$

There is a countable cofinal collection of subgroups $U$ and subgroups $H$, confirming that the adelic Schwartz space is an LF-space [13.8], [13.9]. In particular, a compact subset lies in some limitand $\mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$, by [13.8.5]. Thus, the compactness [7.2.8] implies that the Schwartz functions $\psi \rightarrow \widehat{\varphi}_{x, y}(\psi)$ are inside a compact subset of some $\mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$.

Thus, for $\widehat{\varphi}_{x, y}(\psi) \neq 0$, the finite-prime part $\xi_{\text {fin }}$ of $\xi$ is in some compact $U \subset \mathfrak{n}_{\mathbb{A}_{\text {fin }}} \approx \mathbb{A}_{\text {fin }}$. Thus, $\xi \in \frac{1}{h} \mathfrak{o}_{k}$ for some $0<h \in \mathbb{Z}$, and the collection of infinite-prime parts $\xi_{\infty}$ of such $\xi$ is a lattice $\Lambda^{*} \subset \mathfrak{n}_{\infty} \approx k_{\infty}$. Give the real vectorspace $k_{\infty}$ the real inner product

$$
\left\langle\xi, \xi^{\prime}\right\rangle=\sum_{v \mid \infty} \operatorname{Re}\left(\xi_{v} \cdot \bar{\xi}^{\prime}{ }_{v}\right)
$$

## 7. Discrete decomposition of cuspforms

with the complex conjugation to accommodate complex $k_{v}$. Let $|\cdot|$ be the associated norm. Thus, for $0 \neq \xi_{\infty} \in \Lambda^{*}$, we have $|\xi| \gg 0 \frac{1}{h}$, with an implied constant depending on the ring of integers $\mathfrak{o}$. In these terms, for $F \in \mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H},|F(\psi)|<_{r}(1+|\psi|)^{-\ell}$ for all $\ell>0$. Thus, for fixed Siegel set $\mathfrak{S}$ and $\varphi$ in fixed compact $E$, the compactness [7.2.9] gives a uniform implied constant depending only on $\ell$ so that

$$
\left|\widehat{\varphi}_{x, y}(\psi)\right|<_{\ell}(1+|\psi|)^{-\ell} \quad(\text { for all } x, y \in \mathfrak{S}, \text { for all } \varphi \in E)
$$

For $m_{y}=z m^{\prime} a_{t}$ with $z \in Z^{+}, m^{\prime} \in C_{M} \subset M^{1}$, and $t>0$, our parametrization of the archimedean split component $A^{+}$gives $\delta\left(m_{y}\right)=t$, and

$$
\left|\psi^{m_{y}}\right|<_{C_{M}} t^{1 / r_{o}} \cdot|\psi|=\delta\left(m_{y}\right)^{1 / r_{o}} \cdot|\psi|
$$

where $r_{o}$ is the number of archimedean completions of $k$ modulo complex conjugation. Thus, with fixed Siegel set,

$$
\delta\left(m_{y}\right)\left|\widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right)\right| \lll \delta\left(m_{y}\right) \cdot\left(1+\delta\left(m_{y}\right)^{1 / r_{o}} \cdot|\psi|\right)^{-\ell} \quad(\text { for all } x, y \in \mathfrak{S})
$$

[7.2.11] Claim: For fixed Siegel set $\mathfrak{S}$ and $r \gg 1$, the kernel $K_{\varphi}(x, y)$ with its $0^{\text {th }}$ Fourier component removed satisfies

$$
\left|K_{\varphi}(x, y)-\widehat{\varphi}_{x, y}\left(\psi_{0}\right)\right|<_{q}|\eta(y)|^{-q} \quad(\text { for } x, y \in \mathfrak{S})
$$

Proof: First, we claim that, with fixed lattice $\Lambda^{*} \subset k_{\infty}$ obtained by projecting $k \cap(k \infty \times U)$ to $k_{\infty}$, there is an implied constant such that

$$
\left.\eta(y)=\delta\left(m_{y}\right) \ll\left|\psi^{m_{y}}\right|^{r_{o}} \quad \text { (for all } y \in \mathfrak{S}, \text { for all } 0 \neq \psi \in k \cap U\right)
$$

Again, $\eta(y)=t$ for $y=n a_{t} k$ with $t>0$, and $\left|\psi^{m_{y}}\right| \ll \mathfrak{S} t^{1 / r_{o}} \cdot|\psi|$. Since the norms of non-zero elements of $\Lambda^{*}$ have a positive inf,

$$
\left|\psi^{m_{y}}\right|^{r_{o}} \gg_{\mathfrak{S}}\left(\eta(y)^{1 / r_{o}} \cdot|\psi|\right)^{r_{o}} \geq \eta(y) \cdot \inf _{0 \neq \lambda \in \Lambda^{*}}|\lambda|^{r_{o}}
$$

Since $\widehat{\varphi}_{x, y}$ is a Schwartz function, $\left|\widehat{\varphi}_{x, y}(\psi)\right| \ll \ell(1+|\psi|)^{-\ell}$ for every $\ell>0$. By the comparison of $\eta(y)$ to $\left|\psi^{m_{y}}\right|$, for $0 \neq \psi \in k \cap\left(k_{\infty} \times U\right)$,

$$
\begin{aligned}
\delta\left(m_{y}\right) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-\ell} & =\eta(y) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-\ell}=\eta(y) \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-(q+1) \cdot r_{o}-\left(\ell-(q+1) \cdot r_{o}\right)} \\
& \ll \eta(y) \cdot \eta(y)^{-(q+1)} \cdot\left(1+\left|\psi^{m_{y}}\right|\right)^{-\left(\ell-(q+1) \cdot r_{o}\right)}
\end{aligned}
$$

For $\ell$ sufficiently large depending on $q$, the sum of this over $0 \neq \psi$ converges, giving the claim.
[7.2.12] Claim: Cuspforms $f$ ignore the trivial Fourier component of $K_{\varphi}(x, y)$ :

$$
(\varphi \cdot f)(y)=\int_{\exp (\Lambda) \backslash G}\left(K_{\varphi}(x, y)-\widehat{\varphi}_{x, y}\left(\psi_{0}\right)\right) \cdot f(x) d x
$$

(Direct computation, identical to [7.1.11].)
The proof of [7.2.9] is almost complete. From above, for $y$ in a fixed Siegel set $\mathfrak{S}$ and for fixed test function $\varphi$, there is a compact $C_{M} \subset A^{+}$such that, for $\varphi\left(y^{-1} n x\right)$ to be nonzero, the Iwasawa $A^{+}$-component $m_{x}$ of $x$ must lie in $m_{y} \cdot C_{M}$. Thus,

$$
\left\{x \in \mathfrak{S}: \varphi\left(y^{-1} n x\right) \neq 0 \text { for some } n \in N\right\} \subset m_{y} C_{M} \cdot K
$$

Combining this with the estimate just obtained, for cuspform $f$,

$$
|(\varphi \cdot f)(y)|<_{r}|\eta(y)|^{-r} \cdot \int_{Z+G_{k} \backslash G_{k}\left(C \cdot m_{y} C_{M} \cdot K_{\mathbb{A}}\right)}|f(x)| d x
$$

By Cauchy-Schwarz-Bunyakowsky,

$$
\left(\int_{Z^{+} G_{k} \backslash G_{k}\left(C \cdot m_{y} C_{M} \cdot K_{\mathbb{A}}\right)}|f(x)| d x\right)^{2} \leq \int_{Z^{+}+G_{k} \backslash G_{k}\left(C \cdot m_{y} C_{M} \cdot K_{\mathbb{A}}\right)} 1 d x \quad \cdot \int_{Z^{+} G_{k} \backslash G_{k}\left(C \cdot m_{y} C_{M} \cdot K_{\mathbb{A}}\right)}|f(x)|^{2} d x \ll|f|_{L^{2}}^{2}
$$

This gives the desired decay, proving theorem [7.2.9].
We are getting closer to the compactness of the operators $f \rightarrow \varphi \cdot f$ on cuspforms $f$. Again, a collection $E$ of continuous functions on $G_{\mathbb{A}}$ or $Z^{+} G_{k} \backslash G_{\mathbb{A}}$ is (uniformly) equicontinuous when, given $\varepsilon>0$, there is a neighborhood $U$ of 1 in $G_{\mathbb{A}}$ such that

$$
|f(x)-f(y)|<\varepsilon \quad\left(\text { for all } f \in E, \text { for all } x \in G_{\mathbb{A}}, \text { for all } y \in x \cdot U\right)
$$

As in the previous section, for general reasons, we have
[7.2.13] Lemma: Let $\mathfrak{g}_{v}$ be the Lie algebra of $G_{v}$ for archimedean $v$. For $X \in \mathfrak{g}_{v}$, the left-derivative map

$$
C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \longrightarrow C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \quad \text { by } \quad \varphi \longrightarrow\left(\left.g \rightarrow \frac{d}{d t}\right|_{t=0} \varphi\left(e^{-t X} g\right)\right)
$$

is continuous. (Same proof as [7.1.12].)
[7.2.14] Corollary: For a compact set $E$ of test functions on $G_{\mathbb{A}}$, for a compact $C_{\mathfrak{g}}$ in $\mathfrak{g}=\mathfrak{g}_{v}$, and for $f$ ranging over cuspforms in the unit ball in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, there is a uniform implied constant such that

$$
\left.\left|\frac{d}{d t}\right|_{t=0}(\varphi \cdot f)\left(g e^{t X}\right) \right\rvert\, \ll 1 \quad \text { (for all } g \in G_{\mathbb{A}}, \text { for all } \varphi \in E \text {, for all } X \in C_{\mathfrak{g}} \text { ) }
$$

Proof: As in the proof of [7.1.13], differentiation of $\varphi \cdot f$ can be rewritten as a differentiation of $\varphi$, followed by action of the resulting function on $f$, by changing variables:

$$
\begin{gathered}
(X \cdot \varphi \cdot f)(x)=\left.\frac{d}{d t}\right|_{t=0} \int_{G} \varphi(y) f\left(x e^{t X} y\right) d y=\left.\frac{d}{d t}\right|_{t=0} \int_{G} \varphi\left(e^{-t X} y\right) f(x y) d y \\
\left.=\int_{G}\left(\left.\frac{d}{d t}\right|_{t=0} \varphi\left(e^{-t X} y\right)\right) f(x y) d y \quad \text { (replacing } y \text { by } e^{-t X} y\right)
\end{gathered}
$$

justifying interchange of differentiation and integration by the continuity of differentiation on test functions $\varphi$, and Gelfand-Pettis integral properties [14.1]. That is, $X \cdot \varphi \cdot f=\left(X^{\text {left }} \varphi\right) \cdot f$ with $X^{\text {left }}$ the left action. Since $\mathfrak{g}_{v}$ is a finite-dimensional real vector-space and the action is linear in $X$, this gives the continuity in $X \in \mathfrak{g}_{v}$. Thus, the collection of test functions $X^{\text {left }} \varphi$ with $X \in C_{\mathfrak{g}}$ and $\varphi \in E$ is again compact in $C_{c}^{\infty}(G)$. Thus, by the bound of theorem [7.2.10],

$$
\left.|(X \cdot \varphi \cdot f)(y)|=\mid\left(X^{\text {left }} \varphi\right) \cdot f\right)\left.(y)\left|<_{r} \eta(y)^{-r} \cdot\right| f\right|_{L^{2}} \quad\left(\text { for all } y \in \mathfrak{S}, X \in C_{\mathfrak{g}}, \varphi \in E\right)
$$

For large-enough Siegel set to cover the quotient, and any $r>0$, this gives

$$
\sup |X \cdot \varphi \cdot f(y)|<_{r}|f|_{L^{2}}
$$

proving uniform boundedness for $|f|_{L^{2}} \leq 1$.
The smoothing property of $f \rightarrow \varphi \cdot f$ as in [14.5] assures that each $\varphi \cdot f$ is smooth. Smoothness of $\varphi$ at finite places is uniform, since that of $\varphi$ is: with $\varphi$ left-invariant by compact open subgroup $K^{\prime} \subset G_{\mathbb{A}_{\text {fin }}}$, for $h \in K^{\prime}$,

$$
(\varphi \cdot f)(g \cdot h)=\int_{G} \varphi(x) f(g h x) d x=\int_{G} \varphi\left(h^{-1} x\right) f(g x) d x=\int_{G} \varphi(x) f(g x) d x=(\varphi \cdot f)(g)
$$

## 7. Discrete decomposition of cuspforms

A uniform bound on derivatives at archimedean places will imply uniform continuity:
[7.2.15] Lemma: Let $F$ be a smooth function on $G$, (uniformly) right $K^{\prime}$-invariant for some compact open subgroup $K^{\prime} \subset G_{\mathbb{A}_{\text {fin }}}$, with a uniform pointwise bound on all $X \cdot F$ with $X$ in a compact neighborhood $C_{\mathfrak{g}}$ of 0 in $\mathfrak{g}_{\infty}$, namely,

$$
|(X \cdot F)(x)| \leq B \quad\left(\text { for all } x \in G_{\mathbb{A}}, \text { all } X \in C_{\mathfrak{g}}\right)
$$

Then $F$ is uniformly continuous: for every $\varepsilon>0$ there is a neighborhood $U$ of 1 in $G_{\mathbb{A}}$ such that $|F(x)-F(y)|<\varepsilon$ for all $x \in G_{\mathbb{A}}$ and $y \in x U$.
Proof: The only issue is at archimedean places. Let $V$ be a small enough open containing 1 such that $V$ is contained in $\exp \left(C_{\mathfrak{g}}\right) \cdot K^{\prime}$, and that the exponential map on the archimedean part is injective to $V$. Let $y=x \cdot e^{s X}$ for $X \in C_{\mathfrak{g}}$ and $0 \leq s \leq 1$. By hypothesis, the function $h(t)=F\left(x \cdot e^{s X}\right)$ has

$$
h^{\prime}(s)=\left.\frac{d}{d t}\right|_{t=0} h(s+t)=F\left(x \cdot e^{s X} \cdot e^{t X}\right)
$$

which is bounded by $B$. From the mean value theorem, $\left|F\left(x \cdot e^{t X}\right)-F(x)\right| \leq t \cdot B$. Thus, for all $|t|<B \cdot \varepsilon$ we have the desired inequality.
[7.2.16] Corollary: For a compact set $E$ of test functions on $G_{\mathbb{A}}$, and for $f$ ranging over cuspforms in the unit ball in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, the family of images $\varphi \cdot f$ is (uniformly) equicontinuous on $G$.

We are almost done with the proof that $f \rightarrow \varphi \cdot f$ is compact on cuspforms. We again have a compactness lemma vaguely reminiscent of Arzela-Ascoli:
[7.2.17] Lemma: Let $E$ be a equicontinuous, uniformly bounded, set of functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. Then $E$ has compact closure in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. (Same proof as [7.1.16].)

Finally, we prove the theorem [7.2.1]. To summarize: the asymptotics of the kernels prove pointwise boundedness of the image of the unit ball $B$ of $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, and consideration of derivatives proves equicontinuity of the image of $B$. The faux-Arzela-Ascoli compactness lemma above proves compactness of the closure of $\{\varphi \cdot B: \varphi \in E\}$. Being integrated versions of right translations, these operators stabilize the subspace of cuspforms, as the latter is defined by a left integral condition. Thus, $\varphi$ maps the unit ball to a set with compact closure, so is a compact operator.

Adjoints are easily computed: letting $\langle$,$\rangle be the inner product on L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$,

$$
\begin{aligned}
\langle\varphi \cdot f, F\rangle & =\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \int_{G_{\mathbb{A}}} \varphi(x) f(y x) \bar{F}(y) d x d y=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \int_{G_{\mathbb{A}}} \varphi(x) f(y) \bar{F}\left(y x^{-1}\right) d x d y \\
& =\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} \int_{G_{\mathbb{A}}} f(y) \varphi\left(x^{-1}\right) \bar{F}(y x) d x d y=\int_{Z^{+} G_{k} \backslash G_{\mathbb{A}}} f(y) \overline{\varphi^{\vee} \cdot F}(y) d y
\end{aligned}
$$

where $\varphi^{\vee}(x)=\bar{\varphi}\left(x^{-1}\right)$ as suggested by the computation. The space of test functions is stable under the operation $\varphi \rightarrow \varphi^{\vee}$.

The non-degeneracy is [14.1.5], finishing the proof of theorem [7.2.1].
Proof: Now we can prove Corollary [7.2.2], that $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ decomposes as (the closure of) a direct sum of irreducible representations of $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, each occurring with finite multiplicity.
[7.2.18] Claim: Let $A$ be an adjoint-stable algebra of compact operators on a Hilbert space $H$, non-degenerate in the sense that for every non-zero $v \in H$ there is $a \in A$ with $a \cdot v \neq 0$. Then $H$ is (the completion of) an orthogonal direct sum of closed $A$-irreducible subspaces, and each isomorphism class of $A$-irreducible $V$ occurs with finite multiplicity.
Proof: To prove that $H$ is (the completion of) a direct sum of (closed) $A$-irreducibles, reduce to the case that $H$ has no proper irreducible $A$-subspaces, by replacing $H$ by the orthogonal complement to the sum of all irreducible $A$-subspaces. By non-degenerateness, there is a non-zero self-adjoint operator $T$ in $A$, since for a non-zero operator $S$ in $A$, either $S+S^{*}$ or $S-S^{*}$ is non-zero (and $S+S^{*}$ and ( $S-S^{*}$ )/i are self-adjoint).

From the spectral theorem for self-adjoint compact operators [9.A.6], there is a non-zero eigenvalue $\lambda$ of the (non-zero) self-adjoint compact operator $T$ on $H$, and the $\lambda$-eigenspace is finite-dimensional. Among
all $A$-stable closed subspaces choose one, $W$, such that the $\lambda$-eigenspace $W_{\lambda}$ is of minimal positive (finite) dimension. Let $w$ be a non-zero vector in $W_{\lambda}$. The closure of $A \cdot w$ is a closed subspace of $W$, and we claim that it is irreducible. Suppose that $\operatorname{closure}(A \cdot w)=X \oplus Y$ is a decomposition into mutually orthogonal, closed $A$-stable subspaces. With $w=w_{X}+w_{Y}$ the corresponding decomposition,

$$
\lambda w_{X}+\lambda w_{Y}=\lambda w=T w=T\left(w_{X}+w_{Y}\right)
$$

By the orthogonality and stability,

$$
\lambda w_{X}=T w_{X} \quad \text { and } \quad \lambda w_{Y}=T w_{Y}
$$

By the minimality of the $\lambda$-eigenspace in $W$, either $w_{X}=0$ or $w_{Y}=0$. That is, $\lambda w=T w \subset X$ or $\lambda w=T w \subset Y$. That is, since $\lambda \neq 0$, either $w \in X$ or $w \in Y$. Thus, either $A \cdot w \subset X$ or $A \cdot w \subset Y$, and likewise for the closures. But this implies that one or the other of $X, Y$ is 0 . This proves the irreducibility of the closure of $A \cdot w$, contradicting the assumption that $H$ had no irreducible $A$-subspaces.

For finite multiplicities: an irreducible $V$ is non-degenerate, since otherwise the subspace annihilated by all $a \in A$ would be a proper subspace. Thus, there is an operator $T \in A$ compact and self-adjoint on $H$, and non-zero on $V$. If the orthogonal direct sum $V \oplus \ldots \oplus V$ of $n$ copies appeared inside $H$ for arbitrarily large $n$, this would give $T$ infinite multiplicities of non-zero eigenvalues on $H$, contradicting the spectral theorem.

Then the proofs of [7.2.3] and [7.2.4] are special cases, where the algebra $\mathcal{H}$ of compact operators is designed to be commutative, so the notion of irreducibles simplifies to simultaneous eigenspace. Gelfand's criterion [2.4.5] and the $p$-adic and archimedean Cartan decompositions from [2.1] and [1.2] show that this Hecke algebra $\mathcal{H}$ is commutative, as in [2.4] and in the proof of [7.1.2]. For non-zero $f$ in a simultaneous eigenspace $V_{\chi}$ for $\mathcal{H}$, by non-degeneracy there is a test function $\varphi$ such that $\varphi \cdot f=\chi(\varphi) \cdot f$ and $\chi(\varphi) \neq 0$. With $\eta=\varphi / \chi(\phi)$, then $\eta \cdot f=f$. By smoothing, as in [14.5] for example, $f$ is smooth.

Just as in the proof of [7.1.3], the fact that the Casimir operators $\Omega_{v}$ on archimedean $G_{v}$ commute with left and right translation implies that $\Omega_{v}$ commutes with the action of $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, by integrating. On right $K_{v^{-}}$ invariant functions, $\Omega_{v}$ is the invariant Laplacian $\Delta_{v}$. Thus, each $\Delta_{v}$ stabilizes the simultaneous eigenspaces $V_{\chi}$ of $\mathcal{H}$, all of which are finite-dimensional, consisting of smooth functions. The restriction of $\Delta_{v}$ to $V_{\chi}$ is still symmetric, so by finite-dimensional spectral theory $V_{\chi}$ has a basis of $\Delta_{v}$-eigenfunctions.

As in the previous section, the implications of the following corollaries will be apparent in chapter 8:
[7.2.19] Corollary: The space of right $K_{\mathbb{A}}$-invariant $L^{2}$ cuspforms has an orthonormal basis of cuspforms $f$ such that there is a test function $\varphi \in C_{c}^{\infty}\left(K_{\mathbb{A}} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ such that $\varphi \cdot f=f$. Such a cuspform $f$ is smooth and of rapid decay in the sense that, on a standard Siegel set $\mathfrak{S}$, for every $q>0$,

$$
|f(g)|<_{q} \quad \eta(g)^{-q} \quad(\text { for all } g \in \mathfrak{S})
$$

Proof: The proof is identical to that of [7.1.20].
The proof of the following more general case is subtler than the previous:
[7.2.20] Corollary: The space of $L^{2}$ cuspforms has an orthonormal basis of cuspforms $f$ such that there is a test function $\varphi \in C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ such that $\varphi \cdot f=f$. Such a cuspform $f$ is of rapid decay in the sense that, on a standard Siegel set $\mathfrak{S}$, for every $q>0$,

$$
|f(g)|<_{q} \quad \eta(g)^{-q} \quad(\text { for all } g \in \mathfrak{S})
$$

Proof: The irreducible modules over $\mathcal{H}=C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ appearing in the space of $L^{2}$ cuspforms are merely finitedimensional, each occurring with finite multiplicity. Let $f$ be in a copy $V$ of an irreducible module for $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$.

Since $V$ is irreducible, it has no proper, topologically closed $\mathcal{H}$-stable subspace. Since $V$ is finitedimensional, all vector subspaces are closed. Thus, $\mathcal{H} \cdot f=V$. In particular, there is a test function $\varphi$ such that $\varphi \cdot f=f$. Then [7.2.20] shows that $f=\varphi \cdot f$ is of rapid decay in Siegel sets. Smoothness follows as in [14.5] and [14.6].
[7.2.21] Remark: The irreducibles for adjoint-closed rings of compact operators need not be finitedimensional, even though non-zero eigenspaces of self-adjoint compact operators are finite-dimensional. It is not so easy to give examples of such infinite-dimensional irreducibles. Perhaps the simplest examples would be the irreducibles in the decomposition of $L^{2}(\Gamma \backslash G)$ for compact quotients [7.B].

$$
7.3 Z^{+} G L_{r}(k) \backslash G L_{r}(\mathbb{A})
$$

The only new ingredient beyond the previous two sections is treatment of more complicated asymptotics for $G=G L_{r}$ for $r \geq 3$. The statements of results, and the proof mechanisms, are essentially identical to the previous section:
[7.3.1] Theorem: $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ acts on square-integrable cuspforms $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ by compact operators. The collection of such operators is closed under adjoints, and is non-degenerate in the sense that for every $f \in L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ there is $\varphi \in C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ such that $\varphi \cdot f \neq 0$. (Proof below.)
[7.3.2] Corollary: The space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ decomposes discretely with finite multiplicities into irreducibles for $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$. (Proof below.)

As in the previous sections, without further information about the irreducible representations of these noncommutative convolution algebras, or of the groups $G_{v}$ and $G_{\mathbb{A}}$, corollaries on commutative Hecke algebras are more immediately informative. Fix $K^{\prime}=K_{\infty} \prod_{v<\infty} K_{v}^{\prime}$, a compact subgroup of $G_{\mathbb{A}}$, with $K_{v}^{\prime}$ equal to $K_{v}=G L_{r}\left(\mathfrak{o}_{v}\right)$ for almost all $v$. The finite primes $v$ for which $K_{v}^{\prime}$ is $K_{v}$ are good primes, while the finite $v$ for which $K_{v}^{\prime}$ is strictly smaller than $K_{v}$ are bad primes. Let $S$ be the set of bad finite primes, of course depending on $K^{\prime}$. With notation differing from chapter 3 , a suitable spherical Hecke algebra, depending on $K^{\prime}$, that does not attempt to do anything with bad primes is the collection $\mathcal{H}$ of left and right $K^{\prime}$-invariant test functions $\varphi$ on $G_{\mathbb{A}}$ which vanish at $g \in G_{\mathbb{A}}$ unless the $v^{t h}$ component $g_{v}$ is in $K_{v}^{\prime}$ for every $v \in S$. Gelfand's criterion [2.4.5] and the $p$-adic and archimedean Cartan decompositions [3.2] show that $\mathcal{H}$ is commutative.
[7.3.3] Corollary: The space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)^{K^{\prime}}$ of right $K^{\prime}$-invariant cuspforms has an orthonormal basis of simultaneous eigenfunctions for the spherical Hecke algebra $\mathcal{H}$ attached to $K^{\prime}$, with each eigenspace finite-dimensional. The simultaneous eigenfunctions are smooth. (Proof below.)
[7.3.4] Corollary: The space $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$ of right $K_{\mathbb{A}}$-invariant square-integrable cuspforms decomposes into simultaneous eigenspaces for operators in the maximal spherical Hecke algebra $C_{c}^{\infty}\left(K_{\mathbb{A}} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$, with finite multiplicities. The simultaneous eigenfunctions are smooth. (Proof below.)
[7.3.5] Corollary: There is an orthonormal Hilbert-space basis for the space of $K_{\mathbb{A}}$-invariant squareintegrable cuspforms consisting of simultaneous eigenfunctions for the invariant Laplacians on the archimedean factors $G_{v}$. (Proof below.)

As in the previous section, the discussion needs a slightly refined version of Siegel set. Let $\Phi^{o}$ be the collection of positive simple roots (composed with norms), namely, the characters on diagonal matrices given by

$$
\alpha_{j}\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right)=\left|\frac{m_{j}}{m_{j+1}}\right| \quad(\text { for } 1 \leq j \leq r-1, \text { with idele norm })
$$

Let $B=P^{\text {min }}=P^{1,1, \ldots, 1}$ be the minimal parabolic. Put

$$
M^{1}=\left\{\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right) \in M_{\mathbb{A}}^{B}:\left|m_{j}\right|=1, \text { for all } j\right\}
$$

and

$$
A^{+}=\left\{a_{t}=\left(\begin{array}{ccc}
t_{1} & & \\
& \ddots & \\
& & t_{r}
\end{array}\right): \text { all } t_{j}>0, \text { diagonally in archimedean } \prod_{v \mid \infty} M_{v}^{B}\right\}
$$

Thus, for $m^{\prime} \in M^{1}, \alpha_{j}\left(m^{\prime} a_{t}\right)=\left(t_{j} / t_{j+1}\right)^{r_{o}}$ where $r_{o}$ is the number of isomorphism classes of archimedean completions of $k$. For $\tau>0$, let

$$
A_{\tau}=\left\{a \in A^{+}: \alpha(a) \geq \tau, \text { for all } \alpha \in \Phi^{o}\right\}
$$

For $\tau>0$, compact $C_{N} \subset N_{\mathbb{A}}^{B}$, compact $C_{M} \subset M^{1}$, and compact subgroup $K_{\mathbb{A}}$ of $G_{\mathbb{A}}$, the corresponding Siegel set attached to the minimal parabolic $B$ is

$$
\mathfrak{S}=\mathfrak{S}\left(C_{N}, C_{M}, \tau\right)=C_{N} \cdot C_{M} \cdot A_{\tau} \cdot K_{\mathbb{A}}
$$

## Garrett: Modern Analysis of Automorphic Forms

The compactness [2.A] of $\mathbb{J}^{1} / k^{\times}$and reduction theory [3.3] showed that for sufficiently large compact $C_{N}$ and $C_{M}$ and sufficiently small $\tau>0$ the corresponding Siegel set surjects to $Z^{+} G_{k} \backslash G_{\mathbb{A}}$.

In the following, unadorned $N, M$ will be $N=N^{B}$ and $M=M^{B}$, and the unipotent radical and standard Levi components for other parabolics $P$ will be $N^{P}$ and $M^{P}$.
[7.3.6] Claim: Fix a Siegel set $\mathfrak{S}=\mathfrak{S}\left(C_{N}, C_{M}, \tau\right)$ with compact $C_{N} \subset N_{\mathbb{A}}$, compact $C_{M} \subset M^{1}$, and $\tau>0$. Write $x \in \mathfrak{S}$ as $x=n_{x} m_{x} k_{x}$ with $n_{x} \in C_{N}, m_{x}=m_{x}^{\prime} a_{t}$ with $m_{x}^{\prime} \in C_{M}$ and $a_{t} \in A_{\tau}$, and $k_{x} \in K_{\mathbb{A}}$. Then there is a compact subset $C^{\prime}$ of $N_{\mathbb{A}}$ such that $x \in m_{x} \cdot C^{\prime} \cdot K_{\mathbb{A}}$.
Proof: Rewrite the Iwasawa decomposition of $x \in \mathfrak{S}$ as

$$
\begin{gathered}
x=n_{x} \cdot m_{x} \cdot k_{x} \in C_{N} \cdot m_{x} \cdot K_{\mathbb{A}}=m_{x} \cdot m_{x}^{-1} C_{N} m_{x} \cdot K_{\mathbb{A}} \\
\subset m_{x} \cdot a_{t}^{-1}\left(\left(m_{x}^{\prime}\right)^{-1} C_{N} m_{x}^{\prime}\right) a_{t} \cdot K_{\mathbb{A}} \subset m_{x} \cdot a_{t}^{-1}\left(C_{M}^{-1} C_{N} C_{M}\right) a_{t} \cdot K_{\mathbb{A}}
\end{gathered}
$$

Being a continuous image of a compact set, $D=C_{M}^{-1} C_{N} C_{M}$ is a compact subset of $N_{\mathbb{A}}$. Now we claim that because $a_{t} \in A_{\tau}$, the union of all conjugates $a_{t}^{-1} D a_{t}$ is contained in a compact set. Indeed, for $u=\left\{u_{v}\right\} \in N_{\mathbb{A}}$, since $a_{t}$ is purely archimedean, conjugation by $a_{t}$ does not alter the non-archimedean components $u_{v}$ of $u$. At archimedean places $v$, for $i<j$, the $i j^{t h}$ entry of $a_{t}^{-1}\left(u_{v}\right) a_{t}$ is $t_{i}^{-1} t_{j}$ times the $i j^{t h}$ entry of $u_{v}$, and

$$
\frac{t_{j}}{t_{i}}=\frac{t_{j}}{t_{j-1}} \cdot \frac{t_{j-1}}{t_{j-2}} \cdot \ldots \cdot \frac{t_{i+1}}{t_{i}}=\left(\chi_{i}\left(a_{t}\right) \chi_{i+1}\left(a_{t}\right) \ldots \chi_{j-1}\left(a_{t}\right)\right)^{-1 / r_{o}} \leq\left(\tau^{j-i-1}\right)^{-1 / r_{o}}
$$

where $r_{o}$ is the number of isomorphism classes of archimedean completions of $k$. Those entries are bounded on $D$, so the entries of the conjugate are uniformly bounded for all $a_{t} \in A_{\tau}$. Thus, $\bigcup_{a \in A_{\tau}} a^{-1} D a$ is contained in a compact subset $C^{\prime}$ of $N_{\mathbb{A}}$, as claimed.

For strictly upper-triangular square matrices $x$ with entries in any field of characteristic zero, the series for the matrix exponential $e^{x}=\exp (x)=\sum_{\ell \geq 0} x^{\ell} / \ell!$ is finite. Thus, the Lie algebra $\mathfrak{n}^{P}$ of the unipotent radical $N^{P}$ for any standard parabolic $P$ has a purely algebraic sense:

$$
\mathfrak{n}^{P}=\left\{n-\operatorname{by}-n x: \exp (x) \in N^{P}\right\}
$$

For example, for the minimal parabolic $B=P^{\text {min }}$, the Lie algebra $\mathfrak{n}=\mathfrak{n}^{B}$ is all upper-triangular matrices with zeros on the diagonal. For maximal proper $P=P^{i, r-i}$,

$$
\mathfrak{n}^{P}=\left\{\left(\begin{array}{cc}
0 & u \\
0 & 0
\end{array}\right)\right\} \quad(\text { where } u \text { is } i \text {-by- }(r-i))
$$

In the latter case, because $N^{P}$ is abelian, the exponential map is an isomorphism $\exp : \mathfrak{n}^{P} \longrightarrow N^{P}$. For all parabolics, the discrete additive subgroup $\mathfrak{n}_{k}^{P} \subset \mathfrak{n}_{\mathbb{A}}^{P}$ exponentiates to $N_{k}^{P}$. For test function $\varphi \in C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, we can wind up the integral for $f \rightarrow \varphi \cdot f$ along the unipotent radical $N_{k}$ of the standard minimal parabolic $B:$ for $y \in G_{\mathbb{A}}$,

$$
\begin{gathered}
(\varphi \cdot f)(y)=\int_{G_{\mathbb{A}}} \varphi(x) f(y x) d x=\int_{G_{\mathbb{A}}} \varphi\left(y^{-1} x\right) f(x) d x=\int_{N_{k} \backslash G_{\mathbb{A}}}\left(\sum_{\gamma \in \exp \left(\mathfrak{n}_{k}^{P}\right)} \varphi\left(y^{-1} \gamma x\right)\right) f(x) d x \\
=\int_{N_{k} \backslash G_{\mathbb{A}}}\left(\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)\right) f(x) d x
\end{gathered}
$$

The kernel function for this wound-up form of the operator is the latter left- $N_{k}$-invariant inner sum:

$$
K_{\varphi}(x, y)=\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)
$$

As just above, write $B$-Iwasawa decompositions for $x \in \mathfrak{S}$ as $x=n_{x} m_{x} k_{x}$ with $n_{x} \in C_{N}, m_{x}=m_{x}^{\prime} a_{t}$ with $m_{x}^{\prime} \in C_{M}$ and $a_{t} \in A_{\tau}$, and $k_{x} \in K_{\mathbb{A}}$.

## 7. Discrete decomposition of cuspforms

[7.3.7] Claim: For a fixed Siegel set $\mathfrak{S}$, fixed compact $E \subset C_{c}^{\infty}(G)$, there is compact $C_{M}^{\prime} \subset M_{\mathbb{A}}$ such that if there exist $n \in N_{\mathbb{A}}$ and $\varphi \in E$ with $\varphi\left(y^{-1} \cdot n \cdot x\right) \neq 0$ for any $x, y \in \mathfrak{S}$, then $m_{x} \in m_{y} \cdot C_{M}^{\prime}$. That is, $K_{\varphi}(x, y)=0$ for all $x, y \in \mathfrak{S}$ and all $\varphi \in E$ unless $m_{x} \in m_{y} \cdot C_{M}^{\prime}$.
Proof: From the previous claim, there is compact $D \subset N_{\mathbb{A}}$ such that $\left(m_{y}\right)^{-1} y \in D \cdot K_{\mathbb{A}}$. A compact set of test functions has a common compact support $C_{G}$, because a compact set is bounded in the topological vector space sense, and a bounded subset of an LF-space such as $C_{c}^{\infty}(G)$ lies in some Fréchet limitand, by [13.8.5]. Non-vanishing of $\varphi\left(y^{-1} n x\right)$ implies $y^{-1} n x \in C_{G}$, so

$$
\left.n x \in y \cdot C_{G} \subset m_{y} \cdot D \cdot K_{\mathbb{A}} \cdot C_{G} \subset m_{y} \cdot C_{G}^{\prime} \quad \quad \text { (with } C_{G}^{\prime}=D K_{\mathbb{A}} C_{G}=\text { compact }\right)
$$

That is,

$$
C_{G}^{\prime} \ni m_{y}^{-1} \cdot n x=m_{y}^{-1} \cdot n n_{x} \cdot m_{x} \cdot k_{x}=\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x} \cdot k_{x}
$$

That is,

$$
\left(m_{y}^{-1} n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x} \in C_{G}^{\prime} \cdot K_{\mathbb{A}}=\text { compact }
$$

Since $M$ normalizes $N$, the element $m_{y}^{-1} n n_{x} m_{y}$ is in $N$. Since $N_{\mathbb{A}} \cap M_{\mathbb{A}}=\{1\}$, the multiplication map $N_{\mathbb{A}} \times M_{\mathbb{A}} \rightarrow B_{\mathbb{A}}$ is a homeomorphism. Thus, for the product $\left.\left(m_{y}^{-1}\right) n n_{x} m_{y}\right) \cdot m_{y}^{-1} m_{x}$ to lie in a compact set in $G_{\mathbb{A}}$ requires that its $N$-component lies in a compact set in $N_{\mathbb{A}}$ and its $M$-component lies in a compact set in $M_{\mathbb{A}}$. Thus, there is compact $C_{M}^{\prime} \subset M_{\mathbb{A}}$ such that $m_{y}^{-1} m_{x} \in C_{M}^{\prime}$, as claimed.
[7.3.8] Corollary: For fixed Siegel set $\mathfrak{S}$ and fixed compact $E \subset C_{c}^{\infty}(G)$, there is a compact $C_{M}^{\prime} \subset M_{\mathbb{A}}$ such that, if $\varphi\left(y^{-1} n x\right) \neq 0$ for some $x, y \in \mathfrak{S}$, some $n \in N_{\mathbb{A}}$ and some $\varphi \in E$, then $m_{y}^{-1} x \in C_{M}^{\prime}$.
Proof: By [7.1.5], there is a compact $C_{G}$ in $G$ such that $x \in m_{x} \cdot C_{G}$. By [7.1.6], there is a compact $C_{M}^{\prime}$ in $M_{\mathbb{A}}$ such that $m_{x} \in m_{y} C_{M}^{\prime}$. Thus, $x \in m_{x} C_{G} \subset m_{y} C_{M}^{\prime} C_{G}^{\prime}$, rearranging to give the claim.

As in the previous section, the notion of Schwartz function on an archimedean vectorspace such as $\mathfrak{n}_{\infty}=\oplus_{v \mid \infty} \mathfrak{n}_{v}$ is as in [13.7], and on non-archimedean vectorspaces $\mathfrak{n}_{v}$, Schwartz functions are the same as test functions, namely, locally constant, compactly supported. Similarly, on the finite-adeles part of an adelic vectorspace, Schwartz functions are simply test functions, that is, locally constant, compactly supported. Then Schwartz functions on adelic vector spaces $\mathfrak{n}_{\mathbb{A}}$ are finite sums $\sum_{i} f_{\infty, i} \otimes f_{o, i}$ where the functions $f_{\infty, i}$ are Schwartz functions on the archimedean part, and the functions $f_{o, i}$ are Schwartz/test functions on the non-archimedean part. Topologies on such spaces are as in [6.2], [6.3], and as simpler examples in [13.7], [13.8], and [13.9].
[7.3.9] Claim: With $\omega_{y}=y^{-1} m_{y}$ and $\omega_{x, y}=m_{y}^{-1} x$, the functions $\nu \longrightarrow \varphi_{x, y}(\nu)=\varphi\left(\omega_{y} \cdot \exp (\nu) \cdot \omega_{x, y}\right)$ for $x, y$ in fixed Siegel set, $\varphi\left(y^{-1} n x\right) \neq 0, \varphi \in E$, constitute a compact subset of the Schwartz space $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$.
Proof: The left and right translation actions $G_{\mathbb{A}} \times G_{\mathbb{A}} \times C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ are continuous, by [6.4]. With fixed Siegel set $\mathfrak{S}$, by [7.2.6], [7.2.7], and [7.1.7], $\left\{\omega_{y}: y \in \mathfrak{S}\right\}$ and $\left\{\omega_{x, y}: x \in \mathfrak{S}, y \in \mathfrak{S}\right\}$ are compact. This gives compactness of the image of

$$
\left\{\omega_{y}: y\right\} \times\left\{\omega_{x, y}: x, y\right\} \times E \longrightarrow C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \quad\left(x, y \text { in fixed Siegel set, } \varphi\left(y^{-1} n x\right) \neq 0\right)
$$

Since $N_{\mathbb{A}}$ is closed in $G_{\mathbb{A}}$, the restriction map $C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \rightarrow C_{c}^{\infty}\left(N_{\mathbb{A}}\right) \approx C_{c}^{\infty}\left(\mathfrak{n}_{\mathbb{A}}\right)$ is continuous. We are fortunate that test functions are characterized by compact support together with purely local smoothness properties, so that $C_{c}^{\infty}\left(N_{\mathbb{A}}\right) \approx C_{c}^{\infty}\left(\mathfrak{n}_{\mathbb{A}}\right)$. From [13.9.3], the inclusion $C_{c}^{\infty}\left(\mathfrak{n}_{\mathbb{A}}\right) \subset \mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$ of test functions to Schwartz functions is continuous, giving compactness of the image.

Poisson summation for the lattice $\mathfrak{n}_{k} \subset \mathfrak{n}_{\mathbb{A}}$ gives

$$
\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)=\sum_{\psi \in\left(\mathfrak{n}_{k}\right)^{*}} \int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}(\nu) \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu
$$

with suitably normalized measure, where $\left(\mathfrak{n}_{k}\right)^{*}$ is the collection of $\mathbb{C}^{\times}$-valued characters on $\mathfrak{n}_{\mathbb{A}}$ trivial on the lattice $\mathfrak{n}_{k}$. As in appendix [7.A], we can identify the dual $\left(\mathfrak{n}_{\mathbb{A}}\right)^{*}$ with $\mathfrak{n}_{\mathbb{A}}$ and the dual lattice $\left(\mathfrak{n}_{k}\right)^{*}$ with $\mathfrak{n}_{k}$.

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One reasonable identification is as follows. Make an $\mathbb{A}$-valued pairing on $\mathfrak{n}_{\mathbb{A}}$ by $\langle\nu, \xi\rangle=\sum_{i<j} \nu_{i j} \xi_{i j}$, and for fixed non-trivial character $\psi_{1}$ of $\mathbb{A} / k$ put

$$
\psi_{\xi}(\nu)=\psi_{1}(\langle\nu, \xi\rangle) \quad\left(\text { with } \xi, \nu \text { in } \mathfrak{n}_{\mathbb{A}}\right)
$$

Thus,

$$
\sum_{\nu \in \mathfrak{n}_{k}} \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right)=\sum_{\xi \in k} \int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{\xi}(\nu) \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu
$$

We have

$$
\left.\int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{\xi}(\nu) \cdot \varphi\left(y^{-1} \cdot \exp (\nu) \cdot x\right) d \nu=\int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{\xi}(\nu) \cdot \varphi\left(y^{-1} m_{y} \cdot \exp \left(m_{y}^{-1}\right) \nu m_{y}\right) \cdot m_{y}^{-1} x\right) d \nu
$$

Replacing $\nu$ by $m_{y} \nu m_{y}^{-1}$ and letting $\varphi_{x, y}(\nu)=\varphi\left(y^{-1} m_{y} \cdot \nu \cdot m_{y}^{-1} x\right)$,

$$
K_{\varphi}(x, y)=\delta_{B}\left(m_{y}\right) \sum_{\xi \in k} \widehat{\varphi}_{x, y}\left(\psi_{\xi}^{m_{y}}\right)
$$

where $\widehat{\varphi}_{x, y}\left(\psi_{\xi}^{m_{y}}\right)$ is the Fourier transform $\widehat{\varphi}_{x, y}$ of $\varphi_{x, y}$ along $\mathfrak{n}_{\mathbb{A}}$, evaluated at $\psi_{\xi}$, where $\psi_{\xi}^{m_{y}}(\nu)=$ $\psi_{\xi}\left(m_{y} \nu\left(m_{y}\right)^{-1}\right)$, and where $\delta_{B}$ is the modular function of $B_{\mathbb{A}}$.
[7.3.10] Theorem: For $y$ in a fixed Siegel set attached to the minimal parabolic $B$, and for $\varphi$ in a fixed compact $E \subset C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, for every $q>0$ there is an implied constant depending only on such that, for every $L^{2}$ cuspform $f$,

$$
|(\varphi \cdot f)(y)|<_{q}\left(\inf _{\alpha \in \Phi^{o}} \alpha\left(m_{y}\right)\right)^{-q} \cdot|f|_{L^{2}}
$$

Proof: We need several preliminary results:
[7.3.11] Claim: Let $P$ be a standard maximal proper parabolic. For every character $\psi$ of $\mathfrak{n}_{\mathbb{A}}$ trivial on $\mathfrak{n}_{\mathbb{A}}^{P}$, the Fourier component $\widehat{\varphi}_{x, y}(\psi)$ is left $N_{\mathbb{A}}^{P}$-invariant in $x$, and therefore integrates to 0 against every cuspform.
Proof: For $n \in N_{\mathbb{A}}^{P}$, replace $x$ by $n x$ in the original integral defining $\widehat{\varphi}_{x, y}(\psi)$, obtaining

$$
\widehat{\varphi}_{n x, y}\left(\psi_{0}\right)=\int_{\mathfrak{n}_{\mathbb{A}}^{P}} \bar{\psi}_{0}(\nu) \cdot \varphi\left(y^{-1} \exp (\nu) \cdot n x\right) d \nu=\int_{\mathfrak{n}_{\mathbb{A}}^{P}} \bar{\psi}_{0}(\nu) \cdot \varphi\left(y^{-1} \exp \left(\nu+\nu^{\prime}\right) \cdot x\right) d \nu
$$

where $\nu^{\prime}$ is a continuous function of $\nu$ determined by the obvious $\exp \left(\nu+\nu^{\prime}\right)=\exp (\nu) \cdot n$. That $\nu^{\prime}$ is in the subalgebra $\mathfrak{n}^{P}$ rather than merely in $\mathfrak{n}$ follows from a computation in block decompositions of the appropriate size:

$$
\exp (\nu) \cdot n=\exp \left(\begin{array}{cc}
\nu_{11} & \nu_{12} \\
0 & \nu_{22}
\end{array}\right) \cdot\left(\begin{array}{cc}
1 & v \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
e^{\nu_{11}} & b \\
0 & e^{\nu_{22}}
\end{array}\right) \cdot\left(\begin{array}{cc}
1 & v \\
0 & 1
\end{array}\right)=\left(\begin{array}{cc}
e^{\nu_{11}} & e^{\nu_{11}} \cdot v+b \\
0 & e^{\nu_{22}}
\end{array}\right)
$$

for some block $b$. This is still of the form $\exp \left(\begin{array}{cc}\nu_{11} & \nu_{12}^{\prime} \\ 0 & \nu_{22}\end{array}\right)$ for some $\nu_{12}^{\prime}$ depending on $\nu$ and $n$, and we take

$$
\nu^{\prime}=\left(\begin{array}{cc}
0 & \nu_{12}^{\prime} \\
0 & 0
\end{array}\right) \in \mathfrak{n}_{\alpha} \quad \text { (in suitable blocks) }
$$

Replacing $\nu$ by $\nu-\nu^{\prime}$ in the integral gives

$$
\int_{\mathfrak{n}_{\mathbb{A}}} \bar{\psi}_{0}(\nu) \cdot \varphi\left(y^{-1} \exp (\nu) \cdot n x\right) d \nu=\psi_{0}\left(\nu^{\prime}\right) \cdot \widehat{\varphi}_{x, y}\left(\psi_{0}\right)=\widehat{\varphi}_{x, y}\left(\psi_{0}\right)
$$

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proving the left $N_{\mathbb{A}}^{P}$-invariance in $x$. The corresponding integral is

$$
\begin{gathered}
\int_{N_{k} \backslash G_{\mathbb{A}}} \widehat{\varphi}_{x, y}\left(\psi_{0}^{m_{y}}\right) \cdot f(x) d x=\int_{N_{k} N_{\mathbb{A}}^{P} \backslash G_{\mathbb{A}}} \int_{N_{k}^{P} \backslash N_{\mathbb{A}}^{P}} \widehat{\varphi}_{n x, y}\left(\psi_{0}^{m_{y}}\right) \cdot f(n x) d n d x \\
=\int_{N_{k} N_{\mathbb{A}}^{P} \backslash G_{\mathbb{A}}} \widehat{\varphi}_{x, y}\left(\psi_{0}^{m_{y}}\right) \cdot\left(\int_{N_{k}^{P} \backslash N_{\mathbb{A}}^{P}} f(n x) d n\right) d x
\end{gathered}
$$

and the inner integral is 0 because $f$ is a cuspform.
Thus, for cuspforms $f$,

$$
(\varphi \cdot f)(y)=\delta_{B}\left(m_{y}\right) \sum_{\psi \in \mathfrak{n}_{k}}^{*} \int_{N_{k} \backslash G_{\mathbb{A}}} \widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right) f(x) d x
$$

where the $\sum^{*}$ is to mean that the sum omits $\psi \in \mathfrak{n}_{k}$ such that $\left.\psi\right|_{\mathfrak{n}_{\mathbb{A}}}=1$ identically for some maximal standard proper parabolic $P$.

Returning to the proof of theorem [7.3.10], Fourier transform is a continuous map of $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$ to itself, so Fourier transform maps the compact set of [7.3.9] to a compact set of Schwartz functions
$\left\{\widehat{\varphi}_{x, y}: x, y \in \mathscr{S}, \varphi \in E\right\} \subset \mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$. The adelic Schwartz space $\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)$ is an LF-space, a strict colimit of Fréchet spaces, characterized as a countable ascending union of Fréchet subspaces described by restricting support at finite primes and by requiring uniform local constancy at finite primes. That is, for $U$ (large) compact open (additive) subgroup of the finite-adele part $\mathfrak{n}_{\mathbb{A}_{\text {fin }}}$ of $\mathfrak{n}_{\mathbb{A}}$, and for $H$ a (small) compact open (additive) subgroup of $\mathfrak{n}_{\mathbb{A}_{\text {fin }}}$, let $\mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$ be the space of $H$-invariant Schwartz functions supported on $\mathfrak{n}_{\mathbb{A}_{\infty}} \times U$. Then

$$
\mathscr{S}\left(\mathfrak{n}_{\mathbb{A}}\right)=\bigcup_{H, U} \mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}=\operatorname{colim}_{H, U} \mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H} \quad \quad \text { (as } H \text { shrinks and } U \text { grows) }
$$

There is a countable cofinal collection of subgroups $U$ and subgroups $H$, certifying the LF-space structure [13.8], [13.9]. In particular, a compact subset lies in some limitand $\mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$, by [13.8.5]. Thus, the compactness [7.2.8] implies that all the Schwartz functions $\psi \rightarrow \widehat{\varphi}_{x, y}(\psi)$ lie in a compact subset of some $\mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$.

Thus, for $\widehat{\varphi}_{x, y}\left(\psi_{\xi}\right) \neq 0$, the finite-prime part $\xi_{\text {fin }}$ of $\xi$ is in some compact $U \subset \mathfrak{n}_{\mathbb{A}_{\text {fin }}}$. Thus, $\xi \in \frac{1}{h} \mathfrak{n}_{\mathfrak{o}}$ for some $0<h \in \mathbb{Z} \ell$, where $\mathfrak{n}_{\mathfrak{o}}$ is the collection of elements of $\mathfrak{n}_{k}$ with entries in the ring of algebraic integers $\mathfrak{o}$ of $k$. The collection of infinite-prime parts $\xi_{\infty}$ of such $\xi$ is a lattice $\Lambda^{*} \subset \mathfrak{n}_{\infty}$. Give the finite-dimensional $\mathbb{R}$-vectorspace $\mathfrak{n}_{\infty}$ a positive-definite inner product

$$
\left\langle\xi, \xi^{\prime}\right\rangle_{\mathfrak{n}}=\sum_{v \mid \infty} \operatorname{Re}\left(\left(\xi_{v}\right)_{i j} \cdot{\overline{\left(\xi_{v}^{\prime}\right)}}_{i j}\right)
$$

with complex conjugation to accommodate complex $k_{v}$. Let $|\cdot|_{\mathfrak{n}}$ be the associated norm, and write $\left|\psi_{\xi}\right|_{\mathfrak{n}}=\left|\xi_{\infty}\right|_{\mathfrak{n}}$. There is the lower bound $\left|\xi_{\infty}\right|_{\mathfrak{n}} \geq h$ for $0 \neq \xi_{\infty} \in \Lambda^{*}$. In these terms, for a fixed compact subset $E^{\prime} \subset \mathscr{S}\left(\mathfrak{n}_{\infty} \times U\right)^{H}$, for each $\ell>0$, there is an uniform implied constant depending on $\ell$, not on $\varphi^{\prime} \in E^{\prime}$, such that

$$
\left|\varphi^{\prime}\left(\psi_{\xi}\right)\right|_{\mathfrak{n}} \ll \ell\left(1+\left|\psi_{\xi}\right|_{\mathfrak{n}}\right)^{-\ell}=\left(1+\left|\xi_{\infty}\right|_{\mathfrak{n}}\right)^{-\ell} \quad\left(\text { for all } \xi \in \mathfrak{n}_{k}, \text { for all } \varphi^{\prime} \in E^{\prime}\right)
$$

[7.3.12] Lemma: For fixed $\mathfrak{S}$ and $\Lambda^{*}$, there is a uniform implied constant such that, for every $\psi_{\xi} \in \Lambda^{*}$ not vanishing identically on any $\mathfrak{n}_{\mathbb{A}}^{P}$, for every $y=n_{y} m_{y} k_{y} \in \mathfrak{S}$, and for every $\alpha \in \Phi^{o}$,

$$
\alpha(y) \ll\left|\psi_{\xi}^{m_{y}}\right|_{\mathfrak{n}}
$$

Proof: As above, for given $\Lambda^{*}$, there is a lower bound $b>0$ such that for all $\xi_{\infty} \in \Lambda^{*},\left|\xi_{i j}\right| \geq b$ for all indices $i j$ with $\xi_{i j} \neq 0$. Given $\alpha=\alpha_{i} \in \Phi^{o}$, take $P=P^{i-1, r-i+1}$. The non-zero entries of elements of $\mathfrak{n}^{P}$ are at

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$i^{\prime} j^{\prime}$ with $i^{\prime} \leq i$ and $j^{\prime} \geq i+1$. Thus, the condition that $\psi_{\xi}$ restricted to $n_{\mathbb{A}}^{P}$ is not identically 1 requires that there are indices $i^{\prime} \leq i$ and $j^{\prime} \geq i+1$ such that the $i^{\prime} j^{\prime t h}$ component $\xi_{i^{\prime} j^{\prime}}$ is non-zero. With such $i^{\prime} j^{\prime}$,

$$
\left.\left|\psi_{\xi}^{m_{y}}\right|_{\mathfrak{n}}=\left|m_{y} \xi m_{y}^{-1}\right|_{\mathfrak{n}} \geq\left|\left(m_{y} \xi m_{y}^{-1}\right)_{i^{\prime} j^{\prime}}\right|=\mid\left(m_{y}\right)_{i^{\prime}} \cdot \xi_{i^{\prime} j^{\prime}} \cdot\left(m_{y}\right)_{j^{\prime}}^{-1}\right) \mid
$$

where the latter two norms are on the real vector space $k_{\infty}$. Since $m_{y}=m_{y}^{\prime} a_{t}$ with $m_{y}^{\prime}$ in the compact $C_{M}$, there is a uniform implied constant, independent of $i, i^{\prime}, j^{\prime}, \xi$, and $y$, such that

$$
\left.\left.\mid\left(m_{y}\right)_{i^{\prime}} \cdot \xi_{i^{\prime} j^{\prime}} \cdot\left(m_{y}\right)_{j^{\prime}}^{-1}\right)|\gg|\left(a_{t}\right)_{i^{\prime}} \cdot \xi_{i^{\prime} j^{\prime}} \cdot\left(a_{t}\right)_{j^{\prime}}^{-1}\right) \left.\left|=\frac{t_{i^{\prime}}}{t_{j^{\prime}}} \cdot\right| \xi_{i^{\prime} j^{\prime}} \right\rvert\, \geq \frac{t_{i^{\prime}}}{t_{j^{\prime}}} \cdot b
$$

Every character $a_{t} \rightarrow t_{i^{\prime}} / t_{j^{\prime}}$ with $i^{\prime}<j^{\prime}$ is a product of non-negative powers of the simple positive characters $a_{t} \rightarrow t_{\ell} / t_{\ell+1}$ for $1 \leq \ell<r$ :

$$
\frac{t_{i^{\prime}}}{t_{j^{\prime}}}=\prod_{i^{\prime} \leq \ell<j^{\prime}} \frac{t_{\ell}}{t_{\ell+1}}
$$

For $i^{\prime} \leq i<j^{\prime}$, the exponent of $\alpha_{i}$ in such an expression is 1 . Thus, for $y=n_{y} m_{y}^{\prime} a_{t} k_{y} \in \mathfrak{S}=C_{N} C_{M} A_{\tau} K_{\mathbb{A}}$,

$$
\frac{t_{i^{\prime}}}{t_{j^{\prime}}} \geq \frac{t_{i}}{t_{i+1}} \cdot \prod_{i^{\prime} \leq \ell<j^{\prime}, \ell \neq i} \frac{t_{\ell}}{t_{\ell+1}} \geq \frac{t_{i}}{t_{i+1}} \cdot \prod_{i^{\prime} \leq \ell<j^{\prime}, \ell \neq i} \tau \geq \frac{t_{i}}{t_{i+1}} \cdot \tau^{r} \gg_{\mathfrak{S}} \frac{t_{i}}{t_{i+1}}
$$

Thus,

$$
\left|\psi_{\xi}^{m_{y}}\right|_{\mathfrak{n}} \gg \frac{t_{i}}{t_{i+1}}=\alpha_{i}\left(m_{y}\right) \quad(\text { for } y \in \mathfrak{S})
$$

as was claimed.
[7.3.13] Corollary: For fixed $\mathfrak{S}$ and $\Lambda^{*}$, there is a uniform implied constant such that, for every $\psi_{\xi} \in \Lambda^{*}$ not vanishing identically on any $\mathfrak{n}_{\mathbb{A}}^{P}$, for every $y=n_{y} m_{y} k_{y} \in \mathfrak{S}$,

$$
\delta_{B}\left(m_{y}\right) \ll\left|\psi_{\xi}^{m_{y}}\right|^{r(r-1) / 2}
$$

Proof: $\delta_{B}$ is the product of all characters $a_{t} \rightarrow t_{i} / t_{j}$ with $i<j$. Apply the lemma.
Thus, for any $q, \ell>0$, for any $\alpha \in \Phi^{o}$, the part of the kernel $K_{\varphi}(x, y)$ that interacts with cuspforms has an estimate

$$
\begin{aligned}
& \delta_{B}\left(m_{y}\right) \sum_{\psi \in \mathfrak{n}_{k}}^{*} \widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right) \ll \ell \sum_{\psi \in \mathfrak{n}_{k}}^{*}\left|\psi_{\xi}^{m_{y}}\right|^{r(r-1) / 2} \cdot\left(1+\left|\psi_{\xi}\right|\right)^{-\ell} \\
\ll & \sum_{\psi \in \mathfrak{n}_{k}}^{*} \alpha\left(m_{y}\right)^{-q}\left|\psi_{\xi}^{m_{y}}\right|^{r(r-1) / 2+q-\ell}=\alpha\left(m_{y}\right)^{-q} \sum_{\psi \in \mathfrak{n}_{k}}^{*}\left|\psi_{\xi}^{m_{y}}\right|^{r(r-1) / 2+q-\ell}
\end{aligned}
$$

For given $q$, for $\ell$ sufficiently large, the sum $\sum_{\psi}^{*}|\psi|^{r(r-1) / 2+q-\ell}$ is convergent, so

$$
\delta_{B}\left(m_{y}\right) \sum_{\psi \in \mathfrak{n}_{k}}^{*} \widehat{\varphi}_{x, y}\left(\psi^{m_{y}}\right)<_{q} \alpha\left(m_{y}\right)^{-q} \quad\left(\text { for every } \alpha \in \Phi^{o}\right)
$$

From this estimate, for cuspform $f$,

$$
|(\varphi \cdot f)(y)|<_{q}\left(\inf _{\alpha \in \Phi^{o}} \alpha\left(m_{y}\right)\right)^{-q} \cdot \int_{Z^{+} G_{k} \backslash G_{k} \mathfrak{S}}|f(x)| d x
$$

and by Cauchy-Schwarz-Bunyakowsky

$$
\int_{Z^{+} G_{k} \backslash G_{k} \mathfrak{S}}|f(x)| d x \leq\left(\int_{Z^{+} G_{k} \backslash G_{k} \mathfrak{S}} 1 d x\right)^{1 / 2} \cdot\left(\int_{Z^{+} G_{k} \backslash G_{k} \mathfrak{S}}|f(x)|^{2} d x\right)^{1 / 2} \leq \operatorname{meas}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right) \cdot|f|_{L^{2}}
$$

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That is, at last, for $y$ in a fixed Siegel set,

$$
|(\varphi \cdot f)(y)|<_{q}\left(\inf _{\alpha \in \Phi^{o}} \alpha\left(m_{y}\right)\right)^{-q} \cdot|f|_{L^{2}}
$$

This is the decay property of $\varphi \cdot f$ asserted in theorem [7.3.10].
The remainder of the arguments for theorem [7.3.1] and corollaries [7.3.2]-[7.3.5] is essentially identical to that for theorem [7.2.1] and corollaries [7.2.2]-[7.2.5]. We review the points of the argument.

As earlier, for general reasons, we have
[7.3.14] Lemma: Let $\mathfrak{g}_{v}$ be the Lie algebra of $G_{v}$ for archimedean $v$. For $X \in \mathfrak{g}_{v}$, the left-derivative map

$$
C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \longrightarrow C_{c}^{\infty}\left(G_{\mathbb{A}}\right) \quad \text { by } \quad \varphi \longrightarrow\left(\left.g \rightarrow \frac{d}{d t}\right|_{t=0} \varphi\left(e^{-t X} g\right)\right)
$$

is continuous. (Proof as [7.1.12].)
[7.3.15] Corollary: For a compact set $E$ of test functions on $G_{\mathbb{A}}$, for a compact $C_{\mathfrak{g}}$ in $\mathfrak{g}=\mathfrak{g}_{v}$, and for $f$ ranging over cuspforms in the unit ball in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, there is a uniform implied constant such that

$$
\left.\left|\frac{d}{d t}\right|_{t=0}(\varphi \cdot f)\left(g e^{t X}\right) \right\rvert\, \ll 1 \quad \text { (for all } g \in G_{\mathbb{A}}, \text { for all } \varphi \in E \text {, for all } X \in C_{\mathfrak{g}} \text { ) }
$$

(Proof as [7.2.14].)
The smoothing property of $f \rightarrow \varphi \cdot f$ as in [14.5] assures that each $\varphi \cdot f$ is smooth. Smoothness of $\varphi$ at finite places is uniform, since that of $\varphi$ is: with $\varphi$ left-invariant by compact open subgroup $K^{\prime} \subset G_{\mathbb{A}_{\text {fin }}}$, for $h \in K^{\prime}$,

$$
(\varphi \cdot f)(g \cdot h)=\int_{G} \varphi(x) f(g h x) d x=\int_{G} \varphi\left(h^{-1} x\right) f(g x) d x=\int_{G} \varphi(x) f(g x) d x=(\varphi \cdot f)(g)
$$

A uniform bound on derivatives at archimedean places will imply uniform continuity:
[7.3.16] Lemma: Let $F$ be a smooth function on $G_{\mathbb{A}}$, (uniformly) right $K^{\prime}$-invariant for some compact open subgroup $K^{\prime} \subset G_{\mathbb{A}_{\text {fin }}}$, with a uniform pointwise bound on all $X \cdot F$ with $X$ in a compact neighborhood $C_{\mathfrak{g}}$ of 0 in $\mathfrak{g}_{\infty}$, namely,

$$
|(X \cdot F)(x)| \leq B \quad\left(\text { for all } x \in G_{\mathbb{A}}, \text { all } X \in C_{\mathfrak{g}}\right)
$$

Then $F$ is uniformly continuous: for every $\varepsilon>0$ there is a neighborhood $U$ of 1 in $G_{\mathbb{A}}$ such that $|F(x)-F(y)|<\varepsilon$ for all $x \in G_{\mathbb{A}}$ and $y \in x U$. (Proof as [7.2.15].)
[7.3.17] Corollary: For a compact set $E$ of test functions on $G_{\mathbb{A}}$, and for $f$ ranging over cuspforms in the unit ball in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, the family of images $\varphi \cdot f$ is (uniformly) equicontinuous on $G$.
///
Again, a compactness lemma reminiscent of Arzela-Ascoli:
[7.3.18] Lemma: Let $E$ be a equicontinuous, uniformly bounded, set of functions on $Z^{+} G_{k} \backslash G_{\mathbb{A}}$. Then $E$ has compact closure in $L^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$. (Proof as [7.1.16].)

Finally, we prove the theorem [7.3.1]. To summarize: the asymptotics of the kernels prove pointwise boundedness of the image of the unit ball $B$ of $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$, and consideration of derivatives proves equicontinuity of the image of $B$. The faux-Arzela-Ascoli compactness lemma proves compactness of the closure of $\{\varphi \cdot B: \varphi \in E\}$. Being integrated versions of right translations, these operators stabilize the subspace of cuspforms, as the latter is defined by a left integral condition. Thus, $\varphi$ maps the unit ball to a set with compact closure, so is a compact operator.

As in earlier examples, by direct computation, the adjoint of $f \rightarrow \varphi \cdot f$ is $f \rightarrow \varphi^{\vee} \cdot f$, where $\varphi^{\vee}(x)=\bar{\varphi}\left(x^{-1}\right)$. The space of test functions is stable under the operation $\varphi \rightarrow \varphi^{\vee}$.

Again, the general non-degeneracy result is [14.1.5], finishing the proof of [7.3.1].
The proof of corollary [7.3.2], that $L_{o}^{2}\left(Z^{+} G_{k} \backslash G_{\mathbb{A}}\right)$ decomposes as (the closure of) a direct sum of irreducible representations of $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, each occurring with finite multiplicity, is identical to the proof of [7.2.2].

The proofs of [7.3.3] and [7.3.4] are special cases, where the algebra $\mathcal{H}$ of compact operators is designed to be commutative, so the notion of irreducibles simplifies to simultaneous eigenspace. Gelfand's criterion [2.4.5] and the $p$-adic and archimedean Cartan decompositions from [3.2] show that this Hecke algebra $\mathcal{H}$ is commutative. For non-zero $f$ in a simultaneous eigenspace $V_{\chi}$ for $\mathcal{H}$, by non-degeneracy there is a test function $\varphi$ such that $\varphi \cdot f=\chi(\varphi) \cdot f$ and $\chi(\varphi) \neq 0$. With $\eta=\varphi / \chi(\phi)$, then $\eta \cdot f=f$. By smoothing, as in [14.5] for example, $f$ is smooth.

Just as in the proof of [7.1.3], the fact that the Casimir operators $\Omega_{v}$ on archimedean $G_{v}$ commute with left and right translation implies that $\Omega_{v}$ commutes with the action of $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$, by integrating. On right $K_{v^{-}}$ invariant functions, $\Omega_{v}$ is the invariant Laplacian $\Delta_{v}$. Thus, each $\Delta_{v}$ stabilizes the simultaneous eigenspaces $V_{\chi}$ of $\mathcal{H}$, all of which are finite-dimensional, consisting of smooth functions. The restriction of $\Delta_{v}$ to $V_{\chi}$ is still symmetric, so by finite-dimensional spectral theory $V_{\chi}$ has a basis of $\Delta_{v}$-eigenfunctions. ///
[7.3.19] Corollary: The space of $L^{2}$ cuspforms has an orthonormal basis of cuspforms $f$ such that there is a test function $\varphi \in C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ such that $\varphi \cdot f=f$. Such a cuspform $f$ is smooth and of rapid decay in the sense that, given a standard Siegel set $\mathfrak{S}$ and $q>0$,

$$
|f(g)|<_{q}\left(\inf _{\alpha \in \Phi^{o}} \alpha\left(m_{g}\right)\right)^{-q}
$$

Proof: The proof is the same as that of [7.2.19]: From above, the irreducibles modules over $\mathcal{H}=C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$ appearing in the space of $L^{2}$ cuspforms are finite-dimensional, each occurring with finite multiplicity. Let $f$ be in a copy $V$ of an irreducible module for $C_{c}^{\infty}\left(G_{\mathbb{A}}\right)$.

Since $V$ is irreducible, it has no proper, topologically closed $\mathcal{H}$-stable subspace. Since $V$ is finitedimensional, all vector subspaces are closed. Thus, $\mathcal{H} \cdot f=V$. In particular, there is a test function $\varphi$ such that $\varphi \cdot f=f$. Then [7.3.10] applies to $f=\varphi \cdot f$. Smoothness follows as in [14.5] and [14.6].

## 7.A Appendix: dualities

For an abelian topological group $G$ and $\mathbb{T}$ the unit circle in $\mathbb{C}$, the unitary dual of $G$ is

$$
\widehat{G}=\operatorname{Hom}^{o}\left(G, S^{1}\right)=\{\text { continuous group homomorphisms } G \rightarrow \mathbb{T}\}
$$

Pointwise multiplication makes $\widehat{G}$ an abelian group. A reasonable topology ${ }^{[58]}$ on $\widehat{G}$ is the compact-open topology, with a sub-basis of opens

$$
U=U_{C, E}=\{f \in \widehat{G}: f(C) \subset E\} \quad(\text { for compact } C \text { in } G, \text { open } E \text { in } \mathbb{T} \text { ) }
$$

From [7.A.4], the compact-open topology makes $\widehat{G}$ a abelian (locally-compact, Hausdorf) topological group, [7.A.1] Claim: The unitary dual of a compact abelian group is discrete. The unitary dual of a discrete abelian group is compact.
Proof: Let $G$ be compact. Let $E$ be a small-enough open in $\mathbb{T}$ so that $E$ contains no non-trivial subgroups of $G$. Noting that $G$ itself is compact, let $U \subset \widehat{G}$ be the open

$$
U=\{f \in \widehat{G}: f(G) \subset E\}
$$

Since $E$ is small, $f(G)=\{1\}$. That is, $f$ is the trivial homomorphism. This proves discreteness of $\widehat{G}$. For $G$ discrete, every group homomorphism to $\mathbb{T}$ is continuous. The space of all functions $G \rightarrow \mathbb{T}$ is the cartesian

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product of copies of $\mathbb{T}$ indexed by $G$. By Tychonoff's theorem, with the product topology, this product is compact. Indeed, for discrete $X$, the compact-open topology on the space $C^{o}(X, Y)$ of continuous functions from $X \rightarrow Y$ is the product topology on copies of $Y$ indexed by $X$. The subset of functions $f$ satisfying the group homomorphism condition

$$
f(g h)=f(g) \cdot f(h) \quad(\text { for } g, h \in G)
$$

is closed, since the group multiplication $f(g) \times f(h) \rightarrow f(g) \cdot f(h)$ in $\mathbb{T}$ is continuous. Since the product is also Hausdorff, $\widehat{G}$ is also compact.
[7.A.2] Claim: Local fields $k_{v}$ are self-dual, as are the adeles of a number field $k: \mathbb{A}^{\vee} \approx \mathbb{A}$.
Proof: For compact totally disconnected $G$, since $\mathbb{C}^{\times}$contains no small subgroups [2.4.3], every element of $G^{\vee}$ has image in roots of unity in $\mathbb{C}^{\times}$, which can be identified with $\mathbb{Q} / \mathbb{Z}$. Thus, for compact totally disconnected $G$,

$$
G^{\vee} \approx \operatorname{Hom}^{o}(G, \mathbb{Q} / \mathbb{Z}) \quad \text { (continuous homomorphisms) }
$$

where $\mathbb{Q} / \mathbb{Z}=\operatorname{colim} \frac{1}{N} \mathbb{Z} / \mathbb{Z}$ is discrete. As a topological group, $\mathbb{Z}_{p}=\lim \mathbb{Z} / p^{\ell} \mathbb{Z}$. It is also useful to observe that $\mathbb{Z}_{p}$ is a limit of the corresponding quotients of itself, namely,

$$
\mathbb{Z}_{p} \approx \lim \mathbb{Z}_{p} / p^{\ell} \mathbb{Z}_{p}
$$

Indeed, more generally, every abelian totally disconnected topological group $G$ has the property that

$$
G \approx \lim _{K} G / K
$$

where $K$ ranges over compact open subgroups of $G$. Also, as a topological group,

$$
\mathbb{Q}_{p}=\bigcup \frac{1}{p^{\ell}} \mathbb{Z}_{p}=\operatorname{colim} \frac{1}{p^{\ell}} \mathbb{Z}_{p}
$$

Because of the no small subgroups property [2.4.3] of the unit circle in $\mathbb{C}^{\times}$, every continuous element of $\mathbb{Z}_{p}^{\vee}$ factors through some limitand

$$
\mathbb{Z}_{p} / p^{\ell} \mathbb{Z}_{p} \approx \mathbb{Z} / p^{\ell} \mathbb{Z}
$$

Thus,

$$
\mathbb{Z}_{p}^{\vee}=\operatorname{colim}\left(\mathbb{Z}_{p} / p^{\ell} \mathbb{Z}_{p}\right)^{\vee}=\operatorname{colim} \frac{1}{p^{\ell}} \mathbb{Z}_{p} / \mathbb{Z}_{p}
$$

since $\frac{1}{p^{\ell}} \mathbb{Z}_{p} / \mathbb{Z}_{p}$ is the dual to $\mathbb{Z}_{p} / p^{\ell} \mathbb{Z}_{p}$ under the pairing

$$
\frac{1}{p^{\ell}} \mathbb{Z}_{p} / \mathbb{Z}_{p} \times \mathbb{Z}_{p} / p^{\ell} \mathbb{Z}_{p} \approx \frac{1}{p^{\ell}} \mathbb{Z} / \mathbb{Z} \times \mathbb{Z} / p^{\ell} \mathbb{Z} \ni\left(\frac{x}{p^{\ell}}+\mathbb{Z}\right) \times\left(y+p^{\ell} \mathbb{Z}\right) \longrightarrow x y+\mathbb{Z} \in \mathbb{Q} / \mathbb{Z}
$$

The transition maps in the colimit expression for $\mathbb{Z}_{p}^{\vee}$ are inclusions, so

$$
\mathbb{Z}_{p}^{\vee}=\operatorname{colim} \frac{1}{p^{\ell}} \mathbb{Z}_{p} / \mathbb{Z}_{p} \approx\left(\operatorname{colim} \frac{1}{p^{\ell}} \mathbb{Z}_{p}\right) / \mathbb{Z}_{p} \approx \mathbb{Q}_{p} / \mathbb{Z}_{p}
$$

Thus,

$$
\mathbb{Q}_{p}^{\vee}=\left(\operatorname{colim} \frac{1}{p^{\ell}} \mathbb{Z}_{p}\right)^{\vee}=\lim \left(\frac{1}{p^{\ell}} \mathbb{Z}_{p}^{\vee}\right.
$$

As a topological group, $\frac{1}{p^{\ell}} \mathbb{Z}_{p} \approx \mathbb{Z}_{p}$ by multiplying by $p^{\ell}$, so the dual of $\frac{1}{p^{\ell}} \mathbb{Z}_{p}$ is isomorphic to $\mathbb{Z}_{p}^{\vee} \approx \mathbb{Q}_{p} / \mathbb{Z}_{p}$. However, the inclusions for varying $\ell$ are not the identity map, so for compatibility take

$$
\left(\frac{1}{p^{\ell}} \mathbb{Z}_{p}\right)^{\vee}=\mathbb{Q}_{p} / p^{\ell} \mathbb{Z}_{p}
$$

Thus,

$$
\mathbb{Q}_{p}^{\vee}=\lim \mathbb{Q}_{p} / p^{\ell} \mathbb{Z}_{p} \approx \mathbb{Q}_{p}
$$

because, again, any abelian totally disconnected group is the projective limit of its quotients by compact open subgroups. The same argument applies to $\widehat{\mathbb{Z}}=\lim \mathbb{Z} / N \mathbb{Z}$ and finite adeles $\mathbb{A}_{\text {fin }}=\operatorname{colim} \frac{1}{N} \widehat{\mathbb{Z}}$, proving the self-duality of $\mathbb{A}_{\text {fin }}$. Fourier inversion asserts the self-duality of $\mathbb{R}$ and $\mathbb{C}$, giving the self-duality of $\mathbb{A}$. The same argument applies over an arbitrary finite extension $k_{v}$ of $\mathbb{Q}_{p}$, but now the pairing is composed with the local trace from $k_{v}$ to $\mathbb{Q}_{p}$ and the dual lattice to the local integers $\mathfrak{o}_{v}$ is (by definition) the inverse different.
[7.A.3] Claim: The unitary dual $(\mathbb{A} / k)^{\wedge}$ of the compact quotient $\mathbb{A} / k$ is isomorphic to $k$. In particular, given any non-trivial character $\psi$ on $\mathbb{A} / k$, all characters on $\mathbb{A} / k$ are of the form $x \rightarrow \psi(\alpha \cdot x)$ for some $\alpha \in k$. Proof: Because $\mathbb{A} / k$ is compact, $(\mathbb{A} / k)^{\wedge}$ is discrete. Since multiplication by elements of $k$ respects cosets $x+k$ in $\mathbb{A} / k$, the unitary dual has a $k$-vectorspace structure given by

$$
(\alpha \cdot \psi)(x)=\psi(\alpha \cdot x) \quad(\text { for } \alpha \in k, x \in \mathbb{A} / k)
$$

There is no topological issue in this $k$-vectorspace structure, because $(\mathbb{A} / k)^{\wedge}$ is discrete. The quotient map $\mathbb{A} \rightarrow \mathbb{A} / k$ gives a natural injection $(\mathbb{A} / k)^{\wedge} \rightarrow \widehat{\mathbb{A}}$.

Given non-trivial $\psi \in(\mathbb{A} / k)^{\wedge}$, the $k$-vectorspace $k \cdot \psi$ inside $(\mathbb{A} / k)^{\wedge}$ injects to a copy of $k \cdot \psi$ inside $\widehat{\mathbb{A}} \approx \mathbb{A}$. Assuming for a moment that the image in $\mathbb{A}$ is essentially the same as the diagonal copy of $k$, the quotient $(\mathbb{A} / k)^{\wedge} / k$ injects to the compact $\mathbb{A} / k$. The topology of $(\mathbb{A} / k)^{\wedge}$ is discrete, and the quotient $(\mathbb{A} / k)^{\wedge} / k$ is still discrete. Since all these maps are continuous group homomorphisms, the image of $(\mathbb{A} / k)^{\wedge} / k$ in $\mathbb{A} / k$ is a discrete subgroup of a compact group, so is finite. Since $(\mathbb{A} / k)^{\wedge}$ is a $k$-vectorspace, the quotient $(\mathbb{A} / k)^{\wedge} / k$ must be a singleton. This proves that $(\mathbb{A} / k)^{\wedge} \approx k$, granting that the image of $k \cdot \psi$ in $\mathbb{A} \approx \widehat{\mathbb{A}}$ is the usual diagonal copy.

To see how $k \cdot \psi$ is imbedded in $\mathbb{A} \approx \widehat{\mathbb{A}}$, fix non-trivial $\psi$ on $\mathbb{A} / k$, and let $\psi$ be the induced character on $\mathbb{A}$. The self-duality of $\mathbb{A}$ is that the action of $\mathbb{A}$ on $\widehat{\mathbb{A}}$ by $(x \cdot \psi)(y)=\psi(x y)$ gives an isomorphism. The subgroup $x \cdot \psi$ with $x \in k$ is certainly the usual diagonal copy.

For completeness, we prove
[7.A.4] Claim: The unitary dual $\widehat{G}$ of an abelian (locally compact, Hausdorff) topological group is an abelian (locally compact, Hausdorff) topological group.
[7.A.5] Remark: We do not prove the local compactness in general. The important special cases, that the dual of discrete is compact, and vice-versa, give the local compactness of the duals in those cases.
Proof: That the unitary dual is abelian is immediate, since the multiplication is pointwise by values, and the target group $\mathbb{T}$ is abelian. First, verify that the topology is invariant. That is, given a sub-basis open

$$
U(C, E)=\{f \in \widehat{G}: f(c) \in E, \text { for all } c \in C\} \quad \text { (with } C \text { compact in } G, E \text { open in } \mathbb{T} \text { ) }
$$

and given $f_{o} \in \widehat{G}$, show that $f_{o} \cdot U(C, E)$ is open. This is not completely trivial, as $f_{o} \cdot U(C, E)$ is not obviously of the form $U\left(C^{\prime}, E^{\prime}\right)$ :

$$
f_{o} \cdot U(C, E)=\left\{f \in \widehat{G}: f(c) \in f_{o}(c) \cdot E, \text { for all } c \in C\right\}
$$

To show that $f_{o} \cdot U(C, E)$ is open, we show that every point is contained in a finite intersection of the basic opens, with that intersection contained in $f_{o} \cdot U(C, E)$.

Fix $f \in f_{o} \cdot U(C, E)$. Since $f_{o}^{-1}(c) f(c) \in E$, each $c \in C$ has a neighborhood $N_{c}$ such that $f_{o}^{-1}\left(N_{c}\right) \cdot f\left(N_{c}\right) \subset$ $E$. Shrink each $N_{c}$ to have compact closure $\bar{N}_{c}$, and so that $f_{o}^{-1}\left(\bar{N}_{c}\right) \cdot f\left(\bar{N}_{c}\right) \subset E$. By compactness of $C$, it has a finite subcover $N_{i}=N_{c_{i}}$. Thus,

$$
f\left(\bar{N}_{i}\right) \subset f_{o}\left(c^{\prime}\right) \cdot E \quad\left(\text { for all } i, \text { for all } c^{\prime} \in \bar{N}_{i}\right)
$$

From the result of the following subsection, an intersection of a compact family of opens is open, so

$$
E_{i}=\bigcap_{c^{\prime} \in \bar{N}_{i}} f_{o}\left(c^{\prime}\right) \cdot E=\text { open }
$$

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This open $E_{i}$ is non-empty, since it contains $f\left(\bar{N}_{i}\right)$. Thus,

$$
f \in \bigcap_{i} U\left(\bar{N}_{i}, E_{i}\right) \quad \text { (a finite intersection) }
$$

On the other hand, with $c_{i}$ and $\bar{N}_{i}$ determined by $f$, take

$$
f^{\prime} \in \bigcap_{i} U\left(\bar{N}_{i}, E_{i}\right)
$$

Then

$$
f^{\prime}\left(\bar{N}_{i}\right) \subset f_{o}(c) \cdot E \quad\left(\text { for all } c \in \bar{N}_{i}\right)
$$

In particular,

$$
f^{\prime}(c) \in f_{o}(c) \cdot E \quad\left(\text { for all } c \in \bar{N}_{i}\right)
$$

Since the sets $\bar{N}_{i}$ cover $C$, we have $f^{\prime} \in f_{o} \cdot U(C, E)$. That is,

$$
\bigcap_{i} U\left(\bar{N}_{i}, E_{i}\right) \subset f_{o} \cdot U(C, E)
$$

This proves that the translate $f_{o} \cdot U(C, E)$ is open, in the compact-open topology. That is, the compact-open topology is translation-invariant.

Now we prove the fact needed above, that compact intersections of opens are open, in the following sense. Let $H$ be a topological group, Hausdorff, but not necessarily locally compact. We claim that

$$
\bigcap_{k \in K} k \cdot U=\text { open } \quad \text { (for } U \subset H \text { open, and } K \subset H \text { compact) }
$$

For $u \in k \cdot U$ for all $k$, by the continuity of inversion and the group operation, there are neighborhoods $U_{k}$ of $u$ and $V_{k}$ of $k$ such that

$$
V_{k}^{-1} \cdot U_{k} \subset U
$$

Let $V_{i}=V_{k_{i}}$ be a finite subcover of $K$, and put $U_{i}=U_{k_{i}}$. Thus, for $k \in V_{i}$,

$$
k^{-1} \cdot U_{i} \subset U \quad\left(\text { for } k \in V_{i}\right)
$$

Thus,

$$
k^{-1} \cdot \bigcap_{i} U_{i} \subset U \quad(\text { for all } k \in K)
$$

Since finite intersections of opens are open, the intersection of the $U_{i}$, each containing $u$, is an open neighborhood of $u$. That is, the intersection of the translates $k \cdot U$ is open. This proves the claim.

Next, show that the pointwise multiplication operation

$$
\left(f_{1} \cdot f_{2}\right)(x)=f_{1}(x) \cdot f_{2}(x) \quad\left(\text { for } f_{i} \in \widehat{G} \text { and } x \in G\right)
$$

in $\widehat{G}$ is continuous in the compact-open topology. Given a sub-basis neighborhood $U(C, E)$ of $f_{1} \cdot f_{2}$, the already-demonstrated invariance of the topology implies that $\left(f_{1} f_{2}\right)^{-1} U(C, E)$ is open, and is a neighborhood of the trivial character. Thus, without loss of generality, take $f_{1}=f$ and $f_{2}=f^{-1}$. Given a sub-basic neighborhood $U(C, E)$ of the trivial character in $\widehat{G}$, show that there are neighborhoods $U_{1}$ of $f$ and $U_{2}$ of $f^{-1}$ such that $U_{1} \cdot U_{2} \subset U(C, E)$. For $U(C, E)$ to be a neighborhood of the trivial character means exactly that $1 \in E$ (and $C \neq \phi)$.

Let $E^{\prime}$ be an open neighborhood of 1 such that $E^{\prime} \cdot E^{\prime} \subset E$. For

$$
f^{\prime} \in\left(f \cdot U\left(C, E^{\prime}\right)\right) \cdot\left(f^{-1} \cdot U\left(C, E^{\prime}\right)\right)
$$

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we have

$$
f^{\prime}(c) \in\left(f(c) \cdot E^{\prime}\right) \cdot\left(f^{-1}(c) \cdot E^{\prime}\right)=E^{\prime} \cdot E^{\prime} \subset E \quad(\text { for all } c \in C)
$$

That is,

$$
\left(f \cdot U\left(C, E^{\prime}\right)\right) \cdot\left(f^{-1} \cdot U\left(C, E^{\prime}\right)\right) \subset U(C, E)
$$

This proves continuity of multiplication. Continuity of inversion is similar.
Finally, Hausdorffness: take $f_{1} \neq f_{2}$ in $\widehat{G}$. For some $g \in G, f_{1}(g) \neq f_{2}(g)$. Since the target $\mathbb{T}$ is Hausdorff, there are opens $E_{1} \ni f_{1}(g)$ and $E_{2} \ni f_{2}(g)$ with $E_{1} \cap E_{2}=\phi$. Since the source $G$ is Hausdorff, the singleton $\{g\}$ is compact. Thus, $f_{i} \in U\left(\{g\}, E_{i}\right)$, and these opens are disjoint. This completes the discussion of the unitary dual.

## 7.B Appendix: compact quotients $\Gamma \backslash G$

Here we see the enormous simplification in the argument for discrete decomposition when $\Gamma \backslash G$ is compact, so that Gelfand's cuspform condition is vacuously met. Let $G$ be any unimodular topological group, and $\Gamma$ a discrete subgroup so that $\Gamma \backslash G$ is compact. As in [6.1] and [6.2], $G$ acts continuously on $L^{2}(\Gamma \backslash G)$ by right translation, in fact by unitary operators, and $\eta \in C_{c}^{o}(G)$ acts continuously by the integral operators

$$
(\eta \cdot f)(x)=\int_{G} \eta(g) f(x g) d g
$$

[7.3.1] Theorem: $C_{c}^{o}(G)$ acts on $L^{2}(\Gamma \backslash G)$ by Hilbert-Schmidt (hence, compact) operators. The collection of such operators is closed under adjoints, and is non-degenerate in the sense that for every $f \in L^{2}(\Gamma \backslash G)$ there is $\eta \in C_{c}^{\infty}(K \backslash G / K)$ such that $\eta \cdot f \neq 0$.
Proof: Just as in the more complicated arguments concerning cuspforms, first rearrange:

$$
\begin{aligned}
& (\eta \cdot f)(x)=\int_{G} \eta\left(x^{-1} g\right) f(g) d g=\int_{\Gamma \backslash G} \sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma g\right) f(\gamma g) d g \\
& =\int_{\Gamma \backslash G} \sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma g\right) f(g) d g=\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma g\right)\right) \cdot f(g) d g
\end{aligned}
$$

Thus, with Schwartz kernel $K_{\eta}(x, y)=\sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma y\right)$ (see [Schwartz 1950]),

$$
(\eta \cdot f)(x)=\int_{\Gamma \backslash G} K_{\eta}(x, y) f(y) d y
$$

Since $\Gamma$ is discrete in $G$, for $x, y$ in a fixed compact subset of $G$, the sum for $K_{\eta}(x, y)$ is finite, so $K_{\eta}$ is continuous on $\Gamma \backslash G \times \Gamma \backslash G$. Unlike the more general situation, since $\Gamma \backslash G$ is compact, $K_{\eta} \in L^{2}(\Gamma \backslash G \times \Gamma \backslash G)$. Thus, $f \rightarrow \eta \cdot f$ is Hilbert-Schmidt (see [9.A]), and therefore a compact operator.

The adjoint of $f \rightarrow \eta \cdot f$ is easily expressed by changing variables:

$$
\begin{aligned}
& \langle\eta f, F\rangle_{L^{2}(\Gamma \backslash G)}=\int_{\Gamma \backslash G} \int_{G} \eta(g) f(x g) \overline{F(x)} d g d x=\int_{G} \int_{\Gamma \backslash G} \eta(g) f(x g) \overline{F(x)} d x d g \\
= & \left.\int_{G} \int_{\Gamma \backslash G} f(x) \overline{\bar{\eta}(g) F\left(x g^{-1}\right)} d x d g=\langle f, \check{\eta} \cdot F\rangle_{L^{2}(\Gamma \backslash G)} \quad \quad \text { (where } \check{\eta}(g)=\overline{\eta\left(g^{-1}\right)}\right)
\end{aligned}
$$

The non-degeneracy is [14.1.5].
A representation of $G$ on a topological vector space $V$ is a continuous map $G \times V \rightarrow V$ that sends $g \in G$ to (continuous) linear maps on $V$. The representation $V$ is irreducible when there are no (topologically) closed $G$-stable subspaces except $\{0\}$ and $V$ itself. Homomorphisms $\varphi: V \rightarrow W$ of $G$-representations are continuous linear maps commuting with the action of $G: \varphi(g \cdot v)=g \cdot \varphi(v)$ for all $g \in G$ and $v \in V$. The
multiplicity of an irreducible $V$ in another representation $W$ of $G$ is $\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{G}(V, W)$, as elaborated in [9.D.14]. The same terminology applies to the integral-operator action of $C_{c}^{o}(G)$ on a $G$-representation.
[7.3.2] Corollary: $L^{2}(\Gamma \backslash G)$ is (the completion of) an orthogonal direct sum of irreducible $C_{c}^{o}(G)$ subrepresentations, each occurring with finite multiplicity.
Proof: Given the theorem, this is mostly just [7.2.18].
[7.3.3] Corollary: $L^{2}(\Gamma \backslash G)$ is (the completion of) an orthogonal direct sum of irreducible unitary representations of $G$, each occurring with finite multiplicity.
Proof: By [14.1.6] and [14.1.7], $C_{c}^{o}(G)$-subrepresentations of the $G$-representation $L^{2}(\Gamma \backslash G)$ are $G$ subrepresentations, and $C_{c}^{o}(G)$-irreducibility implies $G$-irreducibility.
[7.3.4] Remark: Discrete subgroups $\Gamma$ of $G=S L_{2}(\mathbb{R})$ with compact quotient $\Gamma \backslash G$ can be obtained in several ways. A purely analytical device is the uniformization theorem, asserting that every compact, connected Riemann surface of genus $\geq 2$ is such a quotient. More number-theoretic examples are obtained by taking quaternion division algebras $B$ over $\mathbb{Q}$ (that is, $\mathbb{Q}$-four-dimensional simple division algebras with $\mathbb{Q}$ in the center) split over $\mathbb{R}$, that is, so that $B \otimes_{\mathbb{Q}} \mathbb{R}$ is isomorphic to the 2 -by- 2 matrix algebra $M_{2}(\mathbb{R})$. For a maximal finite- $\mathbb{Z}$-module-rank subring $\mathfrak{o}$ of $B$, imbed $\mathfrak{o}^{\times}$into $G L_{2}(\mathbb{R})$, and let $\Gamma=\mathfrak{o}^{\times} \cap S L_{2}(\mathbb{R})$. Then $\Gamma \backslash G$ is compact, by an argument resembling that in [2.A] for the compactness of $\mathbb{J}^{1} / k^{\times}$. Similarly, for a number field $k$ and quaternion division algebra $B$ over $k$, with $B \otimes_{k} k_{v} \approx M_{2}(\mathbb{R})$ for exactly one archimedean place $v$ of $k$, and $B \otimes_{k} k_{v^{\prime}} \approx \mathbb{H}$ for all other archimedean places, for maximal subring $\mathfrak{o}$, let $\Gamma$ be the intersection of $S L_{2}(\mathbb{R})$ with the projection of $\mathfrak{o}^{\times}$to $G L_{2}\left(k_{v}\right) \approx G L_{2}(\mathbb{R})$. Then $\Gamma \backslash S L_{2}(\mathbb{R})$ is compact. The corresponding compact quotients $\Gamma \backslash \mathfrak{H}$ are Shimura curves.

## 8. Moderate growth functions, theory of the constant term

1. The four small examples
2. $G L_{2}(\mathbb{A})$
3. $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$
4. Moderate growth of convergent Eisenstein series
5. Integral operators on cuspidal-data Eisenstein series

Appendix A: continuity of bilinear maps
From [7.1.20], [7.2.20], and [7.3.19], there is an orthonormal basis for cuspforms $f$ such that there are test functions $\varphi$ with $\varphi \cdot f=f$, and these cuspforms are of rapid decay in Siegel sets. This is a special case of the idea that the asymptotic behavior of moderate growth automorphic forms $f$ in Siegel sets is dominated by their constant terms, under the hypothesis that there is a test function $\varphi$ such that $\varphi \cdot f=f$. Eisenstein series are of moderate growth, even after meromorphic continuation, and under reasonable hypotheses on the data used to form them, meet the condition $\varphi \cdot f=f$ for suitable test function. Thus, although Eisenstein series are not in $L^{2}$, they admit good asymptotic approximations by their constant terms.

The underlying mechanism for the results of this chapter is essentially the fundamental theorem of calculus. Thus, these results are essentially archimedean. Thus, the second section indicates how to reduce the $G L_{2}(\mathbb{A})$ example to the four simple examples, and the third section treats only the simplest purely archimedean version of $G L_{r}$.

### 8.1 The four small examples

First, consider the four small examples from chapter 1 , with $G, \Gamma, P, M, N, A^{+}, K$ as there. For this section, write Iwasawa decompositions as $x=n_{x} a_{x} k_{x}$ with $a_{x} \in A^{+}$. The height function $\eta$ is

$$
\eta(n a k)=\eta(a)=\delta_{P}(a)=t^{2 r} \quad\left(\text { for } n \in N, a=\left(\begin{array}{cc}
t & 0 \\
0 & 1 / t
\end{array}\right) \in A^{+}, \text {and } k \in K\right)
$$

where $r=1,2,3,4$ in the respective examples, and $\delta_{P}$ is the modular function. A left $N \cap \Gamma$-invariant function $\mathbb{C}$-valued $f$ on $G$ or on $G / K$ is of moderate growth of exponent $\lambda \in \mathbb{R}$ on a fixed standard Siegel set $\mathfrak{S}$ when

$$
|f(x)| \ll \mathfrak{S} \eta(x)^{\lambda} \quad(\text { for } x \in \mathfrak{S})
$$

A left $N \cap \Gamma$-invariant function $\mathbb{C}$-valued $f$ on $G$ or $G / K$ is of moderate growth of exponent $\lambda$ (on standard Siegel sets) when it is of moderate growth of exponent $\lambda$ on every standard Siegel set. Say $f$ is of moderate growth if it is of moderate growth of exponent $\lambda$ for some $\lambda$. Since constant terms

$$
c_{P} f(x)=\int_{(\Gamma \cap N) \backslash N} f(n x) d n
$$

of automorphic forms $f$ are merely $N \cap \Gamma$-invariant, the notion of moderate growth needs to be more broadly applicable than just to $\Gamma$-invariant functions. A left $N \cap \Gamma$-invariant function $\mathbb{C}$-valued $f$ on $G$ or $G / K$ is of rapid decay on standard Siegel sets when

$$
|f(x)| \ll_{\lambda, \mathfrak{S}} \eta(x)^{\lambda} \quad(\text { for } x \in \mathfrak{S}, \text { for every } \mathfrak{S}, \text { for every } \lambda \in \mathbb{R})
$$

We can see directly that, for $f$ of moderate growth of exponent $\lambda$, the constant term $c_{P} f$ is also of moderate growth of exponent $\lambda$, and for $f$ of rapid decay, the constant term $c_{P} f$ is also of rapid decay: on a fixed standard Siegel set $\mathfrak{S}$,

$$
\left|c_{P} f(x)\right| \leq \int_{(N \cap \Gamma) \backslash N}|f(n x)| d n \ll_{\mathfrak{S}} \int_{(N \cap \Gamma) \backslash N} \eta(n x)^{\lambda} d n=\int_{(N \cap \Gamma) \backslash N} \eta(x)^{\lambda} d n=\eta(x)^{\lambda}
$$

For $f \in C^{o}(G)$ and for $\varphi \in \mathcal{D}(G)$, as usual

$$
(\varphi \cdot f)(x)=\int_{G} \varphi(g) f(x g) d g \quad(\text { for } \varphi \in \mathcal{D}(G))
$$

This converges at least as a $\mathbb{C}$-valued integral, for each $x \in G$.
[8.1.1] Theorem: For $f$ on $(N \cap \Gamma) \backslash G / K$ of moderate growth, if there is $\varphi \in \mathcal{D}(G)$ with $\varphi \cdot f=f$, then $f-c_{P} f$ is of rapid decay on Siegel sets.
Proof: The proof of the theorem is in several stages. First, the action of test functions does preserve moderate growth of a given exponent:
[8.1.2] Claim: Let $f$ be a function on $(\Gamma \cap N) \backslash G$ of moderate growth of exponent $\lambda$ in standard Siegel sets. Then, for every test function $\varphi \in \mathcal{D}(G)$, the function $\varphi \cdot f$ is also of moderate growth of exponent $\lambda$ in standard Siegel sets.
Proof: Let compact $C \subset G$ contain the support of $\varphi$. Without loss of generality, replace $C$ by $C \cdot K$ to make $C$ right $K$-stable. Fix a Siegel set $\mathfrak{S}=C_{N} A_{\tau} K$ with compact $C_{N} \subset N$ and $A_{\tau}=\left\{a \in A^{+}: \eta(a) \geq \tau\right\}$. For $x \in \mathfrak{S}$,

$$
(\varphi \cdot f)(x)=\int_{G} f(x y) \varphi(y) d y=\int_{G} f(y) \varphi\left(x^{-1} y\right) d y=\int_{x C} f(y) \varphi\left(x^{-1} y\right) d y
$$

[8.1.3] Lemma: Given compact $C \subset G$, there is compact $C_{A} \subset A^{+}$such that $a_{y} \in a_{x} C_{A}$ for $y \in x C$.
Proof: (of lemma) For right $K$-stable $C, C \subset\left(N A^{+} \cap C\right) \cdot K$, since in Iwasawa coordinates $p k \in C$ with $p \in N A^{+}$and $k \in K$ implies that $p=(p k) \cdot k^{-1}$ is also in $C$. Since $N \cdot A^{+} \approx N \times A^{+}$is a topological product, there are compacts $C_{N}^{\prime} \subset N, C_{A} \subset A^{+}$so that $K \cdot C \subset C_{N}^{\prime} \cdot C_{A} \cdot K$. Then

$$
x C \subset N a_{x} K \cdot C \subset N a_{x} \cdot C_{N}^{\prime} C_{A} K \subset N a_{x} \cdot N C_{A} K \subset N \cdot\left(a_{x} N a_{x}^{-1}\right) \cdot\left(a_{x} C_{A}\right) \cdot K \subset N \cdot\left(a_{x} C_{A}\right) \cdot K
$$

That is, $a_{y} \in a_{x} C_{A}$ for $y \in x C$, as claimed.
Returning to the proof of the claim, for $x \in \mathfrak{S}$, for $y$ in the support $x C$ of the integral,

$$
\eta(y)=\eta\left(a_{y}\right) \in \eta\left(a_{x} C_{A}\right)=\eta\left(a_{x}\right) \cdot \eta\left(C_{A}\right)
$$

Let

$$
\mu=\inf _{a \in C_{A}} \eta(a) \quad \sigma=\sup _{a \in C_{A}} \eta(a)^{\lambda}
$$

Take compact $C_{N}^{\prime \prime} \subset N$ large enough to surject to $(N \cap \Gamma) \backslash N$. Then, up to adjustment by $N \cap \Gamma$, for $x \in \mathfrak{S}$, $y \in x C$ implies that $y$ is in $\mathfrak{S}^{\prime}=C_{N^{\prime \prime}} A_{\mu \cdot \tau} K$. Invoking the moderate growth of $f$ on $\mathfrak{S}^{\prime}$,

$$
\begin{gathered}
|\varphi \cdot f(x)| \leq \sup |\varphi| \int_{x C}|f(y)| d y<_{f, \mathfrak{S}} \sup |\varphi| \cdot \int_{x C} \eta(y)^{\lambda} d y \ll \sup |\varphi| \cdot \sigma \cdot \int_{x C} \eta(x)^{\lambda} d y \\
=\sup |\varphi| \cdot \sigma \cdot \eta(x)^{\lambda} \cdot \operatorname{meas}(C)<_{\varphi} \eta(x)^{\lambda} \quad(\text { for } x \in \mathfrak{S})
\end{gathered}
$$

This argument applies to every standard Siegel set $\mathfrak{S}$, giving the moderate growth of $f$.
[8.1.4] Claim: For $f$ of moderate growth of exponent $\lambda$, if $\varphi \cdot f=f$ for some $\varphi \in \mathcal{D}(G)$, then $f$ is smooth, and is of uniform moderate growth of exponent $\lambda$, in the sense that for any $L$ in the universal enveloping algebra $U \mathfrak{g}$ of the Lie algebra $\mathfrak{g}$ of $G$, the derivative $L f$ is of moderate growth with exponent $\lambda$ on standard Siegel sets.
Proof: The image $\varphi \cdot f$ is smooth, by [14.5]. The point is that the left- $G$-invariant differential operators attached to $X$ in the Lie algebra of $G$ by

$$
X f(x)=\left.\frac{\partial}{\partial s}\right|_{s=0} f\left(x \cdot e^{s X}\right)
$$

can be absorbed into the action of varying $\varphi \in \mathcal{D}(G)$ on $f$ :

$$
X(\varphi \cdot f)(x)=\left.\frac{\partial}{\partial s}\right|_{s=0} \int_{G} f\left(x \cdot e^{s X} g\right) \varphi(g) d g=\left.\frac{\partial}{\partial s}\right|_{s=0} \int_{G} f(g) \varphi\left(e^{-s X} x^{-1} g\right) d g
$$

by replacing $g$ by $e^{-s X} x^{-1} g$. This is

$$
\left.\int_{G} f(x) \frac{\partial}{\partial s}\right|_{s=0} \varphi\left(e^{-s X} x^{-1} g\right) d g=\int_{G} f(x) X^{\mathrm{left}} \varphi\left(x^{-1} g\right) d g
$$

where $X^{\text {left }}$ is the (right-G-invariant) differential operator on the left naturally attached to $X$ via the left translation action. The interchange of differentiation and integration is justified by Gelfand-Pettis [14.1], observing that the integral is compactly supported, continuous, and takes values in a quasi-complete locally convex topological vector space on which differentiation is a continuous linear map. Using $\varphi \cdot f=f$,

$$
X f(x)=X(\varphi \cdot f)(x)=\left(\left(X^{\mathrm{left}} \varphi\right) \cdot f\right)(x)
$$

which is of moderate growth of exponent $\lambda$, by the previous claim. By induction on the degree of the differential operator $L, L f$ is of moderate growth of exponent $\lambda$.

The key bootstrapping property is the following:
[8.1.5] Claim: For $f$ smooth and left $(N \cap \Gamma)$-invariant, of uniform moderate growth of exponent $\lambda$ on standard Siegel sets, on a standard Siegel set $\mathfrak{S}$

$$
\left|\left(f-c_{P} f\right)(x)\right|<_{\mathfrak{G}} \eta(x)^{\lambda-1}
$$

Proof: For notational simplicity, we first carry out the argument for the example with $G=S L_{2}(\mathbb{R})$. Normalizing the measure of $(\Gamma \cap N) \backslash N$ to be 1,

$$
\left(c_{P} f-f\right)(x)=\int_{(\Gamma \cap N) \backslash N} f(n x)-f(x) d n=\int_{0 \leq t \leq 1} f\left(e^{t X} \cdot x\right)-f(x) d t
$$

where $X=\left(\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right)$ in the Lie algebra of $N$. By the fundamental theorem of calculus,
$f\left(e^{t X} \cdot x\right)-f(x)=\left.\int_{0}^{t} \frac{\partial}{\partial u}\right|_{u=0} f\left(e^{(u+s) X} \cdot x\right) d s=\int_{0}^{t} f\left(e^{s X} \cdot x \cdot e^{u x^{-1} X x}\right) d s=\int_{0}^{t}\left(x^{-1} X x \cdot f\right)\left(e^{s X} \cdot x\right) d s$
Let $x=n_{x} a_{x} k_{x}$ with $n_{x} \in N, a_{x} \in A^{+}, k_{x} \in K$. Then

$$
x^{-1} X x=\left(k_{x}^{-1} a_{x}^{-1} n_{x}^{-1}\right) X\left(n_{x} a_{x} k_{x}\right)=\left(k_{x}^{-1} a_{x}^{-1}\right) X\left(a_{x} k_{x}\right)
$$

Further,

$$
a_{x}^{-1} X a_{x}=\eta\left(a_{x}\right)^{-1} \cdot X
$$

Then

$$
x^{-1} X x=\left(k_{x}^{-1} a_{x}^{-1}\right) X\left(k_{x}^{-1} a_{x}^{-1}\right)=\eta\left(a_{x}\right)^{-1} k_{x}^{-1} X k_{x}=\eta\left(a_{x}\right)^{-1} \cdot \sum_{i} c_{i}\left(k_{x}\right) X_{i}
$$

where the $c_{i}$ are continuous functions (depending upon $X$ ) on $K$ and $\left\{X_{i}\right\}$ is a basis for the Lie algebra of $G$. Since the $c_{i}$ are continuous on the compact set $K$, they have a uniform bound $c$ in absolute value. Altogether,

$$
\begin{gathered}
\left(c_{P} f-f\right)(x)=\int_{0 \leq t \leq 1} \int_{0 \leq s \leq t} \eta\left(a_{x}\right)^{-1} \cdot\left(-\sum_{i} c_{i}\left(k_{x}\right) X_{i}\right) f\left(e^{s X} \cdot x\right) d s d t \\
=\eta\left(a_{x}\right)^{-1} \cdot \sum_{i} c_{i}\left(k_{x}\right) \int_{0 \leq t \leq 1} \int_{0 \leq s \leq t}\left(-X_{i} f\right)\left(e^{s X} \cdot x\right) d s d t=\eta\left(a_{x}\right)^{-1} \cdot \sum_{i} c_{i}\left(\theta_{x}\right) \int_{0 \leq t \leq 1}\left(-X_{i} f\right)\left(e^{t X} \cdot x\right) d t \\
=\eta\left(a_{x}\right)^{-1} \cdot \sum_{i} c_{i}\left(\theta_{x}\right) \cdot c_{P}\left(-X_{i} f\right)(x)
\end{gathered}
$$

so

$$
\begin{aligned}
&\left|\left(c_{P} f-f\right)(x)\right| \leq \eta\left(a_{x}\right)^{-1} \sum_{i}\left|c_{i}\left(\theta_{x}\right)\right| \cdot c_{P}\left|\left(-X_{i} f\right)(x)\right| \leq \eta\left(a_{x}\right)^{-1} \cdot \sum_{i} c \cdot c_{P}\left|\left(-X_{i} f\right)(x)\right| \\
& \ll \mathfrak{S} \eta\left(a_{x}\right)^{-1} \cdot \sum_{i} c \cdot \eta(x)^{\lambda} \ll \eta\left(a_{x}\right)^{\lambda-1}
\end{aligned}
$$

as claimed, for $G=S L_{2}(\mathbb{R})$. For the other three small examples, since $N$ is $r=2,3,4$-dimensional, respectively, commensurately more differentiations are needed, but the pattern is the same:

$$
\left(c_{P} f-f\right)(x)=\int_{(\Gamma \cap N) \backslash N} f(n x)-f(x) d n=\int_{0 \leq t_{1} \leq 1} \ldots \int_{0 \leq t_{r} \leq 1} f\left(e^{t_{1} X_{1}+\ldots+t_{r} X_{r}} \cdot x\right)-f(x) d t_{1} \ldots d t_{r}
$$

where $X_{1}, \ldots, X_{r}$ is a suitable basis for the Lie algebra of $N$. By the fundamental theorem of calculus,

$$
\begin{gathered}
f\left(e^{t_{1} X_{1}+\ldots+t_{r} X_{r}} \cdot x\right)-f(x) \\
=\left.\left.\int_{0}^{t_{1}} \cdots \int_{0}^{t_{r}} \frac{\partial}{\partial u_{1}}\right|_{u_{1}=0} \cdots \frac{\partial}{\partial u_{r}}\right|_{u_{r}=0} f\left(e^{\left(u_{1}+s_{1}\right) X_{1}+\ldots+\left(u_{r}+s_{r}\right) X_{r}} \cdot x\right)-f(x) d s_{1} \ldots d s_{r} \\
=\int_{0}^{t_{1}} \ldots \int_{0}^{t_{r}} f\left(e^{s_{1} X_{1}+\ldots+s_{r} X_{r}} \cdot x \cdot e^{u_{1} x^{-1} X_{1} x+\ldots+u_{r} x^{-1} X_{r} x}\right) d s_{1} \ldots d s_{r} \\
=\int_{0}^{t_{1}} \cdots \int_{0}^{t_{r}}\left(x^{-1} X_{1} x \cdots x^{-1} X_{r} x \cdot f\right)\left(e^{s_{1} X_{1}+\ldots+s_{r} X_{r}} \cdot x\right) d s_{1} \ldots d s_{r}
\end{gathered}
$$

As in the simpler version above,

$$
\begin{aligned}
x^{-1} X_{j} x & =\left(k_{x}^{-1} a_{x}^{-1}\right) X_{j}\left(k_{x}^{-1} a_{x}^{-1}\right)=\eta\left(a_{x}\right)^{-1 / r} k_{x}^{-1} X_{j} k_{x} \\
& =\eta\left(a_{x}\right)^{-1 / r} \sum_{i} c_{i j}\left(k_{x}\right) X_{i} \quad\left(\text { with } \eta\left(\begin{array}{cc}
t & 0 \\
0 & 1 / t
\end{array}\right)=t^{2 r}\right)
\end{aligned}
$$

for some continuous functions $c_{i j}$ on $K$. Thus, again, these functions have a bound $c$, and

$$
\begin{aligned}
&\left|\int_{0}^{t_{1}} \cdots \int_{0}^{t_{r}}\left(x^{-1} X_{1} x \cdots x^{-1} X_{r} x \cdot f\right)\left(e^{s_{1} X_{1}+\ldots+s_{r} X_{r}} \cdot x\right) d s_{1} \ldots d s_{r}\right| \\
& \ll \eta(x)^{-1} \sum_{i_{1}, \ldots, i_{r}} \int_{0}^{t_{1}} \cdots \int_{0}^{t_{r}}\left(X_{i_{1}} \cdots X_{i_{r}} f\right)\left(e^{s_{1} X_{1}+\ldots+s_{r} X_{r}} \cdot x\right) d s_{1} \ldots d s_{r} \\
& \ll \eta(x)^{-1} \sum_{i_{1}, \ldots, i_{r}} c_{P}\left|X_{i_{1}} \cdots X_{i_{r}} f(x)\right| \ll \eta(x)^{-1} \cdot \eta(x)^{\lambda}
\end{aligned}
$$

giving the claim.
Now finish the proof of the theorem. Take $\varphi \cdot f=f$ on $(\Gamma \cap N) \backslash G$ of moderate growth of exponent $\lambda$ in Siegel sets. Then $\alpha f$ is of moderate growth of exponent $\lambda$ for all $\alpha \in U \mathfrak{g}$, and $f-c_{P} f$ is of moderate growth of exponent $\lambda-1$. Then $\alpha\left(f-c_{P} f\right)$ is of moderate growth of exponent $\lambda-1$ for all $\alpha \in U \mathfrak{g}$, and $\left(f-c_{P} f\right)-c_{P}\left(f-c_{P} f\right)$ is of moderate growth of exponent $(\lambda-1)-1$. But

$$
\left(f-c_{P} f\right)-c_{P}\left(f-c_{P} f\right)=f-c_{P} f-c_{P} f+c_{P} c_{P} f=f-c_{P} f-c_{P} f+c_{P} f=f-c_{P} f
$$

By induction, $f-c_{P} f$ is of moderate growth of exponent $\lambda-\ell$ for every $\ell \in \mathbb{Z}$.
[8.1.6] Remark: In fact, the above arguments apply to $f$ eventually having the property $f=\varphi \cdot f$, in the sense that this property holds in a region where $\eta\left(a_{x}\right)$ is sufficiently large.

## 8.2 $G L_{2}(\mathbb{A})$

The purely archimedean argument of the previous section applies to the archimedean local factors $G_{v}$ of groups $G L_{2}(\mathbb{A})$. Again, the fundamental device is the fundamental theorem of calculus. For simplicity, we consider only trivial central characters.

A function $f$ on $Z_{\mathbb{A}} P_{k} \backslash G_{\mathbb{A}} / K$ is of moderate growth of exponent $\lambda$ on standard Siegel sets $\mathfrak{S}$ when

$$
|f(n m k)| \ll{ }_{f, \mathfrak{S}}\left|m_{1} / m_{2}\right|^{\lambda} \quad\left(\text { where } n \in N_{\mathbb{A}}, m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) \in M_{\mathbb{A}}, k \in K_{\mathbb{A}}\right)
$$

Such a function is of rapid decay (on Siegel sets) when it is of moderate growth of exponent $\lambda$ for all $\lambda \in \mathbb{R}$, with implied constant allowed to depend on $\lambda$. We may suppress the phrase on Siegel sets, but this is implied throughout.
[8.2.1] Theorem: For $f$ on $P_{k} \backslash G_{\mathbb{A}} / K$ of moderate growth, $K_{\mathbb{A}}$-finite, if there is $\varphi \in \mathcal{D}\left(G_{\mathbb{A}}\right)$ with $\varphi \cdot f=f$, then $f-c_{P} f$ is of rapid decay.
Proof: The proof is completely parallel to that of the previous sections, so we merely outline it, highlighting differences and adaptations. First, the action of test functions preserves moderate growth of a given exponent:
[8.2.2] Claim: Let $f$ be a function on $Z_{\mathbb{A}} P_{k} \backslash G$ of moderate growth of exponent $\lambda$ in standard Siegel sets. Then, for every test function $\varphi \in \mathcal{D}(G)$, the function $\varphi \cdot f$ is also of moderate growth of exponent $\lambda$ in standard Siegel sets. (Same proof as [8.1.2].)
[8.2.3] Claim: For $f$ of moderate growth of exponent $\lambda$, if $\varphi \cdot f=f$ for some $\varphi \in \mathcal{D}(G)$, then $f$ is smooth, and is of uniform moderate growth of exponent $\lambda$, in the sense that for any $L$ in the universal enveloping algebra $U \mathfrak{g}$ of the Lie algebra $\mathfrak{g}$ of $G$, the derivative $L f$ is of moderate growth with exponent $\lambda$ on standard Siegel sets. (Same proof as [8.1.4].)

Again, we have the bootstrapping property
[8.2.4] Claim: For $f$ smooth and left $(N \cap \Gamma)$-invariant, of uniform moderate growth of exponent $\lambda$ on standard Siegel sets, on a standard Siegel set $\mathfrak{S}$

$$
\left.\left|\left(f-c_{P} f\right)(n m k)\right|<_{\mathfrak{S}}\left|m_{1} / m_{2}\right|^{\lambda-1} \quad \text { (with } m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right)
$$

(Same proof as [8.1.5], using the fundamental theorem of calculus.)
The proof of the theorem is finished up as follows. Take $\varphi \cdot f=f$ on $Z_{\mathbb{A}} G_{k} \backslash G_{\mathbb{A}}$ and of moderate growth of exponent $\lambda$ in Siegel sets. Then $\alpha f$ is of moderate growth of exponent $\lambda$ for all $\alpha \in U \mathfrak{g}$, and $f-c_{P} f$ is of moderate growth of exponent $\lambda-1$. Then $\alpha\left(f-c_{P} f\right)$ is of moderate growth of exponent $\lambda-1$ for all $\alpha \in U \mathfrak{g}$, and $\left(f-c_{P} f\right)-c_{P}\left(f-c_{P} f\right)$ is of moderate growth of exponent $(\lambda-1)-1$. But

$$
\left(f-c_{P} f\right)-c_{P}\left(f-c_{P} f\right)=f-c_{P} f-c_{P} f+c_{P} c_{P} f=f-c_{P} f-c_{P} f+c_{P} f=f-c_{P} f
$$

By induction, $f-c_{P} f$ is of moderate growth of exponent $\lambda-\ell$ for every $\ell \in \mathbb{Z}$.

$$
8.3 S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

For these larger examples, for simplicity continue the archimedean aspects are emphasized. The general case is a superposition of copies of the archimedean, as in the previous section. The proofs for $S L_{r}(\mathbb{Z})$ are in essence mild extensions of those of [6.1], repeated for intelligibility, with appropriate modifications.

Let $G=S L_{r}(\mathbb{R}), \Gamma=S L_{r}(\mathbb{Z})$, and

$$
A=\left\{\left(\begin{array}{lll}
* & & \\
& \ddots & \\
& & *
\end{array}\right) \in S L_{r}(\mathbb{R})\right\}
$$

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and let $A^{+}$be the connected component of the identity in $A$, namely, diagonal matrices with positive entries. Let $B$ be the standard minimal parabolic (Borel subgroup) of upper triangular matrices, so $A$ is its standard Levi component. Let $N^{B}$ be the unipotent radical of $B$, namely, the upper-triangular unipotent matrices. With $K=S O_{n}(\mathbb{R})$, an Iwasawa decomposition of $G$ is $G=N^{B} \cdot A^{+} \cdot K$. The function $g \rightarrow a_{g}$ defined by expressing $g=n a_{g} k$ with $n \in N^{B}, a_{g} \in A^{+}, k \in K$ is well-defined.

Let $\log : A^{+} \rightarrow \mathfrak{a}$ be the inverse of the Lie exponential map $\alpha \rightarrow e^{\alpha}$ from the Lie algebra $\mathfrak{a}$ of $A^{+}$to $A^{+}$ itself. For $\lambda$ in the space of characters $\mathfrak{a}^{*}$ of $\mathfrak{a}$, write

$$
a^{\lambda}=e^{\lambda(\log a)}
$$

The roots of $A^{+}$or $\mathfrak{a}$ on the Lie algebra $\mathfrak{g}$ of $G$ are the characters $\lambda$ such that the $\lambda$-eigenspace

$$
\mathfrak{g}_{\lambda}=\left\{x \in \mathfrak{g}: a x a^{-1}=a^{\lambda} \cdot x, \text { for all } a \in A^{+}\right\}
$$

is non-zero, and then the eigenspace is called the $\lambda$-rootspace. The non-trivial roots are

$$
\chi_{i j}\left(\begin{array}{cccc}
m_{1} & & & \\
& m_{2} & & \\
& & \ddots & \\
& & & m_{n}
\end{array}\right)=\frac{m_{i}}{m_{j}} \quad(\text { for } i \neq j)
$$

As in [3.3], [3.10], [3,12], for $G=S L_{r}$ the standard simple roots are $\chi_{i, i+1}$. A left $N^{B} \cap \Gamma$-invariant $\mathbb{C}$-valued function $f$ on $G$ is of moderate growth of exponent $\lambda$ on a fixed standard Siegel set

$$
\mathfrak{S}=\mathfrak{S}_{t}=\left\{x \in G: a_{x}^{\alpha} \geq t \text { for all simple roots } \alpha\right\}
$$

when $|f(g)| \ll \mathfrak{S} a_{g}^{\lambda}$ for $g \in \mathfrak{S}$. Such $f$ is of rapid decay on a standard Siegel set $\mathfrak{S}$ when $|f(g)| \ll \mathfrak{S}^{\prime}, \lambda a_{g}^{\lambda}$ for all characters $\lambda \in \mathfrak{a}^{*}$.

Recall that, for standard parabolic $P$ with unipotent radical $N^{P}$, for left $\left(\Gamma \cap N^{P}\right)$-invariant $f$ on $G$, the constant term of $f$ along $P$ is

$$
c_{P} f(x)=\int_{\left(N^{P} \cap \Gamma\right) \backslash N^{P}} f(n x) d n
$$

Left-invariance under a larger subgroup of $P$ than just $\Gamma \cap N^{P}$ is inherited by the constant term, since $N^{P}$ is normal in $P$. In particular, left $\left(\Gamma \cap N^{B}\right)$-invariance is inherited. For $f$ left $\Gamma \cap N^{B}$-invariant, of moderate growth of exponent $\lambda$ in standard Siegel sets, for maximal proper standard parabolic $P$, the constant term $c_{P} f$ is also of moderate growth of exponent $\lambda$ in standard Siegel sets: normalizing the measure of $\left(N^{P} \cap \Gamma\right) \backslash N^{P}$ to be 1,

$$
\left|c_{P} f(x)\right| \leq \int_{\left(N^{P} \cap \Gamma\right) \backslash N^{P}}|f(n x)| d n \ll f, \mathfrak{S} \int_{\left(N^{P} \cap \Gamma\right) \backslash N^{P}} a_{n x}^{\lambda} d n=\int_{\left(N^{P} \cap \Gamma\right) \backslash N^{P}} a_{x}^{\lambda} d n=a_{x}^{\lambda}
$$

Similarly, for $f$ of rapid decay in standard Siegel sets, the constant term $c_{P} f$ is also of rapid decay.
Standard maximal proper parabolics $P$ have the convenient feature that their unipotent radicals $N^{P}$ are abelian. Further, the standard maximal proper parabolics are in bijection with simple roots, as follows. For two roots $\lambda, \mu$, write $\lambda \geq \mu$ if $\lambda-\mu$ is a linear combination of simple roots with non-negative coefficients. In $G=S L_{r}$, we have $\chi_{i j} \geq \chi_{i^{\prime} j^{\prime}}$ if and only if $i \leq i^{\prime}$ and $j \geq j^{\prime}$. Then the maximal proper parabolic associated to a simple root $\alpha$ is specified by saying that the Lie algebra $\mathfrak{n}$ of its unipotent radical $N=N^{P}$ is the sum of root-spaces $\mathfrak{g}_{\beta}$ for $\beta \geq \alpha$. In terms of matrix entries, for $\alpha=\chi_{i, i+1}$, the parabolic has diagonal blocks of sizes $i \times i$ and $(r-i) \times(r-i)$, and the unipotent radical $N$ consists of an $i \times(r-i)$ block in the upper right.

For $f \in C^{o}(G)$, write

$$
(\varphi \cdot f)(x)=\int_{G} \varphi(g) f(x g) d g \quad(\text { for } \varphi \in \mathcal{D}(G))
$$

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This converges at least as a $\mathbb{C}$-valued integral, for each $x \in G$.
[8.3.1] Theorem: Let $f$ be a left $N^{B} \cap \Gamma$-invariant $\mathbb{C}$-valued function on $G$, of moderate growth on every standard Siegel set. Suppose that there is $\varphi \in \mathcal{D}(G)$ such that $\varphi \cdot f=f$. Then for each simple root $\alpha$ and associated standard maximal parabolic $P, f-c_{P} f$ is of rapid decay in the direction $\alpha$, in the sense that $f-c_{P} f \ll_{\mathfrak{S}} a^{-\ell \cdot \alpha}$ for all $\ell \in \mathbb{Z}$.
[8.3.2] Corollary: For a left $N^{B} \cap \Gamma$-invariant $\mathbb{C}$-valued function $f$ on $G$, of moderate growth on every standard Siegel set, with $\varphi \in \mathcal{D}(G)$ such that $\varphi \cdot f=f$, if $c_{P} f=0$ for every standard maximal compact $P$, then $f$ is of rapid decay in standard Siegel sets.

Proof: The proof of the theorem is in stages. First, show that the action of $\mathcal{D}(G)$ preserves moderate growth on standard Siegel sets:
[8.3.3] Claim: For any $\varphi \in \mathcal{D}(G)$, if $f$ on $\left(N^{B} \cap \Gamma\right) \backslash G$ is of moderate growth of exponent $\lambda$ on standard Siegel sets, then $\varphi \cdot f$ is of moderate growth of exponent $\lambda$ on standard Siegel sets.
Proof: (of claim) For a compact set $C$ containing the support of $\varphi$,

$$
\varphi \cdot f(x)=\int_{G} f(x g) \varphi(g) d g=\int_{G} f(g) \varphi\left(x^{-1} g\right) d g=\int_{x C} f(g) \varphi\left(x^{-1} g\right) d g
$$

The proof of the following is identical to the proof of [8.1.3]:
[8.3.4] Lemma: Let $C$ be a compact set in $G, x \in G$. Then there is a compact subset $C_{A}$ of $A^{+}$such that $y \in x C$ implies $a_{y} \in a_{x} \cdot C_{A}$.
Take $x$ in a standard Siegel set

$$
\mathfrak{S}=\mathfrak{S}_{t}=\left\{x \in G: a_{x}^{\alpha} \geq t \text { for all simple roots } \alpha\right\}
$$

For $y$ in the support $x C$ of the integral, for simple root $\alpha$,

$$
a_{y}^{\alpha}=\in\left\{\left(a_{x} \cdot a\right)^{\alpha}: a \in C_{A}\right\}=a_{x}^{\alpha} \cdot\left\{a^{\alpha}: a \in C_{A}\right\}
$$

Let

$$
\mu=\inf _{\alpha} \inf _{a \in C_{A}} a^{\alpha} \quad \sigma=\sup _{a \in C_{A}} a^{\lambda}
$$

Take compact $C_{N}^{\prime \prime} \subset N^{B}$ large enough to surject to $\left(N^{B} \cap \Gamma\right) \backslash N^{B}$. Then, up to adjustment by $N^{B} \cap \Gamma$, $x \in \mathfrak{S}$ and $y \in x C$ implies that $y$ is in the Siegel set $\mathfrak{S}^{\prime}=C_{N^{\prime \prime}} A_{\mu \cdot t K}$. Invoking the moderate growth of $f$ on $\mathfrak{S}^{\prime}$,

$$
\begin{gathered}
|\varphi \cdot f(x)| \leq \sup |\varphi| \int_{x C}|f(y)| d y \ll{ }_{f, \mathfrak{S}} \sup |\varphi| \cdot \int_{x C} \eta(y)^{\lambda} d y \ll \sup |\varphi| \cdot \sigma \cdot \int_{x C} \eta(x)^{\lambda} d y \\
=\sup |\varphi| \cdot \sigma \cdot \eta(x)^{\lambda} \cdot \operatorname{meas}(C)<_{\varphi} a_{x}^{\lambda} \quad(\text { for } x \in \mathfrak{S})
\end{gathered}
$$

This argument applies for every standard Siegel set $\mathfrak{S}$, giving the moderate growth of $f$.
[8.3.5] Claim: For $f$ of moderate growth of exponent $\lambda$, if $\varphi \cdot f=f$ for some $\varphi \in \mathcal{D}(G)$, then $f$ is smooth, and is of uniform moderate growth of exponent $\lambda$, in the sense that for any $L$ in the universal enveloping algebra $U \mathfrak{g}$ of the Lie algebra $\mathfrak{g}$ of $G$, the derivative $L f$ is of moderate growth with exponent $\lambda$ on standard Siegel sets.
Proof: The image $\varphi \cdot f$ is smooth, by [14.5]. The key mechanism is that the left- $G$-invariant differential operators $X$ acting on the right attached to the right regular representation of $G$, arising from $X$ in the Lie algebra of $G$ by

$$
X f(x)=\left.\frac{\partial}{\partial s}\right|_{s=0} f\left(x \cdot e^{s X}\right)
$$

interact intelligibly with the action of $\varphi \in \mathcal{D}(G)$ on $f$, as follows.

$$
X(\varphi \cdot f)(x)=\left.\frac{\partial}{\partial s}\right|_{s=0} \int_{G} f\left(x \cdot e^{s X} g\right) \varphi(g) d g=\left.\frac{\partial}{\partial s}\right|_{s=0} \int_{G} f(g) \varphi\left(e^{-s X} x^{-1} g\right) d g
$$

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by replacing $g$ by $e^{-s X} x^{-1} g$. This is

$$
\left.\int_{G} f(x) \frac{\partial}{\partial s}\right|_{s=0} \varphi\left(e^{-s X} x^{-1} g\right) d g=\int_{G} f(x) X^{\mathrm{left}} \varphi\left(x^{-1} g\right) d g
$$

where $X^{\text {left }}$ is the (right- $G$-invariant) differential operator attached to $X$ via the left regular representation. [59] Thus, since $\varphi \cdot f=f$,

$$
X f(x)=X(\varphi \cdot f)(x)=\left(\left(X^{\mathrm{left}} \varphi\right) \cdot f\right)(x)
$$

which is of moderate growth of exponent $\lambda$, by the previous. By induction on the degree of the differential operator $L, L f$ is of moderate growth of exponent $\lambda$.
[8.3.6] Claim: Let $P$ be a maximal (proper) parabolic attached to simple root $\alpha$. Let $f$ be smooth and left $(N \cap \Gamma)$-invariant. Suppose that for all $Y \in \mathfrak{g}$ the (right) Lie derivative $Y f$ is of moderate growth of exponent $\lambda$ in Siegel sets. Then

$$
\left|\left(f-c_{P} f\right)(x)\right| \ll a_{x}^{\lambda-\alpha}
$$

Proof: Normalizing the measure of $\left(N^{P} \cap \Gamma\right) \backslash N^{P}$ to be 1,

$$
\left(f-c_{P} f\right)(x)=\int_{\left(N^{P} \cap \Gamma\right) \backslash N^{P}} f(n x)-f(x) d n=\int_{[0,1]^{k}} f\left(e^{t_{1} X_{1}+\cdots+t_{k} X_{k}} \cdot x\right)-f(x) d t_{1} \ldots d t_{k}
$$

where $X_{1}, \ldots, X_{k}$ is a basis for the Lie algebra of $N$ so that

$$
\left\{t_{1} X_{1}+\cdots+t_{k} X_{k}: 0 \leq t_{i} \leq 1,1 \leq i \leq k\right\}
$$

maps bijectively to $\left(N^{P} \cap \Gamma\right) \backslash N^{P}$, using the abelian-ness to see that this is easily possible. By the fundamental theorem of calculus, for $X$ in the Lie algebra,

$$
f\left(e^{t X} \cdot x\right)-f(x)=\left.\int_{0}^{t} \frac{\partial}{\partial r}\right|_{r=0} f\left(e^{(r+s) X} \cdot x\right) d s=\int_{0}^{t}-X^{\mathrm{left}} f\left(e^{s X} \cdot x\right) d s
$$

where $X^{\text {left }}$ is the natural right- $G$-invariant operator attached to $X$. For $X$ in the $\beta$ rootspace $\mathfrak{g}_{\beta}$ in the Lie algebra $\mathfrak{n}^{P}$ of $N^{P}$, writing $\operatorname{Ad}(g)(X)$ for $g X g^{-1}$,

$$
\operatorname{Ad}\left(a_{x}^{-1}\right)(X)=a_{x}^{-\beta} \cdot X
$$

and

$$
\operatorname{Ad}\left(\theta_{x}^{-1} a_{x}^{-1}\right)(X)=a_{x}^{-\beta} \cdot \operatorname{Ad}\left(\theta_{x}^{-1}\right)(X)=a_{x}^{-\beta} \cdot \sum_{1 \leq i \leq k} c_{i}\left(\theta_{x}\right) Y_{i}
$$

where the $c_{i}$ are continuous functions (depending upon $X$ ) on $K$ and $\left\{Y_{i}\right\}$ is a basis for the Lie algebra of $G$. Since the $c_{i}$ are continuous on the compact $K$, they have a uniform bound $c$ (depending on $X$ ). Altogether,

$$
\int_{0 \leq t \leq 1} f\left(e^{Y+t X} \cdot x\right)-f\left(e^{Y} \cdot x\right) d t=\beta\left(a_{x}\right)^{-1} \cdot \sum_{1 \leq i \leq k} c_{i}\left(\theta_{x}\right) \int_{0 \leq t \leq 1} \int_{0 \leq s \leq t}\left(-Y_{i} f\right)\left(e^{Y+s X} \cdot x\right) d s d t
$$

We need
[8.3.7] Lemma: For $x$ in a fixed standard Siegel set $\mathfrak{S}=\mathfrak{S}_{t}$,

$$
a_{x}^{-\beta} \ll \mathfrak{S} a_{x}^{-\alpha} \quad\left(\text { for roots } \beta \text { with } \mathfrak{g}_{\beta} \subset \mathfrak{n}^{P}\right)
$$

[59] As usual, the interchange of differentiation and integration is justified by observing that the integral is compactly supported, continuous, and takes values in a quasi-complete locally convex topological vector space on which differentiation is a continuous linear map. See [14.1] and [14.2].

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Proof: Since $\mathfrak{g}_{\beta} \subset \mathfrak{n}^{P}$, in an expression $\beta=\sum_{j} c_{j} \alpha_{j}$ for $\beta$ in terms of the simple roots $\alpha_{1}, \ldots, \alpha_{n-1}$, all coefficients are non-negative, and the coefficient of $\alpha$ is 1 . Thus, letting $\alpha=\alpha_{j_{o}}$,

$$
a_{x}^{-\beta}=a_{x}^{-\alpha} \cdot \prod_{j \neq j_{o}} a_{x}^{-c_{j} \alpha_{j}} \leq a_{x}^{-\alpha} \cdot \prod_{j \neq j_{o}} t^{-c_{j}}<_{t} a_{x}^{-\alpha}
$$

because of the inequalities $\alpha_{j}\left(a_{x}\right) \geq t$ characterizing the Siegel set $\mathfrak{S}=\mathfrak{S}_{t}$.
///
Continuing the proof of the claim, using the exponent $\lambda$ of moderate growth of all of the functions $Y_{i} f$,

$$
\int_{0 \leq t \leq 1} f\left(e^{Y+t X} \cdot x\right)-f\left(e^{Y} \cdot x\right) d t=O\left(a_{x}^{\lambda-\alpha}\right)
$$

or,

$$
\int_{0 \leq t \leq 1} f\left(e^{t_{1} X_{1}+\cdots+t_{i} X_{i}} \cdot x\right)-f\left(e^{t_{1} X_{1}+\cdots+t_{i-1} X_{i-1}} \cdot x\right) d t_{i}=O\left(a_{x}^{\lambda-\alpha}\right)
$$

Integrating in $d t_{1}, \ldots, d t_{i-1}$ and in $d t_{i+1}, \ldots, d t_{k}$ over copies of $[0,1]$ gives the same estimate for the $k$-fold integral:

$$
\int_{[0,1]^{k}} f\left(e^{t_{1} X_{1}+\cdots+t_{i} X_{i}} \cdot x\right)-f\left(e^{t_{1} X_{1}+\cdots+t_{i-1} X_{i-1}} \cdot x\right) d t_{1} \ldots d t_{k}=O\left(a_{x}^{\lambda-\alpha}\right)
$$

This is the assertion.
Continuing in this context, returning to the proof of the theorem, the previous claim shows that since every $X f$ is of exponent $\lambda, f-c_{P} f$ is of exponent $\lambda-\alpha$. The uniform moderate growth assures that every $X\left(f-c_{f}\right)$ is of exponent $\lambda-\alpha$, as well. Applying the last claim again,

$$
\left(X f-X c_{P} f\right)-c_{P}\left(X f-X c_{P} f\right)=X f-X c_{P} f=X\left(f-c_{P} f\right)
$$

is of exponent $\lambda-2 \cdot \alpha$, beginning an induction which proves the theorem.
[8.3.8] Corollary: For $f=\varphi \cdot f$ of moderate growth, if $c_{P} f=0$ for all maximal proper parabolics $P$, then $f$ is of rapid decay in standard Siegel sets $\mathfrak{S}$.

### 8.4 Moderate growth of convergent Eisenstein series

Now that the importance of the moderate growth property is clearer, we use an even simpler version of the approximation methods of [3.11] to prove moderate growth of convergent Eisenstein series, at least for parameters sufficiently far into the region of convergence. Chapter 11 will use this partial result to prove that Eisenstein series meromorphically continue as functions of moderate growth, therefore are of moderate growth everywhere.

We carry out the argument in detail for the maximal proper parabolic cuspidal-data Eisenstein series $E_{s, f}^{P}$ on $S L_{r}(\mathbb{Z}) \backslash S L_{r}(\mathbb{R}) / S O_{n}(\mathbb{R})$. After that proof, we indicate how to obtain a corresponding result more generally.
[8.4.1] Theorem: With maximal proper parabolic $P$ in $S L_{r}$ and cuspidal data $f=f_{1} \otimes f_{2}$ on $M^{P}$ with $f_{j}$ cuspforms in a strong sense on the factors of $M^{P}$, the Eisenstein series $E_{s, f}^{P}$ is of moderate growth on Siegel sets, and uniformly so for $s$ in compacts.
Proof: From chapter 7, strong-sense cuspforms are bounded, so as in [3.11] it suffices to treat potentially degenerate Eisenstein series $E_{s}^{P}=\sum_{\gamma \in P_{k} \backslash G L_{r}(k)} \varphi_{\sigma}^{P} \circ \gamma$ formed from

$$
\varphi_{\sigma}^{P}(n m k)=\left|\operatorname{det} m_{1}\right|^{\sigma_{1}} \cdot\left|\operatorname{det} m_{2}\right|^{\sigma_{2}} \quad\left(\text { with } n \in N^{P}, k \in K, \text { and } m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) \in M^{P}\right)
$$

where $P=P^{r_{1}, r_{2}}$ with $r_{1}+r_{2}=r, \sigma=\left(\operatorname{sig}_{1}, \sigma_{2}\right)$. For convergence, it suffices to take $\sigma_{1} \gg \sigma_{2}$. Allowing nontrivial central character in the notation and computations helps avoid some needless concern over artifactual details.

We reduce to the case of parabolics of the form $P^{r-1,1}$. Given $P=P^{q, r} \subset G L_{p+q}$, let $\rho: G L_{p+q} \rightarrow G L_{N}$ where $N=\binom{q+r}{r}$, by acting on $\wedge^{r}\left(k^{q+r}\right)$. The parabolic $P^{q, r}$ maps to the stabilizer $P^{\prime}$ in $G L_{N}$ of the line generated by $v_{o}=e_{q+1} \wedge \ldots \wedge e_{q+r}$ in $\wedge^{r}\left(k^{q+r}\right)$. Certainly

$$
\left(e_{q+1} \wedge \ldots \wedge e_{q+r}\right) \cdot \wedge^{r}\left(\begin{array}{cc}
1_{q} & 0 \\
0 & m_{2}
\end{array}\right)=\left(e_{q+1} \wedge \ldots \wedge e_{q+r}\right) \cdot \operatorname{det} m_{2}
$$

for $m_{2} \in G L_{r}$, and $\operatorname{det}\left(\wedge^{r} g\right)=(\operatorname{det} g)^{r}$ for $g \in G L_{q+r}$. Thus, for diagonal $m_{1} \in G L_{q}$ and $m_{2} \in G L_{r}$,

$$
\begin{aligned}
\varphi_{\sigma_{1}, \sigma_{2}}^{P^{q, r}}\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) & =\left|\operatorname{det} m_{1}\right|^{\sigma_{1}} \cdot\left|\operatorname{det} m_{2}\right|^{\sigma_{2}}=\left|\operatorname{det}\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right|^{\sigma_{1}} \cdot\left|\operatorname{det} m_{2}\right|^{\sigma_{2}-\sigma_{1}} \\
& =\left|\operatorname{det}\left(\wedge^{r}\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right)\right|^{\sigma_{1} \cdot r /\binom{q+r}{r}} \cdot\left|\wedge^{r} m_{2}\right|^{\sigma_{2}-\sigma_{1}}
\end{aligned}
$$

which in turn can be written in the corresponding form $\varphi_{\tau_{1}, \tau_{2}}^{P^{\prime}}$ on $G L_{N}$ for suitable $\tau_{1}, \tau_{2}$, whose precise form is inessential. Since $\rho\left(P^{q, r}\right) \subset P^{\prime}$, there is the immediate domination for $g \in G L_{q+r}$ :

$$
\sum_{\gamma \in P_{k}^{q, r} \backslash G L_{q+r}(k)} \varphi_{\sigma_{1}, \sigma_{2}}^{P^{q, r}}(\gamma \cdot g) \leq \sum_{\gamma^{\prime} \in P_{k}^{\prime} \backslash G L_{N}(k)} \varphi_{\tau_{1}, \tau_{2}}^{P^{\prime}}\left(\gamma^{\prime} \cdot \rho(g)\right)
$$

Thus, it suffices to prove
[8.4.2] Claim: Let $|\cdot|$ be the usual norm on $\mathbb{R}^{r}$.

$$
\Phi_{\sigma}(g)=\sum_{0 \neq v \in \mathbb{Z}^{r}} \frac{1}{|v \cdot g|^{2 \sigma}} \quad \text { (with } 1 \ll \sigma \in \mathbb{R} \text { ) }
$$

is bounded on Siegel sets.
Proof: Indeed, up to powers of absolute value of the determinant, the indicated sum dominates the sum for a degenerate Eisenstein series attached to the $P^{r-1,1}$ parabolic in $G L_{r}$. We can take Siegel sets to be of the form

$$
\mathfrak{S}_{t, C_{N}}=\left\{n m k: n \in C_{N} \subset N^{B}, m=\left(\begin{array}{ccc}
m_{1} & & \\
& \ddots & \\
& & m_{r}
\end{array}\right) \in M^{B}, k \in K,\left|m_{j} / m_{j+1}\right| \geq t\right\}
$$

where $C_{N}$ is compact and $t>0$. With operator norm $\|\cdot\|$, letting $g=n m k$ with $n \in C_{N}$, we have

$$
|v \cdot g|=|v \cdot n m k|=|v \cdot n m|=\left|v \cdot m \cdot m^{-1} n m\right|
$$

From

$$
|v \cdot m|=\left|v \cdot m \cdot m^{-1} n m \cdot m^{-1} n^{-1} m\right| \leq\left|v \cdot m \cdot m^{-1} n m\right| \cdot\left\|m^{-1} n^{-1} m\right\|=\left|v \cdot m \cdot m^{-1} n m\right| \cdot\left\|n^{-1}\right\|
$$

for $\sigma>0$ we have

$$
\frac{1}{|v \cdot n m k|^{2 \sigma}} \leq \frac{1}{|v \cdot m|^{2 \sigma}} \cdot\left\|n^{-1}\right\|^{2 \sigma} \leq \frac{1}{|v \cdot m|^{2 \sigma}} \cdot \sup _{n \in C_{N}}\left\|n^{-1}\right\|^{2 \sigma}<_{C_{N}} \frac{1}{|v \cdot m|^{2 \sigma}}
$$

since the operator norm is continuous and $C_{N}$ is compact. Thus,

$$
\Phi_{\sigma}(n m k)<_{C_{N}} \sum_{0 \neq v \in \mathbb{Z}^{r}} \frac{1}{|v \cdot m|^{2 \sigma}}=\sum_{0 \neq v \in \mathbb{Z}^{r}} \frac{1}{\left(\left(v_{1} m_{1}\right)^{2}+\ldots+\left(v_{r} m_{r}\right)^{2}\right)^{\sigma}}
$$

Without loss of generality, take $m_{r}=1$. With $\alpha_{j}=\left|m_{j} / m_{j+1}\right| \geq t>0$, this is

$$
\sum_{0 \neq v \in \mathbb{Z}^{r}} \frac{1}{\left(\left(v_{1} \alpha_{1} \alpha_{2} \cdots \alpha_{r-1} m_{1}\right)^{2}+\ldots+\left(v_{r-2} \alpha_{r-2} \alpha_{r-1}\right)^{2}+\left(v_{r-1} \alpha_{r-1}\right)^{2}+v_{r}^{2}\right)^{\sigma}} \leq \frac{1}{\min (t, 1)^{2 \sigma}} \sum_{0 \neq v \in \mathbb{Z}^{r}} \frac{1}{|v|^{2 \sigma}}
$$

This is finite for $2 \sigma>r$, and proves that $\Phi_{\sigma}$ is bounded on Siegel sets.
Powers of determinants are certainly of uniformly moderate growth on Siegel sets, so the theorem is proven.

The same type of argument gives a more general result, as follows.
From chapter 7, strong-sense cuspforms are bounded, so as in [3.11] it suffices to treat potentially degenerate Eisenstein series $E_{s}^{P}=\sum_{\gamma \in P_{k} \backslash G L_{r}(k)} \varphi_{\sigma}^{P} \circ \gamma$ formed from

$$
\varphi_{\sigma}^{P}(n m k)=\left|\operatorname{det} m_{1}\right|^{\sigma_{1}} \cdot\left|\operatorname{det} m_{2}\right|^{\sigma_{2}} \cdot \ldots \quad\left(\text { with } n \in N^{P}, k \in K, \text { and } m=\left(\begin{array}{ccc}
m_{1} & & \\
& m_{2} & \\
& & \ddots .
\end{array}\right) \in M^{P}\right)
$$

where $P=P^{r_{1}, r_{2}, \ldots}$ with $r_{1}+r_{2}+\ldots=r, \sigma=\left(\operatorname{sig}_{1}, \sigma_{2}, \ldots\right)$, and $m_{j} \in G L_{r_{j}}$. For easy convergence, it suffices to take $\sigma_{1} \gg \sigma_{2} \gg \ldots$.

In Iwasawa coordinates for the minimal parabolic $B$, such $\varphi_{\sigma}^{P}$ can readily be expressed as a product of analogous functions attached to maximal proper parabolics. For example, given $\sigma_{1} \gg \sigma_{2} \gg \sigma_{3}$,

$$
\left|\operatorname{det} m_{1}\right|^{\sigma_{1}} \cdot\left|\operatorname{det} m_{2}\right|^{\sigma_{2}} \cdot\left|\operatorname{det} m_{3}\right|^{\sigma_{3}}=\left(\left|\operatorname{det} m_{1} \cdot \operatorname{det} m_{2}\right|^{a} \cdot\left|\operatorname{det} m_{3}\right|^{b}\right) \cdot\left(\left|\operatorname{det} m_{1}\right|^{a^{\prime}} \cdot\left|\operatorname{det} m_{2} \cdot \operatorname{det} m_{3}\right|^{b^{\prime}}\right)
$$

is the requirement

$$
a+b=\sigma_{1} \quad a+b^{\prime}=\sigma_{2} \quad a^{\prime}+b^{\prime}=\sigma_{3}
$$

This is readily satisfied by taking $b^{\prime}$ large negative, so that the $a^{\prime}$ determined from the last equation satisfies $a^{\prime} \gg b^{\prime}$, then $a$ determined from the second equation and $b$ determined from the first satisfy $a \gg b$. Thus, we write

$$
\varphi_{\sigma}^{P}=\prod_{j} \varphi_{\tau_{1}^{j}, \tau_{2}^{j}}^{P^{j, r-j}}
$$

where $j$ indexes maximal proper parabolics and $\tau_{1}^{j} \gg \tau_{2}^{j}$ for all $j$. There is an immediate domination

$$
\sum_{\gamma \in P_{k} \backslash G L_{r}(k)} \varphi_{\sigma}^{P} \circ \gamma \leq \prod_{j} \sum_{\gamma \in P_{k}^{j, r-j} \backslash G L_{r}(k)} \varphi_{\tau_{1}^{j}, \tau_{2}^{j}}^{P^{j, r-j}} \circ \gamma
$$

Certainly moderate growth is preserved by products. Thus, it suffices to prove moderate growth for degenerate Eisenstein series attached to maximal proper parabolics.

A similar, complementary device reduces to the case of groundfield $k=\mathbb{Q}$, thus essentially reducing to the case explicitly treated.

### 8.5 Integral operators on cuspidal-data Eisenstein series

Having seen the significance of the property $\varphi \cdot f=f$ for some test function $\varphi$, the analogous property [7.2.20] for cuspforms $f_{1}, f_{2}$ on $G L_{r_{1}} \times G L_{r_{2}} \subset G L_{r_{1}+r_{2}}$ can be invoked to prove a similar property for Eisenstein series $E_{s, f}^{P}$ with $f=f_{1} \otimes f_{2}$ and $P=P^{r_{1}, r_{2}}$.

As in the computation [3.11.9] of the general form of constant terms, we need to assume a multiplicity-one property of $f_{1}$ and $f_{2}$, namely, that each is the unique cuspform on respective $S L_{r_{j}}(\mathbb{Z}) \backslash S L_{r_{j}}(R) / S O\left(r_{j}, \mathbb{R}\right)$ with given Laplacian eigenvalue, up to constant multiples.

Let $M^{1}=S L_{r_{1}} \times S L_{r_{2}} \subset M^{P}$, and

$$
\varphi_{s, f}\left(n m^{\prime} z_{y} k\right)=y^{s} \cdot f\left(m^{\prime}\right) \quad\left(\text { with } n \in N^{P}, m^{\prime} \in M^{1}, k \in K\right)
$$

and with

$$
z_{y}=\left(\begin{array}{cc}
y^{\frac{1}{r_{1} r_{2}}} \cdot 1_{r_{1}} & 0 \\
0 & 1_{r_{2}}
\end{array}\right)
$$

Let $E_{s, f}=\sum_{\gamma \in(\Gamma \cap P) \backslash \Gamma} \varphi_{s, f} \circ \gamma$, for $\operatorname{Re}(s) \gg 1$ for convergence.
[8.5.1] Claim: For every $\eta \in C_{c}^{\infty}(K \backslash G / K)$, there is an entire $\mathbb{C}$-valued function $s \rightarrow \mu_{s, f}(\eta)$ such that $\eta \cdot \varphi_{s, f}=\mu(\eta) \cdot \varphi_{s, f}$. At least in $\operatorname{Re}(s)>1$, similarly, $\eta \cdot E_{s, f}=\mu_{s, f} \cdot E_{s, f}$. Given $s, f$, there is $\eta$ such that $\mu_{s, f}(\eta)$ is not identically 0 .
Proof: In the region of convergence, $E_{s, f}$ is a sum of left translates of $\varphi_{s, f}$, so it suffices to prove the eigenfunction property of $\varphi_{s, f}$. The eigenfunction property for meromorphically-continued $E_{s, f}$ will follow by the identity principle from complex analysis.

Certainly the right action of such $\eta$ preserves left $N$-invariance. Since $\eta$ is $K$-bi-invariant, it preserves right $K$-invariant functions, as well. Computing directly, using an Iwasawa decomposition $G=P \cdot K=N^{P} \cdot M^{P} \cdot K$, noting that $P \cap K$ is compact, for $m_{o} \in M^{P}$,

$$
\left(\eta \cdot \varphi_{s, f}\right)\left(m_{o}\right)=\int_{G} \eta(h) \varphi_{s, f}\left(m_{o} h\right) d h=\int_{N^{P}} \int_{M^{P}} \int K_{\eta}(n m k) \varphi_{s, f}\left(m_{o} n m k\right) \delta(m)^{-1} d n d m d k
$$

with modular function $\delta$ on $P$. Continuing, this is

$$
\begin{gathered}
\int_{N^{P}} \int_{M^{P}} \int_{K} \eta(n m k) \varphi_{s, f}\left(m_{o} n m_{o}^{-1} \cdot m_{o} m\right) \delta(m)^{-1} d n d m d k \\
\quad=\int_{M^{P}}\left(\delta(m)^{-1} \int_{N^{P}} \eta(n m) d n\right) \cdot \varphi_{s, f}\left(m_{o} \cdot m\right) d m
\end{gathered}
$$

As a function of $m \in M^{P}$, the inner integral $\eta^{\prime}(m)$ is left and right $K \cap M^{P}$-invariant, smooth, and compactly supported. Action of $\eta^{\prime} \in C_{c}^{\infty}\left(\left(K \cap M^{P}\right) \backslash M^{P} /\left(K \cap M^{P}\right)\right)$ on functions on $M^{P}$ commutes with Casimir operators on the factors of $M^{P}$, and preserves central equivariance $u\left(m z_{t}\right)=|t|^{s} \cdot u(m)$. The right action on functions $m^{\prime} a_{y} \rightarrow y^{s} \cdot f\left(m^{\prime}\right)$ with $f=f_{1} \otimes f_{1}$ preserves cuspidality of $f$. Since $f=f_{1} \otimes f_{1}$ is assumed to be the unique cuspform with its eigenvalue, the action of $\eta^{\prime}$ sends $y^{s} \cdot f\left(m^{\prime}\right)$ to a multiple of itself, by scalar $\mu_{s, f}(\eta)$. Since $s \rightarrow \varphi_{s, f}$ is a holomorphic $C^{o}(G)$-valued function (for example), and since the integral giving the action of $\eta$ or $\eta^{\prime}$ exists as a Gelfand-Pettis integral, the function $s \rightarrow \mu_{s, f}(\eta)$ is holomorphic in $s$.

For a sequence $\left\{\eta_{j}\right\}$ of functions forming an approximate identity, $\mu_{s, f}(\eta) \cdot f=\eta_{j} \cdot f \rightarrow f$. Thus, for $f$ not the zero vector, $\mu_{s, f}\left(\eta_{j}\right) \neq 0$ for sufficiently large $j$.

## 8.A Appendix: joint continuity of bilinear maps

This is a corollary of Banach-Steinhaus [13.12.3], useful in removing ambiguities in considering averaged actions, for example, the action $C_{c}^{\infty}(\mathbb{R}) \times V \rightarrow V$ of test functions.
[8.A.1] Corollary: Let $\beta: X \times Y \rightarrow Z$ be a bilinear map on Fréchet spaces $X, Y, Z$, continuous in each variable separately. Then $\beta$ is jointly continuous.
Proof: Fix a neighborhood $N$ of 0 in $Z$ Let $x_{n} \rightarrow x_{o}$ in $X$ and $y_{n} \rightarrow y_{o}$ in $Y$. For each $x \in X$, by continuity in $Y, \beta\left(x, y_{n}\right) \rightarrow \beta\left(x, y_{o}\right)$. Thus, for each $x \in X$, the set of values $\beta\left(x, y_{n}\right)$ is bounded in $Z$. The linear functionals $x \rightarrow \beta\left(x, y_{n}\right)$ are equicontinuous, by Banach-Steinhaus, so there is a neighborhood $U$ of 0 in $X$ so that $b_{n}(U) \subset N$ for all $n$. In the identity

$$
\beta\left(x_{n}, y_{n}\right)-\beta\left(x_{o}, y_{o}\right)=\beta\left(x_{n}-x_{o}, y_{n}\right)+\beta\left(x_{o}, y_{n}-y_{o}\right)
$$

we have $x_{n}-x_{o} \in U$ for large $n$, and $\beta\left(x_{n}-x_{o}, y_{o}\right) \in N$. Also, by continuity in $Y, \beta\left(x_{o}, y_{n}-y_{o}\right) \in N$ for large $n$. Thus, $\beta\left(x_{n}, y_{n}\right)-\beta\left(x_{o}, y_{o}\right) \in N+N$, proving sequential continuity. Since $X \times Y$ is metrizable, sequential continuity implies continuity.
[8.A.2] Corollary: The same conclusion holds when $X$ is an LF-space.
Proof: Continuous linear functionals from an LF-space are exactly given by compatible families of continuous maps from the limitands.

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## 9. Unbounded operators on Hilbert spaces

1. Unbounded symmetric operators on Hilbert spaces
2. Friedrichs self-adjoint extensions of semi-bounded operators
3. Example: incommensurable self-adjoint extensions
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Appendix A: compact operators
Appendix B: closed graph theorem
Appendix C: irreducibles of compact groups
Appendix D: spectral theorem, Schur's lemma, multiplicities
Appendix E: Tietze-Urysohn-Brouwer extension theorem
This is preparation for eigenfunction decompositions of Hilbert spaces by operators closely related to invariant Laplacians.

Amazingly, resolvents $R_{\lambda}=(T-\lambda)^{-1}$ can exist, as everywhere-defined, continuous linear maps on a Hilbert space, even for $T$ unbounded and only densely-defined. Further hypotheses on $T$ are needed, but these hypotheses are met in useful situations occurring in practice. In particular, we need that $T$ is symmetric, in the sense that $\langle T v, w\rangle=\langle v, T w\rangle$ for $v, w$ in the domain $D_{T}$ of $T$, and semi-bounded in the sense that there is a constant $C$ such that either $\langle T v, v\rangle \geq C \cdot\langle v, v\rangle$ for all $v$ in $D_{T}$ or $\langle T v, v\rangle \leq C\langle v, v\rangle$ for all $v$ in $D_{T}$. In that circumstance, $T$ has a self-adjoint Friedrichs extension, with several good features, described explicitly below.

In practice, anticipating that a given unbounded operator is self-adjoint when extended suitably, a simple version of the operator is defined on an easily described, small, dense domain, specifying a symmetric operator. Then a self-adjoint extension is shown to exist, as in Friedrichs' theorem below.

For example, [9.5] gives a simple application, recovering the standard fact that the Hilbert space $L^{2}(\mathbb{T})$ on the circle $\mathbb{T}=\mathbb{R} / 2 \pi \mathbb{Z}$ has an orthogonal Hilbert-space basis of exponentials $e^{i n x}$ with $n \in \mathbb{Z}$, using ideas still applicable to situations lacking analogues of Dirichlet or Fejér kernels. These exponentials are eigenfunctions for the Laplacian $\Delta=d^{2} / d x^{2}$, so it would suffice to show that $L^{2}(\mathbb{T})$ has an orthogonal basis of eigenfunctions for $\Delta$. Two technical issues must be overcome: the most awkward is that $\Delta$ does not $\operatorname{map} L^{2}(\mathbb{T})$ to itself. Second, there is no guarantee that infinite-dimensional Hilbert spaces have Hilbertspace bases of eigenfunctions for a given linear operator. Indeed, reasonable operators on infinite-dimensional spaces may fail to have any eigenvectors. For example, on $L^{2}[a, b]$, the multiplication operator $T f(x)=x \cdot f(x)$ is continuous, possesses the symmetry property $\langle T f, g\rangle=\langle f, T g\rangle$, but has no eigenvectors. That is, the spectrum of operators on infinite-dimensional Hilbert spaces typically includes more than eigenvalues.

Natural operators like $d^{2} / d x^{2}$ on $L^{2}[a, b]$ are not bounded, that is, not continuous operators. Not-necessarily-bounded operators are called unbounded, despite the inconsistency of language.

Self-adjoint operators on Hilbert spaces generally do not give orthogonal Hilbert-space bases of eigenvectors, but an important special class does have a spectral theory imitating finite-dimensional spectral theory: the compact self-adjoint operators, which always give an orthogonal Hilbert-space basis of eigenvectors, as in [9.A].

Genuinely unbounded operators such as $\Delta$ are never continuous, much less compact, but in happy circumstances they may have compact resolvent $(1-\Delta)^{-1}$. Often, this compactness can be proven by a Sobolev imbedding lemma, a Rellich compactness lemma, and Friedrichs' construction, recovering a good spectral theory, as in examples below.

### 9.1 Unbounded symmetric operators on Hilbert spaces

The natural differential operator $\Delta=\frac{d^{2}}{d x^{2}}$ on $\mathbb{R}$ has no sensible definition as mapping all of the Hilbert space $L^{2}(\mathbb{R})$ to itself, whatever else can be said. The possibility of thinking of $\Delta$ as differentiating $L^{2}$ functions distributionally is useful, but sacrifices information if it abandons the $L^{2}$ behavior too completely. There is substantial motivation to accommodate discontinuous (unbounded) linear maps on Hilbert spaces, under some reasonable technical hypotheses, as apply to operators like $\Delta$. At the most cautious, $\Delta$ certainly $\operatorname{maps} C_{c}^{\infty}(\mathbb{R})$ to itself, and, by integration by parts,

$$
\langle\Delta f, g\rangle=\langle f, \Delta g\rangle \quad\left(\text { for both } f, g \in C_{c}^{\infty}(\mathbb{R})\right)
$$

This is symmetry of $\Delta$.
A not-necessarily continuous, that is, not-necessarily bounded, linear operator $T$, defined on a dense subspace $D_{T}$ of a Hilbert space $V$, is called an unbounded operator on $V$, even though it is likely not defined or definable on all of $V$. We consider mostly symmetric unbounded operators $T$, meaning that $\langle T v, w\rangle=\langle v, T w\rangle$ for $v, w$ in the domain $D_{T}$ of $T$. For example, the Laplacian is symmetric on $C_{c}^{\infty}(\mathbb{R})$, by integration by parts.

For unbounded operators on $V$, description of the domain is critical: an unbounded operator $T$ on $V$ is a subspace $D$ of $V$ and a linear map $T: D \longrightarrow V$. Nevertheless, explicit naming of the domain of an unbounded operator is often suppressed, instead writing $T_{1} \subset T_{2}$ when $T_{2}$ is an extension of $T_{1}$, in the sense that the domain of $T_{2}$ contains that of $T_{1}$, and the restriction of $T_{2}$ to the domain of $T_{1}$ agrees with $T_{1}$. Unlike self-adjoint operators on finite-dimensional spaces, and unlike self-adjoint bounded operators on Hilbert spaces, symmetric unbounded operators, even when densely defined, usually need to be extended in order to behave more like self-adjoint operators in finite-dimensional and bounded-operator situations.

An operator $T^{\prime}, D^{\prime}$ is a sub-adjoint to a symmetric operator $T, D$ when

$$
\langle T v, w\rangle=\left\langle v, T^{\prime} w\right\rangle \quad\left(\text { for } v \in D, w \in D^{\prime}\right)
$$

For dense domain $D$, for given $D^{\prime}$ there is at most one $T^{\prime}$ meeting the sub-adjointness condition.
In various useful circumstances there is a unique maximal element, in terms of domain, among all subadjoints to $T$, the adjoint $T^{*}$ of $T$. Uniqueness of a maximal sub-adjoint is proven below for $T$ symmetric. Perhaps surprisingly, we see below that the adjoint $T^{*}$ of a symmetric operator $T$ is not symmetric unless already $T$ is self-adjoint, that is, unless $T=T^{*}$. In particular, existence of adjoints for symmetric, denselydefined operators $T$ does not immediately imply existence of $\left(T^{*}\right)^{*}$. Paraphrasing the notion of symmetry: a densely-defined operator $T$ is symmetric when $T \subset T^{*}$, and self-adjoint when $T=T^{*}$. These comparisons refer to the domains of these not-everywhere-defined operators. In the following claim and its proof, the domain of a map $S$ on $V$ is incorporated in a reference to its graph

$$
\text { graph } S=\{v \oplus S v: v \in \text { domain } S\} \subset V \oplus V
$$

The direct sum $V \oplus V$ is a Hilbert space with natural inner product $\left\langle v \oplus v^{\prime}, w \oplus w^{\prime}\right\rangle=\left\langle v, v^{\prime}\right\rangle+\left\langle w, w^{\prime}\right\rangle$. Define an isometry $U: V \oplus V \rightarrow V \oplus V$ by $v \oplus w \rightarrow-w \oplus v$.
[9.1.1] Claim: Given $T$ with dense domain $D$, there is a unique maximal $T^{*}, D^{*}$ among all sub-adjoints to $T, D$. The adjoint $T^{*}$ is closed, in the sense that its graph is closed in $V \oplus V$. In fact, the adjoint is characterized by its graph, the orthogonal complement in $V \oplus V$ to the image of the graph of $T$ under $U$, namely,

$$
\text { graph } T^{*}=\text { orthogonal complement of } U(\operatorname{graph} T)
$$

Proof: The adjointness condition $\langle T v, w\rangle=\left\langle v, T^{*} w\right\rangle$ for given $w \in V$ is an orthogonality condition

$$
\left\langle w \oplus T^{*} w, U(v \oplus T v)\right\rangle=0 \quad \text { (for all } v \text { in the domain of } T \text { ) }
$$

The graph of any sub-adjoint is a subset of $X=U(\text { graph } T)^{\perp}$. Since $T$ is densely-defined, for given $w \in V$ there is at most one possible value $w^{\prime}$ such that $w \oplus w^{\prime} \in X$, so this orthogonality condition determines

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a well-defined function $T^{*}$ on a subset of $V$, by $T^{*} w=w^{\prime}$ if there exists $w^{\prime} \in V$ such that $w \oplus w^{\prime} \in X$. Linearity of $T^{*}$ is immediate. It is maximal among sub-adjoints to $T$ because the graph of any sub-adjoint is a subset of the graph of $T^{*}$. Orthogonal complements are closed, so $T^{*}$ has a closed graph.
[9.1.2] Corollary: For $T_{1} \subset T_{2}$ with dense domains, $T_{2}^{*} \subset T_{1}^{*}$.
[9.1.3] Corollary: $T \subset T^{* *}$ for densely-defined, symmetric $T$.
Proof: Since $T$ is symmetric, and by uniqueness of the adjoint, $T^{*} \supset T$. In particular, $T^{*}$ is densely defined. Hence, from above, $T^{*}$ has an adjoint $T^{* *}$. The description of the adjoint in terms of orthogonality in $V \oplus V$ shows that $T^{* *} \supset T$.
[9.1.4] Corollary: A densely-defined self-adjoint operator has a closed graph.
Proof: Self-adjointness of densely-defined $T$ includes equality of domains $T=T^{*}$. Again, since the graph of $T^{*}$ is an orthogonal complement, it is closed.

Closed-ness of the graph of a self-adjoint operator is essential in proving existence of resolvents, below. [60]
[9.1.5] Corollary: The adjoint $T^{*}$ of a symmetric densely-defined operator $T$ is also symmetric if and only if $T=T^{*}$.

The use of the term symmetric in this context is potentially misleading. The notation $T=T^{*}$ allows an inattentive reader to forget non-trivial assumptions on the domains of the operators. The equality of domains of $T$ and $T^{*}$ is critical for legitimate computations.
[9.1.6] Proposition: Eigenvalues for symmetric operators $T, D$ are real.
Proof: Suppose $0 \neq v \in D$ and $T v=\lambda v$. Then

$$
\lambda\langle v, v\rangle=\langle\lambda v, v\rangle=\langle T v, v\rangle=\left\langle v, T^{*} v\right\rangle \quad \text { (because } v \in D \subset D^{*} \text { ) }
$$

Because $T^{*}$ agrees with $T$ on $D,\left\langle v, T^{*} v\right\rangle=\langle v, \lambda v\rangle=\bar{\lambda}\langle v, v\rangle$. Thus, $\lambda\langle v, v\rangle=\bar{\lambda}\langle v, v\rangle$.
The resolvent of $T$ is $R_{\lambda}=(T-\lambda)^{-1}$ for $\lambda \in \mathbb{C}$, when this inverse exists as a continuous, everywhere-defined linear operator on $V$.
[9.1.7] Theorem: Let $T$ be self-adjoint with dense domain $D$. For $\lambda \in \mathbb{C}, \lambda \notin \mathbb{R}$, the image $(T-\lambda) D$ is the whole Hilbert space $V$. The resolvent $R_{\lambda}$ exists. For $T$ positive, for $\lambda \notin[0,+\infty)$, the image $(T-\lambda) D$ is the whole space $V$, and $R_{\lambda}$ exists.
Proof: For $\lambda=x+i y$ off the real line and $v$ in the domain of $T$,

$$
\begin{gathered}
|(T-\lambda) v|^{2}=|(T-x) v|^{2}+\langle(T-x) v, i y v\rangle+\langle i y v,(T-x) v\rangle+y^{2}|v|^{2} \\
=|(T-x) v|^{2}-i y\langle(T-x) v, v\rangle+i y\langle v,(T-x) v\rangle+y^{2}|v|^{2}
\end{gathered}
$$

By the symmetry of $T$, and the fact that the domain of $T^{*}$ contains that of $T,\langle v, T v\rangle=\left\langle T^{*} v, v\right\rangle=\langle T v, v\rangle$. Thus,

$$
|(T-\lambda) v|^{2}=|(T-x) v|^{2}+y^{2}|v|^{2} \geq y^{2}|v|^{2}
$$

For $y \neq 0$ and $v \neq 0,(T-\lambda) v \neq 0$, so $T-\lambda$ is injective. We must show that $(T-\lambda) D$ is the whole Hilbert space $V$. If

$$
0=\langle(T-\lambda) v, w\rangle \quad(\text { for all } v \in D)
$$

then the adjoint of $T-\lambda$ can be defined on $w$ simply as $(T-\lambda)^{*} w=0$, since

$$
\langle(T-\lambda) v, w\rangle=0=\langle v, 0\rangle \quad \text { (for all } v \in D)
$$

[60] In general, the graph-closure of an operator need not be the graph of an operator! An operator whose graphclosure is the graph of an operator is close-able. A broader discussion of unbounded operators would consider such issues at greater length, but such a discussion is not necessary here.

## 9. Unbounded operators on Hilbert spaces

Thus, $T^{*}=T$ is defined on $w$, and $T w=\bar{\lambda} w$. For $\lambda$ not real, this implies $w=0$. Thus, $(T-\lambda) D$ is dense in $V$.

Since $T$ is self-adjoint, it is closed, so $T-\lambda$ is closed. The equality

$$
|(T-\lambda) v|^{2}=|(T-x) v|^{2}+y^{2}|v|^{2}
$$

gives

$$
|(T-\lambda) v|^{2} \gg_{y}|v|^{2}
$$

Thus, for fixed $y \neq 0$, the map

$$
F: v \oplus(T-\lambda) v \longrightarrow(T-\lambda) v
$$

respects the metrics, in the sense that

$$
|(T-\lambda) v|^{2} \leq|(T-\lambda) v|^{2}+|v|^{2} \lll y|(T-\lambda) v|^{2} \quad(\text { for fixed } y \neq 0)
$$

The graph of $T-\lambda$ is closed, so is a complete metric subspace of $V \oplus V$. Since $F$ respects the metrics, it preserves completeness. Thus, the metric space $(T-\lambda) D$ is complete, so is a closed subspace of $V$. Since the closed subspace $(T-\lambda) D$ is dense, it is $V$. Thus, for $\lambda \notin \mathbb{R}, R_{\lambda}$ is everywhere-defined. Its norm is bounded by $1 /|\operatorname{Im} \lambda|$, so it is a continuous linear operator on $V$.

Similarly, for $T$ positive, for $\operatorname{Re}(\lambda)<0$,

$$
|(T-\lambda) v|^{2}=|T v|^{2}-\lambda\langle T v, v\rangle-\bar{\lambda}\langle v, T v\rangle+|\lambda|^{2} \cdot|v|^{2}=|T v|^{2}+2|\operatorname{Re} \lambda|\langle T v, v\rangle+|\lambda|^{2} \cdot|v|^{2} \geq|\lambda|^{2} \cdot|v|^{2}
$$

Then the same argument proves the existence of an everywhere-defined inverse $R_{\lambda}=(T-\lambda)^{-1}$, with $\left\|R_{\lambda}\right\|<1 /|\lambda|$ for $\operatorname{Re} \lambda<0$.
[9.1.8] Theorem: (Hilbert) For $T$ self-adjoint, for points $\lambda, \mu$ off the real line, or, for $T$ positive self-adjoint and $\lambda, \mu$ off $[0,+\infty)$,

$$
R_{\lambda}-R_{\mu}=(\lambda-\mu) R_{\lambda} R_{\mu}
$$

For the operator-norm topology, $\lambda \rightarrow R_{\lambda}$ is holomorphic at such points.
Proof: From the previous theorem, for such $T$, $\lambda$, the image $(T-\lambda) D$ is the whole Hilbert space $V$. Applying $R_{\lambda}$ to

$$
1_{V}-(T-\lambda) R_{\mu}=((T-\mu)-(T-\lambda)) R_{\mu}=(\lambda-\mu) R_{\mu}
$$

gives

$$
R_{\lambda}\left(1_{V}-(T-\lambda) R_{\mu}\right)=R_{\lambda}((T-\mu)-(T-\lambda)) R_{\mu}=R_{\lambda}(\lambda-\mu) R_{\mu}
$$

Then

$$
\frac{R_{\lambda}-R_{\mu}}{\lambda-\mu}=R_{\lambda} R_{\mu}
$$

For holomorphy, with $\lambda \rightarrow \mu$,

$$
\frac{R_{\lambda}-R_{\mu}}{\lambda-\mu}-R_{\mu}^{2}=R_{\lambda} R_{\mu}-R_{\mu}^{2}=\left(R_{\lambda}-R_{\mu}\right) R_{\mu}=(\lambda-\mu) R_{\lambda} R_{\mu} R_{\mu}
$$

Taking operator norm, using $\left\|R_{\lambda}\right\| \leq 1 /|\operatorname{Im} \lambda|$ from the previous computation,

$$
\left\|\frac{R_{\lambda}-R_{\mu}}{\lambda-\mu}-R_{\mu}^{2}\right\| \leq \frac{|\lambda-\mu|}{|\operatorname{Im} \lambda| \cdot|\operatorname{Im} \mu|^{2}}
$$

Thus, for $\mu \notin \mathbb{R}$, as $\lambda \rightarrow \mu$, this operator norm goes to 0 , demonstrating the holomorphy.
For positive $T$, the estimate $\left\|R_{\lambda}\right\| \leq 1 /|\lambda|$ for $\operatorname{Re} \lambda \leq 0$ yields holomorphy on the negative real axis by the same argument.

### 9.2 Friedrichs' self-adjoint extensions of semi-bounded operators

Semi-bounded operators are more tractable than general unbounded symmetric operators. A denselydefined symmetric operator $T, D$ is positive (or non-negative), denoted $T \geq 0$, when $\langle T v, v\rangle \geq 0$ for all $v \in D$. All the eigenvalues of a positive operator are non-negative real. Similarly, $T$ is negative when $\langle T v, v\rangle \leq 0$ for all $v$ in the (dense) domain of $T$. Generally, if there is a constant $c \in \mathbb{R}$ such that $\langle T v, v\rangle \geq c \cdot\langle v, v\rangle$ (written $T \geq c$ ), or $\langle T v, v\rangle \leq c \cdot\langle v, v\rangle$ (written $T \leq c$ ), say $T$ is semi-bounded. The following argument for positive operators can easily be adapted to the general semi-bounded situation.

For positive, symmetric $T$ on $V$ with dense domain $D$, define a hermitian form $\langle,\rangle_{1}$ and corresponding norm $|\cdot|_{1}$ by ${ }^{[61]}$

$$
\langle v, w\rangle_{1}=\langle v, w\rangle+\langle T v, w\rangle=\langle v,(1+T) w\rangle=\langle(1+T) v, w\rangle \quad(\text { for } v, w \in D)
$$

The symmetry and positivity of $T$ make $\langle,\rangle_{1}$ positive-definite hermitian on $D$, and $\langle v, w\rangle_{1}$ has sense whenever at least one of $v, w$ is in $D$. Let $V^{1}$ be the Hilbert-space completion of $D$ with respect to the metric $d_{1}$ induced by the norm $|\cdot|_{1}$ on $D$. The completion $V^{1}$ continuously maps to $V$ : for $v_{i}$ a $d_{1}$-Cauchy sequence in $D, v_{i}$ is Cauchy in $V$ in the original topology, since $\left|v_{i}-v_{j}\right| \leq\left|v_{i}-v_{j}\right|_{1}$. For two sequences $v_{i}, w_{j}$ with the same $d_{1}$-limit $v$, the $d_{1}$-limit of $v_{i}-w_{i}$ is 0 , so $\left|v_{i}-w_{i}\right| \leq\left|v_{i}-w_{i}\right|_{1} \rightarrow 0$.
[9.2.1] Theorem: (Friedrichs) A positive, densely-defined, symmetric operator $T$ with domain $D$ dense in Hilbert space $V$ has a positive self-adjoint extension $\widetilde{T}$ with domain $\widetilde{D} \subset V^{1}$, characterized by

$$
\left.\left\langle(1+T) v,(1+\widetilde{T})^{-1} w\right\rangle=\langle v, w\rangle \quad \text { (for } v \in D \text { and } w \in V\right)
$$

The bound $\langle\widetilde{T} v, v\rangle \geq 0$ for $v$ in the domain $\widetilde{D}$ of $\widetilde{T}$ is preserved. The resolvent $(1+\widetilde{T})^{-1}: V \rightarrow V^{1}$ is continuous with the finer topology on $V^{1}$.
Proof: First, let $j$ be the continuous linear map $j: V^{1} \rightarrow V$ obtained by extending by continuity the identity map $D \rightarrow D$, with the source being given the $|\cdot|_{1}$ topology and the target being given the $|\cdot|$ topology. We claim that $j$ is an injection. By construction, $\langle v, w\rangle_{1}=\langle j v, T w\rangle$ for $v \in V^{1}$ and $w \in D$. For $0 \neq v \in V^{1}$, since $D$ is dense in $V^{1}$, there exists $w \in D$ such that $\langle v, w\rangle_{1} \neq 0$. For that $v$,

$$
0 \neq\langle v, w\rangle_{V^{1}}=\langle j v, T w\rangle
$$

Thus, $j v \neq 0$ for $0 \neq v \in V^{1}$, and $j$ is indeed injective. We may identify $V^{1}$ with its image in $V$, noting that $V^{1}$ has a finer topology than that induced from $V$.

For $h \in V$ and $v \in V^{1}$, the functional $\lambda_{h}: v \rightarrow\langle v, h\rangle$ has a bound

$$
\left|\lambda_{h} v\right| \leq|v| \cdot|h| \leq|v|_{1} \cdot|h|
$$

so the norm of the functional $\lambda_{h}$ on $V^{1}$ is at most $|h|$. By Riesz-Fréchet, there is unique $B h$ in the Hilbert space $V^{1}$ with $|B h|_{1} \leq|h|$, such that $\lambda_{h}(v)=\langle v, B h\rangle_{1}$ for $v \in V^{1}$, and then $|B h| \leq|B h|_{1} \leq|h|$. The map $B: V \rightarrow V^{1}$ is verifiably linear. There is an obvious symmetry of $B:$

$$
\langle B v, w\rangle=\lambda_{w}(B v)=\langle B v, B w\rangle_{1}=\overline{\langle B w, B v\rangle_{1}}=\overline{\lambda_{v}(B w)}=\overline{\langle B w, v\rangle}=\langle v, B w\rangle \quad(\text { for } v, w \in V)
$$

Positivity of $B$ is similar:

$$
\langle v, B v\rangle=\lambda_{v}(B v)=\langle B v, B v\rangle_{1} \geq\langle B v, B v\rangle \geq 0
$$

$B$ is injective: for $B w=0$, for all $v \in V^{1}$

$$
0=\langle v, 0\rangle_{1}=\langle v, B w\rangle_{1}=\lambda_{w}(v)=\langle v, w\rangle
$$

[^17]
## 9. Unbounded operators on Hilbert spaces

Since $V^{1}$ is dense in $V$, this gives $w=0$. The image of $B$ is dense in $V^{1}$ : if $w \in V^{1}$ is such that $\langle B v, w\rangle_{1}=\lambda_{v}(w)=0$ for all $v \in V$, taking $v=w$ gives

$$
0=\lambda_{w}(w)=\langle w, B w\rangle_{1}=\langle B w, B w\rangle
$$

and by injectivity $w=0$. Thus, $B: V \rightarrow V^{1} \subset V$ is bounded, symmetric, positive, injective, with dense image. In particular, $B$ is self-adjoint.

Thus, $B$ has a possibly unbounded positive, symmetric inverse $A$. Since $B$ injects $V$ to a dense subset $V^{1}$, necessarily $A$ surjects from its domain (inside $V^{1}$ ) to $V$. We claim that $A$ is self-adjoint. Let $S: V \oplus V \rightarrow V \oplus V$ by $S(v \oplus w)=w \oplus v$. Then graph $A=S$ (graph $B$ ). In computing orthogonal complements $X^{\perp}$, clearly

$$
(S X)^{\perp}=S\left(X^{\perp}\right)
$$

From the obvious $U \circ S=-S \circ U$, compute

$$
\begin{aligned}
\operatorname{graph} A^{*}= & (U \text { graph } A)^{\perp}=(U \circ S \text { graph } B)^{\perp}=(-S \circ U \text { graph } B)^{\perp} \\
& =-S\left((U \text { graph } B)^{\perp}\right)=-\operatorname{graph} A=\operatorname{graph} A
\end{aligned}
$$

since the domain of $B^{*}$ is the domain of $B$. Thus, $A$ is self-adjoint.
We claim that for $v$ in the domain of $A,\langle A v, v\rangle \geq\langle v, v\rangle$. Indeed, letting $v=B w$,

$$
\langle v, A v\rangle=\langle B w, w\rangle=\lambda_{w} B w=\langle B w, B w\rangle_{1} \geq\langle B w, B w\rangle=\langle v, v\rangle
$$

Similarly, with $v^{\prime}=B w^{\prime}$, and $v \in V^{1}$,

$$
\left\langle v, A v^{\prime}\right\rangle=\left\langle v, w^{\prime}\right\rangle=\lambda_{w^{\prime}} v=\left\langle v, B w^{\prime}\right\rangle_{1}=\left\langle v, v^{\prime}\right\rangle_{1} \quad\left(v \in V^{1}, v^{\prime} \text { in the domain of } A\right)
$$

Last, prove that $A$ is an extension of $S=1+T$. On one hand, as above,

$$
\langle v, S w\rangle=\lambda_{S w} v=\langle v, B S w\rangle_{1} \quad(\text { for } v, w \in D)
$$

On the other hand, by definition of $\langle,\rangle_{1}$,

$$
\langle v, S w\rangle=\langle v, w\rangle_{1} \quad(\text { for } v, w \in D)
$$

Thus,

$$
\langle v, w-B S w\rangle_{1}=0 \quad(\text { for all } v, w \in D)
$$

Since $D$ is $d_{1}$-dense in $V^{1}, B S w=w$ for $w \in D$. Thus, $w \in D$ is in the range of $B$, so is in the domain of $A$, and

$$
A w=A(B S w)=S w
$$

Thus, the domain of $A$ contains that of $S$ and extends $S$, so the domain of $A$ is dense in $V^{1}$ in the $d_{1}$-topology. In fact, $B=(1+\widetilde{T})^{-1}$ maps $V \rightarrow V^{1}$ continuously even with the finer $\langle,\rangle_{1}$-topology on $V^{1}$ : the relation $\langle v, B w\rangle_{1}=\langle v, w\rangle$ for $v \in V^{1}$ with $v=B w$ gives

$$
|B w|_{1}^{2}=\langle B w, B w\rangle_{1}=\langle B w, w\rangle \leq|B w| \cdot|w| \leq|B w|_{1} \cdot|w|
$$

The resulting $|B w|_{1} \leq|w|$ is continuity in the finer topology.
Continuity of $(1+\widetilde{T})^{-1}: V \rightarrow V^{1}$ with the finer topology on $V^{1}$ is a useful property of Friedrichs' selfadjoint extensions not shared by the other self-adjoint extensions of a given symmetric operator. It has an important corollary:
[9.2.2] Corollary: When the inclusion $V^{1} \rightarrow V$ is compact, the resolvent $(1+\widetilde{T})^{-1}: V \rightarrow V$ is compact.
Proof: In the notation of the proof of the theorem, $B: V \rightarrow V^{1} \rightarrow V$ is the composition of this continuous map with the injection $V^{1} \rightarrow V$ where $V^{1}$ has the finer topology. The composition of a continuous linear map with a compact operator is compact, so compactness of $V^{1} \rightarrow V$ with the finer topology on $V^{1}$ suffices to prove compactness of the resolvent.

### 9.3 Example: incommensurable self-adjoint extensions

The differential operator $\frac{d^{2}}{d x^{2}}$ on $L^{2}[a, b]$ or $L^{2}(\mathbb{R})$ is a prototypical natural unbounded operator. It is undeniably not continuous in the $L^{2}$ topology: on $L^{2}[0,1]$ the norm of $f(x)=x^{n}$ is $1 / \sqrt{2 n+1}$, and the second derivative of $x^{n}$ is $n(n-1) x^{n-2}$, so the ratio of $L^{2}$ norms $\left|\left(x^{n}\right)^{\prime \prime}\right| /\left|x^{n}\right|$ goes to $+\infty$ as $n \rightarrow+\infty$. Since the operator is unbounded on polynomials, it certainly has no bounded extension to $L^{2}[0,1]$.

Just below, we exhibit a continuum of mutually incomparable self-adjoint extensions of the restriction $T$ of $-\frac{d^{2}}{d x^{2}}$ to smooth functions on $[a, b]$ vanishing at the endpoints. As this will illustrate, it is unreasonable to expect naturally-occurring positive/negative, symmetric operators to have unique self-adjoint extensions. In brief, for unbounded operators arising from differential operators, varying boundary conditions gives mutually incomparable self-adjoint extensions. In that situation, the graph-closure $\bar{T}=T^{* *}$ is not selfadjoint. Equivalently, $T^{*}$ is not symmetric, proven as follows.

In general, the graph-closure of an unbounded operator need not be the graph of an operator, but this potential problem does not exist for a semi-bounded operator $T$, since the Friedrichs self-adjoint extension $\widetilde{T}$ exists, and the graph of $\widetilde{T}$ contains the graph-closure of $T$.

Suppose positive, symmetric, densely-defined $T$ has positive, symmetric extensions $A, B$ admitting no common symmetric extension. Let $\bar{A}=A^{* *}, \bar{B}=B^{* *}$ be the graph-closures of $A, B$. Friedrichs' construction $T \rightarrow \widetilde{T}$ applies to $T, A, B$. The inclusion-reversing property of $S \rightarrow S^{*}$ gives a diagram of extensions, where ascending lines indicate extensions:


Since $T^{*}$ is a common extension of $A, B$, but $A, B$ have no common symmetric extension, $T^{*}$ cannot be symmetric. Thus, any such situation gives an example of non-symmetric adjoints of symmetric operators. Equivalently, $\bar{T}$ cannot be self-adjoint, because its adjoint is $T^{*}$, which cannot be symmetric.

Further, although the graph closures $\bar{A}$ and $\bar{B}$ are (not necessarily proper) extensions of $\bar{T}$, neither of their Friedrichs extensions can be directly comparable to that of $\bar{T}$ without being equal to it, since comparable self-adjoint densely-defined operators are necessarily equal: a densely-defined self-adjoint operator cannot be a proper extension of another such: for $S \subset T$ with $S=S^{*}$ and $T=T^{*}$, the inclusion-reversing property gives $T=T^{*} \subset S^{*}=S$. By hypothesis, $A, B$ have no common symmetric extension, so both equalities cannot hold.

Let $V=L^{2}[a, b], T=-d^{2} / d x^{2}$, with domain

$$
D_{T}=\left\{f \in C_{c}^{\infty}[a, b]: f \text { vanishes at } a, b\right\}
$$

The sign on the second derivative makes $T$ positive: using the boundary conditions, integrating by parts,

$$
\langle T v, v\rangle=-\left\langle v^{\prime \prime}, v\right\rangle=-v^{\prime}(b) \bar{v}(b)+v^{\prime}(a) \bar{v}(a)+\left\langle v^{\prime}, v^{\prime}\right\rangle=\left\langle v^{\prime}, v^{\prime}\right\rangle \geq 0 \quad\left(\text { for } v \in D_{T}\right)
$$

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Integration by parts twice proves symmetry:

$$
\begin{aligned}
& \langle T v, w\rangle=-\left\langle v^{\prime \prime}, w\right\rangle=-v^{\prime}(b) \bar{w}(b)+v^{\prime}(a) \bar{w}(a)+\left\langle v^{\prime}, w^{\prime}\right\rangle=\left\langle v^{\prime}, w^{\prime}\right\rangle \\
& =v(b) \bar{w}^{\prime}(b)-v(a) \bar{w}^{\prime}(a)-\left\langle v, w^{\prime \prime}\right\rangle=\langle v, T w\rangle \quad\left(\text { for } v, w \in D_{T}\right)
\end{aligned}
$$

For each pair $\alpha, \beta$ of complex numbers, an extension $T_{\alpha, \beta}=-d^{2} / d x^{2}$ of $T$ is defined by taking a larger domain, by relaxing the boundary conditions:

$$
D_{\alpha, \beta}=\left\{f \in C_{c}^{\infty}[a, b]: f(a)=\alpha \cdot f(b), f^{\prime}(a)=\beta \cdot f^{\prime}(b)\right\}
$$

Integrating by parts,

$$
\left\langle T_{\alpha, \beta} v, w\right\rangle=v^{\prime}(b) \bar{w}(b) \cdot(1-\beta \bar{\alpha})+v(b) \bar{w}^{\prime}(b) \cdot(1-\alpha \bar{\beta})+\left\langle v, T_{\alpha, \beta} w\right\rangle \quad\left(\text { for } v, w \in D_{\alpha, \beta}\right)
$$

The values $v^{\prime}(b), v(b), w(b)$, and $w^{\prime}(b)$ can be arbitrary, so the extension $T_{\alpha, \beta}$ is symmetric if and only if $\alpha \bar{\beta}=1$, and in that case $T$ is positive, since again

$$
\left\langle T_{\alpha, \beta} v, v\right\rangle=-\left\langle v^{\prime \prime}, v\right\rangle=\left\langle v^{\prime}, v^{\prime}\right\rangle \geq 0 \quad\left(\text { for } \alpha \bar{\beta}=1 \text { and } v \in D_{\alpha, \beta}\right)
$$

For two values $\alpha, \alpha^{\prime}$, taking $\beta=1 / \bar{\alpha}$ and $\beta^{\prime}=1 / \bar{\alpha}^{\prime}$, for the symmetric extensions $T_{\alpha, \beta}$ and $T_{\alpha^{\prime}, \beta^{\prime}}$ to have a common symmetric extension $\widetilde{T}$ requires that the domain of $\widetilde{T}$ include both $D_{\alpha, \beta} \cup D_{\alpha^{\prime}, \beta^{\prime}}$. The integration by parts computation gives

$$
\begin{gathered}
\langle\widetilde{T} v, w\rangle=v^{\prime}(b) \bar{w}(b) \cdot(1-\beta \bar{\alpha})+v(b) \bar{w}^{\prime}(b) \cdot(1-\alpha \bar{\beta})+\left\langle v, T_{\alpha, \beta} w\right\rangle \\
=v^{\prime}(b) \bar{w}(b)\left(1-\beta \bar{\alpha}^{\prime}\right)+v(b) \bar{w}^{\prime}(b) \cdot\left(1-\alpha \bar{\beta}^{\prime}\right)+\langle v, \widetilde{T} w\rangle \quad\left(\text { for } v \in D_{\alpha, \beta}, w \in D_{\alpha^{\prime}, \beta^{\prime}}\right)
\end{gathered}
$$

Thus, the required symmetry $\langle\widetilde{T} v, w\rangle=\langle v, \widetilde{T} w\rangle$ holds only for $\alpha=\alpha^{\prime}$ and $\beta=\beta^{\prime}$. That is, the original operator $T$ has a continuum of distinct symmetric extensions, no two of which admit a common symmetric extension. In particular, no two of these symmetric extensions have a common self-adjoint extension. Yet, each does have at least the Friedrichs positive, self-adjoint extension. Thus, $T$ has infinitely-many distinct positive, self-adjoint extensions.

For example, the two similar boundary-value problems on $L^{2}[0,2 \pi]$

$$
\begin{cases}u^{\prime \prime}=\lambda \cdot u \quad \text { and } \quad u(0)=u(2 \pi), u^{\prime}(0)=u^{\prime}(2 \pi) \\ u^{\prime \prime}=\lambda \cdot u \quad \text { and } \quad u(0)=0=u(2 \pi)\end{cases}
$$

(provably) have eigenfunctions and eigenvalues

$$
\left\{\begin{array}{cll}
1, \sin (n x), \cos (n x) & n=1,2,3, \ldots & \text { eigenvalues } \\
\sin \left(\frac{n x}{2}\right) & n=1,2,3, \ldots & \text { eigenvalues }
\end{array} \frac{1}{4}, 1, \frac{9}{4}, 4, \frac{25}{4}, 9, \frac{49}{4}, \ldots .\right.
$$

Half the eigenfunctions and eigenvalues are common, while the other half of eigenvalues of the first are shifted upward for the second. Both collections of eigenfunctions give orthogonal bases for $L^{2}[0,2 \pi]$. Expressions of the unshared eigenfunctions of one in terms of those of the other are not trivial, despite considerable mythology suggesting the contrary.

### 9.4 Unbounded self-adjoint operators with compact resolvents

The following unsurprising claim and its proof are standard:
[9.4.1] Claim: For a not-necessarily-bounded self-adjoint operator $T$, if $T^{-1}$ exists and is compact, then $(T-\lambda)^{-1}$ exists and is a compact operator for $\lambda$ off a discrete set in $\mathbb{C}$, and is meromorphic in $\lambda$. Further, the spectrum of $T$ and non-zero spectrum of $T^{-1}$ are in the bijection $\lambda \leftrightarrow \lambda^{-1}$.
Proof: The set of eigenvalues or point spectrum of a possibly-unbounded operator $T$ consists of $\lambda \in \mathbb{C}$ such that $T-\lambda$ fails to be injective. The continuous spectrum consists of $\lambda$ with $T-\lambda$ injective and with dense image, but not surjective. Further, for possibly unbounded operators, we require a bounded (=continuous) inverse $(T-\lambda)^{-1}$ on $(T-\lambda) D_{T}$ for $\lambda$ to be in the continuous spectrum. The residual spectrum consists of $\lambda$ with $T-\lambda$ injective, but $(T-\lambda) D_{T}$ not dense.

The description of continuous spectrum simplifies for closed $T$, that is, for $T$ with closed graph: we claim that for $(T-\lambda)^{-1}$ densely defined and continuous, $(T-\lambda) D_{T}$ is the whole space, so $(T-\lambda)^{-1}$ is everywhere defined, so $\lambda$ cannot be in the residual spectrum. Indeed, the continuity gives a constant $C$ such that $|x| \leq C \cdot|(T-\lambda) x|$ for all $x \in D_{T}$. Then $(T-\lambda) x_{i}$ Cauchy implies $x_{i}$ Cauchy, and $T$ closed implies $T\left(\lim x_{i}\right)=\lim T x_{i}$. Thus, $(T-\lambda) D_{T}$ is closed. Then density of $(T-\lambda) D_{T}$ implies it is the whole space.

Now prove that for $T^{-1}$ compact, the resolvent $(T-\lambda)^{-1}$ exists and is compact for $\lambda$ off a discrete set, and is meromorphic in $\lambda$. The non-zero spectrum of the compact self-adjoint operator $T^{-1}$ is point spectrum, from basic spectral theory for such operators, as in [9.A]. We claim that the spectrum of $T$ and non-zero spectrum of $T^{-1}$ are in the obvious bijection $\lambda \leftrightarrow \lambda^{-1}$. From the algebraic identities

$$
T^{-1}-\lambda^{-1}=T^{-1}(\lambda-T) \lambda^{-1} \quad T-\lambda=T\left(\lambda^{-1}-T^{-1}\right) \lambda
$$

failure of either $T-\lambda$ or $T^{-1}-\lambda^{-1}$ to be injective forces the failure of the other, so the point spectra are identical.

For (non-zero) $\lambda^{-1}$ not an eigenvalue of compact $T^{-1}, T^{-1}-\lambda^{-1}$ is injective and has a continuous, everywhere-defined inverse. That $S-\lambda$ is surjective for compact self-adjoint $S$ and $\lambda \neq 0$ not an eigenvalue follows from the spectral theorem for self-adjoint compact operators. For such $\lambda$, inverting the relation $T-\lambda=T\left(\lambda^{-1}-T^{-1}\right) \lambda$ gives

$$
(T-\lambda)^{-1}=\lambda^{-1}\left(\lambda^{-1}-T^{-1}\right)^{-1} T^{-1}
$$

from which $(T-\lambda)^{-1}$ is continuous and everywhere-defined. That is, $\lambda$ is not in the spectrum of $T$. Finally, $\lambda=0$ is not in the spectrum of $T$, because $T^{-1}$ exists and is continuous. This establishes the bijection.

Thus, for $T^{-1}$ compact self-adjoint, the spectrum of $T$ is countable, with no accumulation point in $\mathbb{C}$. Letting $R_{\lambda}=(T-\lambda)^{-1}$, the resolvent relation

$$
R_{\lambda}=\left(R_{\lambda}-R_{0}\right)+R_{0}=(\lambda-0) R_{\lambda} R_{0}+R_{0}=\left(\lambda R_{\lambda}+1\right) \circ R_{0}
$$

expresses $R_{\lambda}$ as the composition of a continuous operator with a compact operator, proving its compactness.

As earlier, continuity is immediate from Hilbert's relation

$$
(T-\lambda)^{-1}(\lambda-\mu)(T-\mu)^{-1}=(T-\lambda)^{-1}((T-\mu)-(T-\lambda))(T-\mu)^{-1}=(T-\lambda)^{-1}-(T-\mu)^{-1}
$$

Then dividing through by $\lambda-\mu$ gives

$$
\frac{(T-\lambda)^{-1}-(T-\mu)^{-1}}{\lambda-\mu}=(T-\lambda)^{-1}(T-\mu)^{-1}
$$

proving differentiability.

### 9.5 Example: $\Delta$ on $L^{2}(\mathbb{T})$ and Sobolev spaces

There are many ways to understand and prove that exponentials $x \rightarrow e^{i \xi x}$ give an orthogonal basis for $L^{2}$ of the circle. Here, we use Friedrichs extensions and compact resolvents.

On the circle $\mathbb{T}=\mathbb{R} / 2 \pi \mathbb{Z}$ or $\mathbb{R} / \mathbb{Z}$, there are no boundary terms in integration by parts, so $\Delta$ has the symmetry

$$
\left.\langle\Delta f, g\rangle=\langle f, \Delta g\rangle \quad \text { (with usual }\langle f, g\rangle=\int_{\mathbb{T}} f \cdot \bar{g}, \text { for } f, g \in C^{\infty}(\mathbb{T})\right)
$$

Friedrichs' self-adjoint extension $\widetilde{\Delta}$ of $\Delta$ is essentially described by the relation

$$
\left\langle(1-\widetilde{\Delta})^{-1} x,(1-\Delta) y\right\rangle=\langle x, y\rangle \quad\left(\text { for } x \in L^{2}(\mathbb{T}) \text { and } y \in C^{\infty}(\mathbb{T})\right)
$$

The compactness of the resolvent $(\widetilde{\Delta}-z)^{-1}$, proven below, and the spectral theorem for compact, self-adjoint operators, yield an orthogonal Hilbert-space basis for $L^{2}(\mathbb{T})$ consisting of $\widetilde{\Delta}$-eigenfunctions. Further, these eigenfunctions will be proven eigenfunctions for $\Delta$ itself.

The compactness of the resolvent will follow from Friedrichs' construction of $\widetilde{\Delta}$ via the continuous linear map $(1-\widetilde{\Delta})^{-1}$, itself a continuous linear map $L^{2}(\mathbb{T}) \longrightarrow H^{1}(\mathbb{T})$, where the Sobolev space $H^{1}(\mathbb{T})$ is the completion of $C^{\infty}(\mathbb{T})$ with respect to the Sobolev norm

$$
|f|_{H^{1}(\mathbb{T})}=\left(|f|_{L^{2}(\mathbb{T})}^{2}+\left|f^{\prime}\right|_{L^{2}(\mathbb{T})}^{2}\right)^{\frac{1}{2}}=\langle(1-\Delta) f, f\rangle^{\frac{1}{2}}
$$

Compactness of the resolvent is Rellich's lemma: the inclusion $H^{1}(\mathbb{T}) \rightarrow L^{2}(\mathbb{T})$ is compact. The map $(1-\widetilde{\Delta})^{-1}$ of $L^{2}(\mathbb{T})$ to itself is compact because it is the composition of the continuous map $(1-\widetilde{\Delta})^{-1}$ : $L^{2}(\mathbb{T}) \rightarrow H^{1}(\mathbb{T})$ and the compact inclusion $H^{1}(\mathbb{T}) \rightarrow L^{2}(\mathbb{T})$.

The eigenfunctions for the extension $\widetilde{\Delta}$ certainly include the $\Delta$-eigenfunctions $\psi_{n}(x)=e^{i n x}$ with $n \in \mathbb{Z}$, but the issue is exactly to show that there are no other eigenfunctions for $\widetilde{\Delta}$ than for $\Delta$. The genuine possibility of exotic eigenfunctions is illustrated in the following section. That is, while we can solve the differential equation $\Delta u=\lambda \cdot u$ on $\mathbb{R}$ and identify $\lambda$ having $2 \pi \mathbb{Z}$-periodic solutions, more must be done to assure that there are no further $\widetilde{\Delta}$-eigenfunctions in the orthogonal basis promised by the spectral theorem. We hope that the natural heuristic, of straightforward solution of the differential equation $\Delta u=\lambda \cdot u$, gives the whole orthogonal basis, but this is exactly the issue.

The $k^{\text {th }}$ Sobolev space $H^{k}(\mathbb{T})$ is the Hilbert space completion of $C^{\infty}(\mathbb{T})$ with respect to $k^{\text {th }}$ Sobolev norm given by

$$
|f|_{H^{k}(\mathbb{T})}^{2}=\left\langle(1-\Delta)^{k} f, f\right\rangle_{L^{2}} \quad\left(\text { for } 0 \leq k \in \mathbb{Z}, \text { for } f \in C^{\infty}(\mathbb{T})\right)
$$

There are other useful, slightly different, expressions for a $k^{\text {th }}$ Sobolev norm, such as

$$
\left(|f|_{L^{2}(\mathbb{T})}^{2}+\left|f^{\prime}\right|_{L^{2}(\mathbb{T})}^{2}+\ldots+\left|f^{(k)}\right|^{2}\right)^{\frac{1}{2}}
$$

The two are comparable: with uniform implied constants, [62]

$$
\left(|f|_{L^{2}(\mathbb{T})}^{2}+\left|f^{\prime}\right|_{L^{2}(\mathbb{T})}^{2}+\ldots+\left|f^{(k)}\right|^{2}\right)^{\frac{1}{2}} \asymp\left\langle(1-\Delta)^{k} f, f\right\rangle^{\frac{1}{2}} \quad(\text { for } 0 \leq k \in \mathbb{Z})
$$

but they are not constant multiples of each other. In fact, precise comparison of constants between the two versions of the Sobolev norms proves irrelevant, as we will see that the exponentials $\psi_{n}(x)=e^{i n x}$ give an orthogonal basis for any/all versions of the Hilbert-space structure.

The Sobolev imbedding theorem below [9.5.4] shows that $H^{k+1}(\mathbb{T})$ is inside the Banach space $C^{k}(\mathbb{T})$ with norm

$$
|f|_{C^{k}(\mathbb{T})}=\sup _{0 \leq i \leq k} \sup _{x \in \mathbb{T}}\left|f^{(i)}(x)\right|
$$

[62] Here $a \asymp b$ means that there are positive, finite constants $c, c^{\prime}$ such that $c \cdot a \leq b \leq c^{\prime} \cdot a$.

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by showing the dominance relation

$$
|f|_{H^{k}(\mathbb{T})} \ll|f|_{C^{k}(\mathbb{T})} \ll|f|_{H^{k+1}(\mathbb{T})} \quad \text { giving } \quad H^{k+1}(\mathbb{T}) \subset C^{k}(\mathbb{T}) \subset H^{k}(\mathbb{T})
$$

The inclusions $C^{k}(\mathbb{T}) \subset H^{k}(\mathbb{T})$ follow from the density of $C^{\infty}(\mathbb{T})$ in every $C^{k}(\mathbb{T})$. Letting $H^{\infty}(\mathbb{T})=\lim _{k} H^{k}(\mathbb{T})$, the intersection $C^{\infty}(T)$ of Banach spaces $C^{k}(\mathbb{T})$ is an intersection of Hilbert spaces

$$
H^{\infty}(\mathbb{T})=\bigcap_{k} H^{k}(\mathbb{T})=\bigcap_{k} H^{k+1}(\mathbb{T}) \subset \bigcap_{k} C^{k}(\mathbb{T})=C^{\infty}(\mathbb{T}) \subset H^{\infty}(\mathbb{T})
$$

For $f \in C^{\infty}(\mathbb{T})$, let $f \rightarrow \bar{f}$ be the value-wise conjugation. Extend this by continuity to $f \rightarrow \bar{f}$ on $L^{2}(\mathbb{T})$, and let $(\Lambda f)(g)=\langle g, \bar{f}\rangle_{L^{2}}$. From the Riesz-Frechet theorem, $\Lambda$ is an isomorphism of $L^{2}(\mathbb{T})$ to itself. For $1 \leq k \in \mathbb{Z}$, let $H^{-k}(\mathbb{T})$ be the Hilbert-space dual of $H^{k}(\mathbb{T})$, not identified with $H^{k}(\mathbb{T})$ itself via Riesz-Frechet. With $i: H^{1}(\mathbb{T}) \rightarrow L^{2}(\mathbb{T})$ the inclusion, let $i^{*}: L^{2}(\mathbb{T})^{*} \rightarrow H^{-1}(\mathbb{T})$ be the adjoint. Similarly, the adjoint of the inclusion $H^{k+1}(\mathbb{T}) \rightarrow H^{k}(\mathbb{T})$ is $H^{-k}(\mathbb{T}) \rightarrow H^{-(k+1)}(\mathbb{T})$. Thus, suppressing the reference to $\mathbb{T}$, we have


Let $H^{-\infty}(\mathbb{T})=\bigcup_{k \geq 0} H^{-k}(\mathbb{T})=\operatorname{colim}_{k} H^{-k}(\mathbb{T})$. Assuming Sobolev imbedding [9.5.4], we have
[9.5.1] Corollary: The dual space $C^{\infty}(\mathbb{T})^{*}$ to $C^{\infty}(\mathbb{T})$ of distributions on $\mathbb{T}$ is $H^{-\infty}(\mathbb{T})$.
Proof: $C^{\infty}(\mathbb{T})=H^{\infty}(\mathbb{T})$ by Sobolev. In general, the dual of a limit is not the colimit of the duals of the limitands, but when the limitands are normed and the transition maps have dense images, as with $H^{k+1}(\mathbb{T}) \rightarrow H^{k}(\mathbb{T}),[13.14 .4]$ shows that $\left(\lim _{k} H^{k}\right)^{*}=\operatorname{colim}_{k}\left(H^{k}\right)^{*}=\operatorname{colim}_{k} H^{-k}$.
[9.5.2] Corollary: For $f$ in the domain of $\widetilde{\Delta}$, the image $\widetilde{\Delta} f$ is the genuine distributional derivative $\Delta f$, and the domain of $\widetilde{\Delta}$ is $H^{2}(\mathbb{T})$.
Proof: For $g \in C^{\infty}(\mathbb{T})$ and $f$ in the domain of $\widetilde{\Delta}$, by the characterization of the Friedrichs extension $f=(\underset{\sim}{1}-\widetilde{\Delta})^{-1} F$ for some $F$ in $L^{2}$. The characterization $\left\langle(1-\widetilde{\Delta})^{-1} F, g\right\rangle_{H^{1}}=\langle F, g\rangle_{L^{2}}$ gives $\langle f, g\rangle_{H^{1}}=$ $\langle(1-\widetilde{\Delta}) f, g\rangle_{L^{2}}$. Then

$$
\langle(1-\widetilde{\Delta}) f, g\rangle_{L^{2}}=\langle f, g\rangle_{H^{1}}=\langle f,(1-\Delta) g\rangle_{L^{2}}==\langle f,(1-\Delta) g\rangle_{H^{-\infty} \times H^{\infty}}=\langle(1-\Delta) f, g\rangle_{H^{-\infty} \times H^{\infty}}
$$

where we restrict the hermitian pairing on $L^{2} \times L^{2}$ to a hermitian pairing $L^{2} \times H^{\infty}$, and then extend it to $H^{-\infty} \times H^{\infty}$. This holds for all $g \in C^{\infty}(\mathbb{T})$, so $(1-\widetilde{\Delta}) f=(1-\Delta) f$ as distributions. Thus, for $f \in H^{1}(\mathbb{T})$ in the domain of $\widetilde{\Delta}, \Delta f \in L^{2}(\mathbb{T})$. Thus, $\left\langle(1-\Delta)^{2} f, f\right\rangle_{L^{2}}$ is finite, so $f \in H^{2}(\mathbb{T})$.
[9.5.3] Corollary: $\widetilde{\Delta}$-eigenvectors are smooth, and are $\Delta$-eigenvectors.
Proof: A $\lambda$-eigenfunction $u$ for $\widetilde{\Delta}$ is in the domain of $\widetilde{\Delta}$, and by the previous $\widetilde{\Delta} u=\Delta u$, so $\Delta u=\lambda u$. Thus, $(1-\Delta) u=(1-\lambda) u$. By design, $(1-\Delta)$ is a continuous map $H^{k+2} \rightarrow H^{k}:$ for $f \in C^{\infty}(\mathbb{T})$,

$$
|(1-\Delta) f|_{H_{k}}^{2}=\left\langle(1-\Delta)^{k}(1-\Delta) f,(1-\Delta) f\right\rangle_{L^{2}}=\left\langle(1-\Delta)^{k+2} f, f\right\rangle_{L^{2}}=|f|_{H^{k+2}}^{2}
$$

Each $H^{k}$ is the completion of $C^{\infty}(\mathbb{T})$, so $(1-\Delta)$ is an isomorphism $H^{k+2} \rightarrow H^{k}$, and $(1-\Delta)^{-1}$ is an isomorphism $H^{k} \rightarrow H^{k+2}$, so

$$
u=(1-\Delta)^{-1}(1-\lambda) u \in(1-\Delta)^{-1} H^{2}(\mathbb{T}) \subset H^{4}(\mathbb{T})
$$

By induction, $u \in H^{\infty}(\mathbb{T})=C^{\infty}(\mathbb{T})$.
Of course, granting the above, to find smooth $\Delta$-eigenfunctions, solve the differential equation

$$
\lambda \cdot u=\widetilde{\Delta} u=\Delta u=u^{\prime \prime}
$$

## 9. Unbounded operators on Hilbert spaces

on $\mathbb{R}$ : for $\lambda \neq 0$, the solutions are linear combinations of $u(x)=e^{ \pm \sqrt{\lambda} \cdot x}$; for $\lambda=0$, the solutions are linear combinations of $u(x)=1$ and $u(x)=x$. The $2 \pi \mathbb{Z}$-periodicity is equivalent to $\lambda \in i \mathbb{Z}$ in the former case, and eliminates $u(x)=x$ in the latter. As usual, uniqueness is proven via the mean-value theorem. This would prove that the exponentials $\left\{e^{i n x}: n \in \mathbb{Z}\right\}$ are a Hilbert-space basis for $L^{2}(\mathbb{T})$.

Now we prove Sobolev imbedding and Rellich compactness.
[9.5.4] Theorem: (Sobolev imbedding) $H^{k+1}(\mathbb{T}) \subset C^{k}(\mathbb{T})$.
Proof: This is the fundamental theorem of calculus and the Cauchy-Schwarz-Bunyakowsky inequality. The case $k=0$ adequately illustrates the causality: prove that the $H^{1}$ norm dominates the $C^{o}$ norm, namely, sup-norm, on $C^{\infty}(\mathbb{T})$. Use coordinates from the real line, with $\mathbb{T}=\mathbb{R} / \mathbb{Z}$. For $0 \leq x \leq y \leq 1$, the difference between maximum and minimum values of $f \in C^{\infty}[0,1]$ is constrained:

$$
|f(y)-f(x)|=\left|\int_{x}^{y} f^{\prime}(t) d t\right| \leq \int_{0}^{1}\left|f^{\prime}(t)\right| d t \leq\left(\int_{0}^{1}\left|f^{\prime}(t)\right|^{2} d t\right)^{1 / 2} \cdot\left(\int_{x}^{y} 1 d t\right)^{1 / 2}=\left|f^{\prime}\right|_{L^{2}} \cdot|x-y|^{\frac{1}{2}}
$$

Let $y \in[0,1]$ be such that $|f(y)|=\min _{x}|f(x)|$. Using the previous inequality,

$$
\begin{gathered}
|f(x)| \leq|f(y)|+|f(x)-f(y)| \leq \int_{0}^{1}|f(t)| d t+|f(x)-f(y)| \\
\leq \int_{0}^{1}|f| \cdot 1+\left|f^{\prime}\right|_{L^{2}} \cdot 1 \leq|f|_{L^{2}}+\left|f^{\prime}\right|_{L^{2}} \ll 2\left(|f|^{2}+\left|f^{\prime}\right|^{2}\right)^{1 / 2}=2|f|_{H^{1}}
\end{gathered}
$$

Thus, on $C_{c}^{\infty}(\mathbb{T})$ the $H^{1}$ norm dominates the sup-norm. Thus, this comparison holds on the $H^{1}$ completion $H^{1}(\mathbb{T})$, and $H^{1}(\mathbb{T}) \subset C^{o}(\mathbb{T})$. The same argument applies to the $H^{1}$-norm and $C^{o}$-norm of successive derivatives.

The space $C^{\infty}(\mathbb{T})$ of smooth functions on $\mathbb{T}$ is the nested intersection of the spaces $C^{k}(\mathbb{T})$, an instance of a (projective) limit of Banach spaces:

$$
C^{\infty}(\mathbb{T})=\bigcap_{k=0}^{\infty} C^{k}(\mathbb{T})=\lim _{k} C^{k}(\mathbb{T})
$$

so it has a uniquely-determined Fréchet space topology, as in [13.5]. Similarly,

$$
H^{\infty}(\mathbb{T})=\bigcap_{k=0}^{\infty} H^{k}(\mathbb{T})=\lim _{k} H^{k}(\mathbb{T})
$$

[9.5.5] Corollary: $C^{\infty}(\mathbb{T})=H^{\infty}(\mathbb{T})$.
Proof: The interlacing property $C^{k+1}(\mathbb{T}) \subset H^{k+1}(\mathbb{T}) \subset C^{k}(\mathbb{T})$ gives (dashed) compatible maps from $H^{\infty}(\mathbb{T})$ to the spaces $C^{k}(\mathbb{T})$ inducing a unique (dotted) map to the limit $C^{\infty}(\mathbb{T})$ :


Oppositely, we obtain a unique (dotted) map of $C^{\infty}(\mathbb{T})$ to $\lim _{k} H^{k}(\mathbb{T})$ :


Thus, the two dotted maps must be mutual inverses.
Rellich's lemma on $\mathbb{T}$ uses some finer details from the discussion just above, namely, the Lipschitz property $|f(x)-f(y)| \ll|x-y|^{\frac{1}{2}}$ for $|f|_{H^{1}} \leq 1$, and the related fact that the map $H^{1}(T) \rightarrow C^{o}(\mathbb{T})$ has operator norm at most 2 .
[9.5.6] Theorem: (Rellich compactness) The inclusion $H^{k+1}(\mathbb{T}) \rightarrow H^{k}(\mathbb{T})$ is compact.
Proof: The causality is adequately illustrated by the $k=0$ case, showing that the unit ball in $H^{1}(\mathbb{T})$ is totally bounded in $L^{2}(\mathbb{T})$. Approximate $f \in H^{1}(\mathbb{T})$ in $L^{2}(\mathbb{T})$ by piecewise-constant functions

$$
F(x)= \begin{cases}c_{1} & \text { for } 0 \leq x<\frac{1}{n} \\ c_{2} & \text { for } \frac{1}{n} \leq x<\frac{2}{n} \\ \cdots & \\ c_{n} & \text { for } \frac{n-1}{n} \leq x \leq 1\end{cases}
$$

The sup norm of $|f|_{H^{1}} \leq 1$ is bounded by 2 , so we only need $c_{i}$ in the range $\left|c_{i}\right| \leq 2$.
Given $\varepsilon>0$, take $N$ large enough such that the disk of radius 2 in $\mathbb{C}$ is covered by $N$ disks of radius less than $\varepsilon$, with centers $C$. Given $f \in H^{1}(\mathbb{T})$ with $|f|_{1} \leq 1$, choose constants $c_{k} \in C$ such that $\left|f(k / n)-c_{k}\right|<\varepsilon$. Then

$$
\left|f(x)-c_{k}\right| \leq\left|f\left(\frac{k}{n}\right)-c_{k}\right|+\left|f(x)-f\left(\frac{k}{n}\right)\right|<\varepsilon+\left|x-\frac{k}{n}\right|^{\frac{1}{2}} \leq \varepsilon+\frac{1}{\sqrt{n}} \quad\left(\text { for } \frac{k}{n} \leq x \leq \frac{k+1}{n}\right)
$$

Then

$$
\int_{0}^{1}|f-F|^{2} \leq \sum_{k=1}^{n} \int_{k / n}^{(k+1) / n}\left(\varepsilon+\frac{1}{\sqrt{n}}\right)^{2} \leq n \cdot \frac{1}{n} \cdot\left(\varepsilon+\frac{1}{\sqrt{n}}\right)^{2}=\left(\varepsilon+\frac{1}{\sqrt{n}}\right)^{2}
$$

For $\varepsilon$ small and $n$ large, this is small. Thus, the image in $L^{2}(\mathbb{T})$ of the unit ball in $H^{1}(\mathbb{T})$ is totally bounded, so has compact closure. This proves that the inclusion $H^{1}(\mathbb{T}) \subset L^{2}(\mathbb{T})$ is compact.

This completes the arguments for Sobolev imbedding and Rellich compactness, which are used to prove that a $\widetilde{\Delta}$-eigenfunction is actually smooth, so is in the natural domain $\mathbb{C}^{\infty}(\mathbb{T})$ of $\Delta$, justifying determination of the orthogonal basis for $L^{2}(\mathbb{T})$ by solving the differential equations $u^{\prime \prime}=\lambda \cdot u$ in classical (non-distributional terms).

Because the exponentials are smooth eigenfunctions for this self-adjoint extension of $\Delta$,
[9.5.7] Corollary: The exponentials $\psi_{n}(x)=e^{i n x}$ are an orthogonal basis for every Sobolev space $H^{k}(\mathbb{T})$. Proof: The exponentials are smooth, so are inside every $H^{k}(\mathbb{T})$. The $H^{k}$ norm of $\psi_{n}$ is $\left(1+n^{2}\right)^{k / 2}$ times the $L^{2}$ norm:

$$
\left|\psi_{n}\right|_{H^{k}}^{2}=\left\langle(1-\Delta)^{k} \psi_{n}, \psi_{n}\right\rangle_{L^{2}}=\left\langle\left(1+n^{2}\right)^{k} \psi_{n}, \psi_{n}\right\rangle_{L^{2}}=\left(1+n^{2}\right)^{k}\left\langle\psi_{n}, \psi_{n}\right\rangle_{L^{2}}
$$

They are mutually orthogonal in every $H^{k}(\mathbb{T})$ :

$$
\left\langle\psi_{m}, \psi_{n}\right\rangle_{H^{k}}=\left\langle(1-\Delta)^{k} \psi_{m}, \psi_{n}\right\rangle_{L^{2}}=\left\langle\left(1+m^{2}\right)^{k} \psi_{m}, \psi_{n}\right\rangle_{L^{2}}=\left(1+m^{2}\right)^{k} \cdot\left\langle\psi_{m}, \psi_{n}\right\rangle_{L^{2}}
$$

For $f \in H^{k}(\mathbb{T}) \subset L^{2}(\mathbb{T})$, orthogonality

$$
\left.\left.0=\left\langle\psi_{n}, f\right\rangle_{H^{k}}=\left\langle(1-\Delta)^{k} \psi_{n}, f\right\rangle_{L^{2}}=\left\langle\left(1+n^{2}\right)\right)^{k} \psi_{n}, f\right\rangle_{L^{2}}=\left(1+n^{2}\right)\right)^{k}\left\langle\psi_{n}, f\right\rangle_{L^{2}}
$$

gives orthogonality in $L^{2}$, so completeness of the exponentials in $L^{2}$ gives completeness in $H^{k}$.
Granting that the exponentials are an orthogonal basis for every $H^{k}(\mathbb{T})$, we have the spectral characterization of $H^{k}(\mathbb{T})$ :
[9.5.8] Corollary: For $f \in C^{\infty}(\mathbb{T})$, for every $0 \leq k \in \mathbb{Z}$,

$$
|f|_{H^{k}}^{2}=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}}\left|\left\langle f, \psi_{n}\right\rangle_{L^{2}}\right|^{2} \cdot\left(1+n^{2}\right)^{k}
$$

## 9. Unbounded operators on Hilbert spaces

The Fourier series $\frac{1}{2 \pi} \sum_{n}\left\langle F, \psi_{n}\right\rangle_{H^{k}} \cdot \psi_{n}$ of a function $F$ in $H^{k}(\mathbb{T})$ converges to $F$ in the $H^{k}$ topology, and

$$
u\left(\frac{1}{2 \pi} \sum_{n}\left\langle F, \psi_{n}\right\rangle_{H^{k}} \cdot \psi_{n}\right)=\frac{1}{2 \pi} \sum_{n}\left\langle F, \psi_{n}\right\rangle_{H^{k}} \cdot u\left(\psi_{n}\right) \quad\left(\text { for } u \in H^{k}(\mathbb{T})^{*}=H^{-k}(\mathbb{T})\right)
$$

Proof: Starting with Plancherel for Hilbert spaces,

$$
|f|_{H^{k}}^{2}=\frac{1}{2 \pi} \sum_{n} \frac{\left|\left\langle f, \psi_{n}\right\rangle_{H^{k}}\right|^{2}}{\left\langle\psi_{n}, \psi_{n}\right\rangle_{H^{k}}}=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}} \frac{\left|\left\langle f,\left(1+n^{2}\right)^{k} \psi_{n}\right\rangle_{L^{2}}\right|^{2}}{\left(1+n^{2}\right)^{k}}=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}}\left|\left\langle f, \psi_{n}\right\rangle_{L^{2}}\right|^{2} \cdot\left(1+n^{2}\right)^{k}
$$

Completing, the same holds for $F \in H^{k}$. In particular, the partial sums of the Fourier series of $F$ converge to $F$ in $H^{k}$. Thus, by the continuity of $u$ in the dual,

$$
u(F)=u\left(\lim _{M, N} \frac{1}{2 \pi} \sum_{-M \leq n \leq N}\left\langle F, \psi_{n}\right\rangle \cdot \psi_{n}\right)=\lim _{M, N} \frac{1}{2 \pi} \sum_{-M \leq n \leq N}\left\langle F, \psi_{n}\right\rangle \cdot u\left(\psi_{n}\right)=\frac{1}{2 \pi} \sum_{n}\left\langle F, \psi_{n}\right\rangle \cdot u\left(\psi_{n}\right)
$$

a convergent series.
The spectral characterization allows extension to a definition of the $s^{t h}$ Sobolev norm for real $s$ :

$$
|f|_{H^{s}}^{2}=\frac{1}{2 \pi} \sum_{n \in \mathbb{Z}}\left|\left\langle f, \psi_{n}\right\rangle_{L^{2}}\right|^{2} \cdot\left(1+n^{2}\right)^{s} \quad \quad\left(\text { for } f \in C^{\infty}(\mathbb{T})\right)
$$

It would be natural to declare that $H^{s}(\mathbb{T})$ is the completion of $C^{\infty}(\mathbb{T})$ with respect to the $s^{t h}$ norm, for all $s \in \mathbb{R}$, but one potential issue is that $H^{-k}(\mathbb{T})$ is already described as the dual of $H^{k}(\mathbb{T})$, so there is consistency to be checked. First, for duality, the $L^{2}$ Plancherel assertion

$$
\langle f, g\rangle_{L^{2}}=\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle \cdot \overline{\left\langle g, \psi_{n}\right\rangle}
$$

can be desymmetrized:
[9.5.9] Corollary: With Sobolev spaces described as completions of $C^{\infty}(\mathbb{T})$, for $f \in H^{s}(\mathbb{T})$ and $g \in H^{-s}(\mathbb{T})$, the pairing $\langle f, g\rangle_{H^{s} \times H^{-s}}$ defined by

$$
\langle f, g\rangle_{H^{s} \times H^{-s}}=\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{H^{s}} \cdot \overline{\left\langle g, \psi_{n}\right\rangle_{H^{-s}}} \quad\left(\text { for } f, g \in C^{\infty}(\mathbb{T})\right)
$$

gives a (conjugate-linear) isomorphism of $H^{-s}$ to the Hilbert-space dual of $H^{s}$, for all real $s$.
Proof: It suffices to prove that this pairing matches the $L^{2}$ pairing for $f, g \in C^{\infty}(\mathbb{T})$, since the smooth functions are dense in every $H^{s}(\mathbb{T})$. Indeed, for $f, g \in C^{\infty}(\mathbb{T})$,

$$
\begin{gathered}
\langle f, g\rangle_{H^{s} \times H^{-s}}=\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{H^{s}} \cdot \overline{\left\langle g, \psi_{n}\right\rangle_{H^{-s}}}=\frac{1}{2 \pi} \sum_{n}\left(1+n^{2}\right)^{s}\left\langle f, \psi_{n}\right\rangle_{L^{2}} \cdot \overline{\left(1+n^{2}\right)^{-s}\left\langle g, \psi_{n}\right\rangle_{L^{2}}} \\
=\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{L^{2}} \cdot \overline{\left\langle g, \psi_{n}\right\rangle_{L^{2}}}=\langle f, g\rangle_{L^{2}}
\end{gathered}
$$

by ordinary Plancherel for $L^{2}$.
Granting that the Sobolev spaces defined as completions of $C^{\infty}(\mathbb{T})$ are in suitable duality,
[9.5.10] Claim: For a distribution $u$, if $\sum_{n}\left|u\left(\psi_{n}\right)\right|^{2} \cdot\left(1+n^{2}\right)^{-s}<\infty$ then $u \in H^{-s}(\mathbb{T})$, and we have the Fourier expansion

$$
u=\frac{1}{2 \pi} \sum_{n} u\left(\psi_{n}\right) \cdot \psi_{n}=\frac{1}{2 \pi} \sum_{n}\left\langle u, \psi_{n}\right\rangle_{H^{s} \times H^{-s}} \cdot \psi_{n} \quad\left(\text { in } H^{s}(\mathbb{T})\right)
$$

Proof: For $f \in H^{s}(\mathbb{T})$,
$u(f)=\lim _{M, N} u\left(\frac{1}{2 \pi} \sum_{-M \leq n \leq N}\left\langle f, \psi_{n}\right\rangle_{H^{s}} \cdot \psi_{n}\right)=\lim _{M, N} \frac{1}{2 \pi} \sum_{-M \leq n \leq N}\left\langle f, \psi_{n}\right\rangle_{H^{s}} \cdot u\left(\psi_{n}\right)=\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{H^{s}} \cdot u\left(\psi_{n}\right)$
and by Cauchy-Schwarz-Bunyakowsky

$$
\begin{aligned}
& \left|\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{H^{s}} \cdot u\left(\psi_{n}\right)\right|^{2}=\left|\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{L^{2}} \cdot\left(1+n^{2}\right)^{s / 2} \cdot \frac{u\left(\psi_{n}\right)}{\left(1+n^{2}\right)^{s / 2}}\right|^{2} \\
\leq & \frac{1}{2 \pi} \sum_{n}\left|\left\langle f, \psi_{n}\right\rangle\right|_{L^{2}}^{2} \cdot\left(1+n^{2}\right)^{s} \cdot \frac{1}{2 \pi} \sum_{n} \frac{\left|u\left(\psi_{n}\right)\right|^{2}}{\left(1+n^{2}\right)^{s}}=|f|_{H^{s}}^{2} \cdot\left|\sum_{n} u\left(\psi_{n}\right) \cdot \psi_{n}\right|_{H^{-s}}^{2}
\end{aligned}
$$

Thus, $u$ is a continuous linear functional on $H^{s}$, so is in the dual $H^{-s}$. Accommodation of complex conjugation by $u\left(\psi_{n}\right)=\left\langle u, \psi_{-n}\right\rangle_{H^{-s} \times H^{s}}$ and replacing $n$ by $-n$ gives the second form of the Fourier expansion.

Sobolev imbedding becomes simpler and sharper on the spectral side:
[9.5.11] Corollary: For $s>k+\frac{1}{2}, H^{s}(\mathbb{T}) \subset C^{k}(\mathbb{T})$.
Proof: As earlier, it suffices to treat $k=0$, and prove that the $C^{k}$-norm is dominated by the $H^{s}$-norm, on $C^{\infty}$-functions $f$. Indeed,

$$
\begin{gathered}
|f|_{C^{\circ}}=\sup _{x}|f(x)|=\sup _{x}\left|\frac{1}{2 \pi} \sum_{n}\left\langle f, \psi_{n}\right\rangle_{L^{2}} \cdot e^{i n x}\right| \leq \frac{1}{2 \pi} \sum_{n}\left|\left\langle f, \psi_{n}\right\rangle_{L^{2}}\right| \\
=\frac{1}{2 \pi} \sum_{n}\left|\left\langle f, \psi_{n}\right\rangle_{L^{2}}\right|\left(1+n^{2}\right)^{s / 2} \cdot\left(1+n^{2}\right)^{-s / 2} \leq \frac{1}{2 \pi}\left(\sum_{n}\left|\left\langle f, \psi_{n}\right\rangle_{L^{2}}\right|^{2}\left(1+n^{2}\right)^{s}\right)^{\frac{1}{2}} \cdot\left(\sum_{n}\left(1+n^{2}\right)^{-s}\right)^{\frac{1}{2}}
\end{gathered}
$$

by Cauchy-Schwarz-Bunyakowsky. For $s>\frac{1}{2}$, the latter sum is finite. Similarly, for general $k$,

$$
|f|_{C_{k}} \ll_{s}|f|_{H^{s}} \quad\left(\text { for any } s>\frac{1}{2}+k, \text { for } f \in C^{\infty}(\mathbb{T})\right)
$$

Thus, the completion $H^{s}(\mathbb{T})$ has a canonical inclusion to $C^{k}(\mathbb{T})$.
Similarly, granting that the exponentials form an orthonormal basis for every $H^{s}$, proof of an extended form of Rellich's lemma becomes easier:
[9.5.12] Corollary: For real $t>s$, the inclusion $H^{t}(\mathbb{T}) \subset H^{s}(\mathbb{T})$ is compact. For $t>s+\frac{1}{2}$, this inclusion is Hilbert-Schmidt, and for $t>s+1$ it is trace-class.
Proof: The inclusion maps one orthogonal basis to another, but the lengths change. That is, ignoring the constant $2 \pi,\left|\psi_{n}\right|_{H^{s}}=\left(1+n^{2}\right)^{s / 2}$. Thus, letting $e_{n}^{s}=\psi_{n} /\left(1+n^{2}\right)^{s / 2} / \sqrt{2 \pi}$ be an orthonormal basis for $H^{s}$, the inclusion map maps $e_{n}^{t} \rightarrow e_{n}^{s} \cdot\left(1+n^{2}\right)^{s-t}$.

Generally, a map $T: V \rightarrow W$ of Hilbert spaces with orthonormal bases $\left\{e_{n}\right\},\left\{f_{n}\right\}$, of the form $T e_{n}=\lambda_{n} f_{n}$, is compact when $\lambda_{n} \rightarrow 0$, is Hilbert-Schmidt when $\sum_{n}\left|\lambda_{n}\right|^{2}<\infty$, and is trace-class when $\sum_{n}\left|\lambda_{n}\right|<\infty$.
[9.5.13] Corollary: $H^{\infty}(\mathbb{T})=C^{\infty}(\mathbb{T})$ is nuclear Fréchet, in the sense that it is a (projective) limit of Hilbert spaces $V_{n}$, with transition maps $V_{n} \rightarrow V_{n-1}$ Hilbert-Schmidt. ///

Sobolev imbedding and Rellich compactness on $\mathbb{T}^{n}$ can be proven either by reducing to the case of a single circle $\mathbb{T}$, or by repeating analogous arguments directly on $\mathbb{T}^{n}$. Sobolev norms and spaces $H^{s}\left(\mathbb{T}^{n}\right)$ can be defined for real $s$. The index-shift in the Sobolev imbedding is easy to understand in the spectral form:
[9.5.14] Theorem: (Sobolev imbedding) $H^{s}\left(\mathbb{T}^{n}\right) \subset C^{k}\left(\mathbb{T}^{n}\right)$ for $s>k+\frac{n}{2}$.
Proof: Index the characters on $\mathbb{T}^{n}$ by $\psi_{\xi}(x)=e^{i \xi \cdot x}$, where $\xi \cdot x$ is the usual pairing on $\mathbb{R}^{n} \times \mathbb{R}^{n}$. For $k=0$,

$$
\begin{gathered}
|f|_{C^{o}}=\sup _{x}|f(x)|=\sup _{x}\left|\frac{1}{(2 \pi)^{n}} \sum_{\xi \in \mathbb{Z}^{n}}\left\langle f, \psi_{\xi}\right\rangle_{L^{2}} \cdot \psi_{\xi}(x)\right| \leq \frac{1}{(2 \pi)^{n}} \sum_{\xi}\left|\left\langle f, \psi_{\xi}\right\rangle_{L^{2}}\right| \\
=\frac{1}{(2 \pi)^{n}} \sum_{\xi}\left|\left\langle f, \psi_{\xi}\right\rangle_{L^{2}}\right|\left(1+|\xi|^{2}\right)^{s / 2} \cdot\left(1+|\xi|^{2}\right)^{-s / 2} \leq \frac{1}{(2 \pi)^{n}}\left(\sum_{\xi}\left|\left\langle f, \psi_{\xi}\right\rangle_{L^{2}}\right|^{2}\left(1+|\xi|^{2}\right)^{s}\right)^{\frac{1}{2}} \cdot\left(\sum_{\xi}\left(1+|\xi|^{2}\right)^{-s}\right)^{\frac{1}{2}}
\end{gathered}
$$

by Cauchy-Schwarz-Bunyakowsky, where $|\xi|=(\xi \cdot \xi)^{\frac{1}{2}}$. For $s>\frac{n}{2}$, the latter sum is finite.
[9.5.15] Theorem: (Rellich compactness) $H^{t}\left(\mathbb{T}^{n}\right) \subset H^{s}\left(\mathbb{T}^{n}\right)$ is compact for $t>s$, Hilbert-Schmidt for $t>s+\frac{n}{2}$, and trace-class for $t>s+n$.
[9.5.16] Corollary: $H^{\infty}\left(\mathbb{T}^{n}\right)=C^{\infty}\left(\mathbb{T}^{n}\right)$ is nuclear Fréchet.
The relevance of the nuclearity property is the Schwartz kernel theorem [Schwartz 1950] for this situation: every continuous linear map $T: C^{\infty}\left(\mathbb{T}^{m}\right) \rightarrow C^{\infty}\left(\mathbb{T}^{n}\right)^{*}$ from smooth functions to distributions has a Schwartz kernel $K(x, y) \in C^{\infty}\left(\mathbb{T}^{m+n}\right)^{*}$, meaning that

$$
(T \varphi)(\psi)=K(\varphi \otimes \psi) \quad\left(\text { for } \varphi \in C^{\infty}\left(\mathbb{T}^{m}\right) \text { and } \psi \in C^{\infty}\left(\mathbb{T}^{n}\right)\right)
$$

### 9.6 Exotic eigenfunctions on $\mathbb{T}$

In this example, the exotic eigenfunctions are not truly exotic, but do illustrate the possibility that eigenfunctions for a Friedrichs self-adjoint extension of a slight restriction of an otherwise natural differential operator like $\Delta$ may fail to be smooth.

Let $\delta$ be the Dirac delta distribution at 0 in $\mathbb{T} \approx \mathbb{R} / 2 \pi \mathbb{Z}$. The Sobolev imbedding $H^{1}(\mathbb{T}) \subset C^{o}(\mathbb{T})$ essentially shows that $\delta \in H^{-s}$ for every $s>\frac{1}{2}$ :

$$
\left|\sum_{n} \delta \psi_{n} \cdot \psi_{n}\right|_{H^{-s}}^{2}=\sum_{n} \frac{\left|\delta \psi_{n}\right|^{2}}{\left(1+n^{2}\right)^{s}}=\sum_{n} \frac{1}{\left(1+n^{2}\right)^{s}}<+\infty \quad\left(\text { for } s>\frac{1}{2}\right)
$$

Thus, ker $\delta$ is a closed subspace of $H^{1}(\mathbb{T})$. Let $S$ be the restriction of $\Delta$ to the space $D_{S}$ of smooth functions $f$ such that $\delta f=0$. A functional on $L^{2}(\mathbb{T})$ is continuous if and only if its kernel is closed, so $D_{S}$ is dense in $L^{2}(\mathbb{T})$. Let $\widetilde{S}$ be its Friedrichs extension.
[9.6.1] Claim: The domain of $\widetilde{S}$ is

$$
\left\{f \in H^{1}(\mathbb{T}): \Delta f \in L^{2}(\mathbb{T})+\mathbb{C} \cdot \delta, \text { and } \delta f=0\right\}
$$

and $\widetilde{S} f=u$ if and only if $\Delta f=u+c \cdot \delta$ for some $c \in \mathbb{C}$, and $\delta f=0$. In particular, $f \in H^{1}(\mathbb{T})$ such that $(\Delta-\lambda) f=\delta$ and $\delta f=0$ is an eigenfunction for $\widetilde{S}$.
Proof: Since $\underset{\widetilde{S}}{S}$ is non-positive, by the Friedrichs characterization, a function in the domain of $\widetilde{S}$ is expressible as $f=(1-\widetilde{S})^{-1} F$ for some $F \in L^{2}(\mathbb{T})$. The $H^{1}$ completion of the domain $D_{S}$ is the $H^{1}$-closed subspace $\left.\operatorname{ker} \delta\right|_{H^{1}}$. The characterization $\left\langle(1-\widetilde{S})^{-1} F, g\right\rangle_{\operatorname{ker} \delta \mid H^{1}}=\langle F, g\rangle_{L^{2}}$ gives $\langle f, g\rangle_{H^{1}}=\langle(1-\widetilde{S}) f, g\rangle_{L^{2}}$. Then

$$
\langle(1-\widetilde{S}) f, g\rangle_{L^{2}}=\langle f, g\rangle_{H^{1}}=\langle f,(1-\Delta) g\rangle_{L^{2}}==\langle f,(1-\Delta) g\rangle_{H^{-\infty} \times H^{\infty}}=\langle(1-\Delta) f, g\rangle_{H^{-\infty} \times H^{\infty}}
$$

where we restrict the hermitian pairing on $L^{2} \times L^{2}$ to a hermitian pairing $L^{2} \times H^{\infty}$, and then extend it to $H^{-\infty} \times H^{\infty}$. Thus, $\widetilde{S}-\Delta=0$ on $D_{S}$. This does not quite imply that $\widetilde{S}-\Delta=0$ as a distribution, since $D_{S}$ is the kernel of $\delta$ on $C^{\infty}(\mathbb{T})$. Recall [13.14.5] that a continuous linear functional vanishing on the kernel of a (continuous) linear function is a scalar multiple of it. Thus, for $f$ in the domain of $\widetilde{S}$, there is a constant $c_{f}$ such that

$$
\langle(\widetilde{S}-\Delta) f, g\rangle_{H^{-\infty} \times H^{\infty}}=c_{f} \cdot \delta g \quad \quad\left(\text { for all } g \in C^{\infty}(\mathbb{T})\right)
$$

which gives $(\widetilde{S}-\Delta) f=c_{f} \cdot \delta$. Thus, $\widetilde{S} f=u$ if and only if $\Delta f+c_{f} \delta=u$, which is $\Delta f=u-c_{f} \delta$.

On the other hand, if $\Delta f=u+c \cdot \delta$ and $\delta f=0$, then $(1-\Delta) f=f-u-c \delta$, and the Friedrichs characterization of $f=(1-\widetilde{S})^{-1}(f-u)$ is satisfied: essentially running the earlier computation in reverse, for $g \in D_{S}$,

$$
\begin{aligned}
&\langle f, g\rangle_{H^{1}}=\langle f,(1-\Delta) g\rangle_{L^{2}}=\langle f,(1-\Delta) g\rangle_{H^{-\infty} \times H^{\infty}}=\langle(1-\Delta) f, g\rangle_{H^{-\infty} \times H^{\infty}}=\langle f-u-c \delta, g\rangle_{H^{-\infty} \times H^{\infty}} \\
&=\langle f-u, g\rangle_{H^{-\infty} \times H^{\infty}}-c\langle\delta, g\rangle_{H^{-\infty} \times H^{\infty}}=\langle f-u, g\rangle_{H^{-\infty} \times H^{\infty}}+0=\langle f-u, g\rangle_{L^{2}}
\end{aligned}
$$

The condition $\delta f \underset{\widetilde{S}}{=} 0$ is necessary for $f$ to be in the $H^{1}$-closure of the original domain $D_{S}$. If $\delta f \neq 0$, then $f$ cannot be $(1-\widetilde{S})^{-1} F$ for any $F \in L^{2}(\mathbb{T})$.
///
Explicitly, while the $\Delta$-eigenfunctions $u_{n}(x)=\sin n x$ for $n=0,1,2, \ldots$ satisfy $\delta u_{n}=0$, the $\Delta$ eigenfunctions $\cos n x$ for $n>0$ do not, so cannot be $\widetilde{S}$-eigenfunctions. In effect, they are replaced by exotic $\widetilde{S}$-eigenfunctions, essentially $\sin (n x / 2)$ with $n=1,3,5,7, \ldots$ made $2 \pi$-periodic by force. That is, $v_{n}(x)=\sin (n x / 2)$ for $x \in[0,2 \pi]$, and made $2 \pi$-periodic. That is, given $x \in \mathbb{R}$, for $\ell \in \mathbb{Z}$ such that $x-\ell \cdot 2 \pi \in[0,2 \pi)$, put $[x]=x-\ell \cdot 2 \pi$, and $v_{n}(x)=\sin (n[x] / 2)$. These are mildly exotic since the forced $2 \pi$-periodicity gives their graphs corners at $2 \pi \mathbb{Z}$ on $\mathbb{R} / 2 \pi \mathbb{Z}$, and

$$
\Delta v_{n}=-\left(\frac{n}{2}\right)^{2} \cdot v_{n}+n \cdot \delta
$$

The differential equation $(\Delta-\lambda) u=\delta$ can be solved by division, using Fourier expansions of distributions on $\mathbb{T}$ : let $u=\frac{1}{2 \pi} \sum_{n} c_{n} \psi_{n}$,

$$
\frac{1}{2 \pi} \sum_{n}\left(-n^{2}-\lambda\right) \psi_{n}=(\Delta-\lambda) u=\delta=\frac{1}{2 \pi} \sum_{n} \psi_{n} \quad \quad\left(\text { convergent in } H^{-1}(\mathbb{T})\right)
$$

Thus,

$$
u=\sum_{n} \frac{\psi_{n}}{-n^{2}-\lambda} \quad\left(\text { convergent in } H^{1}(\mathbb{T})\right)
$$

For $\lambda \neq-n^{2}$ this has a unique solution $u \in H^{1}$. The condition $\delta u=0$ is

$$
0=\delta u=\frac{1}{2 \pi} \sum_{n} \frac{\delta \psi_{n}}{-n^{2}-\lambda}=\frac{1}{2 \pi} \sum_{n} \frac{1}{-n^{2}-\lambda}
$$

It is perhaps not obvious that this has solutions exactly for $\lambda=-(n / 2)^{2}$ for $n=1,3,5, \ldots$.

### 9.7 Example: usual Sobolev spaces on $\mathbb{R}$

In contrast to $\Delta$ on $\mathbb{T}$, there are no square-integrable $\Delta$-eigenfunctions on $\mathbb{R}$ : the eigenfunction condition $u^{\prime \prime}=\lambda \cdot u$ is an explicitly solvable constant-coefficient differential equation, all whose solutions are linear combinations of $e^{ \pm \sqrt{\lambda} \cdot x}$, and none of these is square-integrable on $\mathbb{R}$. There cannot be an orthogonal basis for $L^{2}(\mathbb{R})$ consisting of $\Delta$-eigenfunctions, although Fourier transform and inversion

$$
\widehat{f}(\xi)=\int_{\mathbb{R}} e^{-i \xi x} f(x) d x \quad f(x)=\frac{1}{2 \pi} \int_{\mathbb{R}} e^{i \xi x} \widehat{f}(\xi) d \xi \quad(\text { for } f \in \mathscr{S}(\mathbb{R}))
$$

do express functions as superpositions of $\Delta$-eigenfunctions. For $0 \leq k \in \mathbb{Z}$, the usual Sobolev spaces are completions of $C_{c}^{\infty}(\mathbb{R})$ with respect to the the $k^{t h}$ Sobolev norm given by

$$
|f|_{H^{k}}^{2}=\int_{\mathbb{R}}(1-\Delta)^{k} f(x) \cdot \bar{f}(x) d x \quad\left(\text { for } f \in C_{c}^{\infty}(\mathbb{R})\right)
$$

By the Plancherel theorem for the Fourier transform, this is

$$
|f|_{H^{k}}^{2}=\int_{\mathbb{R}}(1-\Delta)^{k} f(x) \cdot \bar{f}(x) d x=\frac{1}{2 \pi} \int_{\mathbb{R}}\left(1+\xi^{2}\right)^{k}|\widehat{f}(\xi)|^{2} d \xi \quad \quad\left(\text { for } f \in C_{c}^{\infty}(\mathbb{R})\right)
$$

This also gives an $s^{t h}$ Sobolev norm for all real $s$ :

$$
|f|_{H^{s}}^{2}=\frac{1}{2 \pi} \int_{\mathbb{R}}\left(1+\xi^{2}\right)^{s}|\widehat{f}(\xi)|^{2} d \xi \quad \quad\left(\text { for } s \in \mathbb{R}, \text { for } f \in C_{c}^{\infty}(\mathbb{R})\right)
$$

The same sort of arguments as for $\mathbb{T}$ prove Sobolev imbedding here:
[9.7.1] Claim: $H^{s}(\mathbb{R}) \subset C_{c}^{k}(\mathbb{R})$ for any $s>\frac{1}{2}+k$, because the semi-norm $\nu_{k}(f)=\sup _{0 \leq i \leq k} \sup _{x \in \mathbb{R}}\left|f^{(i)}(x)\right|$ is dominated by the $H^{s}$-norm for $s>k+\frac{1}{2}$.
Proof: The semi-norm comparison implies that the $H^{s}$-completion of $C_{c}^{\infty}(\mathbb{R})$ is contained in the $\nu_{k^{-}}$ completion, which consists of $C^{k}$-functions whose $k$ derivatives all vanish at $\infty$. First, the case $k=0$ illustrates the key idea. By Cauchy-Schwarz-Bunyakowsky, for $f \in C_{c}^{\infty}(\mathbb{R})$, by Fourier inversion,

$$
\begin{gathered}
\sup _{x \in \mathbb{R}}|f(x)|=\sup _{x \in \mathbb{R}}\left|\frac{1}{2 \pi} \int_{\mathbb{R}} \widehat{f}(\xi) \cdot e^{i \xi x} d \xi\right| \leq \frac{1}{2 \pi} \int_{\mathbb{R}}|\widehat{f}(\xi)| d \xi=\frac{1}{2 \pi} \int_{\mathbb{R}}|\widehat{f}(\xi)| \cdot\left(1+\xi^{2}\right)^{s / 2} \cdot \frac{1}{\left(1+\xi^{2}\right)^{s / 2}} d \xi \\
\leq\left(\frac{1}{2 \pi} \int_{\mathbb{R}}|\widehat{f}(\xi)|^{2} \cdot\left(1+\xi^{2}\right)^{s} d \xi\right)^{\frac{1}{2}} \cdot\left(\frac{1}{2 \pi} \int_{\mathbb{R}} \frac{1}{\left(1+\xi^{2}\right)^{s}} d \xi\right)^{\frac{1}{2}}<_{s}|f|_{H^{s}}
\end{gathered}
$$

since for any $s>\frac{1}{2}$ the last integral is finite. For $k \geq 0$, use Gelfand-Pettis corollaries [14.3] to justify moving the differentiation inside the integral:

$$
\begin{gathered}
\sup _{x \in \mathbb{R}}\left|f^{(k)}(x)\right|=\sup _{x \in \mathbb{R}}\left|\frac{\partial^{k}}{\partial x^{k}} \frac{1}{2 \pi} \int_{\mathbb{R}} \widehat{f}(\xi) \cdot e^{i \xi x} d \xi\right|=\sup _{x \in \mathbb{R}}\left|\frac{1}{2 \pi} \int_{\mathbb{R}} \frac{\partial^{k}}{\partial x^{k}} \widehat{f}(\xi) \cdot e^{i \xi x} d \xi\right| \\
=\sup _{x \in \mathbb{R}}\left|\frac{1}{2 \pi} \int_{\mathbb{R}}(i \xi)^{k} \widehat{f}(\xi) \cdot e^{i \xi x} d \xi\right| \leq \frac{1}{2 \pi} \int_{\mathbb{R}}\left(1+\xi^{2}\right)^{k / 2} \cdot|\widehat{f}(\xi)| d \xi \\
=\frac{1}{2 \pi} \int_{\mathbb{R}}|\widehat{f}(\xi)| \cdot\left(1+\xi^{2}\right)^{s / 2} \cdot \frac{1}{\left(1+\xi^{2}\right)^{(s-k) / 2}} d \xi \\
\left.\leq\left(\frac{1}{2 \pi} \int_{\mathbb{R}}|\widehat{f}(\xi)|^{2} \cdot\left(1+\xi^{2}\right)^{s} d \xi\right)^{\frac{1}{2}} \cdot\left(\frac{1}{2 \pi} \int_{\mathbb{R}} \frac{1}{\left(1+\xi^{2}\right)^{s-k}} d \xi\right)^{\frac{1}{2}} \ll s s|f|_{H^{s}} \quad \quad \text { (for } s-k>\frac{1}{2}\right)
\end{gathered}
$$

again because the latter integral is convergent for $s-k>\frac{1}{2}$.
However, there is no Rellich compactness here: the inclusion $H^{1}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$ is not compact, and the Friedrichs extension $\widetilde{\Delta}$ of the restriction of $\Delta$ to $\mathscr{S}(\mathbb{R})$ does not have compact resolvent. The non-compactness of $H^{1}(\mathbb{R}) \subset L^{2}(\mathbb{R})$ follows easily from the spectral characterization, together with Plancherel for Fourier transform on $L^{2}(\mathbb{R})$ : letting $V^{1}$ be the image of $H^{1}(\mathbb{R})$ under Fourier transform, that is, the Hilbert space of $f$ on $\mathbb{R}$ such that $\int_{\mathbb{R}}|f(x)|^{2} \cdot\left(1+x^{2}\right) d x<\infty$, we have a commutative diagram, where vertical maps are isometries given by Fourier transform or inversion,

where the map on the bottom left is multiplication by $1 /\left(1+x^{2}\right)$. If the inclusion $H^{1}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$ were compact, then the composition


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would be compact, so the multiplication operator $L^{2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$ by multiplication by $1 /\left(1+x^{2}\right)$ would be compact. This operator is continuous and self-adjoint, but has no eigenvectors, and has continuous spectrum $[0,1]$, so cannot be compact.

### 9.8 Example: discrete spectrum of $-\Delta+x^{2}$ on $L^{2}(\mathbb{R})$

To obtain an operator on $\mathbb{R}$ related to $\Delta$ but with compact resolvent, add the confining potential $x^{2}$, construed as a multiplication operator, obtaining a Schrödinger operator

$$
S=-\Delta+x^{2}=-\frac{d^{2}}{d x^{2}}+x^{2} \quad S f(x)=-f^{\prime \prime}(x)+x^{2} \cdot f(x)
$$

This operator is also called the quantum harmonic oscillator. We will see that the resolvent $\widetilde{S}^{-1}$ of the Friedrichs extension $\widetilde{S}$ of $S$ is compact, so $\widetilde{S}$ has entirely discrete spectrum.

The eigenfunctions for $S$ are somewhat less well-known than those for $\Delta$, the latter easy to obtain from solving the constant-coefficient equation $u^{\prime \prime}=\lambda u$. The standard device to obtain eigenfunctions is as follows. The relevant Dirac operator here ${ }^{[63]}$ is

$$
\mathbb{D}=i \frac{\partial}{\partial x} \quad \text { so that } \quad \mathbb{D}^{2}=-\Delta
$$

the factorization

$$
-\Delta+x^{2}=(\mathbb{D}-i x)(\mathbb{D}+i x)+[i x, \mathbb{D}]=(\mathbb{D}-i x)(\mathbb{D}+i x)+1 \quad(\text { with }[i x, \mathbb{D}]=i x \circ \mathbb{D}-\mathbb{D} \circ i x)
$$

allows determination of many $S$-eigenfunctions, although proof that all are produced requires some effort. The eigenfunctions will be smooth, but not compactly supported, so it is not optimal to declare the natural domain of the operator to be $C_{c}^{\infty}(\mathbb{R})$. Instead, we take the Schwartz functions $\mathscr{S}(\mathbb{R})$, as in [13.7], to be the domain.

The raising and lowering operators are

$$
R=\text { raising }=\mathbb{D}-i x \quad L=\text { lowering }=\mathbb{D}+i x
$$

[9.8.1] Claim: The operator $S=-\Delta+x^{2}$ satisfies $S \geq 1$, in the sense that $\langle S f, f\rangle \geq\langle f, f\rangle$ for $f \in \mathscr{S}(\mathbb{R})$. Proof: This follows from the Dirac factorization:

$$
\begin{gathered}
\langle S f, f\rangle=\langle((\mathbb{D}-i x)(\mathbb{D}+i x)+1) f, f\rangle=\langle(\mathbb{D}-i x)(\mathbb{D}+i x) f, f\rangle+\langle f, f\rangle \\
=\langle(\mathbb{D}+i x) f,(\mathbb{D}+i x) f\rangle+\langle f, f\rangle \geq\langle f, f\rangle
\end{gathered}
$$

from the integration-by-parts fact $(\mathbb{D}-i x)^{*}=\mathbb{D}+i x$ on Schwartz functions.
Rather than attempting a direct solution of the differential equation $S u=\lambda u$, special features are exploited. First, a smooth function $u$ annihilated by $\mathbb{D}+i x$ will be an eigenfunction for $S$ with eigenvalue 1:

$$
S u=((\mathbb{D}-i x)(\mathbb{D}+i x)+1) u=(\mathbb{D}-i x) 0+u=1 \cdot u \quad(\text { for }(\mathbb{D}+i x) u=0)
$$

Dividing through by $i$, the equation $(\mathbb{D}+i x) u=0$ is

$$
\left(\frac{\partial}{\partial x}+x\right) u=0
$$

[63] Conveniently, the Dirac operator in this situation has complex coefficients. In two dimensions, Dirac operators have Hamiltonian quaternion coefficients, a special case of the general situation, that Dirac operators have coefficients in Clifford algebras.

## 9. Unbounded operators on Hilbert spaces

That is, $u^{\prime}=-x u$ or $u^{\prime} / u=-x$, so $\log u=-x^{2} / 2+C$ for arbitrary constant $C$. With $C=1$

$$
u(x)=e^{-x^{2} / 2}
$$

Conveniently, this is in $L^{2}(\mathbb{R})$, and in fact is in the Schwartz space [13.7] $\mathscr{S}(\mathbb{R})$ on $\mathbb{R}$. The alternative factorization

$$
S=-\Delta+x^{2}=(\mathbb{D}+i x)(\mathbb{D}-i x)-[i x, \mathbb{D}]=(\mathbb{D}+i x)(\mathbb{D}-i x)-1
$$

does also lead to an eigenfunction $u(x)=e^{x^{2} / 2}$, but this grows too fast. It is unreasonable to expect such luck in general, but here the raising and lower operators map $S$-eigenfunctions to other eigenfunctions: for $S u=\lambda u$, noting that $S=R L+1=L R-1$,

$$
\begin{aligned}
S(R u) & =(R L+1)(R u)=R L R u+R u=R(L R) u+R u=R(L R-1) u+2 R u \\
& =R S u+2 R u=R \lambda u+2 R u=(\lambda+2) \cdot R u \quad(\text { for } S u=\lambda u)
\end{aligned}
$$

Similarly, $S(L u)=(\lambda-2) \cdot L u$. Many eigenfunctions are produced by application of $R^{n}$ to $u_{1}(x)=e^{-x^{2} / 2}$ :

$$
R^{n} e^{-x^{2} / 2}=(2 n+1)-\text { eigenfunction for }-\Delta+x^{2}
$$

Repeated application of $R$ to $e^{-x^{2} / 2}$ produces polynomial [64] multiples of $e^{-x^{2} / 2}$

$$
R^{n} e^{-x^{2} / 2}=H_{n}(x) \cdot e^{-x^{2} / 2} \quad\left(\text { with polynomial } H_{n}(x) \text { of degree } n\right)
$$

The commutation relation shows that application of $L R^{n} u$ is just a multiple of $R^{n-1} u$, so application of $L$ to the eigenfunctions $R^{n} u$ produces nothing new.

We can almost prove that the functions $R^{n} u$ are all the square-integrable eigenfunctions. We saw above that $\left\langle\left(-\Delta+x^{2}\right) f, f\right\rangle \geq|f|_{L^{2}}^{2}$, so an $L^{2}$ eigenfunction has real eigenvalue $\lambda \geq 1$. Granting that repeated application of $L$ to a $\lambda$-eigenfunction $u$ stays in $L^{2}(\mathbb{R})$, the function $L^{n} u$ has eigenvalue $\lambda-2 n$, and the requirement $\lambda-2 n \geq 1$ on $L^{2}(\mathbb{R})$ implies that $L^{n} u=0$ for some $n$. Then $L\left(L^{n-1} u\right)=0$, but we already have shown that the only $L^{2}(\mathbb{R})$-function in the kernel of $L$ is $u_{1}(x)=e^{-x^{2} / 2}$.

To make this discussion a proof requires some preparation, since in general a Friedrichs extension can have eigenvectors outside the original domain, as in [9.6].
[9.8.2] Sobolev norms associated to the Schrödinger operator A Friedrichs extension $\widetilde{S}$ requires specification of a domain for $S$. The space $C_{c}^{\infty}(\mathbb{R})$ of test functions is universally reasonable, but we have already seen the not-compactly-supported eigenfunctions for the differential operator $S$. Happily, those eigenfunctions are in the Schwartz space $\mathscr{S}(\mathbb{R})$, confirming specification of $\mathscr{S}(\mathbb{R})$ as the domain of $S$.

There is a hierarchy of Sobolev-like norms

$$
|f|_{\mathfrak{B} \ell}=\left\langle\left(-\Delta+x^{2}\right)^{\ell} f, f\right\rangle_{L^{2}(\mathbb{R})}^{\frac{1}{2}} \quad(\text { for } f \in \mathscr{S}(\mathbb{R}))
$$

with corresponding Hilbert-space completions

$$
\mathfrak{B}^{\ell}=\text { completion of } \mathscr{S}(\mathbb{R}) \text { with respect to }|f|_{\mathfrak{B}^{\ell}}
$$

and $\mathfrak{B}^{0}=L^{2}(\mathbb{R})$. The Friedrichs extension $\widetilde{S}$ is characterized via its resolvent $\widetilde{S}^{-1}$, the resolvent characterized by

$$
\left\langle\widetilde{S}^{-1} f, S g\right\rangle=\langle f, g\rangle \quad\left(\text { for } f \in L^{2}(\mathbb{R}) \text { and } g \in \mathscr{S}(\mathbb{R})\right)
$$

[64] The polynomials $H_{n}$ are the Hermite polynomials, but everything needed about them can be proven from this spectral viewpoint.

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and $\widetilde{S}^{-1}$ maps $L^{2}(\mathbb{R})$ continuously to $\mathfrak{B}^{1}$. Thus, an eigenfunction $u$ for $\widetilde{S}$ is in $\mathfrak{B}^{\infty}=\bigcap_{\ell} \mathfrak{B}^{\ell}=\lim _{\ell} \mathfrak{B}^{\ell}$. We will see that

$$
\mathfrak{B}^{\infty}=\mathscr{S}(\mathbb{R})
$$

In particular, $\widetilde{S}$-eigenfunctions are in the presumed-natural domain $\mathscr{S}(\mathbb{R})$ of $S$, so evaluation of $\widetilde{S}$ on them is evaluation of $S$. Thus, $\widetilde{S}$-eigenfunctions are $S$-eigenfunctions. Further, repeated application of the lowering operator stabilizes $\mathscr{S}(\mathbb{R})$, so the near-proof above becomes a proof that all eigenfunctions in $L^{2}(\mathbb{R})$ are of the form $R^{n} e^{-x^{2} / 2}$.

To prove that these eigenfunctions are a Hilbert space basis for $L^{2}(\mathbb{R})$, we will prove that the resolvent is compact, so the eigenfunctions for the resolvent form an orthogonal Hilbert-space basis, and these are eigenfunctions for $\widetilde{S}$ itself, and then for $S$. That is, there is an orthogonal basis for $L^{2}(\mathbb{R})$ consisting of $S$-eigenfunctions, all obtained as

$$
(2 n+1)-\text { eigenfunction }=R^{n} e^{-x^{2} / 2}=\left(i \frac{\partial}{\partial x}-i x\right)^{n} e^{-x^{2} / 2}
$$

On $\mathbb{R}$, the compactness result depends on both smoothness and decay properties of the functions, in contrast to $\mathbb{T}$, where smoothness was the only issue.
[9.8.3] Theorem: (Rellich compactness) The injection $\mathfrak{B}^{\ell+1} \rightarrow \mathfrak{B}^{\ell}$ is compact.
Proof: The mechanism is well-illustrated by the $\ell=0$ case. We show compactness of $\mathfrak{B}^{1} \rightarrow L^{2}(\mathbb{R})$ by showing total boundedness [14.7.1] of the image of the unit ball. Let $\varphi$ be a smooth cut-off function, with

$$
\varphi_{N}(x)=\left\{\begin{array}{cl}
1 & (\text { for }|x| \leq N) \\
\text { smooth, between 0 and } 1 & (\text { for } N \leq|x| \leq N+1) \\
0 & (\text { for }|x| \geq N+1)
\end{array}\right.
$$

The derivatives of $\varphi_{N}$ in $N \leq|x| \leq N+1$ can easily be arranged to be independent of $N$. For $|f|_{1} \leq 1$, write $f=f_{1}+f_{2}$ with

$$
f_{1}=\varphi_{N} \cdot f \quad f_{2}=\left(1-\varphi_{N}\right) \cdot f
$$

The function $f_{1}$ on $[-N-1, N+1]$ can be considered as a function on a circle $\mathbb{T}$, by sticking $\pm(N+1)$ together. Then the Rellich compactness lemma on $\mathbb{T}[9.5 .6]$ shows that the image of the unit ball from $\mathfrak{B}^{1}$ is totally bounded in $L^{2}(\mathbb{T})$, which we can identify with $L^{2}[-N-1, N+1]$. The $L^{2}$ norm of the function $f_{2}$ is directly estimated

$$
\begin{gathered}
\left|f_{2}\right|_{L^{2}(\mathbb{R})}^{2}=\int_{|x| \geq N} \varphi_{N}^{2}(x) \cdot\left|f_{2}(x)\right|^{2} d x \leq \frac{1}{N^{2}} \int_{|x| \geq N}\left|f_{2}(x) \cdot x\right|^{2} d x \\
\leq \frac{1}{N^{2}} \int_{\mathbb{R}} x^{2} f(x) \cdot \bar{f}(x) d x \leq \frac{1}{N^{2}} \int_{\mathbb{R}}\left(-\frac{d^{2}}{d x^{2}}+x^{2}\right) f(x) \cdot \bar{f}(x) d x=\frac{1}{N^{2}}|f|_{1}^{2} \leq \frac{1}{N^{2}}
\end{gathered}
$$

Thus, given $\varepsilon>0$, for $N$ large the tail $f_{2}$ lies within a single $\varepsilon$-ball in $L^{2}(\mathbb{R})$. This proves total boundedness of the image of the unit ball, and compactness.
[9.8.4] Corollary: The Friedrichs extension $\widetilde{S}$ of $S=-\frac{d^{2}}{d x^{2}}+x^{2}$ has compact resolvent.
Proof: The map $\widetilde{S}^{-1}$ of $L^{2}(\mathbb{R})$ to itself is compact because it is the composition of the continuous map $\widetilde{S}^{-1}: L^{2}(\mathbb{R}) \rightarrow \mathfrak{B}^{1}$ and the compact inclusion $\mathfrak{B}^{1} \rightarrow \mathfrak{B}^{0}=L^{2}(\mathbb{R})$.
[9.8.5] Corollary: The spectrum of $S=-\frac{d^{2}}{d x^{2}}+x^{2}$ is discrete. There is an orthonormal basis of $L^{2}(\mathbb{R})$ consisting of eigenfunctions for $\widetilde{S}$.
Proof: Self-adjoint compact operators have discrete spectrum with finite multiplicities for non-zero eigenvalues. From above, these eigenfunctions are exactly the eigenfunctions for the Friedrichs extension $\widetilde{S}$. Since these eigenfunctions give an orthogonal Hilbert-space basis, $\widetilde{S}$ has no further spectrum.

## 9. Unbounded operators on Hilbert spaces

It remains to show that the eigenfunctions are in $\mathscr{S}(\mathbb{R})$, to know that they are eigenfunctions of $S$ itself, rather than only of the extension $\widetilde{S}$.
[9.8.6] Theorem: $\mathfrak{B}^{\infty}=\mathscr{S}(\mathbb{R})$
Proof: It is clear that $\mathscr{S}(\mathbb{R}) \subset \mathfrak{B}^{\infty}$. The issue is the other containment. The Weyl algebra $A=A_{1}$ of operators, generated over $\mathbb{C}$ by the multiplication $x$ and derivative $\partial=d / d x$, is also generated by $R=i \partial-i x$ and $L=i \partial+i x$. The Weyl algebra is filtered by degree in $R$ and $L$ : let $A \leq n$ be the $C$-subspace of $A$ spanned by all non-commuting monomials in $R, L$ of total degree at most $n$, with $A^{\leq 0}=\mathbb{C}$. Note that $R$ and $L$ commute modulo $A^{\leq 0}$ : as operators, $\partial \circ x=1+x \circ \partial$, and the commutation relation is obtained again, by

$$
\begin{gathered}
{[R, L]=R L-L R=(i \partial-i x)(i \partial+i x)-(i \partial+i x)(i \partial-i x)=-(\partial-x)(\partial+x)+(\partial+x)(\partial-x)} \\
=-\left(\partial^{2}-x \partial+\partial x-x^{2}\right)+\left(\partial^{2}+x \partial-\partial x-x^{2}\right)=2(x \partial-\partial x)=-2
\end{gathered}
$$

[9.8.7] Claim: For a monomial $w_{2 n}$ in $R$ and $L$ of degree $2 n$,

$$
\left|\left\langle w_{2 n} \cdot f, f\right\rangle_{L^{2}(\mathbb{R})}\right| \ll n_{n}|f|_{\mathfrak{B}^{n}}^{2} \quad\left(\text { for } f \in C_{c}^{\infty}(\mathbb{R})\right)
$$

Proof: Induction. First,

$$
\langle R L f, f\rangle=\langle(R L+1) f, f\rangle-\langle f, f\rangle \leq\langle(R L+1) f, f\rangle=\langle S f, f\rangle=|f|_{\mathfrak{B}^{1}}^{2}
$$

A similar argument applies to $L R$. For the length-two word $L^{2}$,

$$
\begin{aligned}
\left|\left\langle L^{2} f, f\right\rangle\right| & =|\langle L f, R f\rangle| \leq|L f| \cdot|R f|=\langle L f, L f\rangle^{\frac{1}{2}} \cdot\langle R f, R f\rangle^{\frac{1}{2}} \\
& =\langle R L f, f\rangle^{\frac{1}{2}} \cdot\langle L R f, f\rangle^{\frac{1}{2}} \leq\langle S f, f\rangle=|f|_{\mathfrak{B}^{1}}^{2}
\end{aligned}
$$

A similar argument applies to $R^{2}$, completing the argument for $n=1$. For the induction step, any word $w_{2 n}$ of length $2 n$ is equal to $R^{a} L^{b} \bmod A^{\leq 2 n-2}$ for some $a+b=2 n$, so, by induction,

$$
\left|\left\langle w_{2 n} f, f\right\rangle\right|=\left|\left\langle R^{a} L^{b} f, f\right\rangle\right|+|f|_{\mathfrak{B}^{n-1}}^{2}
$$

In the case that $a \geq 1$ and $b \geq 1$, by induction

$$
\left|\left\langle R^{a} L^{b} f, f\right\rangle\right|=\left|\left\langle R^{a-1} L^{b-1}(L f), L f\right\rangle\right|<_{n}=|L f|_{\mathfrak{B}^{n-1}}^{2}=\left\langle S^{n-1} L f, L f\right\rangle=\left\langle R S^{n-1} L f, f\right\rangle
$$

Since $R S^{n-1} L$ is $S^{n} \bmod A^{\leq 2 n-2}$, by induction

$$
\left\langle R S^{n-1} L f, f\right\rangle \lll{ }_{n}\left\langle S^{n} f, f\right\rangle+|f|_{\mathfrak{B}^{n-1}}^{2}=|f|_{\mathfrak{B}^{n}}^{2}+|f|_{\mathfrak{B}^{n-1}}^{2} \ll|f|_{\mathfrak{B}^{n}}^{2}
$$

In the extreme case $a=0$,

$$
\left\langle L^{2 n} f, f\right\rangle=\left\langle L^{n} f, R^{n} f\right\rangle \leq\left|L^{n} f\right| \cdot\left|R^{n} f\right|=\left\langle L^{n} f, L^{n} f\right\rangle^{\frac{1}{2}} \cdot\left\langle R^{n} f, R^{n} f\right\rangle^{\frac{1}{2}}=\left\langle R^{n} L^{n} f, f\right\rangle^{\frac{1}{2}} \cdot\left\langle L^{n} R^{n} f, f\right\rangle^{\frac{1}{2}}
$$

which brings us back to the previous case. The extreme case $b=0$ is similar.
[9.8.8] Corollary: For a monomial $w_{n}$ in $R$ and $L$ of degree $n$,

$$
\left|\left\langle w_{n} \cdot f, f\right\rangle_{L^{2}(\mathbb{R})}\right| \ll{ }_{n}|f|_{\mathfrak{B}^{n}} \cdot|f|_{L^{2}} \quad\left(\text { for } f \in C_{c}^{\infty}(\mathbb{R})\right)
$$

Proof: By Cauchy-Schwarz-Bunyakowsky and the claim,

$$
\left|\left\langle w_{n} \cdot f, f\right\rangle_{L^{2}}\right| \leq\left|w_{n} f\right|_{L^{2}} \cdot|f|_{L^{2}}=\left\langle w_{n}^{*} w_{n} f, f\right\rangle^{\frac{1}{2}} \cdot|f|_{L^{2}} \leq|f|_{\mathfrak{B}^{n}} \cdot|f|_{L^{2}}
$$

as claimed.

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[9.8.9] Corollary: The limit $\mathfrak{B}^{\infty}=\lim _{k} \mathfrak{B}^{k}$ is contained in $\mathscr{S}(\mathbb{R})$.
Proof: This is the same idea as in [9.7.1] and its proof. Use the density of test functions in $\mathscr{S}(\mathbb{R})$, whose completion in the $\mathfrak{B}^{k}$ norm gives $\mathfrak{B}^{k}$, for every $k$. Thus, test functions are dense in $\mathfrak{B}$. For $f \in C_{c}^{\infty}(\mathbb{R})$, by Fourier inversion,

$$
\begin{gathered}
\sup _{x \in \mathbb{R}}\left|\left(1+x^{2}\right)^{m} f^{(n)}(x)\right|=\sup _{x \in \mathbb{R}}\left|\frac{1}{2 \pi} \int_{\mathbb{R}}\left(\left(1+x^{2}\right)^{m} f^{(n)}\right)^{\wedge}(\xi) \cdot e^{i \xi x} d \xi\right| \\
=\sup _{x \in \mathbb{R}}\left|\frac{1}{2 \pi} \int_{\mathbb{R}}\left((1-\Delta)^{m} \cdot(-i \xi)^{n} \cdot \widehat{f}(\xi)\right) \cdot e^{i \xi x} d \xi\right| \leq \frac{1}{2 \pi} \int_{\mathbb{R}}\left|(1-\Delta)^{m} \cdot \xi^{n} \cdot \widehat{f}(\xi)\right| d \xi \\
=\frac{1}{2 \pi} \int_{\mathbb{R}}\left|(1-\Delta)^{m} \cdot \xi^{n} \cdot \widehat{f}(\xi)\right|\left(1+\xi^{2}\right)^{s / 2} \cdot \frac{1}{\left(1+\xi^{2}\right)^{s / 2}} d \xi \\
\leq \frac{1}{2 \pi}\left|\left(1+\xi^{2}\right)^{s / 2} \cdot(1-\Delta)^{m} \cdot \xi^{n} \cdot \widehat{f}(\xi)\right|_{L^{2}} \cdot\left|\frac{1}{\left(1+\xi^{2}\right)^{s / 2}}\right|_{L^{2}} \\
=\frac{1}{2 \pi}\left|(1-\Delta)^{s / 2} \cdot\left(1+x^{2}\right)^{m} \cdot f^{(n)}\right|_{L^{2}} \cdot\left|\frac{1}{\left(1+\xi^{2}\right)^{s / 2}}\right|_{L^{2}}
\end{gathered}
$$

taking $s \in 2 \mathbb{Z}$, by Plancherel and Cauchy-Schwarz-Bunyakowsky. The $L^{2}$-norm of $1 /\left(1+\xi^{2}\right)^{s / 2}$ is finite for large-enough $s$, and the $L^{2}$-norm of $(1-\Delta)^{s / 2}\left(1+x^{2}\right) f^{(n)}(x)$ is dominated by a finite linear combination of the seminorms $\mu_{w}(f)=|\langle w f, f\rangle|^{\frac{1}{2}}$.

We return to the proof that all eigenfunctions of $\widetilde{S}$ are in $\mathscr{S}(\mathbb{R})$, and, therefore, in the domain of the original operator $S$. An eigenfunction $u$ for $\widetilde{S}$ lies in $\mathfrak{B}^{1}$, by Friedrichs' construction. Friedrichs' extensions preserve semi-boundedness, so $\widetilde{S} \geq 1$, and the inverse $\widetilde{S}^{-1}$ exists as a bounded operator, and is self-adjoint. An eigenvector relation $\widetilde{S} u=\lambda \cdot u$ entails $\lambda \neq 0$, and gives $u=\lambda^{-1} \widetilde{S}^{-1} u$. Any $\widetilde{S}$-eigenfunction $u$ is in the domain of $\widetilde{S}$, inside $\mathfrak{B}^{1}$, by construction. By induction, $u \in \mathfrak{B}^{\infty}=\mathscr{S}(\mathbb{R})$.

That is, again, the heuristic above that appears to determine an orthogonal basis of eigenfunctions for $S$ does succeed: there is an orthogonal basis of eigenfunctions for $\widetilde{S}$, and we have shown that all eigenfunctions for $\widetilde{S}$ are actually in the domain of $S$.

Up to a constant, the $n^{\text {th }}$ Hermite polynomial $H_{n}(x)$ is characterized by

$$
H_{n}(x) \cdot e^{-x^{2} / 2}=R^{n} e^{-x^{2} / 2}=\left(i \frac{\partial}{\partial x}-i x\right)^{n} e^{-x^{2} / 2}
$$

The above discussion shows that $H_{0}, H_{1}, H_{2}, \ldots$ are orthogonal on $\mathbb{R}$ with respect to the weight $e^{-x^{2}}$, and give an orthogonal basis for the weighted $L^{2}$-space

$$
\left\{f: \int_{\mathbb{R}}|f(x)|^{2} \cdot e^{-x^{2}} d x<\infty\right\}
$$

### 9.9 Essential self-adjointness

The simple examples that exhibit non-symmetric adjoints to symmetric operators show that there can be many self-adjoint extensions, incomparable in the partial ordering on operators. While Friedrichs extensions are canonical positive self-adjoint extension $T_{\mathrm{Fr}} \supset T$ of a positive, symmetric, densely-defined operator $T$, we are interested in clarifying the conditions under which symmetric, densely-defined $T$ has a unique self-adjoint extension.

A symmetric, densely-defined operator $T$ with a unique self-adjoint extension is essentially self-adjoint. Although this use of essential is approximately compatible with the colloquial sense of the word, unfortunately there is some risk that its use in this context be mistaken for the more ambiguous colloquial sense.

In brief, for unbounded operators arising from differential operators, imposition of various boundary conditions often gives rise to mutually incomparable self-adjoint extensions. Thus, free-space situations, lacking boundary conditions, are the best candidates for essential self-adjointness.

## 9. Unbounded operators on Hilbert spaces

Since a self-adjoint operator is (graph-) closed, any self-adjoint extension of symmetric $T$ must extend the (graph-) closure $\bar{T}$.

As above, symmetric $T$ has no non-real complex eigenvalues $\lambda$, that is, that $T-\lambda$ is injective on $D_{T}$. This allows definition of an operator $U$ on the image $(T-\lambda) D_{T}$ by

$$
U=(T-\bar{\lambda})(T-\lambda)^{-1} \quad\left(\text { for } \lambda \notin \mathbb{R}, \text { on the image }(T-\lambda) D_{T}\right)
$$

[9.9.1] Claim: The operator $U$ defined on the image $(T-\lambda) D_{T}$ is unitary, in the sense that $\langle U v, U w\rangle=\langle v, w\rangle$ for $v, w$ in the domain of $U$.
Proof: For $v, w$ in the image $(T-\lambda) D_{T}$, let $v^{\prime}=(T-\lambda)^{-1} v$ and $w^{\prime}=(T-\lambda)^{-1} w$. Then

$$
\langle U v, U w\rangle=\left\langle(T-\bar{\lambda}) v^{\prime},(T-\bar{\lambda}) w^{\prime}\right\rangle
$$

while

$$
\langle v, w\rangle=\left\langle(T-\lambda) v^{\prime},(T-\lambda) w^{\prime}\right\rangle
$$

Thus, we want to show that

$$
\left\langle(T-\bar{\lambda}) v^{\prime},(T-\bar{\lambda}) w^{\prime}\right\rangle=\left\langle(T-\lambda) v^{\prime},(T-\lambda) w^{\prime}\right\rangle
$$

This follows from the symmetry of $T$.
[9.9.2] Theorem: For (graph-) closed, symmetric, densely-defined $T$, if for some non-real $\lambda$ both $(T-\lambda) D_{T}$ and $(T-\bar{\lambda}) D_{T}$ are dense, then $T$ is self-adjoint.
Proof: First, claim that for (graph-) closed and symmetric $T$, for non-real $\lambda$ the image $(T-\lambda) D_{T}$ is closed. To see this, let $(T-\lambda) v_{i}$ be Cauchy, with $v_{i}$ in the domain of $T$. By the unitariness of $U$, the sequence

$$
U\left((T-\lambda) v_{i}\right)=(T-\bar{\lambda}) v_{i}
$$

is also Cauchy. Subtracting one sequence from the other, $(\lambda-\bar{\lambda}) v_{i}$ is Cauchy. Since $\lambda \notin \mathbb{R}, v_{i}$ is Cauchy. Similarly, adding the two sequences, $(2 T+\lambda+\bar{\lambda}) v_{i}$ is Cauchy. Because $v_{i}$ is Cauchy, $(\lambda+\bar{\lambda}) v_{i}$ is Cauchy, so $2 T v_{i}$ and $T v_{i}$ are Cauchy. Since the graph of $T$ is closed, the sequence $v_{i} \oplus T v_{i}$ converges to some $v \oplus T v$ in the graph of $T$. Thus, $(T-\lambda) v_{i}$ certainly converges to $(T-\lambda) v$, and verifies the claim that $(T-\lambda) D_{T}$ is closed. By hypothesis, the closed subspaces $(T-\lambda) D_{T}$ and $(T-\lambda) D_{T}$ are also dense, so each is the whole space $V$.

Given $v$ in the domain $D_{T^{*}}$ of the adjoint $T^{*}$, we show that $v \in D_{T}$. Since $(T-\lambda) D_{T}=V$, there is $v^{\prime} \in D_{T}$ such that

$$
(T-\lambda) v^{\prime}=\left(T^{*}-\lambda\right) v
$$

Thus,

$$
\left\langle v^{\prime},\left(T^{*}-\bar{\lambda}\right) w\right\rangle=\left\langle(T-\lambda) v^{\prime}, w\right\rangle=\left\langle\left(T^{*}-\lambda\right) v, w\right\rangle=\langle v,(T-\bar{\lambda}) w\rangle \quad\left(\text { for all } w \in D_{T}\right)
$$

Since $(T-\bar{\lambda}) D_{T}$ is dense, $v^{\prime}=v$. That is, $v \in D_{T}$.
[9.9.3] Corollary: For symmetric, densely-defined $T$, suppose that for some non-real $\lambda$ both $(T-\lambda) D_{T}$ and $(T-\bar{\lambda}) D_{T}$ are dense. Then the closure $\bar{T}$ of $T$ is self-adjoint, and is the unique self-adjoint extension of $T$. In particular, $T$ is essentially self-adjoint.
Proof: The closure $\bar{T}$ extends $T$, and is symmetric for symmetric $T$. Certainly $(\bar{T}-\lambda) D_{\bar{T}}$ contains $(T-\lambda) D$, so when the latter is dense the former is dense. Thus, $\bar{T}$ meets the hypothesis of the theorem, and is self-adjoint.

Any self-adjoint extension $S=S^{*}$ of $T$ is closed, since adjoints are closed. Thus, any self-adjoint extension $S$ of $T$ contains the closure $\bar{T}=T^{* *}$, for topological reasons. Taking adjoints is inclusion-reversing, so, $S=S^{*} \subset\left(T^{* *}\right)^{*}=T^{*}$ from characterization of adjoints in terms of graphs as above. Therefore, $S=T^{* *}=\bar{T}$.
[9.9.4] Claim: For a symmetric, densely-defined operator $T$, density of $(T-\lambda) D_{T}$ is equivalent to the assertion that $T^{*}$ does not have eigenvalue $\bar{\lambda}$.
Proof: From earlier examples, $T^{*}$ need not be symmetric, so the natural argument that eigenvalues of symmetric operators must be real does not apply to $T^{*}$. Apart from that, the argument is the natural one. The density of $(T-\lambda) D_{T}$ implies that $\langle(T-\lambda) v, w\rangle=0$ for all $v \in D_{T}$ if and only if $w=0$. If $\left(T^{*}-\lambda\right) v=0$, then

$$
0=\left\langle\left(T^{*}-\bar{\lambda}\right) v, w\right\rangle=\langle v,(T-\lambda) w\rangle \quad\left(\text { for all } w \in D_{T}\right)
$$

Since $(T-\lambda) D_{T}$ is dense, this implies $w=0$. Conversely, if $(T-\lambda) D_{T}$ were not dense, then its closure would not be the whole space, so would be orthogonal to some $v \neq 0$. Then

$$
0=\langle v,(T-\lambda) w\rangle=\left\langle\left(T^{*}-\bar{\lambda}\right) v, w\right\rangle \quad\left(\text { for every } w \in D_{T}, \text { for } v \in D_{T^{*}}\right)
$$

Thus, since $D_{T}$ is dense, we imagine that it would be consistent to define $T^{*} v=\bar{\lambda} v$. In fact, by [9.1.1], the graph of $T^{*}$ is the orthogonal complement in $V \oplus V$ of the image of the graph of $S$ under the isometry $J$, so $0=\left\langle(T-\bar{\lambda}) D_{T}, v\right\rangle$ implies that the pair $(v, \lambda v)$ is in the graph of the adjoint. Thus, if $(T-\lambda) D_{T}$ were not dense, then $T^{*}$ would have eigenvalue $\bar{\lambda}$.

Thus, we have a variant form of the criterion for the closure of $T$ being self-adjoint:
[9.9.5] Corollary: For symmetric, densely-defined $T$, if for some non-real $\lambda$, neither $\lambda$ nor $\bar{\lambda}$ is an eigenvalue for the adjoint $T^{*}$, then the closure $\bar{T}$ of $T$ is the unique self-adjoint extension of $T$. In particular, $T$ is essentially self-adjoint.

In the situation of the corollary, since $\bar{T}=T^{* *}$, and $T^{* * *}=T^{*}$, in fact $\bar{T}^{*}=T^{*}$.

### 9.10 Example: essentially self-adjoint operator

In fact, [9.5.2] proved that the domain of the Friedrichs extension $\widetilde{S}$ of the restriction $S$ of $\Delta$ to $C^{\infty}(\mathbb{T})$ is $H^{2}(\mathbb{T})$, which says that the graph-closure of $S$ is its Friedrichs extension. Since every self-adjoint extension contains the graph-closure, and distinct self-adjoint extensions are not comparable, the graph-closure must be the only self-adjoint extension. That is, the restriction of $\Delta$ to $C^{\infty}(\mathbb{T})$ is essentially self-adjoint. Nevertheless, we want to practice application of the criterion above.

Let $S=\frac{d^{2}}{d x^{2}}$ on $\mathbb{T}=\mathbb{R} / 2 \pi \mathbb{Z}$ with domain $D_{S}=C^{\infty}(\mathbb{T})$. Since $S$ is non-positive, there is at least one meaningful self-adjoint extension, the Friedrichs extension. However, we want the (graph-) closure $\bar{S}$ of $S$ to be that self-adjoint extension, giving uniqueness in a strong, unambiguous fashion.

We do not directly characterize the domain $D_{S^{*}}$ of $S^{*}$, apart from the fact that it contains the domain of $S$. It is convenient that $S$ stabilizes $D_{S}$. The translation action of $\mathbb{T}$ on functions on $\mathbb{T}$ is $\left(R_{x} f\right)(y)=f(y+x)$. This action is unitary, and gives a (jointly) continuous map $\mathbb{T} \times L^{2}(\mathbb{T}) \longrightarrow L^{2}(\mathbb{T})$. A constant-coefficient differential operator such as $S$ commutes with the translation action, at least on $D_{S}$ : in symbols, $R_{t} \circ T=T \circ R_{t}$ for all $t \in \mathbb{T}$. Indeed, this invariance allows such operators to descend from $\mathbb{R}$ to the quotient $\mathbb{T}$. Certainly $D_{S}$ is stable under translation.
[9.10.1] Claim: The domain $D_{S^{*}}$ of $S^{*}$ is stable under translation.
Proof: Let $J(x \oplus y)=-y \oplus x$ be the usual map on $L^{2}(\mathbb{T}) \oplus L^{2}(\mathbb{T})$. The map $J$ is an isometry with respect to the usual inner product $\left\langle x+x^{\prime}, y+y^{\prime}\right\rangle=\langle x, y\rangle+\left\langle x^{\prime}, y^{\prime}\right\rangle$ on $L^{2}(\mathbb{T}) \oplus L^{2}(\mathbb{T})$. The graph of the adjoint is characterized as the orthogonal complement of the image by $J$ of the graph of $S$. Thus, for $y \oplus S^{*} y$ in the graph of $S^{*}$, for all $x \oplus S x$ in the graph of $S$, because $S$ commutes with $R_{t}$ on $D_{S}$,

$$
\begin{gathered}
\left\langle R_{t} y \oplus R_{t} S^{*} y, J(x \oplus S x)\right\rangle=\left\langle R_{t} y \oplus R_{t} S^{*} y,-S x \oplus x\right\rangle \\
=\left\langle y \oplus S^{*} y,-R_{t}^{-1} S x \oplus R_{t}^{-1} x\right\rangle=\left\langle y \oplus S^{*} y,-S R_{t}^{-1} x \oplus R_{t}^{-1} x\right\rangle=0
\end{gathered}
$$

because $R_{t}^{-1} x \in D_{S}$. Thus, $R_{t} y \in D_{S^{*}}$, as claimed.
The action of $\varphi \in C^{\infty}(\mathbb{T})$ on $L^{2}(\mathbb{T})$ is by the integral operator

$$
R_{\varphi} v=\int_{\mathbb{T}} \varphi(t) \cdot R_{t} v d t
$$

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Since $t \rightarrow \varphi(t) \cdot R_{t} v$ is a compactly-supported, continuous, $L^{2}(\mathbb{T})$-valued function on $\mathbb{T}$, it has a Gelfand-Pettis integral. Further, $D_{S}$ is stable under this action, and the translation action

$$
\mathbb{T} \times C^{\infty}(\mathbb{T}) \longrightarrow C^{\infty}(\mathbb{T})
$$

is continuous with respect to the Fréchet-space topology. For $\varphi \in C^{\infty}(\mathbb{T})$, the corresponding integral operator $R_{\varphi}$ maps $L^{2}(\mathbb{T})$ to $C^{\infty}(\mathbb{R})$, by smoothing of distributions [14.5].
[9.10.2] Claim: The operators $R_{\varphi}$ for $\varphi \in C^{\infty}(\mathbb{T})$ commute with $S^{*}$.
Proof: Since the operators $S$ on $D_{S}$ and $S^{*}$ on $D_{S^{*}}$ are not continuous on $L^{2}(\mathbb{T})$, the properties of GelfandPettis integrals must be used scrupulously. For $\varphi \in C^{\infty}(\mathbb{T}), v \in D_{S^{*}}$, and $w \in D_{S}$, using the commutativity of Gelfand-Pettis integrals with continuous maps, a sensible computation succeeds:

$$
\begin{gathered}
\left\langle R_{\varphi} S^{*} v, w\right\rangle=\left\langle\int_{\mathbb{T}} \varphi(t) R_{t} S^{*} v d t, w\right\rangle=\int_{\mathbb{T}}\left\langle\varphi(t) R_{t} S^{*} v, w\right\rangle d t=\int_{\mathbb{T}} \varphi(t)\left\langle R_{t} S^{*} v, w\right\rangle d t \\
\quad=\int_{\mathbb{T}} \varphi(t)\left\langle S^{*} v, R_{t}^{-1} w\right\rangle d t=\int_{\mathbb{T}} \varphi(t)\left\langle v, S R_{t}^{-1} w\right\rangle d t=\int_{\mathbb{T}} \varphi(t)\left\langle v, R_{t}^{-1} S w\right\rangle d t \\
=\int_{\mathbb{T}} \varphi(t)\left\langle R_{t} v, S w\right\rangle d t=\left\langle\int_{\mathbb{T}} \varphi(t) R_{t} v d t, S w\right\rangle=\left\langle R_{\varphi} v, S w\right\rangle=\left\langle S^{*} R_{\varphi} v, w\right\rangle
\end{gathered}
$$

This is the desired commutativity.
Now we can prove that $S^{*}$ has no non-real eigenvalues, so $S$ meets the hypotheses of the theorem of the previous section, and its closure $\bar{S}$ is the unique self-adjoint extension of $S$. Suppose $v \in D_{S^{*}}$ and $\left(S^{*}-\lambda\right) v=0$. Then, for any $\varphi \in C^{\infty}(\mathbb{T})$,

$$
0=R_{\varphi} \cdot 0=R_{\varphi}\left(S^{*}-\lambda\right) v=\left(S^{*}-\lambda\right) R_{\varphi} v=(S-\lambda) R_{\varphi} v
$$

the last equality because $R_{\varphi}$ maps everything to $C^{\infty}(\mathbb{T})$, on which $S^{*}$ acts by $S$. Although $S^{*}$ is not assured to be symmetric (unless $S^{*}=\bar{S}=S^{* *}$, which is the sought-after essential self-adjointness of $S$ itself!), the operator $S$ is symmetric, so has no non-real eigenvalues, giving $R_{\varphi} v=0$. For given $\varphi$, taking $\varphi$ sufficiently far along in an approximate identity gives $R_{\varphi} v \neq 0$ for $v \neq 0$. Thus, we conclude that $v=0$, and $S^{*}$ has no non-real eigenvalues.

Thus, the closure $\bar{S}$ of $S$ is the unique self-adjoint extension of $S$. Restricted to the graph of $S$, the metric on $L^{2}(\mathbb{T}) \oplus L^{2}(\mathbb{T})$ gives norm-squared

$$
|x \oplus S x|^{2}=|x|^{2}+|S x|^{2} \geq|x|^{2}+|\langle S x, x\rangle|+|S x|^{2}-\frac{1}{2} \cdot(|x|+|S x|)^{2} \geq \frac{1}{2} \cdot\left(|x|^{2}+\langle-S x, x\rangle+|S x|^{2}\right)
$$

since Cauchy-Schwarz-Bunyakowsky and $2 a b \leq(a+b)^{2}$ give

$$
|\langle S x, x\rangle| \leq|S x| \cdot|x| \leq \frac{1}{2} \cdot(|x|+|S x|)^{2}
$$

Thus, the completion $H^{2}(\mathbb{T})$ of $D_{S}$ with respect to the norm attached to the hermitian inner product

$$
\langle x, y\rangle+\langle-S x, y\rangle+\left\langle S^{2} x, y\right\rangle \asymp\langle x, y\rangle_{H^{2}} \quad\left(\text { for } x, y \in D_{S}\right)
$$

is exactly the domain of $\bar{S}$. In the $\langle,\rangle_{H^{2}}$ topology, $S$ is continuous on $D_{S}$, and $\bar{S}$ is the extension by continuity to $H^{2}(\mathbb{T})$.

## 9.A Appendix: compact operators

The spectrum $\sigma(T)$ of a continuous linear operator $T: X \rightarrow X$ on a Hilbert space $X$ is the collection of complex numbers $\lambda$ such that $T-\lambda$ does not have a continuous linear inverse. The discrete spectrum $\sigma_{\text {disc }}(T)$ is the collection of complex numbers $\lambda$ such that $T-\lambda$ fails to be injective. In other words, the discrete spectrum is the collection of eigenvalues. The continuous spectrum $\sigma_{\text {cont }}(T)$ is the collection of complex

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numbers $\lambda$ such that $T-\lambda \cdot 1_{X}$ is injective, does have dense image, but fails to be surjective. The residual spectrum $\sigma_{\text {res }}(T)$ is everything else: neither discrete nor continuous spectrum. That is, the residual spectrum of $T$ is the collection of complex numbers $\lambda$ such that $T-\lambda \cdot 1_{X}$ is injective, and fails to have dense image (so is certainly not surjective).

To see that there are no other possibilities for failure of existence of an inverse, note that the closed graph theorem [9.B.3] implies that a bijective, continuous, linear map $T: X \rightarrow Y$ of Banach spaces has continuous inverse. Indeed, granting that the inverse exists as a linear map, its graph is

$$
\text { graph of } T^{-1}=\{(y, x) \in Y \times X:(x, y) \text { in the graph of } T \subset X \times Y\}
$$

Since the graph of $T$ is closed, the graph of $T^{-1}$ is closed, and by the closed graph theorem $T^{-1}$ is continuous.
As usual, the adjoint $T^{*}$ of a continuous linear map $T: X \rightarrow Y$ from one Hilbert space is defined by

$$
\langle T x, y\rangle_{Y}=\left\langle x, T^{*} y\right\rangle_{Y}
$$

[9.A.1] Claim: An (bounded) normal operator $T: X \rightarrow X$, that is, with $T T^{*}=T^{*} T$, has empty residual spectrum. That is, for $\lambda$ not an eigenvalue, $T-\lambda$ has dense image.
Proof: The adjoint of $T-\lambda$ is $T^{*}-\bar{\lambda}$, so consider $\lambda=0$ to lighten the notation. Suppose that $T$ does not have dense image. Then there is non-zero $z$ such that

$$
\left.0=\langle z, T x\rangle=\left\langle T^{*} z, x\right\rangle \quad \text { (for every } x \in X\right)
$$

Therefore $T^{*} z=0$, and the 0-eigenspace $Z$ of $T^{*}$ is non-zero. Since $T^{*}(T z)=T\left(T^{*} z\right)=T(0)=0$ for $z \in Z$, $T^{*}$ stabilizes $Z$. That is, $Z$ is both $T$ and $T^{*}$-stable. Therefore, $T=\left(T^{*}\right)^{*}$ acts on $Z$ by (the complex conjugate of) 0 , and $T$ has non-trivial 0 -eigenvectors, contradiction.

A set in a topological space is pre-compact when its closure is compact. A linear operator $T: X \rightarrow Y$ on Hilbert spaces is compact when it maps the unit ball in $X$ to a pre-compact set in $Y$. Equivalently, $T$ is compact if and only if it maps bounded sequences in $X$ to sequences in $Y$ with convergent subsequences.
[9.A.2] Claim: An operator-norm limit of compact operators is compact. A compact operator $T: X \rightarrow Y$ with $Y$ a Hilbert space is an operator norm limit of finite rank operators.
Proof: Let $T_{n} \rightarrow T$ in uniform operator norm, with compact $T_{n}$. Given $\varepsilon>0$, let $n$ be sufficiently large such that $\left|T_{n}-T\right|<\varepsilon / 2$. Since $T_{n}(B)$ is pre-compact, there are finitely many $y_{1}, \ldots, y_{t}$ such that for any $x \in B$ there is $i$ such that $\left|T_{n} x-y_{i}\right|<\varepsilon / 2$. By the triangle inequality

$$
\left|T x-y_{i}\right| \leq\left|T x-T_{n} x\right|+\left|T_{n} x-y_{i}\right|<\varepsilon
$$

Thus, $T(B)$ is covered by finitely many balls of radius $\varepsilon$.
A continuous linear operator is of finite rank when its image is finite-dimensional. A finite-rank operator is compact, since all balls are pre-compact in a finite-dimensional Hilbert space.
[9.A.3] Theorem: A compact operator $T: X \rightarrow Y$ with $Y$ a Hilbert space is an operator norm limit of finite rank operators.
Proof: Let $B$ be the closed unit ball in $X$. Since $T(B)$ is pre-compact it is totally bounded, so for given $\varepsilon>0$ cover $T(B)$ by open balls of radius $\varepsilon$ centered at points $y_{1}, \ldots, y_{n}$. Let $p$ be the orthogonal projection to the finite-dimensional subspace $F$ spanned by the $y_{i}$ and define $T_{\varepsilon}=p \circ T$. Note that for any $y \in Y$ and for any $y_{i}$

$$
\left|p(y)-y_{i}\right| \leq\left|y-y_{i}\right|
$$

since $y=p(y)+y^{\prime}$ with $y^{\prime}$ orthogonal to all $y_{i}$. For $x$ in $X$ with $|x| \leq 1$, by construction there is $y_{i}$ such that $\left|T x-y_{i}\right|<\varepsilon$. Then

$$
\left|T x-T_{\varepsilon} x\right| \leq\left|T x-y_{i}\right|+\left|T_{\varepsilon} x-y_{i}\right|<\varepsilon+\varepsilon
$$

Thus, $T_{\varepsilon} \rightarrow T$ in operator norm as $\varepsilon \rightarrow 0$.

## 9. Unbounded operators on Hilbert spaces

Hilbert-Schmidt operators are an important concrete class of compact operators, as is verified in the claim below. Originally Hilbert-Schmidt operators on function spaces $L^{2}(X)$ arose as operators given by integral kernels: for $X$ and $Y \sigma$-finite measure spaces, and for integral kernel $K \in L^{2}(X \times Y)$, the associated Hilbert-Schmidt operator $T: L^{2}(X) \longrightarrow L^{2}(Y)$ is

$$
T f(y)=\int_{X} K(x, y) f(x) d x
$$

By Fubini's theorem and the $\sigma$-finiteness, for orthonormal bases $\varphi_{\alpha}$ for $L^{2}(X)$ and $\psi_{\beta}$ for $L^{2}(Y)$, the collection of functions $\varphi_{\alpha}(x) \psi_{\beta}(y)$ is an orthonormal basis for $L^{2}(X \times Y)$. Thus, for some scalars $c_{i j}$,

$$
K(x, y)=\sum_{i j} c_{i j} \overline{\varphi_{i}}(x) \psi_{j}(y)
$$

Square-integrability is

$$
\sum_{i j}\left|c_{i j}\right|^{2}=|K|_{L^{2}(X \times Y)}^{2}<\infty
$$

The indexing sets may as well be countable, since an uncountable sum of positive reals cannot converge. Given $f \in L^{2}(X)$, the image $T f$ is in $L^{2}(Y)$, since

$$
T f(y)=\sum_{i j} c_{i j}\left\langle f, \varphi_{i}\right\rangle \psi_{j}(y)
$$

has $L^{2}(Y)$ norm easily estimated by

$$
\begin{gathered}
|T f|_{L^{2}(Y)}^{2} \leq \sum_{i j}\left|c_{i j}\right|^{2}\left|\left\langle f, \varphi_{i}\right\rangle\right|^{2}\left|\psi_{j}\right|_{L^{2}(Y)}^{2} \leq|f|_{L^{2}(X)}^{2} \sum_{i j}\left|c_{i j}\right|^{2}\left|\varphi_{i}\right|_{L^{2}(X)}^{2}\left|\psi_{j}\right|_{L^{2}(Y)}^{2} \\
=|f|_{L^{2}(X)}^{2} \sum_{i j}\left|c_{i j}\right|^{2}=|f|_{L^{2}(X)}^{2} \cdot|K|_{L^{2}(X \times Y)}^{2}
\end{gathered}
$$

The adjoint $T^{*}: L^{2}(Y) \rightarrow L^{2}(X)$ has kernel

$$
K^{*}(y, x)=\overline{K(x, y)}
$$

by computing

$$
\langle T f, g\rangle_{L^{2}(Y)}=\int_{Y}\left(\int_{X} K(x, y) f(x) d x\right) \overline{g(y)} d y=\int_{X} f(x) \overline{\left(\int_{Y} \overline{K(x, y)} g(y) d y\right)} d x
$$

The intrinsic characterization of Hilbert-Schmidt operators $V \rightarrow W$ on Hilbert spaces $V, W$ is as the completion of the space of finite-rank operators $V \rightarrow W$ with respect to the Hilbert-Schmidt norm, whose square is

$$
|T|_{\mathrm{HS}}^{2}=\operatorname{tr}\left(T^{*} T\right) \quad\left(\text { for } T: V \rightarrow W \text { and } T^{*}: W^{*} \rightarrow V^{*}\right)
$$

The trace of a finite-rank operator from a Hilbert space to itself can be described in coordinates and then proven independent of the choice of coordinates, or trace can be described intrinsically, obviating need for proof of coordinate-independence. First, in coordinates, for an orthonormal basis $e_{i}$ of $V$, and finite-rank $T: V \rightarrow V$, define

$$
\left.\operatorname{tr}(T)=\sum_{i}\left\langle T e_{i}, e_{i}\right\rangle \quad \text { (with reference to orthonormal basis }\left\{e_{i}\right\}\right)
$$

With this description, one would need to show independence of the orthonormal basis. For the intrinsic description, consider the map from $V \otimes V^{*}$ to finite-rank operators on $V$ induced from the bilinear map

$$
v \times \lambda \longrightarrow(w \rightarrow \lambda(w) \cdot v) \quad\left(\text { for } v \in V \text { and } \lambda \in V^{*}\right)
$$

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Trace is easy to define in these terms $\operatorname{tr}(v \otimes \lambda)=\lambda(v)$, and

$$
\operatorname{tr}\left(\sum_{v, \lambda} v \otimes \lambda\right)=\sum_{v, \lambda} \lambda(v) \quad \text { (finite sums) }
$$

Expression of trace in terms of an orthonormal basis $\left\{e_{j}\right\}$ is easily obtained from the intrinsic form: given a finite-rank operator $T$ and an orthonormal basis $\left\{e_{i}\right\}$, let $\lambda_{i}(v)=\left\langle v, e_{i}\right\rangle$. We claim that $T=\sum_{i} T e_{i} \otimes \lambda_{i}$. Indeed,

$$
\left(\sum_{i} T e_{i} \otimes \lambda_{i}\right)(v)=\sum_{i} T e_{i} \cdot \lambda_{i}(v)=\sum_{i} T e_{i} \cdot\left\langle v, e_{i}\right\rangle=T\left(\sum_{i} e_{i} \cdot\left\langle v, e_{i}\right\rangle\right)=T v
$$

Then the trace is

$$
\operatorname{tr} T=\operatorname{tr}\left(\sum_{i} T e_{i} \otimes \lambda_{i}\right)=\sum_{i} \operatorname{tr}\left(T e_{i} \otimes \lambda_{i}\right)=\sum_{i} \lambda_{i}\left(T e_{i}\right)=\sum_{i}\left\langle T e_{i}, e_{i}\right\rangle
$$

Similarly, adjoints $T^{*}: W \rightarrow V$ of maps $T: V \rightarrow W$ are expressible in these terms: for $v \in V$, let $\lambda_{v} \in V^{*}$ be $\lambda_{v}\left(v^{\prime}\right)=\left\langle v^{\prime}, v\right\rangle$, and for $w \in W$ let $\mu_{w} \in W^{*}$ be $\mu_{w}\left(w^{\prime}\right)=\left\langle w^{\prime}, w\right\rangle$. Then

$$
\left(w \otimes \lambda_{v}\right)^{*}=v \otimes \mu_{w} \quad(\text { for } w \in W \text { and } v \in V)
$$

since

$$
\left\langle\left(w \otimes \lambda_{v}\right) v^{\prime}, w^{\prime}\right\rangle=\left\langle\lambda_{v}\left(v^{\prime}\right) w, w^{\prime}\right\rangle=\left\langle v^{\prime}, v\right\rangle\left\langle w, w^{\prime}\right\rangle=\left\langle v^{\prime},\left\langle w^{\prime}, w\right\rangle \cdot v\right\rangle=\left\langle v^{\prime},\left(v \otimes \mu_{w}\right) w^{\prime}\right\rangle
$$

Since it is defined as a completion, the collection of all Hilbert-Schmidt operators $T: V \rightarrow W$ is a Hilbert space, with the hermitian inner product $\langle S, T\rangle=\operatorname{tr}\left(T^{*} S\right)$.
[9.A.4] Claim: The Hilbert-Schmidt norm $\left\|\|_{\text {HS }}\right.$ dominates the uniform operator norm $\| \|_{\text {op }}$, so HilbertSchmidt operators are compact.
Proof: Given $\varepsilon>0$, let $e_{1}$ be a vector with $\left|e_{1}\right| \leq 1$ such that $\left|T v_{1}\right| \geq|T|_{\text {op }}-\varepsilon$. Extend $\left\{e_{1}\right\}$ to an orthonormal basis $\left\{e_{i}\right\}$. Then

$$
|T|_{\mathrm{op}}^{2}=\sup _{|v| \leq 1}|T v|^{2} \leq\left|T v_{1}\right|^{2}+\varepsilon \leq \varepsilon+\sum_{j}\left|T v_{j}\right|^{2}=|T|_{\mathrm{HS}}^{2}
$$

Thus, Hilbert-Schmidt norm limits of finite-rank operators are operator-norm limits of finite-rank operators, so are compact.

It is already nearly visible that the $L^{2}(X \times Y)$ norm on kernels $K(x, y)$ is the same as the Hilbert-Schmidt norm on corresponding operators $T: V \rightarrow W$, yielding
[9.A.5] Claim: Operators $T: L^{2}(X) \rightarrow L^{2}(Y)$ given by integral kernels $K \in L^{2}(X \times Y)$ are Hilbert-Schmidt, that is, are Hilbert-Schmidt norm limits of finite-rank operators.
Proof: To prove properly that the $L^{2}(X \times Y)$ norm on kernels $K(x, y)$ is the same as the Hilbert-Schmidt norm on corresponding operators $T: V \rightarrow W, T$ should be expressed as a limit of finite-rank operators $T_{n}$ in terms of kernels $K_{n}(x, y)$ which are finite sums of products $\varphi(x) \otimes \psi(y)$. Thus, first claim that

$$
K(x, y)=\sum_{i} \overline{\varphi_{i}}(x) T \varphi_{i}(y) \quad\left(\text { in } L^{2}(X \times Y)\right)
$$

Indeed, the inner product in $L^{2}(X \times Y)$ of the right-hand side against any $\varphi_{i}(x) \psi_{j}(y)$ agrees with the inner product of the latter against $K(x, y)$, and we have assumed $K \in L^{2}(X \times Y)$. With $K=\sum_{i j} c_{i j} \bar{\varphi}_{i} \otimes \psi_{j}$,

$$
T \varphi_{i}=\sum_{j} c_{i j} \psi_{j}
$$

Since $\sum_{i j}\left|c_{i j}\right|^{2}$ converges,

$$
\lim _{i}\left|T \varphi_{i}\right|^{2}=\lim _{i} \sum_{j}\left|c_{i j}\right|^{2}=0
$$

and

$$
\lim _{n} \sum_{i>n}\left|T \varphi_{i}\right|^{2}=\lim _{n} \sum_{i>n}\left|c_{i j}\right|^{2}=0
$$

so the infinite sum $\sum_{i} \bar{\varphi}_{i} \otimes T \varphi_{i}$ converges to $K$ in $L^{2}(X \times Y)$. In particular, the truncations

$$
K_{n}(x, y)=\sum_{1 \leq i \leq n} \bar{\varphi}_{i}(x) T \varphi_{i}(y)
$$

converge to $K(x, y)$ in $L^{2}(X \times Y)$, and give finite-rank operators

$$
T_{n} f(y)=\int_{X} K_{n}(x, y) f(x) d x
$$

We claim that $T_{n} \rightarrow T$ in Hilbert-Schmidt norm. It is convenient to note that by a similar argument $\overline{K(x, y)}=\sum_{i} T^{*} \psi_{i}(x) \bar{\psi}_{i}(y)$. Then

$$
\begin{aligned}
& \left|T-T_{n}\right|_{\mathrm{HS}}^{2}=\operatorname{tr}\left(\left(T-T_{n}\right)^{*} \circ\left(T-T_{n}\right)\right)=\sum_{i, j>n} \operatorname{tr}\left(\left(T^{*} \psi_{i} \otimes \bar{\psi}_{i}\right) \circ\left(\bar{\varphi}_{j} \otimes T \varphi_{j}\right)\right) \\
& \left.=\sum_{i, j>n}\left\langle T^{*} \psi_{i}, \varphi_{j}\right\rangle_{L^{2}(X)} \cdot\left\langle T \varphi_{j}, \psi_{i}\right\rangle_{L^{2}(Y)}=\sum_{i, j>n}\left|c_{i j}\right|^{2} \longrightarrow 0 \quad \text { (as } n \rightarrow \infty\right)
\end{aligned}
$$

since $\sum_{i j}\left|c_{i j}\right|^{2}$ converges. Thus, $T_{n} \rightarrow T$ in Hilbert-Schmidt norm.
Now we come to the spectral theorem for self-adjoint compact operators. Again, the $\lambda$-eigenspace $V_{\lambda}$ of a self-adjoint compact operator $T$ on a Hilbert space $T$ is

$$
V_{\lambda}=\{v \in V: T v=\lambda \cdot v\}
$$

We have already shown that eigenvalues, if any, of self-adjoint $T$ are real.
[9.A.6] Theorem: Let $T$ be a self-adjoint compact operator on a non-zero Hilbert space $V$.

- The completion of $\oplus V_{\lambda}$ is all of $V$. In particular, there is an orthonormal basis of eigenvectors.
- For infinite-dimensional $V, 0$ is the only accumulation point of the set of eigenvalues.
- Every eigenspaces $X_{\lambda}$ for $\lambda \neq 0$ is finite-dimensional. The 0 -eigenspace may be $\{0\}$, finite-dimensional, or infinite-dimensional.
- (Rayleigh-Ritz) One or the other of $\pm|T|_{\mathrm{op}}$ is an eigenvalue of $T$, with operator norm $|\cdot|_{\mathrm{op}}$.

Proof: An alternative expression for the operator norm is needed:
[9.A.7] Lemma: $|T|_{\mathrm{op}}=\sup _{|x| \leq 1}|\langle T x, x\rangle|$ for $T$ a self-adjoint continuous linear operator on a Hilbert space. Proof: Let $s$ be that supremum. By Cauchy-Schwarz-Bunyakowsky, $s \leq|T|_{\text {op }}$. For any $x, y$, by polarization

$$
\begin{gathered}
2|\langle T x, y\rangle+\langle T y, x\rangle|=|\langle T(x+y), x+y\rangle-\langle T(x-y), x-y\rangle| \\
\leq|\langle T(x+y), x+y\rangle|+|\langle T(x-y), x-y\rangle| \leq s|x+y|^{2}+s|x-y|^{2}=2 s\left(|x|^{2}+|y|^{2}\right)
\end{gathered}
$$

With $y=t \cdot T x$ with $t>0$, because $T=T^{*}$,

$$
\langle T x, y\rangle=\langle T x, t \cdot T x\rangle=t \cdot|T x|^{2} \geq 0 \quad(\text { for } y=t \cdot T x \text { with } t>0)
$$

and

$$
\langle T y, x\rangle=\left\langle t \cdot T^{2} x, t \cdot x\right\rangle=t \cdot\langle T x, T x\rangle=t \cdot|T x|^{2} \geq 0 \quad(\text { for } y=t \cdot T x \text { with } t>0)
$$

Thus,

$$
|\langle T x, y\rangle|+|\langle T y, x\rangle|=\langle T x, y\rangle+\langle T y, x\rangle=|\langle T x, y\rangle+\langle T y, x\rangle| \quad(\text { for } y=t \cdot T x \text { with } t>0)
$$

From this, and from the polarization identity divided by 2 ,

$$
\left.|\langle T x, y\rangle|+|\langle T y, x\rangle|=|\langle T x, y\rangle+\langle T y, x\rangle| \leq s\left(|x|^{2}+|y|^{2}\right) \quad \quad \text { (with } y=t \cdot T x\right)
$$

Divide through by $t$ to obtain

$$
|\langle T x, T x\rangle|+\left|\left\langle T^{2} x, x\right\rangle\right| \leq \frac{s}{t} \cdot\left(|x|^{2}+|T x|^{2}\right)
$$

Minimize the right-hand side by taking $t=|x| /|T x|$, and note that $\left\langle T^{2} x, x\right\rangle=\langle T x, T x\rangle$, giving

$$
2|\langle T x, T x\rangle| \leq 2 s \cdot|x| \cdot|T x| \leq 2 s \cdot|x|^{2} \cdot|T|_{\mathrm{op}}
$$

Thus, $|T|_{\mathrm{op}} \leq s$.
The last assertion of the theorem is the starting point of the proof and uses $|T|=\sup _{|x| \leq 1}|\langle T x, x\rangle|$ and the fact that any value $\langle T x, x\rangle$ is real, by self-adjointness. Choose a sequence $\left\{x_{n}\right\}$ so that $\left|x_{n}\right| \leq 1$ and $|\langle T x, x\rangle| \rightarrow|T|$. Replacing it by a subsequence if necessary, the sequence $\langle T x, x\rangle$ of real numbers has a limit $\lambda= \pm|T|$. Then

$$
\begin{gathered}
0 \leq\left|T x_{n}-\lambda x_{n}\right|^{2}=\left\langle T x_{n}-\lambda x_{n}, T x_{n}-\lambda x_{n}\right\rangle=\left|T x_{n}\right|^{2}-2 \lambda\left\langle T x_{n}, x_{n}\right\rangle+\lambda^{2}\left|x_{n}\right|^{2} \\
\leq \lambda^{2}-2 \lambda\left\langle T x_{n}, x_{n}\right\rangle+\lambda^{2}
\end{gathered}
$$

The right-hand side goes to 0 . By compactness of $T$, replace $x_{n}$ by a subsequence so that $T x_{n}$ converges to some vector $y$. The previous inequality shows $\lambda x_{n} \rightarrow y$. For $\lambda=0$, we have $|T|=0$, so $T=0$. For $\lambda \neq 0, \lambda x_{n} \rightarrow y$ implies $x_{n} \longrightarrow \lambda^{-1} y$. For $x=\lambda^{-1} y$, we have $T x-\lambda x$ and $x$ is the desired eigenvector with eigenvalue $\pm|T|$.

Now use induction. The completion $Y$ of the sum of non-zero eigenspaces is $T$-stable. We claim that the orthogonal complement $Z=Y^{\perp}$ is $T$-stable, and the restriction of $T$ to is a compact operator. Indeed, for $z \in Z$ and $y \in Y,\langle T z, y\rangle=\langle z, T y\rangle=0$, proving stability. The unit ball in $Z$ is a subset of the unit ball $B$ in $X$, so has pre-compact image $T B \cap Z$ in $X$. Since $Z$ is closed in $X$, the intersection $T B \cap Z$ of $Z$ with the pre-compact $T B$ is pre-compact, proving $T$ restricted to $Z=Y^{\perp}$ is still compact. Self-adjoint-ness is clear.

By construction, the restriction $T_{1}$ of $T$ to $Z$ has no eigenvalues on $Z$, since any such eigenvalue would also be an eigenvalue of $T$ on $Z$. Unless $Z=\{0\}$ this would contradict the previous argument, which showed that $\pm\left|T_{1}\right|$ is an eigenvalue on a non-zero Hilbert space. Thus, it must be that the completion of the sum of the eigenspaces is all of $X$.

To prove that eigenspaces $V_{\lambda}$ for $\lambda \neq 0$ are finite-dimensional, and that there are only finitely-many eigenvalues $\lambda$ with $|\lambda|>\varepsilon$ for given $\varepsilon>0$, let $B$ be the unit ball in $Y=\sum_{|\lambda|>\varepsilon} X_{\lambda}$. The image of $B$ by $T$ contains the ball of radius $\varepsilon$ in $Y$. Since $T$ is compact, this ball is pre-compact, so $Y$ is finite-dimensional. Since the dimensions of the $X_{\lambda}$ are positive integers, there can be only finitely-many of them with $|\lambda|>\varepsilon$, and each is finite-dimensional. It follows that the only possible accumulation point of the set of eigenvalues is 0 , and, for $X$ infinite-dimensional, 0 must be an accumulation point.
[9.A.8] Corollary: For a self-adjoint compact operator $T: X \rightarrow X$ on a Hilbert space $X$, for $\lambda \neq 0$ not an eigenvalue, $(T-\lambda) X=X$.
Proof: By the spectral theorem, $(T-\lambda)^{-1}$ exists.

## 9.B Appendix: open mapping and closed graph theorems

[9.B.1] Theorem: (Open Mapping Theorem) For a continuous linear surjection $T: X \rightarrow Y$ of Banach spaces, there is $\delta>0$ such that for all $y \in Y$ with $|y|<\delta$ there is $x \in X$ with $|x| \leq 1$ such that $T x=y$. In particular, $T$ is an open map.
[9.B.2] Corollary: A bijective continuous linear map of Banach spaces is an isomorphism.
///
Proof: In the corollary the non-trivial point is that $T$ is open, which is the point of the theorem. The linearity of the inverse is easy.

For every $y \in Y$ there is $x \in X$ so that $T x=y$. For some integer $n$ we have $n>|x|$, so $Y$ is the union of the sets $T B(n)$, with usual open balls

$$
B(n)=\{x \in X:|x|<n\}
$$

By Baire category [15.A], the closure of some one of the sets $T B(n)$ contains a non-empty open ball

$$
V=\left\{y \in Y:\left|y-y_{o}\right|<r\right\}
$$

for some $r>0$ and $y_{o} \in Y$. Since we are in a metric space, the conclusion is that every point of $V$ occurs as the limit of a Cauchy sequence consisting of elements from $T B(n)$. Certainly

$$
\{y \in Y:|y|<r\} \subset\left\{y_{1}-y_{2}: y_{1}, y_{2} \in V\right\}
$$

Thus, every point in the ball $B_{r}^{\prime}$ of radius $r$ centered at 0 in $Y$ is the sum of two limits of Cauchy sequences from $T B(n)$. Thus, surely every point in $B_{r}^{\prime}$ is the limit of a single Cauchy sequence from the image $T B(2 n)$ of the open ball $B(2 n)$ of twice the radius. That is, the closure of $T B(2 n)$ contains the ball $B^{\prime}(r)$.

Using the linearity of $T$, the closure of $T B(\rho)$ contains the ball $B^{\prime}(r \rho / 2 n)$ in $Y$.
Given $|y|<1$, choose $x_{1} \in B(2 n / r)$ so that $\left|y-T x_{1}\right|<\varepsilon$. Choose $x_{2} \in B\left(\varepsilon \cdot \frac{2 n}{r}\right)$ so that $\left|\left(y-T x_{1}\right)-T x_{2}\right|<\varepsilon / 2$. Choose $x_{3} \in B\left(\frac{\varepsilon}{2} \cdot \frac{2 n}{r}\right)$ so that

$$
\left|\left(y-T x_{1}-T x_{2}\right)-T x_{3}\right|<\varepsilon / 2^{2}
$$

Choose $x_{4} \in B\left(\frac{\varepsilon}{2^{2}} \cdot \frac{2 n}{r}\right)$ so that

$$
\left|\left(y-T x_{1}-T x_{2}-T x_{3}\right)-T x_{4}\right|<\varepsilon / 2^{3}
$$

and so on. The sequence

$$
x_{1}, x_{1}+x_{2}, x_{1}+x_{2}+x_{3}, \ldots
$$

is Cauchy in $X$. Since $X$ is complete, the limit $x$ of this sequence exists in $X$, and $T x=y$. We find that

$$
x \in B\left(\frac{2 n}{r}\right)+B\left(\varepsilon \frac{2 n}{r}\right)+B\left(\frac{\varepsilon}{2} \cdot \frac{2 n}{r}\right)+B\left(\frac{\varepsilon}{2^{2}} \cdot \frac{2 n}{r}\right)+\ldots \subset B\left((1+2 \varepsilon) \frac{2 n}{r}\right)
$$

Thus,

$$
T B\left((1+\varepsilon) \frac{2 n}{r}\right) \supset\{y \in Y:|y|<1\}
$$

This proves open-ness at 0 .
It is straightforward to show ${ }^{[65]}$ that a continuous map $f: X \rightarrow Y$ of Hausdorff topological spaces has closed graph

$$
\Gamma_{f}=\{(x, y): f(x)=y\} \subset X \times Y
$$

[65] To show that a continuous map $f: X \rightarrow Y$ of topological spaces with $Y$ Hausdorff has closed graph $\Gamma_{f}$, show the complement is open. Take $(x, y) \notin \Gamma_{f}$. Let $V_{1}$ be a neighborhood of $f(x)$ and $V_{2}$ a neighborhood of $y$ such that $V_{1} \cap V_{2}=\phi$, using Hausdorff-ness. By continuity of $f$, for $x^{\prime}$ in a suitable neighborhood $U$ of $x$, the image $f\left(x^{\prime}\right)$ is inside $V_{1}$. Thus, the neighborhood $U \times V_{2}$ of $(x, y)$ does not meet $\Gamma_{f}$.

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Similarly, a topological space $X$ is Hausdorff if and only if the diagonal $X^{\Delta}=\{(x, x): x \in X\}$ is closed in $X \times X$. ${ }^{[66]}$
[9.B.3] Theorem: (Closed Graph Theorem) A linear map $T: V \rightarrow W$ of Banach spaces is continuous if it has closed graph $\Gamma=\{(v, w): T v=w\}$.
Proof: The direct sum $V \oplus W$ with norm $|v \oplus w|=|v|+|w|$ is a Banach space. Since $\Gamma$ is a closed subspace of $V \oplus W$, it is a Banach space itself with the restriction of this norm. The projection $\pi_{V}: V \oplus W \rightarrow V$ is a continuous linear map. The restriction $\left.\pi_{V}\right|_{\Gamma}$ of $\pi_{V}$ to $\Gamma$ is still continuous, and still surjective, because it $T$ is an everywhere-defined function on $V$. By the open mapping theorem, $\left.\pi_{V}\right|_{\Gamma}$ is open. Thus, the bijection $\left.\pi_{V}\right|_{\Gamma}$ is a homeomorphism. Letting $\pi_{W}: V \oplus W \rightarrow W$ be the projection to $W$,

$$
T=\pi_{W} \circ\left(\left.\pi_{V}\right|_{\Gamma}\right)^{-1}: V \longrightarrow W
$$

expresses $T$ as a composition of continuous functions.

## 9.C Appendix: irreducibles of compact groups

As usual, now specifically for compact topological groups $K$, a representation of $K$ on a quasi-complete, locally convex topological vector space $V$ is a continuous map $K \times V \rightarrow V$ making $K$ act by continuous linear maps on $V$. Such a representation $V$ is (topologically) irreducible when there are no $K$-stable (topologically) closed subspaces of $V$ except $\{0\}$ and $V$ itself. A $K$-homomorphism $\varphi: V \rightarrow W$ of $K$-representations is a continuous linear map which respects the action of $K: \varphi(k \cdot v)=k \cdot \varphi(v)$.
[9.C.1] Claim: Every representation of compact $K$ on a Hilbert space $V$ is isomorphic to a unitary representation of $K$ on $V$. That is, there is another inner product $\langle,\rangle^{\prime}$ on $V$, comparable to the original inner product $\langle$,$\rangle in the sense that there are finite constants 0<c_{1}, c_{2}$ such that

$$
c_{1} \cdot\langle v, v\rangle \leq\langle v, v\rangle^{\prime} \leq c_{2} \cdot\langle v, v\rangle
$$

for all $v \in V$, and such that $\langle k \cdot v, k \cdot v\rangle=\langle v, v\rangle^{\prime}$ for all $v \in V$ and $k \in K$.
Proof: The natural idea to average the original inner product by the action of $K$ succeeds, because $K$ is compact: let

$$
\langle v, v\rangle^{\prime}=\int_{K}\langle k \cdot v, k \cdot v\rangle d k
$$

Since $K$ is compact, for each $v \in V$ the orbit $K \cdot v=\{k \cdot v: k \in K\}$ is compact, so bounded. By BanachSteinhaus (uniform boundedness) [13.12.3], the action of elements $k \in K$ are uniformly equicontinuous: given $\varepsilon>0$, there is $\delta>0$ such that

$$
|v|<\delta \Longrightarrow|k \cdot v|<\varepsilon
$$

With $\varepsilon=1,|v| \leq 1$ implies $|k \cdot v|<\delta^{-1}$. That is, the operator norm of $v \rightarrow k \cdot v$ is at most $\delta^{-1}$ for all $k \in K$. That is, $|k \cdot v| \leq \delta^{-1} \cdot|v|$ for all $k, v$. Replacing $k$ by $k^{-1}$ and $v$ by $k \cdot v$, we similarly have $|v| \leq \delta^{-1}|k \cdot v|$, which gives $\delta|v| \leq|k \cdot v|$. Thus, integrating,

$$
\delta^{2} \cdot|v|^{2} \cdot \operatorname{meas}(K) \leq \int_{K}|k \cdot v|^{2} d k \leq \delta^{-2} \cdot|v|^{2} \cdot \operatorname{meas}(K)
$$

Not every norm arises from an inner product. To see that the new norm-squared $|v|_{\text {new }}^{2}=\int_{K}|k \cdot v|^{2} d k$ does arise from an inner product, it suffices to prove the polarization identity

$$
|v+w|_{\text {new }}^{2}-|v-w|_{\text {new }}^{2}=2|v|_{\text {new }}^{2}+2|w|_{\text {new }}^{2}
$$

[66] To show that closed-ness of the diagonal $X^{\Delta}$ in $X \times X$ implies $X$ is Hausdorff, let $x_{1} \neq x_{2}$ be points in $X$. Then there is a neighborhood $U_{1} \times U_{2}$ of $\left(x_{1}, x_{2}\right)$, with $U_{i}$ a neighborhood of $x_{i}$, not meeting the diagonal. That is, $\left(x, x^{\prime}\right) \in U_{1} \times U_{2}$ implies $x \neq x^{\prime}$. That is, $U_{1} \cap U_{2}=\phi$.

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This follows by integrating the polarization identity for the original norm.
The action of $K$ on $L^{2}(K)$ by right translation is unitary, because the measure is invariant. The continuity of this action follows by a simpler form of the argument of [6.1] and [6.2]. As there, $K$ acts continuously on $C_{c}^{o}(K)$ by right translation. The density of $C_{c}^{o}(K)$ in $L^{2}(K)$, and the domination of the $L^{2}$ norm by the sup-norm, give the continuity.

For $\lambda \in V^{*}$ and $v \in V$, let $c_{v, \lambda}(k)=\lambda(k \cdot v)$ for $k \in K$. The function $c_{v, \lambda}$ on $K$ is a (matrix) coefficient function.
[9.C.2] Claim: Every Hilbert space irreducible $V$ of $K$ has a $K$-homomorphism to $L^{2}(K)$, by the map

$$
v \longrightarrow c_{v, \lambda} \quad\left(\text { for fixed } 0 \neq \lambda \in V^{*}\right)
$$

Proof: Without loss of generality, we can assume $V$ is unitary, by the previous. The function $k \times v \rightarrow k \cdot v$ is a (jointly) continuous function $K \times V \rightarrow V$, by assumption. Composing with $\lambda$ gives a continuous function $K \times V \rightarrow \mathbb{C}$. We claim that $v \rightarrow\left(k \rightarrow c_{v, \lambda}(k)\right)$ is a continuous $C^{o}(K)$-valued function on $V$ : for $\left|v-v^{\prime}\right|<\delta$,

$$
\begin{gathered}
\left|c_{v, \lambda}(k)-c_{v^{\prime}, \lambda}\left(k^{\prime}\right)\right|=\left|\lambda\left(k \cdot v-k^{\prime} \cdot v^{\prime}\right)\right|=\left|\lambda\left(k \cdot v-k^{\prime} \cdot v\right)+\lambda\left(k^{\prime} \cdot v-k^{\prime} \cdot v^{\prime}\right)\right| \\
\leq|\lambda|_{V^{*}} \cdot\left(\left(k \cdot v-\left.k^{\prime} \cdot v\right|_{V}+\left(k^{\prime} \cdot v-\left.k^{\prime} \cdot v^{\prime}\right|_{V}\right)\right.\right.
\end{gathered}
$$

By unitariness, $\left|k^{\prime} \cdot\left(v-v^{\prime}\right)\right|=\left|v-v^{\prime}\right|$. By the continuity of the action of $K$ on $V,\left|k \cdot v-k^{\prime} \cdot v\right|<\varepsilon$ for given $v$ for $k^{\prime}$ sufficiently close to $k$.

To see that $v \rightarrow c_{v, \lambda}$ is a $K$-homomorphism, for $x, y \in K$,

$$
c_{x \cdot v, \lambda}(y)=\lambda(y \cdot(x \cdot v))=\lambda((y \cdot x) \cdot v)=c_{v, \lambda}(y x)
$$

This proves the claim.
As usual [14.1], $\varphi \in C_{c}^{o}(K)$ acts on a $K$-representation $W$ by integral operators

$$
\varphi \cdot w=\int_{K} \varphi(k) k \cdot w d k
$$

Thus, such $W$ becomes a $C_{c}^{o}(K)$-representation, as discussed in somewhat greater generality in [9.D]. Potential issues about multiplicities are clarified in [9.D.14].
[9.C.3] Claim: $L^{2}(K)$ is the completion of an orthogonal direct sum $\oplus_{V} m_{V} \cdot V$ of orthogonal sums $m_{V} \cdot V=\underbrace{V \oplus \ldots \oplus V}_{m_{V}}$ of $C_{c}^{o}(K)$-irreducibles $V$, each occurring with finite multiplicity $m_{V}$.
[9.C.4] Remark: This claim is an extreme case of [7.B]'s treatment of compact $\Gamma \backslash G$, where now $\Gamma=\{1\}$. The argument simplifies, as well. Potential ambiguities about the notion of multiplicity are resolved in [9.D.14].
Proof: On $L^{2}(K)$ this is

$$
(\varphi \cdot f)(x)=\int_{K} \varphi(y) f(x y) d y=\int_{K} \varphi\left(x^{-1} y\right) f(y) d y
$$

The function $x \times y \rightarrow \varphi\left(x^{-1} y\right)$ is continuous on $K \times K$, so is in $L^{2}(K \times K)$ by the compactness of $K$. Thus, $\varphi$ gives a Hilbert-Schmidt operator [9.A.5] on $L^{2}(K)$. The adjoint of the operator given by $\varphi$ is easily determined, and is again in $C_{c}^{o}(K)$. This action is non-degenerate, in the sense that for given $f \in L^{2}(K)$, there is $\varphi \in C_{c}^{o}(K)$ such that $\varphi \cdot f \neq 0$, from [14.1.5]. That is, the ring of operators on $L^{2}$ is adjoint-stable, non-degenerate, and consists of compact operators, so [7.2.18] applies: $L^{2}(K)$ is the completion of a direct sum of irreducible $C_{c}^{o}(K)$-representations, each occurring with finite multiplicity.
[9.C.5] Corollary: The $C_{c}^{o}(K)$-irreducible subrepresentations in $L^{2}(K)$ are exactly the $K$-irreducible subrepresentations. Thus, $L^{2}(K)$ is the completion of an orthogonal direct sum of $K$-irreducibles $V$, each occurring with finite multiplicity $m_{V}$.

Proof: This is a special case of [14.1.6] and [14.1.7]: irreducible $C_{c}^{o}(K)$-subrepresentations of a $K$ representation are irreducible $K$-subrepresentations.
[9.C.6] Remark: With a little more effort, one can prove that $m_{V}=\operatorname{dim}_{\mathbb{C}} V$, and more (for example, Schur inner-product relations), but the assertion of the claim is all we need for our immediate purposes.
[9.C.7] Corollary: All Hilbert-space irreducibles of compact $K$ are finite-dimensional.
Proof: A copy of every $K$-irreducible appears inside $L^{2}(K)$, where all irreducibles are finite-dimensional.
[9.C.8] Corollary: For two compact groups $K_{1}$ and $K_{2}$, the Hilbert-space irreducibles of $K_{1} \times K_{2}$ are tensor products of Hilbert-space irreducibles of $K_{1}$ and of $K_{2}$.
Proof: Let $V$ be an irreducible Hilbert-space representation of $K_{1} \times K_{2}$. From above, without loss of generality, the representation is unitary. From the previous corollary, $V$ is finite-dimensional. Forgetting the action of $K_{2}, V$ is a finite-dimensional representation of $K_{1}$, so is a finite orthogonal direct sum of irreducibles.

For an irreducible $W$ of $K$ appearing in $V$, the $W$-isotype $V^{W}$ of $V$ is the (not necessarily direct) sum of all copies of $W$ in $V$. By [9.D.14], this sum is expressible an orthogonal direct sum. We claim that $K_{2}$ stabilizes $V^{W}$. If not, the orthogonal projection from some image $k_{2} \cdot V^{W}$ to some other isotype $V^{W^{\prime}}$ would be non-zero. But the orthogonal projections to $K_{1}$-isotypes are $K_{1}$-homomorphisms, as are the orthogonal projections to copies of $W$ inside $V^{W}$. The kernel and image of $K_{1}$-homomorphisms $W \rightarrow W^{\prime}$ are subrepresentations, since in finite-dimensional spaces all subspaces are (topologically) closed. Thus, if the kernel is not all of $W$, the map is an injection, so has non-zero image, so is all of $W^{\prime}$, giving an isomorphism $W \rightarrow W^{\prime}$, which is impossible for non-isomorphic irreducibles. Thus, $K_{2}$ stabilizes $V^{W}$. Thus, $K_{1} \times K_{2}$ stabilizes $V^{W}$, so by irreducibility of $V$ this (non-zero) isotype is all of $V$, that is, $V=V^{W}$.

In any case, $\operatorname{Hom}_{K_{1}}(W, V)$ has a $K_{2}$-representation structure given by post-application of the action of $k_{2}$ :

$$
\left(k_{2} \cdot \varphi\right)(w)=k_{2} \cdot \varphi(w)
$$

The map $W \otimes_{\mathbb{C}} \operatorname{Hom}_{K_{1}}(W, V) \longrightarrow V$ by $w \otimes \varphi \rightarrow \varphi(w)$ is a non-zero $K_{1} \times K_{2}$-homomorphism to $V$, so must surject to $V$, by the irreducibility of $V$.

For the converse: let $W_{1}, W_{2}$ be unitary irreducibles of $K_{1}, K_{2}$, and claim that $V=W_{1} \otimes W_{2}$ is an irreducible $K_{1} \times K_{2}$-representation. For a $K_{1} \times K_{2}$-subrepresentation $W \subset V$, the orthogonal projections to $W$ and $W^{\perp}$ are $K_{1} \times K_{2}$-homomorphisms. Thus, if $V$ is reducible, then it has non-scalar $K_{1} \times K_{2^{-}}$ endomorphisms. Proving that any endomorphism $\varphi$ of $V$ is scalar will prove that $V$ is irreducible. For fixed $w_{2} \in W_{2}$ and $\lambda_{2} \in W_{2}^{*}$, we can map $W_{1} \otimes W_{2} \rightarrow W_{1}$ by $w_{1} \otimes w_{2} \longrightarrow \lambda_{2}\left(w_{2}\right) \cdot w_{1}$, and then consider

$$
w_{1} \longrightarrow w_{1} \otimes w_{2} \longrightarrow \varphi\left(w_{1} \otimes w_{2}\right) \longrightarrow W_{1}
$$

This is a $K_{1}$-homomorphism, so is a scalar $c_{w_{2}, \lambda_{2}}$ by (the finite-dimensional version of) Schur's lemma. The $\operatorname{map} W_{2} \rightarrow W_{2}^{* *} \approx W_{2}$ by $w_{2} \rightarrow\left(\lambda_{2} \rightarrow c_{w_{2}, \lambda_{2}}\right)$ is a $K_{2}$-homomorphism, so by Schur's lemma there is a constant $c$ such that $c_{w_{2}, \lambda_{2}}=c \cdot \lambda_{2}\left(w_{2}\right)$. Then, for all $\lambda_{1} \in W_{1}^{*}$ and $\lambda_{2} \in W_{2}^{*}$,

$$
\left(\lambda_{1} \otimes \lambda_{2}\right)\left(\varphi\left(w_{1} \otimes w_{2}\right)\right)=\lambda_{1}\left(c_{w_{2}, \lambda_{2}} \cdot w_{1}\right)=\lambda_{1}\left(w_{1}\right) \cdot c \cdot \lambda_{2}\left(w_{2}\right)=c \cdot\left(\lambda_{1} \otimes \lambda_{2}\right)\left(w_{1} \otimes w_{2}\right)
$$

Thus, any $K_{1} \times K_{2}$-endomorphism $\varphi$ acts by a scalar, so $V$ is irreducible.

## 9.D Appendix: spectral theorem, Schur's lemma, multiplicities

A portion of a spectral theorem for bounded self-adjoint operators on Hilbert spaces is necessary to prove a form of Schur's lemma [9.D.12], itself used to remove ambiguities about multiplicities of irreducible representations [9.D.14].

The present discussion continues in the context of [9.A]. Let $T$ be a continuous self-adjoint linear map $V \rightarrow V$ for a (separable) Hilbert space $V$, with spectrum

$$
\sigma(T)=\left\{\lambda \in \mathbb{C}:(T-\lambda)^{-1} \text { does not exist }\right\}
$$

## 9. Unbounded operators on Hilbert spaces

[9.D.1] Claim: For self-adjoint $T$, the spectrum $\sigma(T)$ is a non-empty compact subset of $\mathbb{R}$.
Proof: First, we show that $T-\lambda$ is invertible for $|\lambda|>|T|_{\mathrm{op}}$. The natural heuristic expands a geometric series:

$$
(T-\lambda)^{-1}=-\lambda^{-1} \cdot\left(1-\frac{T}{\lambda}\right)^{-1}=-\lambda^{-1} \cdot\left(1+\frac{T}{\lambda}+\left(\frac{T}{\lambda}\right)^{2}+\ldots\right)
$$

Since $|T / \lambda|_{\text {op }}<1$, the latter infinite sum does converge in operator norm. Then, just as with geometric series of real or complex numbers, it is easy to check that this infinite sum converges to $(T-\lambda)^{-1}$.

To prove that $\sigma(T)$ is closed, show that $\mu \in \mathbb{C}$ sufficiently close to $\lambda \notin \sigma(T)$ is also not in $\sigma(T)$. Again, this uses geometric series expansions as a natural heuristic to obtain an expression for $(T-\mu)^{-1}$ as a convergent series:

$$
\begin{gathered}
(T-\mu)^{-1}=((T-\lambda)-(\mu-\lambda))^{-1}=\left(1-(\mu-\lambda)(T-\lambda)^{-1}\right) \circ(T-\lambda)^{-1} \\
\left.=\left(1+(\mu-\lambda)(T-\lambda)^{-1}+\left((\mu-\lambda)(T-\lambda)^{-1}\right)^{2}+\ldots\right)\right) \circ(T-\lambda)^{-1}
\end{gathered}
$$

For $|\mu-\lambda|$ small enough that $\left|(\mu-\lambda) \cdot(T-\lambda)^{-1}\right|=|\mu-\lambda| \cdot\left|(T-\lambda)^{-1}\right|_{\mathrm{op}}<1$, the geometric series converges, and is readily checked to give $(T-\mu)^{-1}$.

To show that $\sigma(T) \subset \mathbb{R}$, show that $T-\lambda$ is both injective and surjective for $\lambda \notin \mathbb{R}$. Then the open mapping theorem [9.B.1] shows that the inverse is continuous. For injectivity, note that $\langle T v, v\rangle=\langle v, T v\rangle=\overline{\langle T v, v\rangle}$ implies that $\langle T v, v\rangle$ is real. Then $(T-\lambda) v=0$ with $v \neq 0$ implies $\langle(T-\lambda) v, v\rangle=0$, from which $\lambda \in \mathbb{R}$. For surjectivity, suppose $\langle(T-\lambda) v, w\rangle=0$ for some $w \neq 0$. In particular, $\langle(T-\lambda) w, w\rangle=0$. Again using the fact that $0 \neq\langle T w, w\rangle \in \mathbb{R}$, this would require that $\lambda \in \mathbb{R}$.

Liouville's theorem on bounded entire functions implies that the spectrum of a continuous linear operator on a Hilbert space is not empty, as follows. If a continuous $R_{\lambda}=(T-\lambda)^{-1}$ exists for every complex $\lambda$, then for $0 \neq v \in V, R_{\lambda} v \in V$ is never $0 \in V$. Take $w \in V$ such that $\left\langle R_{\lambda_{o}} v, w\right\rangle \neq 0$ for some $\lambda_{o} \in \mathbb{C}$. Then $f(\lambda)=\left\langle R_{\lambda} v, w\right\rangle$ is a not-identically 0 entire function. At the same time, for large $|\lambda|$, the operator norm of $R_{\lambda}$ is small. Thus, $f(\lambda)$ is small for large $|\lambda|$, and must be identically 0 , by Liouville, contradiction. ///

For a self-adjoint continuous operator $S$ on $V$, write $S \geq 0$ when $\langle S v, v\rangle \geq 0$ for all $v \in V$. For selfadjoint $S, T$, write $S \leq T$ when $T-S \geq 0$. At the outset, with $a \leq-|T|_{\mathrm{op}}$ and $b \geq|T|_{\mathrm{op}}$, we have, $\langle a \cdot v, v\rangle \leq\langle T v, v\rangle \leq\langle b \cdot v, v\rangle$. That is, $a \leq T \leq b$, where the scalars refer to scalar operators on $V$. Here all functions are real-valued, and $C^{o}[a, b]$ refers to real-valued continuous functions on $[a, b]$.
[9.D.2] Theorem: The map $\mathbb{R}[x] \rightarrow \mathbb{R}[T]$ on polynomials given by $f \rightarrow f(T)$ is continuous, where $\mathbb{R}[x]$ has the sup-norm on $[a, b]$ and $\mathbb{R}[T]$ has the uniform operator norm. Thus, by Weierstra $\beta$ approximation, this map extends to a continuous map $C^{o}[a, b] \rightarrow \overline{\mathbb{R}}[T]$, the latter being the operator-norm completion of $\mathbb{R}[T]$. This map factors through $C^{o}(\sigma(T))$ :

$$
C^{o}[a, b] \longrightarrow C^{o}(\sigma(T)) \longrightarrow \overline{\mathbb{R}[T]}
$$

and the map $C^{o}(\sigma(T)) \rightarrow \overline{\mathbb{R}[T]}$ is an isometric isomorphism, where $C^{o}(\sigma(T))$ has sup-norm.
Proof: We claim that for $f \in \mathbb{R}[x]$ with $f(x) \geq 0$ on $[a, b]$, then $f(T) \geq 0$. From the following lemma on polynomials, $f$ is expressible as a finite sum of the form

$$
f=\sum_{i} P_{i}^{2}+(x-a) \sum_{j} Q_{j}^{2}+(b-x) \sum_{k} R_{k}^{2}
$$

for polynomials $P_{i}, Q_{j}, R_{k}$ in $\mathbb{R}[x]$. Incidentally,
[9.D.3] Lemma: For commuting self-adjoint $S, T$ with $T \geq 0$, also $S^{2} T \geq 0$.
Proof: $\left\langle S^{2} T v, v\right\rangle=\left\langle T S v, S^{*} v\right\rangle=\langle T(S v),(S v)\rangle \geq 0$.
Thus, since $a \leq T \leq b$, and all these operators commute (being polynomials in $T$ ), each $P_{i}^{2}(T) \geq 0$, each $(T-a) Q_{j}^{2}(T) \geq 0$, and $(b-T) R_{k}^{2}(T) \geq 0$. Thus, $f(T) \geq 0$, proving the claim.

Since $g(x)=\sup _{[a, b]}|f| \pm f(x) \geq 0$ on $[a, b], \sup _{[a, b]}|f| \pm f(T) \geq 0$. That is, $-\sup _{[a, b]}|f| \leq f(T) \leq$ $\sup _{[a, b]}|f|$, which gives

$$
|f(T)|_{\mathrm{op}}=\sup _{|v| \leq 1}|f(T) v| \leq \sup _{|v| \leq 1}\left|\sup _{[a, b]}\right| f|\cdot| v\left|=\left|\sup _{[a, b]}\right| f\right|
$$

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which is the desired inequality. Thus, we can extend by continuity to the sup-norm closure of $\mathbb{R}[x]$ in $C^{o}[a, b]$, which by Weierstraß is the whole $C^{o}[a, b]$, giving $C^{o}[a, b] \rightarrow \overline{\mathbb{R}}[T]$, the latter being the operator-norm closure of $\mathbb{R}[T]$, with $|f(T)|_{\text {op }} \leq|f|_{C[a, b]}$. Since $\mathbb{R}[x] \rightarrow \mathbb{R}[T]$ is a ring homomorphism, the extension by continuity is also a ring homomorphism.
[9.D.4] Corollary: (Existence of square roots of positive operators) For $T \geq 0$, there is $S \in \overline{\mathbb{R}}[x]$ such that $S \geq 0$ and $S^{2}=T$.
Proof: Since $T \geq 0$, we can take $[a, b]=[0, b]$ in the previous discussion. The function $f(x)=\sqrt{x} \in C^{o}[0, b]$ is non-negative on $[0, b]$, and $f(T)^{2}=f^{2}(T)=T$. Take $S=f(T)$.
[9.D.5] Corollary: (Positivity of products of commuting positive operators) For $S \geq 0$ and $T \geq 0$ with $S T=T S$, also $S T \geq 0$.
Proof: From the previous corollary, there is $R \in \overline{\mathbb{R}}[S]$ such that $R \geq 0$ and $R^{2}=S$. Also, $R$ commutes with $T$, by continuity. Thus,

$$
\langle S T v, v\rangle=\left\langle R^{2} T v, v\right\rangle=\langle R T R v, v\rangle=\langle T R v, R v\rangle \geq 0
$$

because $T \geq 0$.
The kernel $I$ of $C^{o}[a, b] \rightarrow \overline{\mathbb{R}[T]}$ is an ideal in $C^{o}[a, b]$, and is (topologically) closed because $C^{o}[a, b] \rightarrow \overline{\mathbb{R}[T]}$ is continuous. Let $\tau(T) \subset[a, b]$ be the simultaneous zero-set of all the functions in $I$. Shortly, we will see that $\tau(T)=\sigma(T)$, but we cannot use this yet.
[9.D.6] Claim: The restriction map $C^{o}[a, b] \rightarrow C^{o}(\tau(T))$ has kernel $I$. That is, if $\left.f\right|_{\tau(T)}=0$, then $f(T)=0$. More precisely, $f \geq 0$ on $\tau(T)$ if and only if $f(T) \geq 0$.
Proof: It suffices to show that $f(T) \geq 0$ implies $f \geq 0$ on $\tau(T)$. For $f$ not non-negative on $\tau(T)$, there is $x_{o} \in \tau(T)$ where $f\left(x_{o}\right)<0$. Using the continuity of $f$, take a small neighborhood $N$ of $x_{o}$ in $[a, b]$ such that $f(x)<0$ on $N$. Let $g \in C^{o}[a, b]$ be supported inside $N$, non-negative, and strictly positive at $x_{o}$. Then $f g \leq 0$, and $f g\left(x_{o}\right)<0$, so $-f g(T) \geq 0$. But $f(T) \geq 0$ and $g(T) \geq 0$, so by the corollary on positivity of commuting positive operators, $f g(T) \geq 0$. Thus, $f g(T)=0$, so $f g \in I$, and $\left.f g\right|_{\tau(T)}=0$, contradiction. Thus, $f \geq 0$ on $\tau(T)$. Thus, if $f=0$ on $\tau(T)$, both $f \geq 0$ and $-f \geq 0$ on $\tau(T)$, so both $f(T) \geq 0$ and $-f(T) \geq 0$, so $f(T)=0$, and $f \in I$.
[9.D.7] Corollary: $C^{o}[a, b] \rightarrow \overline{\mathbb{R}[T]}$ factors through $C^{o}(\tau(T))$, giving a commutative diagram


The induced map $C^{o}(\tau(T)) \rightarrow \overline{\mathbb{R}[T]}$ is a bijection, and $|f(T)|_{\mathrm{op}} \geq|f|_{C^{o}(\tau(T))}$.
Proof: By the Tietze-Urysohn-Brouwer extension theorem [9.E.1], every continuous function on $\tau(T)$ has an extension to a continuous function on $[a, b]$, with the same sup-norm. This gives the surjectivity of $C^{o}[a, b] \rightarrow C^{o}(\tau(T))$. By the claim, $C^{o}(\tau(T)) \approx C^{o}[a, b] / I$, giving the injectivity to $\overline{\mathbb{R}}[T]$.

Given the positivity, since $|f(T)|_{\mathrm{op}} \pm f(T) \geq 0$, from the previous claim $\mid f(T)_{\mathrm{op}} \pm f(x) \geq 0$ for $x \in \tau(T)$. Thus, $\sup _{x \in \tau(T)}|f(x)| \leq|f(T)|_{\text {op }}$.

Now a refinement of the earlier argument gives the other inequality on norms:
[9.D.8] Corollary: The induced map $C^{o}(\tau(T)) \rightarrow \overline{\mathbb{R}[T]}$ is an isometric isomorphism. That is, the map is a bijection, and $|f(T)|_{\mathrm{op}}=|f|_{C^{o}(\tau(T))}$.
Proof: For $f \geq 0$ on $\tau(T)$, again by Tietze-Urysohn-Brouwer, there is an extension $g \geq 0$ of $f$ to [a,b] with the same sup norm. The first claim of the proof showed that $|f(T)|_{\mathrm{op}} \leq|g|_{C^{\circ}[a, b]}$, so

$$
|f|_{C^{o}(\tau(T))} \leq|f(T)|_{\mathrm{op}} \leq|g|_{C^{o}[a, b]}=|f|_{C^{o}(\tau(T))}
$$

giving the isometry. In particular, for $f_{n}(T)$ a Cauchy sequence in the operator norm (for $f_{n} \in C^{o}(\tau(T))$ ), the sequence $f_{n}$ is Cauchy in $C^{o}(\tau(T))$, so converges to some $f \in C^{o}(\tau(T))$. By the isometry, $f_{n}(T) \rightarrow f(T)$, giving the surjection to the closure.

It remains to show $\tau(T)=\sigma(T)$.

## 9. Unbounded operators on Hilbert spaces

First, we reprove the fact that $\sigma(T) \subset \mathbb{R}$. For $\lambda \in \mathbb{C}$ such that there is no $(T-\lambda)^{-1}$, the polynomial $f(x)=(x-\lambda)(x-\bar{\lambda})$ is non-zero on $\mathbb{R}$, so certainly on $\tau(T)$, so has an inverse $h(x)=1 / g(x) \in C^{o}(\tau(T))$. Then $h(T)(T-\bar{\lambda})$ would be an inverse for $T-\lambda$, contradiction. Thus, $\sigma(T) \subset \mathbb{R}$.

For $\lambda$ real and not in $\tau(T), x-\lambda$ is invertible on $\tau(T)$ with inverse $h \in C^{o}(\tau(T))$, so

$$
h(T) \circ(T-\lambda)=(h \cdot(x-\lambda))(T)=1(T)=1
$$

and similarly $(T-\lambda) \circ h(T)=1$, so $T-\lambda$ is invertible. For $\lambda \in \tau(T)$, for $n>0$, let $\left.f_{n}(x) \in C^{[ } a, b\right]$ be

$$
f_{n}(x)=\left\{\begin{array}{cl}
N & \left(\text { for }|x-\lambda| \leq \frac{1}{N}\right) \\
\frac{1}{|x-\lambda|} & \left(\text { for }|x-\lambda| \geq \frac{1}{N}\right)
\end{array}\right.
$$

Thus, $\left|(x-\lambda) \cdot f_{n}\right|_{C^{o}(\tau(T))} \leq 1$, and $\left.(T-\lambda) f_{n}(T)\right|_{\mathrm{op}} \leq 1$. If $T-\lambda$ had an inverse $S$, then for all $n$ $n \leq\left|f_{n}\right|_{C^{o}(\tau(T))}=\left|f_{n}(T)\right|_{\mathrm{op}}=\left|1 \cdot f_{n}(T)\right|_{\mathrm{op}}=\left|S \cdot(T-\lambda) \cdot f_{n}(T)\right|_{\mathrm{op}} \leq|S|_{\mathrm{op}} \cdot\left|(T-\lambda) \cdot f_{n}(T)\right|_{\mathrm{op}} \leq|S|_{\mathrm{op}}$ This is impossible, so there is no inverse. This proves that $\tau(T)=\sigma(T)$.

Now we prove the peculiar lemma on polynomials:
[9.D.9] Lemma: Let $f \in \mathbb{R}[x]$ be non-negative-valued on a finite interval $[a, b]$. Then $f$ is expressible as a finite sum of the form

$$
f=\sum_{i} P_{i}^{2}+(x-a) \sum_{j} Q_{j}^{2}+(b-x) \sum_{k} R_{k}^{2}
$$

for polynomials $P_{i}, Q_{j}, R_{k}$ in $\mathbb{R}[x]$.
Proof: It suffices to consider monic $f$, since positive constants can be absorbed. Factor $f$ into irreducibles over $\mathbb{R}$, show that each of the linear and quadratic factors can be expressed in the given form, and then show that a product of such expressions can be re-written in the same form.

For quadratic irreducibles with complex-conjugate roots $z, \bar{z}$, by completing the square,

$$
(x-z)(x-\bar{z})=x^{2}-(z+\bar{z}) x+z \bar{z}=\left(x-\frac{z+\bar{z}}{2}\right)^{2}+\left(z \bar{z}-\left(\frac{z+\bar{z}}{2}\right)^{2}\right)
$$

Since

$$
z \bar{z}-\left(\frac{z+\bar{z}}{2}\right)^{2}=z \bar{z}-\frac{1}{4}\left(z^{2}+2 z \bar{z}+\bar{z}^{2}\right)=-\frac{1}{4}(z-\bar{z})^{2}=\left(\frac{z-\bar{z}}{2 i}\right)^{2}=(\operatorname{Im} z)^{2}>0
$$

we have the desired expression for $(x-z)(x-\bar{z})$.
A linear factor $x-\alpha$ with $a<\alpha<b$ must occur to an even power, since otherwise $f(x)$ would take opposite signs on the two sides of $\alpha$, contradicting the positivity of $f$ on $[a, b]$.

A linear factor $x-\alpha$ with $\alpha \leq a$ can be rewritten as

$$
x-\alpha=(x-a)+(a-\alpha)=(x-a) \cdot 1+(a-\alpha)
$$

Since $a-\alpha \geq 0$, it is a square of an element of $\mathbb{R}$, and this gives the desired expression. Similarly, a linear factor $\alpha-x$ with $\alpha \geq b$ can be rewritten as

$$
\alpha-x=(b-x)+(\alpha-b)
$$

Thus, all the factors of $f$ can be written in the desired form. As for products, we can inductively rewrite them by

$$
\begin{aligned}
P^{2} \cdot Q^{2}= & (P Q)^{2} \quad(x-a) P^{2} \cdot Q^{2}=(x-a) \cdot(P Q)^{2} \quad(x-a) P^{2} \cdot(x-a) Q^{2}=((x-a) P Q)^{2} \\
& (b-x) P^{2} \cdot Q^{2}=(b-x) \cdot(P Q)^{2} \quad(b-x) P^{2} \cdot(b-x) Q^{2}=((b-x) P Q)^{2}
\end{aligned}
$$

The only possible issue is the form $(x-a) P^{2} \cdot(b-x) Q^{2}$. By luck,

$$
(x-a)(b-x)=(x-a)(b-x) \cdot \frac{(b-x)+(x-a)}{(b-x)+(x-a)}=\frac{(x-a) \cdot(b-x)^{2}+(b-x) \cdot(x-a)^{2}}{b-a}
$$

which is of the desired form. Iterating these rewritings gives the lemma.
[9.D.10] Corollary: If $\sigma(T)=\{\lambda\}$, then $T$ is the scalar operator $\lambda$.
Proof: Because the function $f(x)=x$ restricted to $\{\lambda\}$ is equal to the restriction of the constant function $g(x)=\lambda$,

$$
T=f(T)=g(T)=\lambda
$$

meaning the scalar operator.
[9.D.11] Remark: Certainly the converse is not true: there easily can be eigenvalues imbedded in continuous spectrum.
[9.D.12] Corollary: (Schur's lemma) Let $R$ be a set of continuous linear operators on a Hilbert space $V$, and suppose $V$ is $R$-irreducible, in the sense that there is no $R$-stable closed subspace of $V$ other than $\{0\}$ and $V$ itself. Let $T$ be a self-adjoint operator commuting with all operators from $R$. Then $T$ is scalar.
Proof: Suppose that $\sigma(T)$ contains at least two distinct points $x_{1}, x_{2}$, and show that $V$ is not $R$-irreducible. Let $f, g$ be continuous functions with disjoint supports, such that $f\left(x_{1}\right)=1$ and $g\left(x_{2}\right)=1$. Thus, $f g=0$, and $f(T) g(T)=g(T) f(T)=0$, but neither $f(T)$ nor $g(T)$ is 0 , because they are not the zero function on $\sigma(T)$. The image $f(T)(V)$ is not 0 , because $f(T) \neq 0$. Also, $f(T)(V)$ is inside the kernel of $g(T)$, because $g(T) f(T)=(g f)(T)=0$. By continuity of $g(T)$, the closure $W$ of $f(T)(V)$ is also inside the kernel of $g(T)$. Since $g(T) \neq 0$, necessarily $W \neq V$.

Since $T$ commutes with all operators in $R, \mathbb{R}[T]$ commutes with $R$, and by continuity of operators in $R$, $\overline{\mathbb{R}[T]}$ commutes with $R$. Thus, $R$ commutes with $f(T)$ and $g(T)$, so for $S \in R$,

$$
S(f(T)(V))=f(T)(S V) \subset f(T)(V)
$$

That is, $R$ stabilizes $f(T)(V)$. By continuity of operators in $R, R$ stabilizes the closure $W$ of $f(T)(V)$. But $W$ is a proper closed subspace of $V$, so $V$ is not $R$-irreducible. Since $\sigma(T) \neq \phi$, it is a singleton $\{\lambda\}$. By the previous corollary, $T$ is the scalar operator $\lambda$.

Suppose that $W$ is another Hilbert space on which $R$ acts, and let

$$
\operatorname{Hom}_{R}(V, W)=\{\mathbb{C} \text {-linear maps } \varphi: V \rightarrow W \text { such that } \varphi(r \cdot v)=r \cdot \varphi(v) \text { for all } r \in R, v \in V\}
$$

In the situation of the previous corollary, let $R$ act on the orthogonal direct sum $V^{n}=\underbrace{V \oplus \ldots \oplus W}_{n}$ in the natural fashion, by

$$
r \cdot\left(v_{1}, \ldots, v_{n}\right)=\left(r v_{1}, \ldots, r v_{n}\right)
$$

[9.D.13] Corollary: $\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{R}\left(V, V^{n}\right)=n$ for $R$-irreducible $V$.
Proof: Let $p_{i}: V^{n} \rightarrow V$ be the projection to the $i^{t h}$ component. For $\varphi \in \operatorname{Hom}_{R}\left(V, V^{n}\right)$, each $p_{i} \circ \varphi: V \rightarrow V$ respects the action of $R$, so by Schur's lemma is scalar. Thus, there are scalars $c_{1}, \ldots, c_{n}$ so that

$$
\varphi(v)=\left(c_{1} \cdot v, c_{2} \cdot v, \ldots, c_{n} \cdot v\right)
$$

as claimed.
For the following, assume that $R$ has an involution $r \rightarrow r^{*}$, and that the action of $R$ on all vector spaces respects this involution: we only consider actions of $R$ on Hilbert spaces with the property that the adjoint of $v \rightarrow r \cdot v$ is $v \rightarrow r^{*} \cdot v$. [67] Also, now we only consider linear maps $V \rightarrow W$ that respect this additional structure on $R$, still referring to these as $R$-homomorphisms.

[^18][9.D.14] Corollary: Suppose that there is an injection in $\operatorname{Hom}_{R}\left(V^{n}, W\right)$. Then
$$
\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{R}(V, W) \geq n
$$

Further, if there is $n o$ injection in $\operatorname{Hom}_{R}\left(V^{n+1}, W\right)$, then

$$
\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{R}(V, W)=n
$$

Proof: Certainly if there is a copy of $V^{n}$ inside $W$, then we can map $V$ to any one of the $n$ summands, respecting the action of $R$. The converse needs Schur's lemma: suppose $\operatorname{dim}_{\mathbb{C}} \operatorname{Hom}_{R}(V, W)=n$. Let $\varphi_{1}, \ldots, \varphi_{n}$ be $n$ linearly independent homomorphisms. The image $\varphi_{1}(V)+\varphi_{2}(V)+\ldots+\varphi_{n}(V)$ need not be an orthogonal direct sum, but we claim that there is another collection of $n$ maps in $\operatorname{Hom}_{R}(V, W)$ that does produce an orthogonal direct sum inside $W$. In effect, this is a version of a Gram-Schmidt process that refers to copies of the irreducible $V$ rather than to individual vectors.

A key point is that, because of the involution $*$, the orthogonal complement $X^{\perp}$ to an $R$-stable subspace $X$ of $W$ is also $R$-stable. Indeed. For $y \in X^{\perp}$,

$$
\langle r \cdot y, x\rangle=\left\langle y, r^{*} \cdot x\right\rangle \in\langle y, X\rangle=\{0\}
$$

This immediately implies that the orthogonal projection $W \rightarrow X$ is an $R$-homomorphism.
Thus, given $\varphi_{1}$ and $\varphi_{2}$, the orthogonal projection $p$ from $\varphi_{2}(V)$ to $\varphi_{1}(V)$ is an $R$-homomorphism. Since $\varphi_{1}, \varphi_{2}$ are non-zero, $\varphi_{1}(V)$ and $\varphi_{2}(V)$ are $R$-irreducible, so $\varphi_{1}$ and $\varphi_{2}$ are $R$-isomorphisms. If the images $\varphi_{1}(V)$ and $\varphi_{2}(V)$ are orthogonal, we are done. If not, the map $p$ is not 0 , so must be an $R$-isomorphism, by $R$-irreducibility. Thus, the composition

$$
V \xrightarrow{\varphi_{1}} \varphi_{2}(V) \xrightarrow{p} \varphi_{1}(V) \xrightarrow{\varphi_{1}^{-1}} V
$$

is an $R$-isomorphism $V \rightarrow V$. By Schur's lemma, it is a non-zero constant map. That is, there is a uniform constant $c$ such that $p\left(\varphi_{2}(v)\right)=c \cdot \varphi_{1}(v)$ for all $v \in V$. That is, $c \cdot \varphi_{1}-p \circ \varphi_{2}=0$ as element of $\operatorname{Hom}_{R}(V, W)$. Then

$$
p \circ\left(c \cdot \varphi_{1}-p \circ \varphi_{2}\right)=c \cdot p \circ \varphi_{1}-p^{2} \circ \varphi_{2}=p \circ \varphi_{1}-p \circ \varphi_{2}=0
$$

so the image $\left(c \cdot \varphi_{1}-p \circ \varphi_{2}\right)(V)$ is orthogonal to $\varphi_{1}(V)$, as desired. Continue by induction to modify all $\varphi_{i}(V)$ to be mutually orthogonal.

## 9.E Appendix: Tietze-Urysohn-Brouwer extension theorem

Granting Urysohn's lemma [9.E.2], the extension result is not difficult:
[9.E.1] Theorem: For $X$ a normal space (meaning that any two disjoint closed sets have disjoint open neighborhoods), closed subset $E \subset X$, every continuous, bounded, real-valued $f$ on $E$ extends to $F$ on $X$ such that $\sup _{X}|F|=\sup _{E}|f|$.
Proof: Without loss of generality, the image of $f$ is contained in [0,1]. Urysohn's lemma [9.E.2] will be repeatedly invoked: given disjoint, closed $B_{n}, C_{n}$ in $X$, there is continuous $g_{n}$ on $X$ taking values in $\left[0, \frac{1}{2}(2 / 3)^{n}\right]$ such that $g_{n}=0$ on $B_{n}$ and $g_{n}=\frac{1}{2}(2 / 3)^{n}$ on $C_{n}$. Specify the subsets $B_{n}, C_{n}(n=1,2, \ldots)$ of $E$ inductively by

$$
B_{1}=\left\{x \in E: f(x) \leq \frac{1}{3}\right\} \quad C_{1}=\left\{x \in E: f(x) \geq \frac{2}{3}\right\}
$$

and

$$
B_{n}=\left\{x \in E: f(x)-\sum_{i=1}^{n-1} g_{i}(x) \leq \frac{2^{n-1}}{3^{n}}\right\} \quad C_{n}=\left\{x \in E: f(x)-\sum_{i=1}^{n-1} g_{i}(x) \geq \frac{2^{n}}{3^{n}}\right\}
$$

These are disjoint closed subsets of $E$, so are closed in $X$. The sum $F=\sum_{i=1}^{\infty} g_{i}$ converges uniformly, so is continuous. On $E, 0 \leq f-F \leq(2 / 3)^{n}$ for all $n$, so $F=f$ on $E$.

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[9.E.2] Theorem: (Urysohn) In a locally compact Hausdorff topological space $X$, given a compact subset $K$ contained in an open set $U$, there is a continuous function $0 \leq f \leq 1$ which is 1 on $K$ and 0 off $U$.
Proof: First, we prove that there is an open set $V$ such that

$$
K \subset V \subset \bar{V} \subset U
$$

For each $x \in K$ let $V_{x}$ be an open neighborhood of $x$ with compact closure. By compactness of $K$, some finite subcollection $V_{x_{1}}, \ldots, V_{x_{n}}$ of these $V_{x}$ cover $K$, so $K$ is contained in the open set $W=\bigcup_{i} V_{x_{i}}$ which has compact closure $\bigcup_{i} \bar{V}_{x_{i}}$ since the union is finite.

Using the compactness again in a similar fashion, for each $x$ in the closed set $X-U$ there is an open $W_{x}$ containing $K$ and a neighborhood $U_{x}$ of $x$ such that $W_{x} \cap U_{x}=\phi$.

Then

$$
\bigcap_{x \in X-U}(X-U) \cap \bar{W} \cap \bar{W}_{x}=\phi
$$

These are compact subsets in a Hausdorff space, so (again from compactness) some finite subcollection has empty intersection, say

$$
(X-U) \cap\left(\bar{W} \cap \bar{W}_{x_{1}} \cap \ldots \cap \bar{W}_{x_{n}}\right)=\phi
$$

That is,

$$
\bar{W} \cap \bar{W}_{x_{1}} \cap \ldots \cap \bar{W}_{x_{n}} \subset U
$$

Thus, the open set

$$
V=W \cap W_{x_{1}} \cap \ldots \cap W_{x_{n}}
$$

meets the requirements.
Using the possibility of inserting an open subset and its closure between any $K \subset U$ with $K$ compact and $U$ open, we inductively create opens $V_{r}$ (with compact closures) indexed by rational numbers $r$ in the interval $0 \leq r \leq 1$ such that, for $r>s$,

$$
K \subset V_{r} \subset \bar{V}_{r} \subset V_{s} \subset \bar{V}_{s} \subset U
$$

From any such configuration of opens we construct the desired continuous function $f$ by

$$
f(x)=\sup \left\{r \text { rational in }[0,1]: x \in V_{r},\right\}=\inf \left\{r \text { rational in }[0,1]: x \in \bar{V}_{r},\right\}
$$

It is not immediate that this sup and inf are the same, but if we grant their equality then we can prove the continuity of this function $f(x)$. Indeed, the sup description expresses $f$ as the supremum of characteristic functions of open sets, so $f$ is at least lower semi-continuous. [68] The inf description expresses $f$ as an infimum of characteristic functions of closed sets so is upper semi-continuous. Thus, $f$ would be continuous.

To finish the argument, we must construct the sets $V_{r}$ and prove equality of the inf and sup descriptions of the function $f$.

To construct the sets $V_{i}$, start by finding $V_{0}$ and $V_{1}$ such that

$$
K \subset V_{1} \subset \bar{V}_{1} \subset V_{0} \subset \bar{V}_{0} \subset U
$$

Fix a well-ordering $r_{1}, r_{2}, \ldots$ of the rationals in the open interval $(0,1)$. Supposing that $V_{r_{1}}, \ldots, v_{r_{n}}$ have been chosen. let $i, j$ be indices in the range $1, \ldots, n$ such that

$$
r_{j}>r_{n+1}>r_{i}
$$

[68] A (real-valued) function $f$ is lower semi-continuous when for all bounds $B$ the set $\{x: f(x)>B\}$ is open. The function $f$ is upper semi-continuous when for all bounds $B$ the set $\{x: f(x)<B\}$ is open. It is easy to show that a sup of lower semi-continuous functions is lower semi-continuous, and an inf of upper semi-continuous functions is upper semi-continuous. As expected, a function both upper and lower semi-continuous is continuous.
and $r_{j}$ is the smallest among $r_{1}, \ldots, r_{n}$ above $r_{n+1}$, while $r_{i}$ is the largest among $r_{1}, \ldots, r_{n}$ below $r_{n+1}$. Using the first observation of this argument, find $V_{r_{n+1}}$ such that

$$
V_{r_{j}} \subset \bar{V}_{r_{j}} \subset V_{r_{n+1}} \subset \bar{V}_{r_{n+1}} \subset V_{r_{i}} \subset \bar{V}_{r_{i}}
$$

This constructs the nested family of opens.
Let $f(x)$ be the sup and $g(x)$ the inf of the characteristic functions above. If $f(x)>g(x)$ then there are $r>s$ such that $x \in V_{r}$ and $x \notin \bar{V}_{s}$. But $r>s$ implies that $V_{r} \subset \bar{V}_{s}$, so this cannot happen. If $g(x)>f(x)$, then there are rationals $r>s$ such that

$$
g(x)>r>s>f(x)
$$

Then $s>f(x)$ implies that $x \notin V_{s}$, and $r<g(x)$ implies $x \in \bar{V}_{r}$. But $V_{r} \subset \bar{V}_{s}$, contradiction. Thus, $f(x)=g(x)$.

## 10. Discrete decomposition of pseudo-cuspforms

1. Compact resolvents in simplest examples
2. Compact resolvents for $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$
3. Density of domains of operators
4. Tail estimates: simplest example
5. Tail estimates: three further small examples
6. Tail estimates: $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$
7. Compact $\mathfrak{B}_{a}^{1} \longrightarrow L_{a}^{2}$ in four simple examples
8. Compact $\mathfrak{B}_{a}^{1} \longrightarrow L_{a}^{2}$ for $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$
9. Compact resolvents and discrete spectrum

Applications of idiosyncracies of Friedrichs self-adjoint extensions of restrictions of Laplace-Beltrami operators are illustrated here, as exploited in [Lax-Phillips 1976] and [ColinDeVerdière 1981/2/3], for example, and as illustrated in the next chapter. This device is essentially archimedean, related to differential operators.

On any one of the four simple unicuspidal quotients $\Gamma \backslash G / K$ of chapter 1 , the space of pseudo-cuspforms $L_{a}^{2}(\Gamma \backslash G / K)$ with cut-off height $a$ is the space of $L^{2}$ functions whose constant terms vanish above height $\eta(g)=a$. The case $a=0$ is the usual space of $L^{2}$ cuspforms. We will show that $L_{a}^{2}(\Gamma \backslash G / K)$ decomposes discretely for $\widetilde{\Delta}_{a}$, the Friedrichs extension [9.2] of the restriction $\Delta_{a}$ of $\Delta$ to the space

$$
D_{a}=C_{c}^{\infty}(\Gamma \backslash G / K) \cap L_{a}^{2}(\Gamma \backslash G / K)
$$

of test functions inside $L_{a}^{2}(\Gamma \backslash G / K)$. The proof proceeds by showing that $\widetilde{\Delta}_{a}$ has compact resolvent $\left(\widetilde{\Delta}_{a}-\lambda\right)^{-1}$ for $\lambda$ off a discrete set, and then verifying the obvious plausible bijection between the spectrum and eigenfunctions of $\widetilde{\Delta}_{a}$ and those of $\left(\widetilde{\Delta}_{a}-\lambda\right)^{-1}$ in [10.7]. Then the spectral theorem for self-adjoint compact operators [10.10] yields an orthonormal basis for $L_{a}^{2}(\Gamma \backslash G / K)$ consisting of eigenfunctions for $\widetilde{\Delta}_{a}$, with finite multiplicities.

Further, in those four examples, for $a \gg 1$, the space of pseudo-cuspforms $L_{a}^{2}(\Gamma \backslash G / K)$ includes not only cuspforms but infinitely-many $\widetilde{\Delta}_{a}$-eigenfunctions which are (necessarily) not $\Delta$-eigenfunctions. Existence of further eigenfunctions in $L_{a}^{2}(\Gamma \backslash G / K)$ is clear from the spectral decomposition of the orthogonal complement of cuspforms $L_{o}^{2}(\Gamma \backslash G / K)$ in $L^{2}(\Gamma \backslash G / K)$, in terms of integrals of Eisenstein series and residues of Eisenstein series, as in [1.12]. In [11.6], we show that all but finitely-many of the the new eigenfunctions are the truncated Eisenstein series whose constant term vanishes at $y=a$.

For the examples $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$, the notion of pseudo-cuspform is more complicated, but the general pattern of the argument is the same. Again, certain truncated Eisenstein series comprise most of the new discrete spectrum.

In all examples, the critical point is the estimate on tails of pseudo-cuspforms, in [10.3], [10.4], [10.5]. This is used to prove a Rellich-type compactness lemma, from which the compactness of the resolvent of the Friedrichs extensions follows. The seeming paradox of this discretization of the continuous spectrum of $\Delta$ on $L^{2}(\Gamma \backslash G / K)$ is essential to [Colin de Verdière 1981]'s proof of meromorphic continuation of Eisenstein series, which we recapitulate in examples in chapter 11.

### 10.1 Compact resolvents in simplest examples

The statements of the theorems are easier for the four simplest examples $\Gamma \backslash G / K$ of chapter 1 . For $a \geq 0$, consider a space of square-integrable pseudo-cuspforms including the space of cuspforms: these are functions in $L^{2}(\Gamma \backslash G / K)$ whose constant terms $c_{P} f$ vanish above height $a$ :

$$
L_{a}^{2}(\Gamma \backslash G / K)=\left\{f \in L^{2}(\Gamma \backslash G / K): c_{P} f(g)=0 \text { for } \eta(g) \geq a\right\}
$$

where the height function is $\eta\left(n m_{y} k\right)=y^{r}$ with $n \in N, m_{y}=\left(\begin{array}{cc}\sqrt{y} & 0 \\ 0 & 1 / \sqrt{y}\end{array}\right), k \in K$, and $r=1,2,3,4$ in the respective examples. As for cuspforms, vanishing of the constant term $c_{P} f$ above height $a$ can be
expressed precisely as orthogonality $\left\langle f, \Psi_{\varphi}\right\rangle=0$ to all pseudo-Eisenstein series $\Psi_{\varphi}$ with test-function data $\varphi \in C_{c}^{\infty}(0, \infty)$ supported on $[a,+\infty)$. Let

$$
D_{a}=C_{c}^{\infty}(\Gamma \backslash G / K) \cap L_{a}^{2}(\Gamma \backslash G / K)
$$

and $\Delta_{a}$ the restriction of $\Delta$ to $D_{a}$. In [10.3] we prove that $D_{a}$ is dense in $L_{a}^{2}(\Gamma \backslash G / K)$, so $\Delta_{a}$ has a Friedrichs extension $\widetilde{\Delta}_{a}$. The main result is:
[10.1.1] Theorem: $\widetilde{\Delta}_{a}$ has compact resolvent. The space $L_{a}^{2}(\Gamma \backslash G / K)$ of square-integrable pseudo-cuspforms with constant term vanishing above height $\eta(g)=a$ has an orthonormal basis of $\widetilde{\Delta}_{a}$-eigenfunctions, and eigenvalues have finite multiplicities. (Proof in [10.7].)

A seeming paradox: Of course, the space $L_{a}^{2}(\Gamma \backslash G / K)$ contains the space $L^{2}(\Gamma \backslash G / K)$ of $L^{2}$ cuspforms, for every $a \geq 0$. For $a \gg 1$, the corresponding space of pseudo-cuspforms it is demonstrably properly larger, containing part of the continuous spectrum for $\Delta$, namely, an infinite-dimensional space of pseudo-Eisenstein series $\Psi_{\varphi}$. For example, take $a^{\prime}<a$, with $a^{\prime}$ still large enough so that the Siegel set

$$
\mathfrak{S}=\mathfrak{S}_{a^{\prime}}=\left\{g \in G / K: \eta(g)>a^{\prime}\right\}
$$

has the property that $\gamma \mathfrak{S} \cap \mathfrak{S} \neq \phi$ implies $\gamma \in \Gamma \cap P$. Then, for any test function $\varphi$ supported on $\left[a^{\prime}, a\right]$, the pseudo-Eisenstein series $\Psi_{\varphi}$ is identically 0 in the region $\eta(g)>a$. From the spectral decomposition [1.12], these pseudo-Eisenstein series are integrals of Eisenstein series. Yet the Friedrichs extension $\widetilde{\Delta}_{a}$ is proven to have entirely discrete spectrum. Evidently, some part of the continuous spectrum of $\Delta$ becomes discrete for $\widetilde{\Delta}_{a}$. That is, some integrals of Eisenstein series $E_{s}$ become $L^{2}$ eigenfunctions for $\widetilde{\Delta}_{a}$ : in the four simple examples:
[10.1.2] Theorem: For $a \gg 1$, truncated Eisenstein series $\wedge^{a} E_{s}$ such that $c_{P} E_{s}(g)=0$ for $\eta(g)=a$ become $\widetilde{\Delta}_{a}$-eigenfunctions. (Proof in [11.6].)

From the theory of the constant term [8.1], the truncation $\wedge^{a} E_{s}$ is in $L^{2}$. However, $\wedge^{a} E_{s}$ is not smooth. The possibility that non-smooth functions can be eigenfunctions for $\widetilde{\Delta}_{a}$ can be understood in terms of the behavior of Friedrichs extensions, and exploited, as in [ColinDeVerdière 1981/2/3]. We give the application to meromorphic continuation of Eisenstein series in the next chapter.

Conversely, in these examples, we will show that
[10.1.3] Theorem: All non-cuspforms with $\widetilde{\Delta}_{a}$-eigenvalues $\lambda_{w}<-1 / 4$ are truncated Eisenstein series $\wedge^{a} E_{s}$ such that $c_{P} E_{s}(g)=0$ for $\eta(g)=a$. (Proof in [11.6].)

### 10.2 Compact resolvents for $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$

Now consider $\Gamma=S L_{r}(\mathbb{Z}), G=S L_{r}(\mathbb{R})$, and $K=S O_{r}(\mathbb{R})$. Again, we will prove that a certain space of square-integrable functions on $\Gamma \backslash G / K$ with all constant terms vanishing beyond fixed heights has purely discrete spectrum with respect to the Friedrichs extension of the restriction of the invariant Laplacian to (test functions in) this space.

Because we have not discussed a sufficiently general form of truncations for automorphic forms on $G L_{r}$, we cannot make as strong a statement as we might like. Namely, we will not prove that the $L^{2}$ closure of a space of test functions $D_{a}$, the initial domain for a restriction $\Delta_{a}$, is as large as we might imagine. Nevertheless, the application to meromorphic continuation of cuspidal-data Eisenstein series in [11.10], [11.12] does not need the strongest density assertions, so we will have a complete proof of that meromorphic continuation.

As in the simpler examples, the proof will proceed by showing that the resolvent of the Friedrichs extension of a restriction of the invariant Laplacian is compact. Specifically, let $A$ be the standard maximal torus consisting of diagonal real matrices, and $A^{+}$its subgroup of positive real diagonal matrices. A standard choice of positive simple roots is

$$
\Phi=\left\{\alpha_{i}(a)=\frac{a_{i}}{a_{i+1}}: i=1, \ldots, r-1\right\} \quad\left(\text { with } a=\left(\begin{array}{ccc}
a_{1} & & \\
& \ddots & \\
& & a_{r}
\end{array}\right)\right)
$$

Let $N^{\text {min }}$ be the unipotent radical of the standard minimal parabolic $P^{\text {min }}$ consisting of upper-triangular elements of $G$. For $g \in G$, let $g=n_{g} m_{g} k_{g}$ be the corresponding Iwasawa decomposition with respect to $P^{\text {min }}$, with $m_{g} \in A^{+}$. By reduction theory [3.3], there is a sufficiently small $t_{o}>0$ such that the standard Siegel set

$$
\mathfrak{S}=\mathfrak{S}_{t_{o}}=\left\{n m k: n \in N^{\min }, m \in A^{+}, k \in K, \alpha(m) \geq t_{o} \text { for all } \alpha \in \Phi\right\}
$$

satisfies $\Gamma \cdot \mathfrak{S}=G$. Fix such $\mathfrak{S}$ for the following discussion. For real $a \gg 1$, specify a subset of $\mathfrak{S}$ by

$$
Y_{a}=\{n m k \in \mathfrak{S}: \alpha(m) \geq a \text { for some } \alpha \in \Phi\}
$$

where again $n \in N^{\min }, m \in A^{+}$, and $k \in K$. Let $\Delta$ be the Casimir operator for $G$ descended to $G / K$ and to $\Gamma \backslash G / K$ as in [4.2], [4.4]. Let $\Delta_{a}$ be the restriction of the invariant Laplace-Beltrami operator $\Delta$ to the domain

$$
D_{a}=\left\{f \in C_{c}^{\infty}(\Gamma \backslash G / K): \text { for } g \in Y_{a}, c_{P} f(g)=0, \text { for all standard parabolics } P\right\}
$$

Let $V_{a}$ be the closure of $D_{a}$ in $L^{2}(\Gamma \backslash G / K)$. As usual [6.5] integration by parts shows that $\Delta_{a}$ is symmetric and non-positive, in the sense that $\langle\Delta f, f\rangle \leq 0$ for test functions $f$. Since $D_{a}$ is dense in $V_{a}$, it has a Friedrichs extension $\widetilde{\Delta}_{a}$, a self-adjoint unbounded operator on $V_{a}$.
[10.2.1] Theorem: $\widetilde{\Delta}_{a}$ has compact resolvent. The space $V_{a}$ has an orthonormal basis of $\widetilde{\Delta}_{a}$-eigenfunctions, and eigenvalues occur with finite multiplicities. (Proof in [10.8], [10.9].)

Define the $\mathfrak{B}^{1}$ norm on $D_{a}$ by

$$
|f|_{\mathfrak{B}^{1}}=\langle(1-\Delta) f, f\rangle_{L^{2}(\Gamma \backslash G / K)}
$$

and let $\mathfrak{B}^{1}$ be the completion of $D_{a}$ with respect to this norm. As in the discussion [9.2] of Friedrichs extensions, we have a natural imbedding $\mathfrak{B}^{1} \subset V_{a}$. As in the simpler examples, for sufficiently high cut-off heights $\eta$, we will see that there must be infinitely-many eigenfunctions for $\widetilde{\Delta}_{a}$ that were not eigenfunctions for $\Delta$, by exhibiting some pseudo-Eisenstein series in the $\mathfrak{B}^{1}$-closure of $D_{a}$. Specifically, we consider pseudoEisenstein series attached to maximal proper parabolics $P=P^{r, r} \subset S L_{2 r}$, with cuspidal data [3.9], with test function data supported just below the cut-off. Via reduction theory [3.3], the $P$-constant term vanishes for $\alpha(a) \geq a$. From [3.9], all other constant terms along standard parabolics are 0 . Similarly, as explicit examples of eigenfunctions for $\widetilde{\Delta}_{a}$ that are not eigenfunctions for $\Delta$, we again find certain truncated Eisenstein series. The simplest case is the following.

Let $P=P^{r, r} \subset S L_{2 r}$, and $f$ cuspidal data on the Levi component $M=M^{P}$. Let $E_{s, f}$ be the corresponding cuspidal-data Eisenstein series as in [3.11], with constant term $c_{P} E_{s, f}$ as in [3.11.9]. Let $A^{P}$ be the center of $M$, and $M^{1}$ the subgroup of $M$ consisting of matrices in $r$-by- $r$ blocks $\left(\begin{array}{ll}a & 0 \\ 0 & d\end{array}\right)$ with $\operatorname{det} a=1=\operatorname{det} d$.
[10.2.2] Theorem: For $s \in \mathbb{C}$ such that $c_{P} E_{s, f}\left(m m_{1}\right)=0$ for $m \in A^{P}$ with $\alpha_{r}(m)=a$ and for all $m_{1} \in M^{1}$, the truncation $\wedge^{a} E_{s, f}$ of $E_{s, f}$ is a $\widetilde{\Delta}_{a}$-eigenfunction in $V_{a}$. (Proof in [11.11].)

### 10.3 Density of domains of operators

For an unbounded operator to have a well-defined adjoint, its domain must be dense in the ambient Hilbert space. Of course, we could shrink the Hilbert space to be the closure of the domain of the operator, but then there would be the issue of determining that closure, apart from other complications. Test functions are dense in $L^{2}(\Gamma \backslash G / K)$ for general reasons [6.1], [14.5], [14.6]: for an approximate identity $\psi_{n}$ in $C_{c}^{\infty}(K \backslash G / K)$, the averaged action images $\psi_{n} \cdot f$ of $f \in L_{a}^{2}(\Gamma \backslash G / K)$ are smooth. However, each such averaging smears out the support of the constant term of $f$ somewhat, depending on the support of $\psi_{n}$. Let $X_{n}$ be a nested sequence of compact subsets of $\Gamma \backslash G / K$ whose union is $\Gamma \backslash G / K$, and $\alpha_{n} \in C_{c}^{\infty}(\Gamma \backslash G / K)$ identically 1 on $X_{n}$, and $0 \leq \alpha_{n}(g) \leq 1$. Thus, smoothly cut off $\psi_{n} \cdot f$ by multiplying by $\alpha_{n}$. Thus, $f_{n}=\left(\psi_{n} \cdot f\right) \cdot \alpha_{n}$ is a sequence in $C_{c}^{\infty}(\Gamma \backslash G / K)$ approaching $f$ in $L^{2}$.

However, it is not as trivial to understand the interaction with constant-term vanishing conditions. In these simple examples, density of $D_{a}$ in $L_{a}^{2}(\Gamma \backslash G / K)$ is relatively easily proven for $a \gg 1$, in which case the natural smooth cutting-off of the constant term near the given height $a$ interacts with constant-term vanishing in a controlled manner.

## 10. Discrete decomposition of pseudo-cuspforms

The condition $a \gg 1$ refers to a great-enough height $a$ so that the standard Siegel set $\mathfrak{S}_{a}$ has the property that $\mathfrak{S}_{a} \cap \gamma \mathfrak{S}_{a} \neq \phi$ implies $\gamma \in \Gamma \cap P$. By reduction theory [1.5] there exists such $a$. For the simplest case of $S L_{2}(\mathbb{Z}) \subset S L_{2}(\mathbb{R})$, the explicit bound $a>1$ suffices, for example. Again, as expected, first approximate $f \in L_{a}^{2}(\Gamma \backslash G / K)$ by functions $f_{n}$ in $C_{c}^{\infty}(\Gamma \backslash G / K)$ by general methods, and then use the condition $a \gg 1$ to consider a family of smooth cut-offs of the constant term near height $a$, with the width of the cut-off region shrinking to 0 :
[10.3.1] Lemma: For $a \gg 1, D_{a}$ is dense in $L_{a}^{2}(\Gamma \backslash G / K)$.
Proof: As just above, we take $a \gg 1$ so that the Siegel set $\mathfrak{S}_{a-\frac{1}{t}}$ meets its translates $\gamma \mathfrak{S}_{a-\frac{1}{t}}$ only for $\gamma \in \Gamma \cap P$, for all sufficiently large $t$. This allows separation of variables in $\mathfrak{S}_{a-\frac{1}{t}}$, since the cylinder $C_{a-\frac{1}{t}}=(\Gamma \cap P) \backslash \mathfrak{S}_{a-\frac{1}{t}}$ injects to $\Gamma \backslash G / K$. Let

$$
|f|_{C_{a-\frac{1}{t}}}^{2}=\int_{C_{a}-\frac{1}{t}}|f(z)|^{2} \frac{d x d y}{y^{r+1}} \leq \int_{\Gamma \backslash G / K}|f(z)|^{2} \frac{d x d y}{y^{r+1}}=|f|_{L^{2}}^{2}
$$

Let $f_{n} \in C_{c}^{\infty}(\Gamma \backslash G / K)$ with $f_{n} \rightarrow f$ in $L^{2}$. Since $f \in L_{a}^{2}(\Gamma \backslash G / K)$, we naturally expect that the constant term is not too far from that of $f$, so that smooth truncations of the constant terms of $f_{n}$ should produce functions also approaching $f$.

Use the Iwasawa coordinates $x, y$ on $G / K$ with $x \in \mathbb{R}^{r}$ and $y>0$ as in [1.3], so the height is $\eta(x, y)=y^{r}$. Let $\beta$ be a smooth function on $\mathbb{R}$ such that $\beta(y)=0$ for $y<-1,0 \leq \beta(y) \leq 1$ for $-1 \leq y \leq 0$, and $\beta(y)=1$ for $y \geq 0$. For $t>1$, put $\beta_{t}(y)=\beta(t(y-a))$, and define a smooth function on $N \backslash G / K$ by

$$
\varphi_{n, t}(x, y)=\left\{\begin{array}{cl}
\beta_{t}\left(y^{r}\right) \cdot c_{P} f_{n}(y) & \left(\text { for } y^{r} \geq a-\frac{1}{t}\right) \\
0 & \left(\text { for } y^{r}<a-\frac{1}{t}\right)
\end{array}\right.
$$

For $t>0$ large enough so that $\mathfrak{S}_{a-\frac{1}{t}}$ does not meet any of its own translates by $\gamma \in \Gamma$ except $\gamma \in \Gamma \cap P$, let $\Psi_{n, t}=\Psi_{\varphi_{n, t}}$ be the pseudo-Eisenstein series made from $\varphi_{n, t}$. The assumption on $t$ assures that in the region $y^{r}>a-\frac{1}{t}$ we have $\Psi_{n, t}=c_{P} \Psi_{n, t}=\varphi_{n, t}$. Thus, as intended, $c_{P}\left(f_{n}-\Psi_{n, t}\right)$ vanishes in $y \geq a$, so $f_{n}-\Psi_{n, t} \in L_{a}^{2}(\Gamma \backslash G / K)$.

By the triangle inequality,

$$
\mid f-\left(f_{n}-\left.\Psi_{n, t}\right|_{L^{2}} \leq\left|f-f_{n}\right|_{L^{2}}+\left|\Psi_{n, t}\right|_{L^{2}}\right.
$$

and $\left|f-f_{n}\right|_{L^{2}} \longrightarrow 0$. Thus, it suffices to show that the $L^{2}$ norm of the pseudo-Eisenstein series $\Psi_{n, t}$ goes to 0 for large $n, t$. Since $a \gg 1$,

$$
\left.\left|\Psi_{n, t}\right|_{L^{2}}=\left|\Psi_{n, t}\right|_{C_{a-\frac{1}{t}}}=\mid \varphi_{n, t}\right)_{C_{a-\frac{1}{t}}}=\left|\beta(t(y-a)) \cdot c_{P} f_{n}\right|_{C_{a-\frac{1}{t}}} \leq\left|c_{P} f_{n}\right|_{C_{a-\frac{1}{t}}}
$$

The cylinder $C_{a-\frac{1}{t}}$ admits a natural action of the product of circle groups $\mathbb{T}^{r}=(\Gamma \cap N) \backslash N$, by translation, inherited from the translation of the $x$-component in coordinates $x, y$. This induces a continuous action of $\mathbb{T}^{r}$ on $L^{2}\left(C_{a-\frac{1}{t}}\right)$ with the norm $|\cdot|_{C_{a-\frac{1}{t}}}$. Thus, the map $F \rightarrow c_{P} F$, is given by a continuous, compactly-supported $L^{2}\left(C_{a-\frac{1}{t}}\right)$-valued integrand, so from [14.1] exists as a Gelfand-Pettis integral. Thus, unsurprisingly, the restriction of $c_{P} f_{n}$ to $C_{a-\frac{1}{t}}$ goes to $c_{P} f$ in $L^{2}\left(C_{a-\frac{1}{t}}\right)$. Since $c_{P} f$ is supported in the range $\eta(g) \leq a$, and the measure of $C_{a}-C_{a-\frac{1}{t}}^{t}$ goes to 0 as $t \rightarrow+\infty$, the $C_{a-\frac{1}{t}}$-norm of $c_{P} f$ goes to 0 as $t \rightarrow+\infty$, since $c_{P} f$ is locally integrable.

Thus, for example, $\Psi_{n, n}$ goes to 0 in $L^{2}$ norm, so the elements $f_{n}-\Psi_{n, n}$ in $L_{a}^{2}$ go to $f$ in $L^{2}$ norm, proving the density of $D_{a}$ in $L_{a}^{2}$.

### 10.4 Tail estimates: simplest example

The computation for $\Gamma=S L_{2}(\mathbb{Z}), G=S L_{2}(\mathbb{R})$, and $K=S O_{2}(\mathbb{R})$ can take advantage of some convenient technical coincidences. Let the $\mathfrak{B}^{1}$ norm be defined on test functions $C_{c}^{\infty}(\Gamma \backslash G / K)$ by

$$
|f|_{\mathfrak{B}^{1}}^{2}=\langle(1-\Delta) f, f\rangle=\langle f, f\rangle+\langle(-\Delta) f, f\rangle
$$

Let $\mathfrak{B}^{1}$ be the completion of $C_{c}^{\infty}(\Gamma \backslash G / K)$ with respect to the $\mathfrak{B}^{1}$ norm. With

$$
D_{a}=C_{c}^{\infty}(\Gamma \backslash G / K) \cap L_{a}^{2}(\Gamma \backslash G / K)
$$

let $\mathfrak{B}_{a}^{1}$ be the $\mathfrak{B}^{1}$-completion of $D_{a}$. Note that while it is clear that $\mathfrak{B}_{a}^{1} \subset \mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$, it is not clear that equality holds. We do not need to address this for the moment. As in [10.1], let $\Delta_{a}$ be the restriction of $\Delta$ to $D_{a}$, and $\widetilde{\Delta}_{a}$ its Friedrichs extension. The Friedrichs extension $\widetilde{\Delta}_{a}$ maps from $L_{a}^{2}(\Gamma \backslash G / K)$ to $\mathfrak{B}_{a}^{1}$. Let $B$ be the unit ball in $\mathfrak{B}_{a}^{1}$. As in all cases, the crucial estimate is
[10.4.1] Claim: Given $\varepsilon>0$, a cut-off $c \geq a$ can be made sufficiently large so that the image of $B$ in $L^{2}(\Gamma \backslash G / K)$ cut off at height $c$ lies in a single $\varepsilon$-ball in $L^{2}(\Gamma \backslash G / K)$. That is, for $f \in \mathfrak{B}_{a}^{1}$,

$$
\lim _{c \rightarrow \infty} \int_{y>c}|f(z)|^{2} \frac{d x d y}{y^{2}} \longrightarrow 0 \quad \quad \text { (uniformly for }|f|_{\mathfrak{B}_{a}^{1}}=|f|_{\mathfrak{B}^{1}} \leq 1 \text { ) }
$$

[10.4.2] Remark: To be careful, we note that the inequality of the claim does not directly address the issue of smooth truncations of $f$ in $\mathfrak{B}_{a}^{1}$ near height $c$, nor whether a collection of smooth truncations $\varphi_{\infty} \cdot f$ of all heights $c \gg a$ can be chosen with $\mathfrak{B}^{1}$-norms uniformly bounded for $f \in B$. These somewhat secondary points are addressed just below in [10.4.3]: nothing surprising happens.
Proof: This computation roughly follows [Lax-Phillips 1976], pages 204-6. To legitimize the following computation, we should be sure that $f \in \mathfrak{B}^{1}$ has first derivatives in an $L^{2}$ sense. This is a local fact, and thus follows from the discussion on tori $\mathbb{T}^{n}$ in [9.5].

Let the Fourier coefficients of $f(x+i y)$ be $\widehat{f}(n)$, functions of $y$. Take $c>a$ so that the $0^{t h}$ Fourier coefficient $\widehat{f}(0)$ vanishes identically. Use Plancherel for the Fourier expansion in $x$, and then elementary inequalities: integrating over the part of $Y_{\infty}$ above $y=c$, letting $\mathcal{F}$ be Fourier transform in $x$,

$$
\begin{gathered}
\iint_{y>c}|f|^{2} \frac{d x d y}{y^{2}} \leq \frac{1}{c^{2}} \iint_{y>c}|f|^{2} d x d y=\frac{1}{c^{2}} \sum_{n \neq 0} \int_{y>c}|\widehat{f}(n)|^{2} d y \\
\leq \frac{1}{c^{2}} \sum_{n \neq 0}(2 \pi n)^{2} \int_{y>c}|\widehat{f}(n)|^{2} d y=\frac{1}{c^{2}} \sum_{n \neq 0} \int_{y>c}\left|\mathcal{F} \frac{\partial f}{\partial x}(n)\right|^{2} d y=\frac{1}{c^{2}} \iint_{y>c}\left|\frac{\partial f}{\partial x}\right|^{2} d x d y \\
=\frac{1}{c^{2}} \iint_{y>c}-\frac{\partial^{2} f}{\partial x^{2}} \cdot \bar{f}(x) d x d y \leq \frac{1}{c^{2}} \iint_{y>c}-\frac{\partial^{2} f}{\partial x^{2}} \cdot \bar{f}(x)-\frac{\partial^{2} f}{\partial y^{2}} \cdot \bar{f}(x) d x d y \\
=\frac{1}{c^{2}} \iint_{y>c}-\Delta f \cdot \bar{f} \frac{d x d y}{y^{2}} \leq \frac{1}{c^{2}} \iint_{\Gamma \backslash G / K}-\Delta f \cdot \bar{f} \frac{d x d y}{y^{2}} \leq \frac{1}{c^{2}} \cdot|f|_{\mathfrak{B}^{1}}^{2} \leq \frac{1}{c^{2}}
\end{gathered}
$$

giving the uniform bound as claimed.
Now we prove the reassuring lemma that the $\mathfrak{B}^{1}$-norms of systematically specified families of smooth tails of functions in $\mathfrak{B}^{1}$ are uniformly dominated by the $\mathfrak{B}^{1}$-norms of the original functions. Let $\varphi$ be a smooth real-valued function on $(0,+\infty)$ with

$$
\left\{\begin{array}{cl}
\varphi(y)=0 & (\text { for } 0<y \leq 1) \\
0 \leq \varphi(y) \leq 1 & (\text { for } 1<y<2) \\
1 \leq \varphi(y) & (\text { for } 1 \leq y)
\end{array}\right.
$$

[10.4.3] Claim: For fixed $\eta$, for $t \geq 1$, the smoothly cut-off tail $f^{[t]}(x+i y)=\varphi\left(\frac{y}{t}\right) \cdot f(x+i y)$ has $\mathfrak{B}^{1}$-norm dominated by that of $f$ itself:

$$
\left|f^{[t]}\right|_{\mathfrak{B}^{1}} \ll \varphi_{\varphi}|f|_{\mathfrak{B}^{1}} \quad \quad \text { (implied constant independent of } f \text { and of } t \geq 1 \text { ) }
$$

Proof: This is essentially elementary. Since $|a+b i|^{2}=a^{2}+b^{2}$ and $\Delta$ has real coefficients, it suffices to treat real-valued $f$. Since $0 \leq \varphi \leq 1$, certainly $|\varphi f|_{L^{2}} \leq|f|_{L^{2}}$. For the other part of the $\mathfrak{B}^{1}$-norm, letting $S^{1} \approx \mathbb{R} / \mathbb{Z}$ be the circle,

$$
\begin{gathered}
\left\langle-\Delta f^{[t]}, f^{[t]}\right\rangle=-\int_{S^{1}} \int_{y \geq t}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) f^{[t]} \cdot f^{[t]} d x d y \\
=-\int_{S^{1}} \int_{y \geq t} \varphi^{2}\left(\frac{y}{t}\right) f_{x x} f+\frac{1}{t^{2}} \varphi^{\prime \prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right) f^{2}+\frac{2}{t} \varphi^{\prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right) f_{y} f+\varphi\left(\frac{y}{t}\right)^{2} f_{y y} f d x d y
\end{gathered}
$$

Some terms are easy to estimate: using the fact that $\varphi^{\prime}$ and $\varphi^{\prime \prime}$ are supported on $[1,2]$,

$$
\begin{gathered}
\int_{S^{1}} \int_{y \geq t}-\varphi\left(\frac{y}{t}\right)^{2} f_{x x} f+\left|\frac{1}{t^{2}} \varphi^{\prime \prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right) f^{2}\right|-\varphi\left(\frac{y}{t}\right)^{2} f_{y y} f d x d y \\
<_{\varphi} \int_{S^{1}} \int_{t \leq y \leq 2 t} \frac{f^{2}}{t^{2}}-\left(f_{x x} f+f_{y y} f\right) d x d y \leq \int_{S^{1}} \int_{t \leq y \leq 2 t} \frac{(2 t)^{2} f^{2}}{t^{2}}-y^{2}\left(f_{x x}+f_{y y}\right) f \frac{d x d y}{y^{2}} \\
\leq 4|f|_{L^{2}}^{2}+\int_{\Gamma \backslash G / K}(-\Delta) f \cdot f \frac{d x d y}{y^{2}} \ll|f|_{\mathfrak{B}^{1}}^{2}
\end{gathered}
$$

with uniform implied constants. Transform the remaining term by integration by parts:

$$
\begin{aligned}
\int_{S^{1}} \int_{y \geq t} \frac{2}{t} \varphi^{\prime}\left(\frac{y}{t}\right) & \varphi\left(\frac{y}{t}\right) f_{y} f d x d y=\int_{S^{1}} \int_{t \leq y \leq 2 t} \frac{1}{t} \varphi^{\prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right) \cdot \frac{\partial}{\partial y}\left(f^{2}\right) d x d y \\
& =\int_{S^{1}} \int_{t \leq y \leq 2 t} \frac{\partial}{\partial y}\left(\frac{1}{t} \varphi^{\prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right)\right) \cdot f^{2} d x d y
\end{aligned}
$$

This is dominated by

$$
\begin{gathered}
\int_{S^{1}} \int_{t \leq y \leq 2 t}\left|\frac{\partial}{\partial y}\left(\frac{1}{t} \varphi^{\prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right)\right)\right| \cdot f^{2} d x d y \leq \int_{S^{1}} \int_{t \leq y \leq 2 t}\left|\frac{\partial}{\partial y}\left(\frac{1}{t} \varphi^{\prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right)\right)\right| \cdot f^{2} \cdot(2 t)^{2} \frac{d x d y}{y^{2}} \\
=4 \int_{S^{1}} \int_{t \leq y \leq 2 t}\left|\varphi^{\prime \prime}\left(\frac{y}{t}\right) \varphi\left(\frac{y}{t}\right)+\varphi^{\prime}\left(\frac{y}{t}\right)^{2}\right| \cdot f^{2} \frac{d x d y}{y^{2}} \ll \varphi|f|_{L^{2}}^{2}
\end{gathered}
$$

with implied constant independent of $f$ and $t \geq 1$.

### 10.5 Tail estimates: three further small examples

Now we see how to adapt the previous argument to the other three examples from chapter 1. Most of the work involves skirting the needless (but convenient) exploitation of coincidences used in that simplest example: the $y^{2}$ in the coordinate expression for the invariant Laplacian and in the invariant measure in the $S L_{2}(\mathbb{R})$ seem to need to cancel to make the computation succeed. Our main point in this section is seeing that that coincidence is irrelevant.

For all four examples from chapter 1, use the coordinates and conventions there. The coordinates are the Iwasawa coordinates $x, y$ with $x \in \mathbb{R}^{r}$ for $r=1,2,3,4$ (with the previous section's example being the case $r=1$ ), and $0<y \in \mathbb{R}$. The invariant Laplacian is

$$
\Delta=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)-(r-1) y \frac{\partial}{\partial y}
$$

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and the invariant measure is $d x d y / y^{r+1}$. The [Lax-Phillips 1976] argument as in [10.4] requires not only that $-\Delta$ itself is non-negative, but that the two natural summands of $-\Delta$ in Iwasawa coordinates are both non-negative. For example, the seemingly-extra term $(r-1) y \frac{\partial}{\partial y}$ is not only harmless, but necessary, possibly contrary to a visual appraisal. The $-y^{2} \frac{\partial^{2}}{\partial x^{2}}$ summand is non-negative, because the partial derivative in $x$ does not interact with either the leading coefficient $y^{2}$ or the denominator $y^{n}$ in the measure.
[10.5.1] Claim:

$$
\int-\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-(r-1) y \frac{\partial}{\partial y}\right) f \cdot \bar{f} \frac{d x d y}{y^{r+1}} \geq 0
$$

Proof: Integrating by parts once in the second-order derivative,

$$
\begin{gathered}
\int-y^{2} \frac{\partial^{2}}{\partial y^{2}} f \cdot \bar{f} \frac{d x d y}{y^{r+1}}=\int-\frac{\partial^{2}}{\partial y^{2}} f \cdot y^{1-r} \bar{f} d x d y=\int \frac{\partial}{\partial y} f \cdot \frac{\partial}{\partial y}\left(y^{1-r} \bar{f}\right) d x d y \\
=\int \frac{\partial}{\partial y} f \cdot\left((1-r) y^{-r} \bar{f}+y^{1-r} \frac{\partial}{\partial y} \bar{f}\right) d x d y
\end{gathered}
$$

The $\frac{\partial}{\partial y} f \cdot(1-r) y^{-r} \bar{f}$ term cancels the corresponding term in the original expression, so

$$
\int-\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-(r-1) y \frac{\partial}{\partial y}\right) f \cdot \bar{f} \frac{d x d y}{y^{r+1}}=\int y \frac{\partial}{\partial y} f \cdot y \frac{\partial}{\partial y} \bar{f} \frac{d x d y}{y^{r+1}}
$$

Thus, for example, with invariant Laplacian $\Delta$,

$$
\int-\Delta f \cdot \bar{f} \frac{d x d y}{y^{r+1}}=\int\left(y \frac{\partial f}{\partial x}\right)^{2}+\left(y \frac{\partial f}{\partial y}\right)^{2} \frac{d x d y}{y^{r+1}}
$$

This is the desired positiviity.
We grant ourselves that the subordinate issue about uniform estimates on families of smooth cut-offs is resolved, as in [10.4.3]. Let $\xi$ run over characters of $(\Gamma \cap N) \backslash N \approx \mathbb{T}^{r}$, and take $c \geq c_{o} \gg 1$. In Iwasawa coordinates $x, y$, write the Fourier expansion in $x$ as

$$
f(x, y)=\sum_{\xi} \widehat{f}(\xi)(y)
$$

Toward the compactness of $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}(\Gamma \backslash G / K)$, the critical point is the tail estimate:
[10.5.2] Claim: For smooth $f$ with support in $y \geq c \geq c_{o} \gg 1$,

$$
\int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}|f|^{2} \frac{d x d y}{y^{r+1}} \ll \frac{1}{c^{2}} \cdot|f|_{\mathfrak{B}^{1}}^{2}
$$

with implied constants independent of $f$ and of $c$.
Proof: By Plancherel in $x$,

$$
\int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}|f|^{2} \frac{d x d y}{y^{r+1}}=\sum_{\xi} \int_{y \geq c}|\widehat{f}(\xi)(y)|^{2} \frac{d y}{y^{r+1}}
$$

When $\widehat{f}(0)(y)=0$ for $y \geq c_{o}$, since $|\xi| \gg 1$ for $\xi \neq 0$,

$$
\sum_{\xi} \int_{y \geq c}|\widehat{f}(\xi)(y)|^{2} \frac{d y}{y^{r+1}} \ll \sum_{\xi} \int_{y \geq c}|\xi|^{2} \cdot|\widehat{f}(\xi)(y)|^{2} \frac{d y}{y^{r+1}}
$$

With $\Delta_{x}$ the Euclidean Laplacian in $x$,

$$
|\xi|^{2} \cdot \widehat{f}(\xi, y)=\frac{1}{4 \pi^{2}}\left(-\Delta_{x} f\right)^{\wedge}(\xi)(y)
$$

Applying this, and going back by Plancherel,

$$
\sum_{\xi} \int_{y \geq c}|\xi|^{2} \cdot|\widehat{f}(\xi)(y)|^{2} \frac{d y}{y^{r+1}} \ll \sum_{\xi} \int_{y \geq c}\left(-\Delta_{x} f\right)^{\wedge}(\xi)(y) \cdot \bar{f}(\xi)(y) \frac{d y}{y^{r+1}}=\int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}-\Delta_{x} f \cdot \bar{f} \frac{d x d y}{y^{r+1}}
$$

Since $y \geq c \geq c_{o} \gg 1$,

$$
\int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}-\Delta_{x} f \cdot \bar{f} \frac{d x d y}{y^{r+1}} \leq \frac{1}{c^{2}} \int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}-y^{2} \Delta_{x} f \cdot \bar{f} \frac{d x d y}{y^{r+1}}
$$

From the positivity result just above,

$$
\int-\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-(r-1) y \frac{\partial}{\partial y}\right) f \cdot \bar{f} \frac{d x d y}{y^{r+1}} \geq 0
$$

so

$$
\int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}-y^{2} \Delta_{x} f \cdot \bar{f} \frac{d x d y}{y^{r+1}} \leq \frac{1}{c^{2}} \int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}\left(-y^{2} \Delta_{x} f-y^{2} \frac{\partial^{2} f}{\partial y^{2}}+(r-1) y \frac{\partial f}{\partial y}\right) \cdot \bar{f} \frac{d x d y}{y^{r+1}}
$$

Thus, for smooth $f$ with support in $y \geq c \geq c_{o}$,

$$
\frac{1}{c^{2}} \int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}|f|^{2} \frac{d x d y}{y^{r+1}} \ll \frac{1}{c^{2}} \int_{(N \cap \Gamma) \backslash N} \int_{y \geq c}-\Delta f \cdot \bar{f} \frac{d x d y}{y^{r+1}} \leq \frac{1}{c^{2}} \cdot|f|_{\mathfrak{B}^{1}}^{2}
$$

as claimed.

### 10.6 Tail estimate: $S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots$

As in the smaller examples, the global automorphic Sobolev space $\mathfrak{B}^{1}$ is the completion of $C_{c}^{\infty}(\Gamma \backslash G / K)$ with respect to the $\mathfrak{B}^{1}$-norm

$$
|f|_{\mathfrak{B}^{1}}^{2}=\int_{\Gamma \backslash G / K}(1-\Delta) f \cdot \bar{f}
$$

Let $\mathfrak{S}$ be a sufficiently large standard Siegel set so that it surjects to the quotient $\Gamma \backslash G$. For $a>0$, the set

$$
X_{a}=\left\{g \in \mathfrak{S}: \alpha\left(m_{g}\right) \leq a \text { for all } \alpha \in \Phi\right\}
$$

has compact image $\Gamma \backslash\left(\Gamma \cdot X_{a}\right)$. For $\alpha$ in the set $\Phi$ of simple roots and $c \geq a$, let

$$
Y_{c}^{\alpha}=\left\{g \in \mathfrak{S}: \alpha\left(m_{g}\right) \geq a\right\}
$$

and $Y_{a}=\bigcup_{\alpha} Y_{a}^{\alpha}$. Certainly

$$
\mathfrak{S}=X_{a} \cup Y_{a}
$$

Let

$$
L_{a}^{2}(\Gamma \backslash G / K)=\left\{f \in L^{2}(\Gamma \backslash G / K): c_{P} f(g)=0, \text { for all } g \in Y_{a}, \text { for all parabolics } P\right\}
$$

and

$$
D_{a}=D \cap L_{a}^{2}(\Gamma \backslash G / K)
$$

It suffices to require the constant-term vanishing just for standard maximal proper parabolics, because $c_{Q \cap P} f=c_{Q} c_{P} f$ for two standard parabolics, and every standard parabolic is an intersection of maximal ones. Let $\mathfrak{B}_{a}^{1}$ be the $\mathfrak{B}^{1}$ completion of $C_{c}^{\infty}(Z \Gamma \backslash G / K) \cap L_{a}^{2}$ in the $\mathfrak{B}^{1}$ norm.

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To eventually show that the injection $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}$ is compact, as in the simpler examples, we will show that the image of the unit ball of $\mathfrak{B}_{a}^{1}$ is totally bounded in $L_{a}^{2}$. The crucial point is an estimate on the tails of functions in the unit ball $B$ in $\mathfrak{B}_{a}^{1}$, as follows.

We grant ourselves a suitable analogue of [10.4.3], that we can control the $\mathfrak{B}^{1}$-norm of smoothly cut-off versions of $f \in B$ when any single simple root becomes large:
[10.6.1] Lemma: Fix a positive simple root $\alpha$. Given $c \geq a+1$, there are real-valued smooth functions $\varphi_{o}$ and $\varphi_{1}$, taking values in $[0,1]$, summing to 1 , such that $\varphi_{1}$ is supported in $Y_{c}^{\alpha}$, and so that there is a bound $C$ uniform in $c \geq a+1$, such that $\left|f \cdot \varphi_{1}\right|_{\mathfrak{B}^{1}} \leq C \cdot|f|_{\mathfrak{B}^{1}}$.

The key point is a bound going to 0 when any simple root $\alpha$ becomes large:
[10.6.2] Claim: For $\alpha \in \Phi$,

$$
\lim _{c \rightarrow+\infty}\left(\sup _{f \in \mathfrak{B}_{a}^{1} \text { and } \operatorname{spt} f \subset Y_{c}^{\alpha}} \frac{|f|_{L^{2}}}{|f|_{\mathfrak{B}^{1}}}\right)=0
$$

Proof: Fix $\alpha=\alpha_{i} \in \Phi$, and $f \in \mathfrak{B}_{c}^{1}$ with support inside $Y_{c}^{\alpha}$ for $c \gg a$. Let $N=N^{\alpha}, P=P^{\alpha}$, and let $M=M^{\alpha}$ be the standard Levi component of $P$. Use exponential coordinates coordinates

$$
n_{x}=\left(\begin{array}{cc}
1_{i} & x \\
0 & 1_{r-i}
\end{array}\right)
$$

In effect, the coordinate $x$ is in the Lie algebra $\mathfrak{n}$ of $N$. Let $\Lambda \subset \mathfrak{n}$ be the lattice which exponentiates to $N \cap P$. Give $\mathfrak{n}$ the natural inner product $\langle$,$\rangle invariant under the (Adjoint) action of M \cap K$ that makes root spaces mutually orthogonal. Fix a non-trivial character $\psi$ on $\mathbb{R} / \mathbb{Z}$. We have the Fourier expansion

$$
\left.f\left(n_{x} m\right)=\sum_{\xi \in \Lambda^{\prime}} \psi\langle x, \xi\rangle \widehat{f}_{\xi}(m) \quad \text { (with } n \in N \text { and } m \in M\right)
$$

where $\Lambda^{\prime}$ is the dual lattice to $\Lambda$ in $\mathfrak{n}$ with respect to $\langle$,$\rangle , and$

$$
\widehat{f}_{\xi}(m)=\int_{\mathfrak{n} / \Lambda} \bar{\psi}\langle x, \xi\rangle f\left(n_{x} m\right) d x
$$

Let $\Delta^{\mathfrak{n}}$ be the flat Laplacian on $\mathfrak{n}$ associated to the inner product $\langle$,$\rangle , normalized so that$

$$
\Delta^{\mathfrak{n}} \psi\langle x, \xi\rangle=-\langle\xi, \xi\rangle \cdot \psi\langle x, \xi\rangle
$$

Let $U=M \cap N^{\min }$, and $M_{\mathfrak{S}}=M \cap \mathfrak{S}$. Abbreviating $A_{u}=\mathrm{Ad} u$,

$$
|f|_{L^{2}}^{2} \leq \int_{\mathfrak{S}}|f|^{2}=\int_{M_{\mathfrak{S}}} \int_{(U \cap \Gamma) \backslash U} \int_{A_{u}^{-1} \Lambda \backslash \mathfrak{n}}\left|f\left(u n_{x} m\right)\right|^{2} d x d u \frac{d m}{\delta(m)}
$$

with Haar measures $d x, d u, d m$, where $\delta$ is the modular function of $P$. Using the Fourier expansion,

$$
f\left(u n_{x} m\right)=f\left(u n_{x} u^{-1} \cdot u m\right)=\sum_{\xi \in \Lambda^{\prime}} \psi\left\langle A_{u} x, \xi\right\rangle \cdot \widehat{f}_{\xi}(u m)=\sum_{\xi \in \Lambda^{\prime}} \psi\left\langle x, A_{u}^{*} \xi\right\rangle \cdot \widehat{f}_{\xi}(u m)
$$

Then

$$
-\Delta^{\mathfrak{n}} f\left(u n_{x} m\right)=\sum_{\xi \in \Lambda^{\prime}}\left\langle A_{u}^{*} \xi, A_{u}^{*} \xi\right\rangle \cdot \psi\left\langle x, A_{u}^{*} \xi\right\rangle \cdot \widehat{f}_{\xi}(u m)
$$

The compact quotient $(U \cap \Gamma) \backslash U$ has a compact set $R$ of representatives in $U$, so there is a uniform lower bound for $0 \neq \xi \in \Lambda^{\prime}$ :

$$
0<b \leq \inf _{u \in R} \inf _{0 \neq \xi \in \Lambda^{\prime}}\left\langle A_{u}^{*} \xi, A_{u}^{*} \xi\right\rangle
$$

By Plancherel applied to the Fourier expansion in $x$, using the hypothesis that $\widehat{f_{0}}=0$ in $X_{a}^{\alpha}$,

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$$
\begin{gathered}
\int_{A_{u}^{-1} \Lambda \backslash \mathfrak{n}}\left|f\left(u n_{x} m\right)\right|^{2} d x=\int_{A_{u}^{-1} \Lambda \backslash \mathfrak{n}}\left|f\left(u n_{x} u^{-1} \cdot u m\right)\right|^{2} d x=\sum_{\xi \in \Lambda^{\prime}}\left|\widehat{f}_{\xi}(u m)\right|^{2} \\
\leq b^{-1} \sum_{\xi \in \Lambda^{\prime}}\left\langle A_{u}^{*} \xi, A_{u}^{*} \xi\right\rangle \cdot\left|\widehat{f}_{\xi}(u m)\right|^{2}=\sum_{\xi \in \Lambda^{\prime}}-{\widehat{\Delta^{\mathfrak{n}}} f_{\xi}(u m) \cdot \bar{f}(u m)}_{=\int_{u^{-1} \Lambda u \backslash \mathfrak{n}}-\Delta^{\mathfrak{n}} f\left(u n_{x} u^{-1} \cdot u m\right) \cdot \bar{f}\left(u n_{x} u^{-1} \cdot u m\right) d x=\int_{A_{u}^{-1} \Lambda \backslash \mathfrak{n}}-\Delta^{\mathfrak{n}} f\left(u n_{x} m\right) \cdot \bar{f}\left(u n_{x} m\right) d x}
\end{gathered}
$$

Thus, for $f$ with $\widehat{f}(0)=0$ on $Y_{a}^{\alpha}$,

$$
|f|_{L^{2}}^{2} \ll \int_{M_{\mathfrak{G}}} \int_{(U \cap \Gamma) \backslash U} \int_{A_{u}^{-1} \Lambda \backslash \mathfrak{n}}-\Delta^{\mathfrak{n}} f\left(u n_{x} m\right) \cdot \bar{f}\left(u n_{x} m\right) d x d u \frac{d m}{\delta(m)}
$$

Next, we compare $\Delta^{\mathfrak{n}}$ to the invariant Laplacian $\Delta$. Let $\mathfrak{g}$ be the Lie algebra of $G_{\mathbb{R}}$, with non-degenerate invariant pairing $\langle u, v\rangle=\operatorname{tr}(u v)$. The Cartan involution $v \rightarrow v^{\theta}=-v^{\top}$ has +1 eigenspace the Lie algebra $\mathfrak{k}$ of $K$, and -1 eigenspace $\mathfrak{s}$, the space of symmetric matrices.

Let $\Phi^{N}$ be the set of positive roots $\beta$ whose root-space $\mathfrak{g}_{\beta}$ appears in $\mathfrak{n}$. For each $\beta \in \Phi^{N}$, take $x_{\beta} \in \mathfrak{g}_{\beta}$ such that $x_{\beta}+x_{\beta}^{\theta} \in \mathfrak{s}, x_{\beta}-x_{\beta}^{\theta} \in \mathfrak{k}$, and $\left\langle x_{\beta}, x_{\beta}^{\theta}\right\rangle=-1$ : for $\beta(a)=a_{i} / a_{j}$ with $i<j, x_{\beta}$ has a single non-zero entry, at the $i j^{t h}$ place. Let

$$
\Omega^{\prime}=-\sum_{\beta \in \Phi^{N}}\left(x_{\beta} x_{\beta}^{\theta}+x_{\beta}^{\theta} x_{\beta}\right) \quad \text { (in the universal enveloping algebra } U \mathfrak{g} \text { ) }
$$

Let $\Omega^{\prime \prime} \in U \mathfrak{g}$ be the Casimir element for the Lie algebra $\mathfrak{m}$ of $M_{\mathbb{R}}$, normalized so that Casimir $\Omega$ for $\mathfrak{g}$ is the $\operatorname{sum} \Omega=\Omega^{\prime}+\Omega^{\prime \prime}$. We rewrite $\Omega^{\prime}$ to fit the Iwasawa coordinates: for each $\beta$,

$$
x_{\beta} x_{\beta}^{\theta}+x_{\beta}^{\theta} x_{\beta}=-2 x_{\beta}^{2}+2 x_{\beta}\left(x_{\beta}+x_{\beta}^{\theta}\right)+\left[x_{\beta}^{\theta}, x_{\beta}\right] \in-2 x_{\beta}^{2}+\left[x_{\beta}^{\theta}, x_{\beta}\right]+\mathfrak{k}
$$

Thus,

$$
\left.\Omega^{\prime}=\sum_{\beta \in \Phi^{N}} 2 x_{\beta}^{2}-\left[x_{\beta}^{\theta}, x_{\beta}\right] \quad \quad \text { (modulo } \mathfrak{k}\right)
$$

The commutators $\left[x_{\beta}^{\theta}, x_{\beta}\right]$ are in $\mathfrak{m}$. In the coordinates $u n_{x} a$ with $U \mathfrak{g}$ acting on the right, $x_{\beta} \in \mathfrak{n}$ is acted on by $a$ before translating $x$, by

$$
u n_{x} a \cdot e^{t x_{\beta}}=u n_{x} \cdot e^{t \beta(a) \cdot x_{\beta}} \cdot a=u n_{x+\beta(a) x_{\beta}} a
$$

That is, $x_{\beta}$ acts by $\beta(a) \cdot \frac{\partial}{\partial x_{\beta}}$.
For two symmetric operators $S, T$ on a not-necessarily-complete inner product space $V$, write $S \leq T$ when

$$
\langle S v, v\rangle \leq\langle T v, v\rangle \quad(\text { for all } v \in V)
$$

Say a symmetric operator $T$ is non-negative when $0 \leq T$. Since $m \in M_{\mathfrak{S}}$, there is an absolute constant so that $\alpha(m) \geq c$ implies $\beta(m) \gg c$. Thus,

$$
-\Delta^{\mathfrak{n}}=-\sum_{\beta \in \Phi^{N}} \frac{\partial^{2}}{\partial x_{\beta}^{2}} \leq \frac{1}{c^{2}} \cdot\left(-\sum_{\beta \in \Phi^{N}} x_{\beta}^{2}\right) \quad \text { (operators on } C_{c}^{\infty}\left(Y_{a}^{\alpha}\right)^{K} \text { ) }
$$

where $C_{c}^{\infty}\left(Y_{a}^{\alpha}\right)$ has the $L^{2}$ inner product. We claim that

$$
\sum_{\beta \in \Phi^{N}}\left[x_{\beta}^{\theta}, x_{\beta}\right]-\Omega^{\prime \prime} \geq 0 \quad \quad \text { (operators on } C_{c}^{\infty}\left(Y_{a}^{\alpha}\right) \text { ) }
$$

From this, it would follow that

$$
-\Delta^{\mathfrak{n}} \ll \frac{1}{c^{2}} \cdot\left(-\sum_{\beta \in \Phi^{N}} x_{\beta}^{2}\right) \leq \frac{1}{c^{2}} \cdot\left(-\sum_{\beta \in \Phi^{N}} x_{\beta}^{2}+\sum_{\beta \in \Phi^{N}}\left[x_{\beta}^{\theta}, x_{\beta}\right]-\Omega^{\prime \prime}\right)=\frac{1}{c^{2}} \cdot(-\Delta)
$$

Then for $f \in \mathfrak{B}_{a}^{1}$ with support in $X_{a}^{\alpha}$ we would have

$$
|f|_{L^{2}}^{2} \ll \int_{\mathfrak{S}}-\Delta^{\mathfrak{n}} f \cdot \bar{f} \ll \frac{1}{c^{2}} \int_{\mathfrak{S}}-\Delta f \cdot \bar{f} \ll \frac{1}{c^{2}} \int_{\Gamma \backslash G}-\Delta f \cdot \bar{f} \ll \frac{1}{c^{2}} \cdot|f|_{\mathfrak{B}^{1}}^{2}
$$

Taking $c$ large makes this small.
To prove the claimed non-negativity of $T=\sum_{\beta \in \Phi^{N}}\left[x_{\beta}^{\theta}, x_{\beta}\right]-\Omega^{\prime \prime}$, exploit the Fourier expansion along $N$ and the fact that $x \in \mathfrak{n}$ does not appear in $T$ : noting that the order of coordinates $n_{x} u$ differs from that above,

$$
\begin{gathered}
\int_{M_{\mathfrak{S}}} \int_{(U \cap \Gamma) \backslash U} \int_{\Lambda \backslash \mathfrak{n}} T f\left(n_{x} u m\right) \bar{f}\left(n_{x} u m\right) d x d u \frac{d m}{\delta(m)} \\
=\int_{M_{\mathfrak{S}}} \int_{(U \cap \Gamma) \backslash U} \int_{\Lambda \backslash \mathfrak{n}} T\left(\sum_{\xi} \psi\langle x, \xi\rangle \widehat{f}_{\xi}(u m)\right) \sum_{\xi^{\prime}} \bar{\psi}\left\langle x, \xi^{\prime}\right\rangle \overline{\widehat{f}}_{\xi}(u m) d x d u \frac{d m}{\delta(m)}
\end{gathered}
$$

Only the diagonal summands survive the integration in $x \in \mathfrak{n}$, and the exponentials cancel, so this is

$$
\int_{M_{\mathfrak{S}}} \int_{(U \cap \Gamma) \backslash U} \sum_{\xi} T \widehat{f}_{\xi}(u m) \cdot \overline{\hat{f}}_{\xi}(u m) d u \frac{d m}{\delta(m)}
$$

Let $F_{\xi}$ be a left- $N$-invariant function taking the same values as $\widehat{f_{\xi}}$ on $U A^{+} K$, defined by

$$
F_{\xi}\left(n_{x} u m k\right)=\widehat{f_{\xi}}(u m k) \quad\left(\text { for } n_{x} \in N, u \in U, m \in M^{+}, k \in K\right)
$$

Since $T$ does not involve $\mathfrak{n}$, and since $F_{\xi}$ is left $N$-invariant,

$$
T \widehat{f_{\xi}}(u m)=T F_{\xi}\left(n_{x} u m\right)=-\Delta F_{\xi}\left(n_{x} u m\right)
$$

and then

$$
\int_{M_{\mathfrak{S}}} \int_{(U \cap \Gamma) \backslash U} \sum_{\xi} T \widehat{f}_{\xi}(u m) \cdot \widehat{\widehat{f}}_{\xi}(u m) d u \frac{d m}{\delta(m)}=\int_{M_{\mathfrak{S}}} \int_{(U \cap \Gamma) \backslash U} \sum_{\xi}-\Delta F_{\xi}(u m) \cdot \bar{F}_{\xi}(u m) d u \frac{d m}{\delta(m)}
$$

The individual summands are not left- $U \cap \Gamma$-invariant. Since $\widehat{f}_{\xi}(\gamma g)=\widehat{f}_{A_{\gamma}^{*} \xi}(g)$ for $\gamma$ normalizing $\mathfrak{n}$, we can group $\xi \in \Lambda^{\prime}$ by $(U \cap \Gamma)$ orbits to obtain $(U \cap \Gamma)$ subsums, and then unwind. Pick a representative $\omega$ for each orbit $[\omega]$, and let $U_{\omega}$ be the isotropy subgroup of $\omega$ in $(U \cap \Gamma)$, so

$$
\begin{aligned}
& \int_{(U \cap \Gamma) \backslash U} \sum_{\xi}-\Delta F_{\xi}(u m) \cdot \bar{F}_{\xi}(u m) d u=\sum_{[\omega]} \int_{(U \cap \Gamma) \backslash U} \sum_{\xi \in[\omega]}-\Delta F_{\xi}(u m) \cdot \bar{F}_{\xi}(u m) d u \\
= & \sum_{\omega} \int_{(U \cap \Gamma) \backslash U} \sum_{\gamma \in U_{\omega} \backslash(U \cap \Gamma)}-\Delta F_{A_{\gamma}^{*} \omega}(u m) \cdot \bar{F}_{A_{\gamma}^{*} \omega}(u m) d u=\sum_{\omega} \int_{U_{\omega} \backslash U}-\Delta F_{\omega}(u m) \cdot \bar{F}_{\omega}(u m) d u
\end{aligned}
$$

Then

$$
\int_{M_{\mathfrak{F}}} \int_{(U \cap \Gamma) \backslash U} \sum_{\xi}-\Delta F_{\xi}(u m) \cdot \bar{F}_{\xi}(u m) d u=\sum_{\omega} \int_{M_{\mathfrak{F}}} \int_{U_{\omega} \backslash U}-\Delta F_{\omega}(u m) \cdot \bar{F}_{\omega}(u m) d u \frac{d a}{\delta(a)}
$$

Since $-\Delta$ is a non-negative operator on functions on every quotient $N U_{\omega} \backslash G / K$ of $G / K$, each double integral is non-negative, proving $T$ is non-negative. This completes the estimate of the tails.

### 10.7 Compact $\mathfrak{B}_{a}^{1} \longrightarrow L_{a}^{2}$ in four simple examples

As remarked above, the discrete decomposition of $L_{a}^{2}(\Gamma \backslash G / K)$, for the Friedrichs extension $\widetilde{\Delta}_{a}$ of a restriction $\Delta_{a}$ of the differential operator $\Delta$ to the dense subspace $D_{a}$ of $L_{a}^{2}(\Gamma \backslash G / K)$, will follow from compactness of the resolvent $\left(1-\widetilde{\Delta}_{a}\right)^{-1}$, which will follow from the compactness of the inclusion $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}(\Gamma \backslash G / K)$, demonstrated just below.
[10.7.1] Theorem: With $a \gg 1$, the inclusion $\mathfrak{B}_{a}^{1}(\Gamma \backslash G / K) \rightarrow L_{a}^{2}(\Gamma \backslash G / K)$ is compact.
Proof: Again, we are roughly following [Lax-Phillips 1976], pages 204-6. The total boundedness criterion for pre-compactness [14.7.1] requires that, given $\varepsilon>0$, the image of the unit ball $B$ in $\mathfrak{B}_{a}^{1}$ in $L_{a}^{2}(\Gamma \backslash G / K)$ can be covered by finitely-many balls of radius $\epsilon$.

The idea is that the usual Rellich compactness lemma, asserting compactness of proper inclusions of Sobolev spaces on (multi-)tori as in [9.5.12] and [9.5.15], reduces the issue to estimates [10.4], [10.5] on the tails. In more detail: let $\mathfrak{S}$ be a fixed Siegel set that surjects to the quotient $\Gamma \backslash G$. Given $c \geq a$, let $Y_{o}$ be the image of $\{g \in \mathfrak{S}: \eta(g) \leq c+1\}$ in $\Gamma \backslash G / K$, and cover it by opens $U_{1}, \ldots, U_{n}$ in $\Gamma \backslash G / K$ with small compact closures, and take one open $U_{\infty}$ covering the image $Y_{\infty}$ of $\eta \geq c$. Compactness of $Y_{o}$ produces a finite sub-cover. Choose a smooth partition of unity $\left\{\varphi_{i}\right\}$ subordinate to the finite subcover and $U_{\infty}$, letting $\varphi_{\infty}$ be a smooth function that is identically 1 for $y \geq c$. That is, $\varphi_{\infty}+\sum_{i} \varphi_{i}=1$, and $\varphi_{i}$ has compact support inside the open $U_{i}$. Note that [10.4.3] showed that we can choose a family of smooth cut-off functions $\varphi_{\infty}$ so that $\varphi_{\infty} \cdot f$ has a uniform $\mathfrak{B}^{1}$ bound in terms of both $f$ and the family.

Maps among function spaces on the compact part $Y_{o}$ behave well for more general reasons, as we see now. Let $\mathfrak{B}_{a}^{1}\left(Y_{o}\right)$ be the closure of $C_{c}^{\infty}\left(Y_{o}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)$ in $\mathfrak{B}_{a}^{1}$, and $L_{a}^{2}\left(Y_{o}\right)$ the closure of $C_{c}^{\infty}\left(Y_{o}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)$ in $L_{a}^{2}(\Gamma \backslash G / K)$.
[10.7.2] Theorem: For $\Gamma=S L_{2}(\mathbb{Z})$ and $G=S L_{2}(\mathbb{R}), \mathfrak{B}_{a}^{1}\left(Y_{o}\right) \rightarrow L_{a}^{2}\left(Y_{o}\right)$ is compact.
Proof: To take advanatage of some fortunate, simplifying (but not strictly necessary) coincidences, we first carry out this part of the argument just for $\Gamma=S L_{2}(\mathbb{Z})$ and $G=S L_{2}(\mathbb{R})$.

For finite $j$, without loss of generality take the opens $U_{j}$ to be small rectangles in the upper half-plane, and the coordinate maps $\psi_{j}$ simply the inclusions. Fix $j$, and let $U_{j}=\left\{z=x+i y: x_{1}<x<x_{2}, y_{1}<y<y_{2}\right\}$. On $U_{j}$, the measure $d x d y / y^{2}$ and the coefficients of the differential operator $\Delta=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)$ differ by bounded amounts from the Euclidean $d x d y$ and $\Delta^{E}=\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}$. Thus, the corresponding $\mathfrak{B}^{1}$ and $L^{2}$ norms are comparable, as follows. Let

$$
L_{a}^{2}\left(U_{j}\right)=\text { closure of } C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K) \text { in } L_{a}^{2}(\Gamma \backslash G / K)
$$

and

$$
\mathfrak{B}_{a}^{1}\left(U_{j}\right)=\text { closure of } C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K) \text { in } \mathfrak{B}_{a}^{1}(\Gamma \backslash G / K)
$$

Letting

$$
|f|_{L^{2}\left(\psi_{j} U_{j}\right)}^{2}=\int_{\psi_{j} U_{j}}|f|^{2} d x d y \quad|f|_{\mathfrak{B}^{1}\left(\psi_{j} U_{j}\right)}^{2}=\int_{\psi_{j} U_{j}}\left(|f|^{2}-\Delta^{E} f \cdot \bar{f}\right) d x d y
$$

there are easy comparisons

$$
\begin{array}{ll}
\frac{1}{y_{2}} \cdot|f|_{L^{2}\left(\psi_{j} U_{j}\right)} \leq|f|_{L^{2}\left(U_{j}\right)} \leq \frac{1}{y_{1}} \cdot|f|_{L^{2}\left(\psi_{j} U_{j}\right)} & \left(\text { for } f \in C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)\right) \\
\frac{1}{y_{2}} \cdot|f|_{\mathfrak{B}^{1}\left(\psi_{j} U_{j}\right)} \leq|f|_{\mathfrak{B}^{1}\left(U_{j}\right)} \leq \frac{1}{y_{1}} \cdot|f|_{L^{2}\left(\psi_{j} U_{j}\right)} & \left(\text { for } f \in C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)\right)
\end{array}
$$

Identification of opposite edges of the rectangles $\psi_{i} U_{i}$ produces a two-torus $T_{j}$, with $L^{2}\left(T_{j}\right)$ and $\mathfrak{B}^{1}\left(T_{j}\right)$ defined from the Euclidean measure and Euclidean Laplacian. The usual Rellich Lemma asserts the compactness of the inclusion $\mathfrak{B}^{1}\left(T_{j}\right) \rightarrow L^{2}\left(T_{j}\right)$. We will repeatedly use
[10.7.3] Lemma: Let $A, B, C, D$ be Hilbert spaces, with a commutative diagram of continuous linear maps

with $B \rightarrow D$ compact, and $S: C \rightarrow D$ with constant $m>0$ such that $|v|_{C} \leq m \cdot|S v|_{D}$ for all $v \in C$. Then $A \rightarrow C$ is also compact.
Proof: (of lemma) Let $X$ be the closed unit ball in $A$, with image $Y$ in $C$. By continuity, the image of $X$ in $B$ is inside a finite-radius ball $Z$. By compactness of $B \rightarrow D$, given $\varepsilon>0$, the image of $Z$ in $D$ is covered by finitely-many $\frac{\varepsilon}{m}$-balls $V_{1}, \ldots, V_{n}$. The condition on $S$ assures that the inverse images $S^{-1}\left(S Y \cap V_{j}\right)$ are contained in $\varepsilon$-balls in $C$. Thus, $Y$ is covered by finitely-many $\varepsilon$-balls in $C$. This holds for every $\varepsilon>0$, so the image $Y$ is pre-compact, and $A \rightarrow C$ is compact.

The lemma applies to the situation


The standard Rellich lemma is that $\mathfrak{B}^{1}\left(T_{j}\right) \rightarrow L^{2}\left(T_{j}\right)$ is compact, and the inclusion $L_{a}^{2}\left(U_{j}\right) \rightarrow L^{2}\left(T_{j}\right)$ satisfies the hypothesis of the lemma with $m=1$, so $\mathfrak{B}_{a}^{1}\left(U_{j}\right) \rightarrow L_{a}^{2}\left(U_{j}\right)$ is compact.

Map $\mathfrak{B}_{a}^{1}\left(Y_{o}\right)$ to $\bigoplus_{j=1}^{n} \mathfrak{B}_{a}^{1}\left(U_{j}\right)$ by $f \rightarrow \bigoplus_{j} \varphi_{j} \cdot f$, and similarly for $L^{2}$. Applying the lemma to

yields the compactness of $\mathfrak{B}_{a}^{1}\left(Y_{o}\right) \rightarrow L_{a}^{2}\left(Y_{o}\right)$.
Returning to the proof of the theorem: let $B$ be the unit ball in $\mathfrak{B}_{a}^{1}(\Gamma \backslash G / K)$. Given $\varepsilon>0$, take $c \geq a$ sufficiently large and smooth cut-off function $\varphi_{\infty}$ such that $\varphi_{\infty} \cdot B$ lies in a single $\varepsilon / 2$-ball in $L_{a}^{2}(\Gamma \backslash G / K)$. By compactness of $\mathfrak{B}_{a}^{1}\left(Y_{o}\right) \rightarrow L_{a}^{2}\left(Y_{o}\right)$, the image of $B$ in

$$
B \longrightarrow\left(1-\varphi_{\infty}\right) \cdot B \subset \mathfrak{B}_{a}^{1}\left(Y_{o}\right) \longrightarrow L_{a}^{2}\left(Y_{o}\right)
$$

is pre-compact, so can be covered by finitely-many $\varepsilon / 2$-balls in $L_{a}^{2}(\Gamma \backslash G / K)$. Thus, $B=\varphi_{\infty} \cdot B+\left(1-\varphi_{\infty}\right) \cdot B$ can be covered by finitely-many $\varepsilon$-balls in $L_{a}^{2}(\Gamma \backslash G / K)$, so is pre-compact there. This proves the compactness of the inclusion $\mathfrak{B}_{a}^{1}(\Gamma \backslash G / K) \rightarrow L_{a}^{2}(\Gamma \backslash G / K)$.
///
Now we give an argument for compactness applicable more generally. To cope more sanely with comparisons of norms, we want a sort of gradient operator $\nabla$ on functions on $\Gamma \backslash G / K$ such that there is an integration by parts property

$$
\int_{\Gamma \backslash G / K}-\Delta f \cdot \bar{f}=\int_{\Gamma \backslash G / K}\langle\nabla f, \nabla f\rangle_{\mathfrak{s}}
$$

with an inner product on the vector space $\mathfrak{s}$ in which $\nabla$ takes values. This can be accomplished quite generally as follows.

Let $\gamma \rightarrow \gamma^{\theta}$ be an involutive automorphism ${ }^{[69]}$ on the Lie algebra of $G$ as $\mathfrak{g}=\mathfrak{s}+\mathfrak{k}$ such that the Lie algebra of $K$ is the +1 eigenspace, and let $\mathfrak{s}$ be the -1 eigenspace. For example, for $G=S L_{2}(\mathbb{R})$ we can

[^19]
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take $\gamma^{\theta}=-\gamma^{\top}$. Let $\langle,\rangle_{\mathfrak{s}}$ be a positive-definite real-valued inner product on $\mathfrak{s}$, invariant under the action of $K$ :

$$
\left\langle k \alpha k^{-1}, k \beta k^{-1}\right\rangle_{\mathfrak{s}}=\langle\alpha, \beta\rangle_{\mathfrak{s}} \quad(\text { for all } \alpha, \beta \in \mathfrak{s} \text { and } k \in K)
$$

In all these examples, $\langle,\rangle_{\mathfrak{s}}$ can be obtained by restricting the trace form ${ }^{[70]}$

$$
\langle\alpha, \beta\rangle_{\text {trace }}=\operatorname{Retr}(\alpha \cdot \beta) \quad(\text { for } \alpha, \beta \in \mathfrak{g})
$$

where $\operatorname{tr}$ is matrix trace. Let $x \in \mathfrak{g}$ act on functions on $G$ or $\Gamma \backslash G$ by differentiating right translation $X_{x}$, as in chapters 4,6 , and subsequently. We need a name for this map, so let $\rho(x)=X_{x}$. It is $K$-equivariant. To describe $\nabla$ independently of coordinates, consider a sequence of $K$-equivariant maps, reminiscent of the analogue in the coordinate-independent description [4.2] of the Casimir operator:

where $\mathfrak{s} \otimes \mathfrak{s}^{*} \rightarrow \operatorname{End}_{\mathbb{R}}(\mathfrak{s})$ is the natural map $(x \otimes \lambda)(y)=\lambda(y) \cdot x$, and where $\mathfrak{s}^{*}$ is identified with $\mathfrak{s}$ via $x \rightarrow\langle-, x\rangle_{\mathfrak{s}}$. Thus,
[10.7.4] Lemma: The image $\nabla$ of the identity automorphism $1_{\mathfrak{s}}$ of $\mathfrak{s}$ is $K$-equivariant.
Thus, for any orthonormal basis $\left\{x_{j}\right\}$ of $\mathfrak{s}$, in coordinates

$$
\nabla=\sum_{j} X_{x_{j}} \cdot x_{j} \quad \nabla f=\sum_{j} X_{x_{j}} f \cdot x_{j} \quad\left(\text { for } f \in C^{\infty}(G)\right)
$$

Since $\nabla$ is right $K$-equivariant, it descends to an operator on functions on $G / K$ and $\Gamma \backslash G / K$.
[10.7.5] Lemma: For $f \in C_{c}^{\infty}(\Gamma \backslash G / K)=C_{c}^{\infty}(\Gamma \backslash G)^{K}$,

$$
\int_{\Gamma \backslash G}-\Delta f \cdot \bar{f}=\int_{\Gamma \backslash G}\langle\nabla f, \nabla f\rangle_{\mathfrak{s}}
$$

Proof: Now write simply $x$ for the operator $X_{x}=\rho(x)$. Let $\left\{\theta_{j}\right\}$ be a basis for $\mathfrak{k}$ such that $\left\langle\theta_{i}, \theta_{j}\right\rangle_{\text {trace }}=-\delta_{i j}$ with Kronecker $\delta$ and the trace pairing. As in [4.2], the Casimir operator is (the image of) $\sum_{j} x_{j}^{2}-\sum_{i} \theta_{i}^{2}$ in the universal enveloping algebra. Thus, on right $K$-invariant functions it is $\sum_{j} x_{j}^{2}$. Thus, integrating by parts,

$$
\begin{gathered}
\int_{\Gamma \backslash G}-\Delta f \cdot \bar{f}=\int_{\Gamma \backslash G}-\sum_{j} x_{j}^{2} f \cdot \bar{f}=\int_{\Gamma \backslash G} \sum_{j} x_{j} f \cdot x_{j} \bar{f} \\
=\int_{\Gamma \backslash G} \sum_{j}\left\langle x_{j} f \cdot x_{j}, x_{j} \bar{f} \cdot x_{j}\right\rangle_{\mathfrak{s}}=\int_{\Gamma \backslash G}\langle\nabla f, \nabla f\rangle_{\mathfrak{s}}
\end{gathered}
$$

as desired.
Using $\nabla$, now we can give a more general proof of compactness, as follows. As above, let $\mathfrak{S}$ be a fixed Siegel set that surjects to the quotient $\Gamma \backslash G / K$. Given $c \geq a$, let $Y_{o}$ be the image of $\{g \in \mathfrak{S}: \eta(g) \leq c+1\}$ in $\Gamma \backslash G / K$, and cover it by opens $U_{1}, \ldots, U_{n}$ in $\Gamma \backslash G / K$ with small compact closures, and take one open $U_{\infty}$ covering the image $Y_{\infty}$ of $\eta \geq c$. Compactness of $Y_{o}$ produces a finite sub-cover. Choose a smooth partition of unity $\left\{\varphi_{i}\right\}$ subordinate to the finite subcover and $U_{\infty}$, letting $\varphi_{\infty}$ be a smooth function identically 1 for $y \geq c$. That is, $\varphi_{\infty}+\sum_{i} \varphi_{i}=1$, and $\varphi_{i}$ has compact support inside the open $U_{i}$. A general version of [10.4.3] shows that we can choose a family of smooth cut-off functions $\varphi_{\infty}$ so that $\varphi_{\infty} \cdot f$ has a uniform $\mathfrak{B}^{1}$ bound in terms of both $f$ and the family. Now we have the somewhat more general version of [10.7.2]: as above,

[^20]let $\mathfrak{B}_{a}^{1}\left(Y_{o}\right)$ be the closure of $C_{c}^{\infty}\left(Y_{o}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)$ in $\mathfrak{B}_{a}^{1}$, and $L_{a}^{2}\left(Y_{o}\right)$ the closure of $C_{c}^{\infty}\left(Y_{o}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)$ in $L_{a}^{2}(\Gamma \backslash G / K)$.
[10.7.6] Theorem: $\mathfrak{B}_{a}^{1}\left(Y_{o}\right) \rightarrow L_{a}^{2}\left(Y_{o}\right)$ is compact.
Proof: As earlier, let $r=1,2,3,4$ in the respective cases. For finite $j$, without loss of generality take the opens $U_{j}$ to be small rectangles in Iwasawa coordinates $x, y$ on $G / K$ with $x \in \mathbb{R}^{r}$ and $y>0$, and the coordinate maps $\psi_{j}$ the inclusions. As in the simplest case, on $U_{j}$, the measure $d x d y / y^{r+1}$ differs by a bounded amount from from the Euclidean invariant measure $d x d y$.

Identifying $\mathfrak{s}$ with the tangent space at every point of $G / K$, on a subset $U$ with compact closure the inner product $\langle,\rangle_{\mathfrak{s}}$ differs from the Euclidean inner product $\langle,\rangle_{E}$ by bounded amounts, simply because continuous functions on compacts are uniformly bounded. Similarly, the coefficients of $\nabla$ differ from those of the Euclidean gradient $\nabla^{E}$ by bounded amounts, for the same reason. Thus, the $\mathfrak{B}^{1}$ and $L^{2}$ norms are comparable to the Euclidean ones on each of the finitely-many $U_{j}$. Specifically, let

$$
L_{a}^{2}\left(U_{j}\right)=\text { closure of } C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K) \text { in } L_{a}^{2}(\Gamma \backslash G / K)
$$

and

$$
\mathfrak{B}_{a}^{1}\left(U_{j}\right)=\text { closure of } C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K) \text { in } \mathfrak{B}_{a}^{1}(\Gamma \backslash G / K)
$$

As earlier in the simplest case, denote the Euclidean versions by

$$
|f|_{L^{2}\left(\psi_{j} U_{j}\right)}^{2}=\int_{\psi_{j} U_{j}}|f|^{2} d x d y \quad|f|_{\mathfrak{B}^{1}\left(\psi_{j} U_{j}\right)}^{2}=\int_{\psi_{j} U_{j}}\left(|f|^{2}+\left\langle\nabla^{E} f, \nabla^{E} \bar{f}\right\rangle_{E}\right) d x d y
$$

Then the comparisons, less explicit than in the proof of [10.7.2], are

$$
|f|_{L^{2}\left(\psi_{j} U_{j}\right)} \ll|f|_{L^{2}\left(U_{j}\right)} \ll|f|_{L^{2}\left(\psi_{j} U_{j}\right)} \quad\left(\text { for } f \in C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)\right)
$$

and

$$
|f|_{\mathfrak{B}^{1}\left(\psi_{j} U_{j}\right)} \ll|f|_{\mathfrak{B}^{1}\left(U_{j}\right)} \ll|f|_{L^{2}\left(\psi_{j} U_{j}\right)} \quad\left(\text { for } f \in C_{c}^{\infty}\left(U_{j}\right) \cap L_{a}^{2}(\Gamma \backslash G / K)\right)
$$

with implied constants uniform in $f$. The rest of the argument proceeds as in [10.7.2]: first, identification of opposite edges of the rectangles $\psi_{i} U_{i}$ produces an $(r+1)$-torus $T_{j}$, with $L^{2}\left(T_{j}\right)$ and $\mathfrak{B}^{1}\left(T_{j}\right)$ defined from the Euclidean measure and Euclidean Laplacian. The usual Rellich Lemma asserts the compactness of the inclusion $\mathfrak{B}^{1}\left(T_{j}\right) \rightarrow L^{2}\left(T_{j}\right)$. Use the lemma [10.7.3] in the situation


As in the proof of [10.7.2], the standard Rellich lemma asserts that $\mathfrak{B}^{1}\left(T_{j}\right) \rightarrow L^{2}\left(T_{j}\right)$ is compact, and the inclusion $L_{a}^{2}\left(U_{j}\right) \rightarrow L^{2}\left(T_{j}\right)$ satisfies the hypothesis of the lemma, so $\mathfrak{B}_{a}^{1}\left(U_{j}\right) \rightarrow L_{a}^{2}\left(U_{j}\right)$ is compact.

As in the proof of [10.7.2], map $\mathfrak{B}_{a}^{1}\left(Y_{o}\right)$ to $\bigoplus_{j=1}^{n} \mathfrak{B}_{a}^{1}\left(U_{j}\right)$ by $f \rightarrow \bigoplus_{j} \varphi_{j} \cdot f$, and similarly for $L^{2}$. Applying the lemma [10.7.3] to

yields the compactness of $\mathfrak{B}_{a}^{1}\left(Y_{o}\right) \rightarrow L_{a}^{2}\left(Y_{o}\right)$.

$$
\text { 10.8 Compact } \mathfrak{B}_{a}^{1} \longrightarrow L_{a}^{2} \text { for } S L_{3}(\mathbb{Z}), S L_{4}(\mathbb{Z}), S L_{5}(\mathbb{Z}), \ldots
$$

Now let $\Gamma=S L_{r}(\mathbb{Z}), G=S L_{r}(\mathbb{R})$, and $K=S O_{r}(\mathbb{R})$. As in the smaller examples, the global automorphic Sobolev space $\mathfrak{B}^{1}$ is the completion of $C_{c}^{\infty}(\Gamma \backslash G / K)$ with respect to the $\mathfrak{B}^{1}$-norm

$$
|f|_{\mathfrak{B}^{1}}^{2}=\int_{\Gamma \backslash G / K}(1-\Delta) f \cdot \bar{f}
$$

For a cut-off height $a \gg 1$, let $\mathfrak{B}_{a}^{1}$ be the $\mathfrak{B}^{1}$ completion of $C_{c}^{\infty}(Z \Gamma \backslash G / K) \cap L_{a}^{2}$ in the $\mathfrak{B}^{1}$ norm. The resolvent of the Friedrichs extension maps continuously from $L_{a}^{2}$ to an automorphic Sobolev space $\mathfrak{B}_{a}^{1}$ with its finer topology. Thus, it suffices to show that the injection $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}$ is compact. As in the simpler examples, to prove this compactness, we will show that the image of the unit ball of $\mathfrak{B}_{a}^{1}$ is totally bounded in $L_{a}^{2}$.
[10.8.1] Theorem: The Friedrichs self-adjoint extension $\widetilde{\Delta}_{a}$ of the restriction of the symmetric operator $\Delta$ to test functions $D_{a}$ in $L_{a}^{2}$ has compact resolvent, thus has purely discrete spectrum.
Proof: First, we grant that we can control smooth cut-off functions:
[10.8.2] Lemma: Fix a positive simple root $\alpha$. Given $\mu \geq \eta(\alpha)+1$, there are smooth functions $\varphi_{a}^{\alpha}$ for $\alpha \in \Phi$ and $\varphi_{a}^{o}$ such that: all these functions are real-valued, taking values between 0 and $1, \varphi^{o}$ is supported in $C_{\mu+1}$ and $\varphi^{\alpha} \mu$ is supported in $X_{a}^{\alpha}$, and $\varphi_{a}^{o}+\sum_{\alpha} \varphi_{a}^{\alpha}=1$. Further, there is a bound $C$ uniform in $\mu \geq \eta(\alpha)+1$, such that $\left|f \cdot \varphi_{a}^{o}\right|_{\mathfrak{B}^{1}} \leq C \cdot|f|_{\mathfrak{B}^{1}}$ and

$$
\left|f \cdot \varphi_{a}^{\alpha}\right|_{\mathfrak{B}^{1}} \leq C \cdot|f|_{\mathfrak{B}^{1}} \quad(\text { for all } \mu \geq \eta(\alpha)+1)
$$

(Proof almost identical to [10.4.3].)
The key point is the estimation of tails as in [10.3] and [10.4]. To prove total boundedness of $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}$, given $\varepsilon>0$, take $\mu \geq \eta(\alpha)+1$ for all $\alpha \in \Phi$, large enough so that $\left|f \cdot \varphi_{a}^{\alpha}\right|_{L^{2}}<\varepsilon$ for all $\alpha \in \Phi$, for all $f \in \mathfrak{B}_{a}^{1}$ with $|f|_{\mathfrak{B}^{1}} \leq 1$. This covers the images $\left\{f \cdot \varphi_{a}^{\alpha}: f \in \mathfrak{B}_{a}^{1}\right\}$ with $\alpha \in \Phi$ with $\operatorname{card}(\Phi)$ open balls in $L^{2}$ of radius $\varepsilon$.

The remaining part $\left\{f \cdot \varphi_{a}^{o}: f \in \mathfrak{B}_{a}^{1}\right\}$ consists of smooth functions supported on the compact $C_{a}$. The latter can be covered by finitely-many coordinate patches $\psi_{i}: U_{i} \rightarrow \mathbb{R}^{d}$. Take smooth cut-off functions $\varphi_{i}$ for this covering. The functions $\left(f \cdot \varphi_{i}\right) \circ \psi_{i}^{-1}$ on $\mathbb{R}^{d}$ have support strictly inside a Euclidean box, whose opposite faces can be identified to form a flat $d$-torus $\mathbb{T}^{d}$. As in the proof of [10.7.6], because continuous functions on compacts are uniformly continuous, the flat gradient and the gradient inherited from $G$ admit uniform comparison on each $\psi\left(U_{i}\right)$, as do the measures, so the $\mathfrak{B}^{1}\left(\mathbb{T}^{d}\right)$-norm of $\left(f \cdot \varphi_{i}\right) \circ \psi_{i}^{-1}$ is uniformly bounded by the $\mathfrak{B}^{1}$-norm. The classical Rellich lemma asserts compactness of $\mathfrak{B}^{1}\left(\mathbb{T}^{d}\right) \rightarrow L^{2}\left(\mathbb{T}^{d}\right)$. By restriction, this gives the compactness of each $\mathfrak{B}^{1} \cdot \varphi_{i} \rightarrow L^{2}$. A finite sum of compact maps is compact, so $\mathfrak{B}^{1} \cdot \varphi_{a}^{o} \rightarrow L^{2}$ is compact. In particular, the image of the unit ball from $\mathfrak{B}^{1}$ admits a cover by finitely-many $\varepsilon$-balls for any $\varepsilon>0$.

Combining these finitely-many $\varepsilon$-balls with the card $(\Phi)$ balls covers the image of $\mathfrak{B}_{a}^{1}$ in $L^{2}$ by finitely-many $\varepsilon$-balls, proving that $\mathfrak{B}_{a}^{1} \rightarrow L^{2}$ is compact.

### 10.9 Compact resolvents and discrete spectrum

The over-all corollary in all these examples:
[10.9.1] Corollary: For $\lambda$ off a discrete set $X$ of points in $\mathbb{C}$, the inverse $\left(\widetilde{\Delta}_{a}-\lambda\right)^{-1}$ exists, is a compact operator, and

$$
\lambda \longrightarrow\left(\left(\widetilde{\Delta}_{a}-\lambda\right)^{-1}: L_{a}^{2}(\Gamma \backslash G / K) \longrightarrow L_{a}^{2}(\Gamma \backslash G / K)\right)
$$

is meromorphic in $\lambda \in \mathbb{C}-X$. The decomposition of $L_{a}^{2}(\Gamma \backslash G / K)$ with respect to $\widetilde{\Delta}_{a}$ is discrete: there is an orthogonal basis of $L_{a}^{2}(\Gamma \backslash G / K)$ consisting of $\widetilde{\Delta}_{a}$-eigenvectors. The eigenvectors of $\widetilde{\Delta}_{a}$ are eigenvectors of $\left(1-\widetilde{\Delta}_{a}\right)^{-1}$, and eigenvalues $\lambda$ of $\widetilde{\Delta}_{a}$ are in bijection with non-zero eigenvalues of $\left(1-\widetilde{\Delta}_{a}\right)^{-1}$ by $\lambda \longleftrightarrow(1-\lambda)^{-1}$. Proof: The previous preparations and [9.4.1].

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## 11. Meromorphic continuation of Eisenstein series

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2. Re-characterization of Friedrichs extensions
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10. Beyond the critical line: $P^{r, r} \subset S L_{2 r}(\mathbb{Z})$
11. Exotic eigenfunctions: $P^{r, r} \subset S L_{2 r}(\mathbb{Z})$
12. Non-self-associate cases

Appendix A: distributions supported on submanifolds
This proof of meromorphic continuation of various Eisenstein series is in part an elaboration of [Colin de Verdière 1981], and parts of [Colin de Verdière 1982,83]. The less-simple examples [11.7-11.12] of cuspidal-data Eisenstein series for maximal proper parabolics in $G L_{r}$ constitute a natural extension.

In the four simplest examples, the compactness [10.7] of the inclusion map of $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}(\Gamma \backslash G / K)$ of pseudocuspforms yields [10.9] the compactness of the resolvent of the Friedrichs self-adjoint extension [9.2] $\widetilde{\Delta}_{a}$ of the restriction of the invariant Laplacian to (a dense subspace of) that subspace, giving its meromorphy. Eisenstein series differ by elementary functions from Eisenstein-series-like functions in the domain of $\widetilde{\Delta}_{a}$, giving the meromorphic continuation of the Eisenstein series.

A noteworthy preliminary result, reminiscent of [Avakumović 1956], [Roelcke 1956], [Selberg 1956], immediately extends Eisenstein series $E_{s}$ up to the critical line $\operatorname{Re}(s)=\frac{1}{2}$. Analytic continuation of the zeta function $\zeta(s)$ to $\operatorname{Re}(s)>0$ is a corollary of this preliminary result, the simplest example of the arguments of [Langlands 1967/76], [Langlands 1971], and [Shahidi 1978] about meromorphic continuation of automorphic $L$-functions.

### 11.1 Up to the critical line: four simple examples

In this section, we consider the four simplest cases of chapter 1 , with Iwasawa coordinates $x, y$ with $x \in \mathbb{R}^{r}$ with $r=1,2,3,4$ respectively, and $y>0$.

Precise discussion of an unbounded operator and its resolvent requires a specified domain [9.1]. Let $\widetilde{\Delta}$ be the Friedrichs extension [9.2] of the restriction of $\Delta$ to $C_{c}^{\infty}(\Gamma \backslash G / K)$. The Friedrichs construction shows that the domain of $\widetilde{\Delta}$ is contained in a Sobolev space:

$$
\text { domain } \widetilde{\Delta} \subset \mathfrak{B}^{1}=\left(\text { completion of } C_{c}^{\infty}(\Gamma \backslash G / K) \text { under }\langle v, w\rangle_{\mathfrak{B}^{1}}=\langle(1-\Delta) v, w\rangle\right)
$$

The domain of $\widetilde{\Delta}$ contains the smaller Sobolev space

$$
\mathfrak{B}^{2}=\left(\text { completion of } C_{c}^{\infty}(\Gamma \backslash G / K) \text { under }\langle v, w\rangle_{\mathfrak{B}^{2}}=\left\langle(1-\Delta)^{2} v, w\right\rangle\right)
$$

As in the previous chapter, the quotient $\Gamma \backslash G / K$ is a union of a compact part $X_{\mathrm{cpt}}$, whose (conceivably complicated) geometry does not matter, and a geometrically simpler non-compact part:

$$
\Gamma \backslash G / K=X_{\mathrm{cpt}} \cup X_{\infty} \quad\left(\text { compact } X_{\mathrm{cpt}}, \text { cusp neighborhood } X_{\infty}\right)
$$

where, with $a \gg 1$, with normalized height function $\eta\left(n m_{y} k\right)=y^{r}$ as in [1.9],

$$
X_{\infty}=\text { image of }\{g \in G / K: \eta(g) \geq a\}=\Gamma_{\infty} \backslash\{g \in G / K: \eta(g) \geq a\} \approx \mathbb{Z}^{r} \backslash \mathbb{R}^{r} \times[a,+\infty)
$$

## 11. Meromorphic continuation of Eisenstein series

Define a smooth cut-off function $\tau$ as usual: fix $a^{\prime \prime}<a^{\prime}$ large enough so that the image of $\{(x, y) \in G / K$ : $\left.y>a^{\prime \prime}\right\}$ in the quotient is in $X_{\infty}$, and let

$$
\tau(g)= \begin{cases}1 & \left(\text { for } \eta(g)>a^{\prime}\right) \\ 0 & \left(\text { for } \eta(g)<a^{\prime \prime}\right)\end{cases}
$$

Form a pseudo-Eisenstein series $h_{s}$ by winding up the smoothly cut-off function $\tau(g) \cdot \eta(g)^{s}$ :

$$
h_{s}(g)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \tau(\gamma g) \cdot \eta(\gamma g)^{s}
$$

Since $\tau$ is supported on $\eta \geq a^{\prime \prime}$ for large $a^{\prime \prime}$, for any $g \in G / K$ there is at most one non-vanishing summand in the expression for $h_{s}$, and convergence is not an issue. Thus, the pseudo-Eisenstein series $h_{s}$ is entire as a function-valued function of $s$. Let

$$
\widetilde{E}_{s}=h_{s}-\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}\left(\Delta-\lambda_{s}\right) h_{s} \quad\left(\text { where } \lambda=r^{2} \cdot s(s-1) \text { with } r=1,2,3,4\right)
$$

[11.1.1] Claim: $\widetilde{E}_{s}-h_{s}$ is a holomorphic $\mathfrak{B}^{1}$-valued function of $s$ for $\operatorname{Re}(s)>\frac{1}{2}$ and $\operatorname{Im}(s) \neq 0$.
Proof: From Friedrichs' construction [9.2], the resolvent $\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}$ exists as an everywhere-defined, continuous operator for $s \in \mathbb{C}$ for $\lambda_{s}$ not a non-positive real number, because of the non-positive-ness of $\Delta$. Further, for $\lambda_{s}$ not a non-positive real, this resolvent is a holomorphic operator-valued function. In fact, for such $\lambda_{s}$, the resolvent $\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}$ injects from $L^{2}(\Gamma \backslash G / K)$ to $\mathfrak{B}^{1}$.
[11.1.2] Remark: The smooth function $\left(\Delta-\lambda_{s}\right) h_{s}$ is supported on the image of $b \leq y \leq b^{\prime}$ in $\Gamma \backslash G / K$, which is compact. Thus, it is in $L^{2}(\Gamma \backslash G / K)$. It might seem $\widetilde{E}_{s}$ vanishes, if it is forgotten that the indicated resolvent maps to the domain of $\widetilde{\Delta}$ inside $L^{2}(\Gamma \backslash G / K)$, and that $h_{s}$ is not in $L^{2}(\Gamma \backslash G / K)$ for $\operatorname{Re}(s)>\frac{1}{2}$. Indeed, since $h_{s}$ is not in $L^{2}(\Gamma \backslash G / K)$ and $\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}\left(\Delta-\lambda_{s}\right) h_{s}$ is in $L^{2}(\Gamma \backslash G / K)$, the difference cannot vanish.
[11.1.3] Theorem: With $\lambda_{s}=r^{2} \cdot s(s-1)$ not non-positive real, $u=\widetilde{E}_{s}-h_{s}$ is the unique element of the domain of $\widetilde{\Delta}$ such that

$$
\left(\widetilde{\Delta}-\lambda_{s}\right) u=-\left(\Delta-\lambda_{s}\right) h_{s}
$$

Thus, $\widetilde{E}_{s}$ is the usual Eisenstein series $E_{s}$ of $[1.9]$ for $\operatorname{Re}(s)>1$, and gives an analytic continuation of $E_{s}-h_{s}$ as $\mathfrak{B}^{1}$-valued function to $\operatorname{Re}(s)>\frac{1}{2}$ with $s \notin\left(\frac{1}{2}, 1\right]$.
Proof: Uniqueness follows from Friedrichs' construction [9.2] and construction of resolvents, because $\widetilde{\Delta}-\lambda_{s}$ is a bijection of its domain to $L^{2}(\Gamma \backslash G / K)$.

On the other hand, for $\operatorname{Re}(s)>\frac{1}{2}$ and $s \notin\left(\frac{1}{2}, 1\right], \widetilde{E}_{s}-h_{s}$ is in $L^{2}(\Gamma \backslash G / K)$, is smooth, and

$$
\begin{aligned}
\Delta\left(\widetilde{E}_{s}-h_{s}\right)= & \left(\Delta-\lambda_{s}\right)\left(\widetilde{E}_{s}-h_{s}\right)+\lambda_{s} \cdot\left(\widetilde{E}_{s}-h_{s}\right)=\left(\Delta-\lambda_{s}\right) h_{s}+\lambda_{s} \cdot\left(\widetilde{E}_{s}-h_{s}\right) \\
& =(\text { smooth }, \text { compactly-supported })+\lambda_{s} \cdot\left(\widetilde{E}_{s}-h_{s}\right)
\end{aligned}
$$

so is in $\mathfrak{B}^{2}$, so certainly in the domain of $\widetilde{\Delta}$. Abbreviating $H_{s}=\left(\Delta-\lambda_{s}\right) h_{s}$, it is legitimate to compute

$$
\left(\widetilde{\Delta}-\lambda_{s}\right)\left(\widetilde{E}_{s}-h_{s}\right)=\left(\widetilde{\Delta}-\lambda_{s}\right)\left(\left(h_{s}-\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1} H_{s}\right)-h_{s}\right)=\left(\widetilde{\Delta}-\lambda_{s}\right)\left(-\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1} H_{s}\right)=-H_{s}
$$

Thus, $\widetilde{E}_{s}-h_{s}$ is a solution. Also, $E_{s}-h_{s}$ is a solution:

$$
\left(\Delta-\lambda_{s}\right)\left(E_{s}-h_{s}\right)=\left(\Delta-\lambda_{s}\right) E_{s}-\left(\Delta-\lambda_{s}\right) h_{s}=0-\left(\Delta-\lambda_{s}\right) h_{s}
$$

By uniqueness, we are done.

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[11.1.4] Remark: Thus, the Eisenstein series $E_{s}$ has an analytic continuation to $\operatorname{Re}(s)>\frac{1}{2}$ and $s \notin\left(\frac{1}{2}, 1\right]$ as an $h_{s}+\mathfrak{B}^{1}$-valued function. Further, the Friedrichs construction gives a bound for the $L^{2}$-norm of $E_{s}-h_{s}$ via an estimate on the operator norm of $\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}$. The $L^{2}$-norm of $\left(\Delta-\lambda_{s}\right) h_{s}$ is not difficult to estimate, since its support is $b \leq y \leq b^{\prime}$ :

$$
\left|\left(\Delta-\lambda_{s}\right) h_{s}\right|_{L^{2}}^{2} \leq \int_{0}^{1} \int_{b}^{b^{\prime}}\left(\left|\Delta h_{s}\right|+\left|\lambda_{s} h_{s}\right|\right)^{2} \frac{d x d y}{y^{2}}{\ll b, b^{\prime}}\left|\lambda_{s}\right|^{2}
$$

Since $\widetilde{\Delta}$ is negative-definite, as in the proof of [9.17], with $\lambda_{s}=a+b i$

$$
\left|\left(\widetilde{\Delta}-\lambda_{s}\right) v\right|^{2}=|(\widetilde{\Delta}-a) v|^{2}-i b\langle(T-a) v, v\rangle+i b\langle v,(T-a) v\rangle+b^{2}|v|^{2} \geq b^{2}|v|^{2}
$$

Thus, the operator norm of the resolvent is estimated by

$$
\left\|\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}\right\| \leq \frac{1}{\operatorname{Im}\left(\lambda_{s}\right)^{2}}=\frac{1}{2 r^{2}\left(\operatorname{Re}(s)-\frac{1}{2}\right) \cdot \operatorname{Im}(s)} \quad\left(\text { for } \operatorname{Re}(s) \sigma>\frac{1}{2}, \operatorname{Im}(s) \neq 0\right)
$$

Thus,

$$
\begin{aligned}
\left|E_{s}-h_{s}\right|_{L^{2}} & \leq\left\|\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1}\right\| \cdot\left|\left(\Delta-\lambda_{s}\right) h_{s}\right|_{L^{2}}<_{b, b^{\prime}} \frac{1}{\left(\operatorname{Re}(s)-\frac{1}{2}\right) \cdot \operatorname{Im}(s)} \cdot\left|\lambda_{s}\right| \\
& =\frac{1}{\left(\operatorname{Re}(s)-\frac{1}{2}\right) \cdot \operatorname{Im}(s)} \cdot\left(\left(\operatorname{Re}(s)-\frac{1}{2}\right)^{2}+\operatorname{Im}(s)^{2}\right)^{\frac{1}{2}}
\end{aligned}
$$

[11.1.5] Remark: From [1.9.4], the Eisenstein series $E_{s}$ has constant term of the form $\eta^{s}+c_{s} \eta^{1-s}$. Thus, the analytic continuation of $E_{s}$ to $\operatorname{Re}(s)>\frac{1}{2}$ analytically continues $c_{s}$ to $\operatorname{Re}(s)>\frac{1}{2}$. In the case $\Gamma=S L_{2}(\mathbb{R})$, since $c_{s}=\xi(2 s-1) / \xi(2 s)$ with $\xi(s)$ the completed zeta-function $\xi(s)=\pi^{-s / 2} \Gamma(s / 2) \zeta(s)$ this yields the analytic continuation of $\zeta(s)$ to $\operatorname{Re}(s)>0$, off the interval $[0,1]$. A similar conclusion holds for $\Gamma=S L_{2}(\mathbb{Z}[i])$ and the zeta function of $\mathbb{Z}[i] \subset \mathbb{Q}(i)$.

### 11.2 Re-characterization of Friedrichs extensions

Friedrichs extensions of restrictions of $\Delta$ admit simple alternative descriptions facilitating finer analysis in terms of distributions. Up to a point, this can be done abstractly, in the same context as the construction of the Friedrichs extension [9.2].

Let $V$ be a Hilbert space with a complex conjugation map $v \rightarrow \bar{v}$, with expected behavior with respect to the hermitian inner product. This gives a complex-linear isomorphism $c: V \rightarrow V^{*}$ of $V$ to its dual $V^{*}$ via Riesz-Fréchet composed with complex conjugation, by $c: v \rightarrow\langle-, \bar{v}\rangle$. Let $S$ be a symmetric operator on $V$ with dense domain $D$, with $\langle S v, v\rangle \geq\langle v, v\rangle$ for $v \in D$. Suppose that $S$ commutes with the conjugation map. Put $\langle x, y\rangle_{V^{1}}=\langle S x, y\rangle$ for $x, y \in D$, and let $V^{1}$ be the completion of $D$ with respect to this norm. The identity map $D \rightarrow D$ induces a continuous injection $j: V^{1} \rightarrow V$ with dense image. This much is the same as in [9.2].

Write $V^{-1}$ for the Hilbert-space dual $\left(V^{1}\right)^{*}$ of $V^{1}$, with hermitian inner product $\langle,\rangle_{V-1}$. Let $j^{*}$ be the adjoint map $j^{*}: V^{*} \rightarrow\left(V^{1}\right)^{*}=V^{-1}$ of $j$, so composition with complex conjugation $c$ gives


There is a continuous linear map $S^{\#}: V^{1} \longrightarrow V^{-1}$, with the respective topologies, given by

$$
S^{\#}(x)(y)=\langle x, \bar{y}\rangle_{V^{1}} \quad\left(\text { for } x, y \in V^{1}\right)
$$

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By Riesz-Fréchet, this map is a topological isomorphism.
[11.2.1] Claim: The restriction of $S^{\#}$ to the domain of the Friedrichs extension $\widetilde{S}$ is $j^{*} \circ c \circ \widetilde{S}$. The domain of $\widetilde{S}$ is

$$
\text { domain } \widetilde{S}=\widetilde{D}=\left\{x \in V^{1}: S^{\#} x \in\left(j^{*} \circ c\right) V\right\}
$$

Proof: By construction [9.2] of the Friedrichs extension $\widetilde{S}$, its domain is exactly $\widetilde{D}=\widetilde{S}^{-1} V$. Thus, for $x=\widetilde{S}^{-1} x^{\prime}$ with $x^{\prime} \in V$, for all $y \in V^{1}$

$$
\left(S^{\#} x\right)(y)=\left(S^{\#} \widetilde{S}^{-1} x^{\prime}\right)(y)=\left\langle\widetilde{S}^{-1} x^{\prime}, \bar{y}\right\rangle_{V^{-1}}=\langle x, \bar{y}\rangle=\left(\left(j^{*} \circ c\right) x^{\prime}\right)(y)=\left(\left(j^{*} \circ c \circ \widetilde{S}\right) x\right)(y)
$$

Thus, the restriction of $S^{\#}$ to the domain $\widetilde{D}$ of $\widetilde{S}$ is essentially $\widetilde{S}$, namely,

$$
\left.S^{\#}\right|_{\widetilde{D}}=\left.\left(j^{*} \circ c \circ \widetilde{S}\right)\right|_{\widetilde{D}}
$$

Thus, $S^{\#}: V^{1} \rightarrow V^{-1}$ extends $\widetilde{S}$. On the other hand, for $S^{\#} x=\left(j^{*} \circ c\right) y$ with $y \in V$, for all $z \in V^{1}$

$$
\langle z, \bar{x}\rangle_{V^{1}}=\left(S^{\#} x\right)(z)=\left(\left(j^{*} \circ c\right) y\right)(z)=(\lambda y)(j z)=\langle j z, \bar{y}\rangle=\left\langle z, \widetilde{S}^{-1} \bar{y}\right\rangle_{V^{1}}
$$

Thus, $\bar{x}=\widetilde{S}^{-1} \bar{y}$. Thus, the domain of $\widetilde{S}$ is as claimed.
Let $\Theta \subset D$ be stable under conjugation, and stable under $S$. For subsequent application, in the simplest examples we are thinking of the collection of pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}(a, \infty)$. Let $V_{\Theta}$ be the orthogonal complement to $\Theta$ in $V$. Let $S_{\Theta}$ be $S$ restricted to $D_{\Theta}=D \cap V_{\Theta}$. The $S$-stability assumption on $\Theta$ gives $S\left(D_{\Theta}\right) \subset V_{\Theta}$. Certainly $D_{\Theta} \subset V^{1} \cap V_{\Theta}$, so the $V^{1}$ closure of $D_{\Theta}$ is a subset of $V^{1} \cap V_{\Theta}$. However, $V^{1}$-density of $D_{\Theta}$ in $V^{1} \cap V_{\Theta}$ equality is not clear in general: we must assume that $D_{\Theta}$ is $V^{1}$-dense in $V^{1} \cap V_{\Theta}$. In the cases of interest, we have proven this under mild hypotheses [10.3]. This density assumption legitimizes the natural sequel: $S_{\Theta}$ with domain $D_{\Theta}$ is densely defined and symmetric on $V_{\Theta}$, so has Friedrichs extension $\widetilde{S}_{\Theta}$, with domain $\widetilde{D}_{\Theta}$.

The extension

$$
\left(S_{\Theta}\right)^{\#}: V^{1} \cap V_{\Theta} \longrightarrow\left(V^{1} \cap V_{\Theta}\right)^{*}
$$

is described by

$$
\left(S_{\Theta}\right)^{\#}(x)(y)=\langle x, y\rangle_{V^{1}} \quad\left(\text { for } x, y \in V^{1} \cap V_{\Theta}\right)
$$

Let

$$
i_{\Theta}: V^{1} \cap V_{\Theta} \longrightarrow V^{1} \quad i_{\Theta}^{*}: V^{-1}=\left(V^{1}\right)^{*} \longrightarrow\left(V^{1} \cap V_{\Theta}\right)^{*}
$$

be the inclusion and its adjoint, fitting into a diagram

[11.2.2] Claim: $\left(S_{\Theta}\right)^{\#}=i_{\Theta}^{*} \circ S^{\#} \circ i_{\Theta}$, and the domain of $\widetilde{S}_{\Theta}$ is

$$
\widetilde{D}_{\Theta}=\left\{x \in V^{1} \cap V_{\Theta}:\left(S^{\#} \circ i_{\Theta}\right) x \in\left(j^{*} \circ c\right) V+\Theta\right\}=\left\{x \in V^{1} \cap V_{\Theta}: S_{\Theta}^{\#} x \in\left(i_{\Theta}^{*} \circ j^{*} \circ c\right) V\right\}
$$

and $\widetilde{S}_{\Theta} x=y$, with $x \in V^{1} \cap V_{\Theta}$ and $y \in V$, if and only if $\left(S^{\#} \circ i_{\Theta}\right) x=\left(j^{*} \circ c\right) y+\theta$ for some $\theta$ in the $V^{-1}$-closure of $\left(j^{*} \circ c\right) \Theta$.
Proof: The assumption of denseness of $D_{\Theta}$ in $V^{1} \cap V_{\Theta}$ legitimizes formation of the Friedrichs extension as an unbounded self-adjoint operator (densely defined) on $V$. For $x, y \in V^{1} \cap V^{\Theta}$

$$
\left(i_{\Theta}^{*} \circ S^{\#} \circ i_{\Theta}\right)(x)(y)=S^{\#}(x)(y)=\left\langle i_{\Theta} x, i_{\Theta} \bar{y}\right\rangle_{V^{1}}=\langle x, \bar{y}\rangle_{V^{1}}=\left(S_{\Theta}\right)^{\#}(x)(y)
$$

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which is the first statement of the claim.
From above, the Friedrichs extension $\widetilde{S}_{\Theta}$ is characterized by

$$
\left\langle z, \widetilde{S}_{\Theta}^{-1} y\right\rangle_{V^{1}}=\langle z, y\rangle \quad\left(\text { for } z \in D_{\Theta} \text { and } y \in V_{\Theta}\right)
$$

Given $S^{\#} x=\left(j^{*} \circ c\right) y+\theta$ with $x \in V^{1} \cap V_{\Theta}, y \in V$, and $\theta$ in the $V^{-1}$ closure of $\left(j^{*} \circ c\right) \Theta$, take $z \in D_{\Theta}$ and compute

$$
\begin{aligned}
& \langle x, \bar{z}\rangle_{V^{1}}=\left(S^{\#} x\right)(z)=\left(\left(j^{*} \circ c\right) y+\theta\right)(z)=\left(j^{*} \bar{y}\right)(z)+\theta(z) \\
& \quad=\langle z, \bar{y}\rangle+0=\left\langle y, \widetilde{S}_{\Theta}^{-1} S \bar{z}\right\rangle=\left\langle\widetilde{S}_{\Theta}^{-1} y, S \bar{z}\right\rangle=\left\langle\widetilde{S}_{\Theta}^{-1} y, \bar{z}\right\rangle_{V^{1}}
\end{aligned}
$$

thus showing that $\widetilde{S}_{\Theta}^{-1} x=y$. On the other hand, $\left(S_{\Theta}\right)^{\#} x=\left(i_{\Theta}^{*} \circ j^{*} \circ c\right) y$ if and only if $\left(S^{\#} \circ i_{\Theta}\right) x=y+\theta$ for some $\theta \in \operatorname{ker} i_{\Theta}^{*}$, and $\operatorname{ker} i_{\Theta}^{*}$ is the closure of $\Theta$ in $V^{-1}$.

### 11.3 Distributional characterization of pseudo-Laplacians

The previous section applies to the pseudo-Laplacians $\widetilde{\Delta}_{a}$ of chapter 10 for $a \gg 1$ large enough so that the density result [10.3.1] legitimizes the discussion. This re-characterization is needed for meromorphic continuation of Eisenstein series beyond the critical line.

Refering to the notation of the previous section, take $V=L^{2}(\Gamma \backslash G / K)$, use the pointwise conjugation $\operatorname{map} c: L^{2}(\Gamma \backslash G / K) \rightarrow L^{2}(\Gamma \backslash G / K)$, let $D=C_{c}^{\infty}(\Gamma \backslash G / K)$, put $S=1-\left.\Delta\right|_{D}$, and let $\Theta=\Theta_{a}$ be the space of pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}(a,+\infty)$ with $a \gg 1$ large enough so that the density lemma [10.3.1] holds. Let $V^{1}=\mathfrak{B}^{1}$ be the completion of $D$ with respect to the norm given by

$$
|f|_{\mathfrak{B}^{1}}^{2}=\int_{\Gamma \backslash G / K}(1-\Delta) f \cdot \bar{f}=\langle(1-\Delta) f, f\rangle
$$

Let $\mathfrak{B}^{-1}$ be the Hilbert space dual of $\mathfrak{B}^{1}$. With inclusion $j: \mathfrak{B}^{1} \rightarrow V$, let $j^{*}$ be its adjoint, and we have a picture

$$
\mathfrak{B}^{1} \xrightarrow{j} V \xrightarrow{j^{*} \circ c} \mathfrak{B}^{-1}
$$

Let $\eta_{a}$ be the functional on $D$ which evaluates constant terms at height $a$.
[11.3.1] Remark: In the present context, we have to prove the following lemma without using any spectral description of $\mathfrak{B}^{1}$ or $\mathfrak{B}^{-1}$, because we are in the process of proving meromorphic continuation of Eisenstein series, which must be done (logically) prior to spectral decompositions. In the following chapter, we can revisit spaces $\mathfrak{B}^{s}$ for $s \in \mathbb{R}$ in a more congenial context, with spectral theory available.

Indeed, the proof of the following lemma uses the already-available spectral theory on multi-tori:
[11.3.2] Lemma: For $a \gg 1$ sufficiently large, $\eta_{a} \in \mathfrak{B}^{-1}$.
Proof: As expected, take $b^{\prime} \gg 1$ large enough so that the standard Siegel set $\mathfrak{S}_{b^{\prime}}$ meets no translate $\gamma \mathfrak{S}_{b^{\prime}}$ with $\gamma \in \Gamma$ unless $\gamma \in N \cap \Gamma$, so that the cylinder $C_{b^{\prime}}=(P \cap \Gamma) \backslash \mathfrak{S}_{b^{\prime}}$ injects to $\Gamma \backslash G / K$. Take $a>b^{\prime}$. Since the support of $\eta_{a}$ is compact and properly inside $\mathfrak{S}_{b^{\prime}}$, there is a test function $\psi$ identically 1 on the support of $\eta_{a}$, and supported inside $\mathfrak{S}_{b^{\prime}}$. Then $\psi \cdot \eta_{a}=\eta_{a}$, in the sense that $\eta_{a}(f)=\eta_{a}(\psi f)$ for all test functions $f$. Thus, it suffices to consider test functions with support in a subset $X=(N \cap \Gamma) \backslash N \times\left(b^{\prime}, b^{\prime \prime}\right)$ of the cylinder $C_{b^{\prime}}=(N \cap \Gamma) \backslash N \times\left(b^{\prime},+\infty\right) \approx(\mathbb{Z} \backslash \mathbb{R})^{r} \times\left(b^{\prime}, b^{\prime \prime}\right)$, with $b^{\prime \prime}<+\infty$.

Identifying the endpoints of the finite interval $\left(b^{\prime}, b^{\prime \prime}\right) \subset\left[b^{\prime}, b^{\prime \prime}\right]$ identifies it with another circle, thus imbedding $X \subset \mathbb{T}^{r+1}$. As in the proofs of [10.7.2] and [10.7.6], the $\mathfrak{B}^{1}$ and $L^{2}$ norms on $X$ are uniformly comparable to those on $\mathbb{T}^{r+1}$ descended from the Euclidean versions. Thus, to prove $\eta_{a} \in \mathfrak{B}^{-1}$, it suffices to prove that the functional $\theta$ given by integration along $\mathbb{T}^{r} \times\{0\}$ inside $\mathbb{T}^{r+1}$ is in the corresponding $\mathfrak{B}^{-1}$ space there. The advantage is that we can use Fourier series, since the spectral theory of $\mathbb{T}$ and $\mathbb{T}^{n}$ is already available, as in [9.5], especially [9.5.9]. That is, parametrizing $\mathbb{T}^{r+1}$ as $\mathbb{Z}^{r+1} \backslash \mathbb{R}^{r+1}$, let $\psi_{\xi}$ be $\psi(x)=e^{2 \pi i \xi \cdot x}$ for $\xi, x \in \mathbb{R}$ and $\xi \cdot x$ the usual inner product on $\mathbb{R}^{r+1}$. Letting $\xi=\left(\xi_{1}, \ldots, \xi_{r+1}\right)$, the Fourier coefficients of $\theta$ are

$$
\widehat{\theta}(\xi)=\theta\left(\bar{\psi}_{\xi}\right)=\int_{\mathbb{T}^{r} \times\{0\}} \psi_{\xi}(x) d x= \begin{cases}0 & \left(\text { for }\left(\xi_{1}, \ldots, \xi_{r}\right) \neq(0, \ldots, 0)\right) \\ 1 & \left(\text { for }\left(\xi_{1}, \ldots, \xi_{r}\right)=(0, \ldots, 0)\right)\end{cases}
$$

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Thus, the $s^{\text {th }}$ Sobolev norm of $\theta$ is

$$
\sum_{\xi \in \mathbb{Z}^{r+1}}|\widehat{\theta}(\xi)|^{2} \cdot\left(1+|\xi|^{2}\right)^{s}=\sum_{\xi_{r+1} \in \mathbb{Z}} 1 \cdot\left(1+\left|\xi_{r+1}\right|^{2}\right)^{s}
$$

which is finite for $\operatorname{Re}(s)<-\frac{1}{2}$. Certainly it is finite for $s=-1$, giving the desired conclusion.
In the previous lemma, on $\mathbb{T}^{r+1}, \theta$ is certainly the suitable Sobolev space limit of its finite subsums, which are smooth. This pulls back to an assertion that $\eta_{a}$ is in the $\mathfrak{B}^{1}$ closure of test functions. We need a stronger assertion in order to use the re-characterization of the previous section:
[11.3.3] Lemma: $\eta_{a}$ is in the $\mathfrak{B}^{-1}$-closure of $\Theta$.
Proof: Again, by the previous lemma, $\eta_{a}$ is a $\mathfrak{B}^{-1}$-limit of a sequence $\left\{f_{n}\right\}$ of test functions on $\Gamma \backslash G / K$ or on the cylinder $C_{b^{\prime}}$. Much as in [10.3], we want to show that suitable smooth truncations of the $f_{n}$, to put them into $\Theta$, still converge to $\eta_{a}$ in $\mathfrak{B}^{-1}$. As in the previous proof, using $a \gg 1$, we can convert the question to one on $\mathbb{T}^{r+1}$ or on $\mathbb{T}^{r} \times \mathbb{R}$. Further, since nothing is happening in the first $r$ coordinates, it suffices to consider prove the following claim on $\mathbb{R}$.

That is, in the standard Sobolev spaces $H^{s}$ on $\mathbb{R}[9.7]$, we claim that the standard Dirac $\delta$ on $\mathbb{R}$ is an $H^{-1}$ limit of a sequence of test functions supported in $[0,+\infty)$. Let $u$ be a test function on $\mathbb{R}$ which is 0 in $(-\infty, 0]$, is non-negative with integral 1 on $[0,+\infty)$. For $n=1,2,3, \ldots$, let $u_{n}(t)=n \cdot u(n t)$. We claim that $u_{n} \rightarrow \delta$ in $H^{-1}$. Taking Fourier transforms,

$$
\widehat{u_{n}}(\xi)=\int_{\mathbb{R}} e^{-2 \pi i \xi t} n \cdot u(n t) d t=\int_{\mathbb{R}} e^{-2 \pi i \xi t / n} u(t) d t=\widehat{u}(\xi / n)
$$

The Fourier transform of $\delta$ is 1 , since $\delta\left(t \rightarrow e^{2} \pi i \xi t\right)=1$ for all $\xi \in \mathbb{R}$. The function $\widehat{u}$ is still a Schwartz function. We want to show that, as $n \rightarrow+\infty$,

$$
\int_{\mathbb{R}}|\widehat{u}(\xi / n)-1|^{2} \cdot\left(1+\xi^{2}\right)^{-1} d \xi \longrightarrow 0
$$

Certainly $\widehat{u}$ is bounded, so, given $\varepsilon>0$, there is $N \gg 1$ such that for all $n$

$$
\int_{|\xi| \geq N}|\widehat{u}(\xi / n)-1|^{2} \cdot\left(1+\xi^{2}\right)^{-1} d \xi<\varepsilon
$$

By the differentiability of $\widehat{u}$,

$$
\widehat{u}(\xi / n)=\widehat{u}(0)+(\xi / n) \cdot \widehat{u}^{\prime}\left(t_{o}\right) \quad\left(\text { for some } t_{o} \text { between } 0 \text { and } \xi / n\right)
$$

Since the integral of $u$ is $1, \widehat{u}(0)=1$. The derivative $\widehat{u}^{\prime}$ is continuous, so has a bound $B$ on $[-1,1]$. For $|\xi| \leq N$, take $n$ large enough so that $|\xi / n|<\varepsilon \leq 1$. Then

$$
\begin{gathered}
\int_{|\xi| \leq N}|\widehat{u}(\xi / n)-1|^{2} \cdot\left(1+\xi^{2}\right)^{-1} d \xi=\int_{|\xi| \leq N}\left|(\xi / n) \cdot \widehat{u}^{\prime}\left(t_{o}\right)\right|^{2} \cdot\left(1+\xi^{2}\right)^{-1} d \xi \\
\quad \leq \int_{|\xi| \leq N} \varepsilon^{2} \cdot B^{2} \cdot\left(1+\xi^{2}\right)^{-1} d \xi \leq \varepsilon^{2} \cdot B^{2} \int_{\mathbb{R}}\left(1+\xi^{2}\right)^{-1} d \xi \ll \varepsilon
\end{gathered}
$$

Thus, in the spectral-side description of the topology on $H^{-1}$, we have the desired convergence.
In the four simplest cases, we have
[11.3.4] Corollary: $\widetilde{\Delta} u=f$ for $f \in L_{a}^{2}(\Gamma \backslash G / K)$ if and only if $u \in \mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$, and $\Delta u=f+c \cdot \eta_{a}$ for some constant $c$.
[11.3.5] Remark: In particular, the proof mechanisms just above show that $u \in \mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$ implies that the constant term is in the Euclidean Sobolev space $H^{1}(\mathbb{R})$ as a function of the coordinate $y$. By Sobolev imbedding [9.5.4], [9.5.11], 9.5.14], this implies continuity of the constant term, so vanishing in $\eta>a$ implies

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$\eta_{a} u=0$. Conversely, if $u \in \mathfrak{B}^{1}$ and $\eta_{a} u=0$, we could truncate $u$ at height $a$ without disturbing the condition $u \in \mathfrak{B}^{1}$, to put $\wedge^{a} u$ in $\mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$. In fact, after we have the meromorphic continuation of Eisenstein series in hand, and once we have a spectral form of global automorphic Sobolev spaces $\mathfrak{B}^{s}$, one can easily prove that the conditions $(\Delta-\lambda) u=\eta_{a}, u \in \mathfrak{B}^{1}$, and $\eta_{a} u=0$ imply $\eta_{b^{\prime}} u=0$ for all $b^{\prime} \geq a$.
[11.3.6] Remark: For $\lambda_{w}$ not the eigenvalue of a cuspform, the homogeneous equation $\left(\Delta-\lambda_{w}\right) u=0$ has no non-zero solution, so the constant $c$ must be non-zero for non-zero $u$.
Proof: Use the characterization [11.2.2]. The previous lemma shows that $\eta_{a}$ is in the $\mathfrak{B}^{-1}$ closure $\Theta_{-1}$ of $\Theta=\Theta_{a}$. Using $a \gg 1$, we must show that the intersection of that closure with the image $\Delta \mathfrak{B}^{1}$ is at most $\mathbb{C} \cdot \eta_{a}$.

On one hand, because $a \gg 1, \Theta_{-1}$ consists of distributions which, on a Siegel set $\mathfrak{S}_{b^{\prime}}$ with $b^{\prime}$ just slightly less than $a$, have support inside $\mathfrak{S}_{a} \subset \mathfrak{S}_{b^{\prime}}$. On the cylinder $C_{b^{\prime}}=\Gamma_{\infty} \backslash \mathfrak{S}_{b^{\prime}}$, the product of circles $(N \cap \Gamma) \backslash N \approx \mathbb{T}^{r}$ acts by translations, descending to the quotient from $G / K$. By reduction theory, the restrictions to $C_{a^{\prime}}$ of every pseudo-Eisenstein series $\Psi_{\varphi}$ with $\varphi \in C_{c}^{\infty}[a, \infty)$ are invariant under $(N \cap \Gamma) \backslash N$, so anything in the $\mathfrak{B}^{-1}$ closure is likewise invariant.

On the other hand, consider the possible images of $\mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$ by $\Delta$. Certainly $D \cap V^{\Theta}$ consists of functions with constant term vanishing in $\eta \geq a$, and taking $\mathfrak{B}^{1}$ completion preserves this property. Since $\Theta_{-1}$ is $(N \cap \Gamma) \backslash N$-invariant and the Laplacian commutes with the group action, it suffices to look at $(N \cap \Gamma) \backslash N$-integral averages restricted to the cylinder $C_{b^{\prime}}$. Such an integral is a restriction of the constant term $c_{P} v$ to $C_{b^{\prime}}$, and vanishes in $\eta>a$.

Thus, the intersection of possible images by $\widetilde{\Delta}_{a}$ with $\Theta_{-1}$ consists of $(N \cap \Gamma) \backslash N$-invariant distributions in $\mathfrak{B}^{-1}$ supported on $Z=\{\eta \leq a\} \cap\{\eta \geq a\} \approx(N \cap \Gamma) \backslash N$. By [11.A], such distributions are obtained as compositions of derivatives transverse to $Z$ composed with a distribution supported on $Z$. By uniqueness of invariant distributions [14.4], the only $(N \cap \Gamma) \backslash N$-invariant distribution on $Z \approx(N \cap \Gamma) \backslash N$ is (a scalar multiple of) integration on $(N \cap \Gamma) \backslash N$.

Certainly $\eta_{a}$ itself is among these functionals. No higher-order derivative (composed with $\eta_{a}$ ) gives a functional in $\mathfrak{B}^{-1}$, as is visible already on $\mathbb{R}$ : computing the $s^{\text {th }}$ Sobolev norm of the $n^{\text {th }}$ derivative $\delta^{(n)}$ of the Euclidean Dirac $\delta$,

$$
\left|\delta^{(n)}\right|_{H^{s}}^{2}=\int_{\mathbb{R}}\left|\widehat{\delta^{(n)}}(\xi)\right|^{2} \cdot\left(1+\xi^{2}\right)^{s} d \xi=\int_{\mathbb{R}}\left|(-2 \pi i \xi)^{n}\right|^{2} \cdot\left(1+\xi^{2}\right)^{s} d \xi
$$

This is finite only for $s<-\left(\frac{1}{2}+n\right)$.

### 11.4 Key density lemma: simple cases

Similar to the description of $E_{s}$ as $\widetilde{E}_{s}$ above in [11.1], but with $\widetilde{\Delta}_{a}$ in place of $\widetilde{\Delta}$, with the pseudo-Eisenstein series $h_{s}$ formed from the smooth cut-off $\tau \cdot \eta^{s}$ of $\eta^{s}$ as in [11.1], put

$$
\widetilde{E}_{a, s}=h_{s}-\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)^{-1}\left(\Delta-\lambda_{s}\right) h_{s}
$$

Since $\left(\Delta-\lambda_{s}\right) h_{s}$ is compactly supported, it is in $L_{a}^{2}(\Gamma \backslash G / K)$. For $\lambda_{s}$ not a non-positive real, $\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)^{-1}$ is a bijection of $L^{2}(\Gamma \backslash G / K)_{a}$ to the domain of $\widetilde{\Delta}_{a}$, so $u=\widetilde{E}_{s, a}-h_{s}$ is the unique element of the domain of $\widetilde{\Delta}_{a}$ satisfying

$$
\left(\widetilde{\Delta}_{a}-\lambda_{s}\right) u=-\left(\Delta-\lambda_{s}\right) h_{s}
$$

Since the pseudo-Eisenstein series $h_{s}$ is entire, the meromorphy of the resolvent $\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)^{-1}$ [10.9] yields the meromorphy of $\widetilde{E}_{a, s}-h_{s}$ as $\mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$-valued function.

Recall that, with $D=C_{c}^{\infty}(\Gamma \backslash G / K)$ and $D_{a}=D \cap L^{2}(\Gamma \backslash G / K)$, $\mathfrak{B}^{1}$ is the $\mathfrak{B}^{1}$-norm completion of $D$, while $\mathfrak{B}_{a}^{1}$ is the $\mathfrak{B}^{1}$-completion of $D_{a}$. We are counting on $\widetilde{E}_{a, s}-h_{s}$ to be in $\mathfrak{B}_{a}^{1}$. This depends upon the assumption $a \gg 1$, and it is critical to verify the following:
[11.4.1] Lemma: For $a \gg 1, \mathfrak{B}_{a}^{1}=\mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$. That is, for $a \gg 1, D_{a}$ is $\mathfrak{B}^{1}$-dense in $\mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$. Proof: The containment $\mathfrak{B}_{a}^{1} \subset \mathfrak{B}^{1} \cap L^{2}(\Gamma \backslash G / K)$ is immediate. For the opposite containment, given a sequence $\left\{\Psi_{\varphi_{i}} \in D\right\}$ of pseudo-Eisenstein series converging to $f \in \mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$ in the $\mathfrak{B}^{1}$-topology, we

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must produce a sequence of pseudo-Eisenstein series in $D_{a}$ converging to $f$ in the topology of $\mathfrak{B}^{1}$. We will do so by smooth cut-offs of the constant terms of the $\Psi_{\varphi_{i}}$. Since the limit $f$ of the $\Psi_{\varphi_{i}}$ has constant term vanishing above height $y=a$ and is in $L_{a}^{2}(\Gamma \backslash G / K)$, that part of the constant terms of the $\Psi_{\varphi_{i}}$ becomes small. More precisely, we proceed as follows.

Let $F$ be a smooth real-valued function on $\mathbb{R}$ with $F(t)=0$ for $t<-1,0 \leq F(t) \leq 1$ for $-1 \leq t \leq 0$, and $F(t)=1$ for $t \geq 0$. For $\varepsilon>0$, let $F_{\varepsilon}(t)=F((t-a) / \varepsilon)$. Fix real $b$ with $a>b>1$. Given $\Psi_{\varphi_{i}} \rightarrow f \in L_{a}^{2}(\Gamma \backslash G / K)$, the $b$-tail of the constant term of $\Psi_{\varphi_{i}}$ is $\tau_{i}(g)=c_{P} \Psi_{\varphi_{i}}(g)$ for $\eta(g) \geq a^{\prime}$, and $\tau_{i}(g)=0$ for $0<\eta(g) \leq a^{\prime \prime}$. By design, $\Psi_{\varphi_{i}}-\Psi_{F_{\varepsilon} \cdot \tau_{i}} \in D_{a}$ for small $\varepsilon$. We will show that, as $i \rightarrow+\infty$, for $\varepsilon_{i}$ sufficiently small depending on $i$, the $\mathfrak{B}^{1}$-norms of $\Psi_{F_{\varepsilon_{i}} \cdot \tau_{i}}$ go to 0 , and $\Psi_{\varphi_{i}}-\Psi_{F_{\varepsilon_{i}} \cdot \tau_{i}} \rightarrow f$ in the $\mathfrak{B}^{1}$-norm.

Let $\mathfrak{S}=\mathfrak{S}_{b}$ with $b \gg 1$, and put $C_{b}=(N \cap \Gamma) \backslash \mathfrak{S}$. The cylinders $C_{b}$ admit natural actions of the compact abelian group $(N \cap \Gamma) \backslash N$, by translation. For $b \gg 1$, by reduction theory [1.5], the further quotient $(\Gamma \cap M) \backslash C_{b}$ injects to its image in $\Gamma \backslash G / K$. Conveniently, $\Gamma \cap M$ is finite in these examples, so, for $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$, letting

$$
|f|_{\mathfrak{B}^{1}\left(C_{b}\right)}^{2}=\int_{C_{b}}|f(z)|^{2}-\Delta f \cdot \bar{f}
$$

we have

$$
|f|_{\mathfrak{B}^{1}\left(C_{b}\right)}^{2} \ll \int_{\Gamma \backslash G / K}|f(z)|^{2}-\Delta f \cdot \bar{f}
$$

For each $b>1$, let $\mathfrak{B}^{1}\left(C_{b}\right)$ be the completion of $C_{c}^{\infty}(\Gamma \backslash G / K)$ with respect to the semi-norm $|\cdot|_{\mathfrak{B}^{1}\left(C_{b}\right)}$ (with collapsing since the $\mathfrak{B}^{1}$-norm ignores function values outside $C_{b}$ ).

As usual, we have a continuous action of $(N \cap \Gamma) \backslash N$ on $\mathfrak{B}^{1}\left(C_{b}\right)$. Thus, the map $u \rightarrow c_{P} u$ gives continuous maps of the spaces $\mathfrak{B}^{1}\left(C_{b}\right)$ to themselves. Thus, $c_{P} \Psi_{\varphi_{i}}$ goes to $c_{P} f$ in $\mathfrak{B}^{1}\left(C_{b}\right)$, and $c_{P} \Psi_{\varphi_{i}} \rightarrow c_{P} f=0$ in $\mathfrak{B}^{1}\left(C_{a}\right)$.

To have a Leibniz rule, write the norms as energy norms by integrating by parts: for $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$, put

$$
|f|_{\mathfrak{B}^{1}}^{2}=|f|_{L^{2}(\Gamma \backslash G / K)}^{2}+\left||\nabla f|_{\mathfrak{s}}^{2}\right|_{L^{2}(\Gamma \backslash G / K)}
$$

where $\nabla$ is the left $G$-invariant, right $K$-equivariant tangent-space-valued gradient on $G$, as in [10.7]. Thus, $\nabla$ descends to $G / K$ and to $\Gamma \backslash G / K$, and $|\cdot|_{\mathfrak{s}}$ is a $K$-invariant norm on the tangent space(s). Explicitly, as in [10.7], for an involutive automorphism $\theta$ of the Lie algebra $\mathfrak{g}$ with the Lie algebra $\mathfrak{k}$ of $K$ the +1 -eigenspace, the -1 -eigenspace $\mathfrak{s}$ can be identified with the tangent space at every point of $G / K$, via left translation of the exponential map: for $\beta \in \mathfrak{s}$, the associated left $G$-invariant differential operator $X_{\beta}$ is

$$
\left(X_{\beta} f\right)(g)=\left.\frac{\partial}{\partial t}\right|_{t=0} f\left(g e^{t \cdot \beta}\right)
$$

It is easy to describe $\nabla$ in coordinates, even though it is independent of coordinates: for an orthonormal basis $\left\{\sigma_{i}\right\}$ for $\mathfrak{s}$,

$$
\nabla f(g)=\sum_{i} X_{\sigma_{i}} f(g) \cdot \sigma_{i} \in \mathfrak{s} \otimes_{\mathbb{R}} \mathbb{C}
$$

Let $|\cdot|_{\mathfrak{s}}$ be the $K$-invariant norm on $\mathfrak{s}$. The essential property is the integration by parts identity

$$
\int_{\Gamma \backslash G / K}\left\langle\nabla F_{1}, \nabla F_{2}\right\rangle_{\mathfrak{s}}=\int_{\Gamma \backslash G / K}-\Delta F_{1} \cdot \bar{F}_{2}
$$

for $F_{1}, F_{2} \in C_{c}^{\infty}(\Gamma \backslash G / K)$. Thus, extending $\nabla$ by continuity in the $\mathfrak{B}^{1}$ topology, $\nabla F$ exists (in an $L^{2}$ sense) for $F \in \mathfrak{B}^{1}\left(C_{b}\right)$. Likewise,

$$
|f|_{\mathfrak{B}^{1}\left(C_{b}\right)}^{2}=|f|_{L^{2}\left(C_{b}\right)}^{2}+\left||\nabla f|_{\mathfrak{s}}^{2}\right|_{L^{2}\left(C_{b}\right)} \ll|f|_{L^{2}(\Gamma \backslash G / K)}^{2}+\left||\nabla f|_{\mathfrak{s}}^{2}\right|_{L^{2}(\Gamma \backslash G / K)}
$$

Since $a \gg 1, \Psi_{F_{\varepsilon} \cdot \tau_{i}}$ is just $F_{\varepsilon} \cdot \tau_{i}$ on $\mathfrak{S}_{b}$, and the support of $F_{\varepsilon} \cdot \tau_{i}$ is inside the image of the cylinder $C_{a-\varepsilon}$. The map $C_{a-\varepsilon} \rightarrow \Gamma \backslash G / K$ is (uniformly) finite-to-one, so

$$
\left|\Psi_{F_{\varepsilon} \cdot \tau_{i}}\right|_{\mathfrak{B}^{1}} \ll\left|F_{\varepsilon} \cdot \tau_{i}\right|_{\mathfrak{B}^{1}\left(C_{a-\varepsilon}\right)} \leq\left|\left(F_{\varepsilon}-1\right) \cdot \tau_{i}\right|_{\mathfrak{B}^{1}\left(C_{a-\varepsilon}\right)}+\left|\tau_{i}-c_{P} f\right|_{\mathfrak{B}^{1}\left(C_{a-\varepsilon}\right)}+\left|c_{P} f\right|_{\mathfrak{B}^{1}\left(C_{a-\varepsilon}\right)}
$$

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by the triangle inequality. The middle summand goes to 0 : from above, by design,

$$
\left|\tau_{i}-c_{P} f\right|_{\mathfrak{B}^{1}\left(C_{a-\varepsilon}\right)} \ll\left|c_{P} \Psi_{\varphi_{i}}-c_{P} f\right|_{\mathfrak{B}^{1}} \ll\left|\Psi_{\varphi_{i}}-f\right|_{\mathfrak{B}^{1}} \longrightarrow 0
$$

The first and third summands require somewhat more care. Estimate

$$
\begin{gathered}
\left|\left(F_{\varepsilon}-1\right) \cdot \tau_{i}\right|_{\mathfrak{B}^{1}\left(C_{a-\varepsilon}\right)}^{2}=\int_{C_{a-\varepsilon}}\left|\left(F_{\varepsilon}-1\right) \tau_{i}\right|^{2}+\left|\nabla\left(F_{\varepsilon}-1\right) \tau_{i}\right|_{\mathfrak{s}}^{2} \\
\leq \int_{C_{a-\varepsilon}}\left|F_{\varepsilon}-1\right|^{2} \cdot\left(\left|\tau_{i}\right|^{2}+\left|\nabla \tau_{i}\right|_{\mathfrak{s}}^{2}\right)+\int_{C_{a-\varepsilon}}\left|\nabla F_{\varepsilon}\right|_{\mathfrak{s}}^{2} \cdot\left|\tau_{i}\right|^{2}+\int_{C_{a-\varepsilon}} 2\left|F_{\varepsilon}\right| \cdot\left|\nabla F_{\varepsilon}\right|_{\mathfrak{s}} \cdot\left|\tau_{i}\right| \cdot\left|\nabla \tau_{i}\right|_{\mathfrak{s}}
\end{gathered}
$$

The first summand in the latter expression goes to 0 as $\varepsilon \rightarrow 0^{+}$because $F_{\varepsilon}-1=0$ when $y \geq a$, and $\tau_{i}$ and $\left|\nabla \tau_{i}\right|_{\mathfrak{s}}$ are continuous.

We can take the orthonormal basis $\left\{\sigma_{i}\right\}$ for $\mathfrak{s}$ to have $\sigma_{1}=\left(\begin{array}{ll}1 & 0 \\ 0 & 1\end{array}\right)$, and $\sigma_{i}=\left(\begin{array}{cc}0 & * \\ * & 0\end{array}\right)$ for $i \geq 2$. Thus, $\sigma_{i} \in \mathfrak{n}+\mathfrak{k}$ for $i \geq 2$, and in terms of the Iwasawa coordinates $x=\left(x_{1}, \ldots, x_{r}\right) \in \mathbb{R}^{r}$ and $0<y \in \mathbb{R}$, for a smooth function $\varphi$ on $N \backslash G / K$, only the $\sigma_{1}$ component is non-zero:

$$
\nabla \varphi=y \frac{\partial \varphi}{\partial y} \cdot \sigma_{1}
$$

Thus,

$$
\left|\nabla F_{\varepsilon}(x+i y)\right|_{\mathfrak{s}}=\left|\frac{1}{\varepsilon} \cdot y g^{\prime}((y-a) / \varepsilon) \cdot h\right|_{\mathfrak{s}}=\frac{1}{\varepsilon} \cdot\left|y g^{\prime}((y-a) / \varepsilon)\right|<_{F} \frac{1}{\varepsilon}
$$

Similarly, since $\tau_{i}$ is a function of $y$ independent of $x, \nabla \tau_{i}=y \tau_{i}^{\prime}(y) \cdot h$. The fundamental theorem of calculus and Cauchy-Schwarz-Bunyakowsky recover an easy instance of a Sobolev inequality:

$$
\left|\tau_{i}(a-v)\right|=\left|0-\int_{0}^{v} \tau_{i}^{\prime}(a-v) d v\right| \leq\left(\int_{0}^{v}\left|\tau_{i}^{\prime}(a-v)\right|^{2} d v\right)^{\frac{1}{2}} \cdot\left(\int_{0}^{v} 1^{2} d v\right)^{\frac{1}{2}} \leq o(1) \cdot \sqrt{v}
$$

with Landau's little-o notation, since $\tau_{i}^{\prime}$ is locally $L^{2}$. Thus,

$$
\begin{gathered}
\int_{C_{a-\varepsilon}}\left|F_{\varepsilon}\right| \cdot\left|\nabla F_{\varepsilon}\right|_{\mathfrak{s}} \cdot\left|\tau_{i}\right| \cdot\left|\nabla \tau_{i}\right|_{\mathfrak{s}} \leq \frac{1}{\varepsilon} \cdot o(1) \cdot \sqrt{\varepsilon} \cdot \int_{0}^{\varepsilon}\left|\nabla \tau_{i}\right|_{\mathfrak{s}} \\
\leq \frac{1}{\varepsilon} \cdot o(1) \cdot \sqrt{\varepsilon} \cdot\left(\int_{0}^{\varepsilon}\left|\tau_{i}^{\prime}\right|^{2}\right)^{\frac{1}{2}} \cdot\left(\int_{0}^{\varepsilon} 1^{2}\right)^{\frac{1}{2}}<_{\tau_{i}} \frac{1}{\varepsilon} \cdot o(1) \cdot \sqrt{\varepsilon} \cdot \sqrt{\varepsilon}=o(1)
\end{gathered}
$$

That is, the summand $\int_{C_{a-\varepsilon}}\left|F_{\varepsilon}\right| \cdot\left|\nabla F_{\varepsilon}\right|_{\mathfrak{s}} \cdot\left|\tau_{i}\right| \cdot\left|\nabla \tau_{i}\right|_{\mathfrak{s}}$ goes to 0 . By the same estimates,

$$
\int_{C_{a-\varepsilon}}\left|\nabla F_{\varepsilon}\right|_{\mathfrak{s}}^{2} \cdot\left|\tau_{i}\right|^{2} \ll \frac{1}{\varepsilon^{2}} \int_{0}^{\varepsilon}(o(1) \cdot \sqrt{v})^{2} d v=\frac{1}{\varepsilon^{2}} \cdot o(1) \cdot \frac{\varepsilon^{2}}{2} \longrightarrow 0
$$

Thus, taking the $\varepsilon_{i}$ sufficiently small, the smooth truncations $\Psi_{\phi_{i}}-\Psi_{F_{\varepsilon_{i}} \cdot \tau_{i}}$ are in $D \cap L_{a}^{2}(\Gamma \backslash G / K)$, and still converge to $f$ in $\mathfrak{B}^{1}$.

### 11.5 Beyond the critical line: four simple examples

We return to the continuation argument. Since $\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)^{-1}$ maps $\left(\Delta-\lambda_{s}\right) h_{s}$ to a function with constant term vanishing above $\eta=a$, above $\eta=a$ the constant term of $\widetilde{E}_{a, s}$ is that of $h_{s}$, namely, $\eta^{s}$. More generally, evaluate $\widetilde{\Delta}_{a}-\lambda_{s}$ distributionally by application of $\Delta-\lambda_{s}$ : for some constant $C_{s}$,

$$
-\left(\Delta-\lambda_{s}\right) h_{s}=\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)\left(\widetilde{E}_{a, s}-h_{s}\right)=\left(\Delta-\lambda_{s}\right)\left(\widetilde{E}_{a, s}-h_{s}\right)+C_{s} \cdot \eta_{a} \quad \text { (as distributions) }
$$

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Everything else in the latter equation is meromorphic in $s$, so $C_{s}$ must be, as well. Thus, rearranging,

$$
\left(\Delta-\lambda_{s}\right) \widetilde{E}_{a, s}=-C_{s} \cdot \eta_{a} \quad \text { (as distributions) }
$$

Since $\Delta$ is $G$-invariant, it commutes with the constant-term map, and the distribution $\left(\Delta-\lambda_{s}\right) c_{P} \widetilde{E}_{a, s}$ is 0 away from $\eta=a$. The distributional differential equation

$$
\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-(r-1) y \frac{\partial}{\partial y}-\lambda_{s}\right) u=0 \quad\left(\text { on } 0<y^{r}=\eta<a\right)
$$

has solutions exactly of the form $A_{s} \eta^{s}+B_{s} \eta^{1-s}$ for constants $A_{s}, B_{s}$, so $c_{P} \widetilde{E}_{a, s}$ must be of this form in $0<\eta<a$. Since $\widetilde{E}_{a, s}$ is meromorphic in $s$, so are $A_{s}, B_{s}$. In summary,

$$
c_{P} \widetilde{E}_{a, s}= \begin{cases}\eta^{s} & (\text { for } \eta>a) \\ A_{s} \eta^{s}+B_{s} \eta^{1-s} & (\text { for } 0<\eta<a)\end{cases}
$$

By construction, $h_{s}$ is smooth, and $\left(\widetilde{\Delta}-\lambda_{s}\right)^{-1} f \in \mathfrak{B}^{1}$ for all $f \in L_{a}^{2}(\Gamma \backslash G / K)$. Thus, $\widetilde{E}_{a, s}$ is locally in $\mathfrak{B}^{1}$ in the sense that $\psi \cdot \widetilde{E}_{a, s}$ is in $\mathfrak{B}^{1}$ for any smooth cut-off $\psi \in C_{c}^{\infty}(\Gamma \backslash G / K)$. In particular, taking $\psi$ with support near $\eta=a$ and identically 1 on a neighborhood of the set where $\eta=a$, since $(N \cap \Gamma) \backslash N$ acts continuously on $\mathfrak{B}^{1}$, the constant term $c_{P}\left(\psi \cdot \widetilde{E}_{a, s}\right)$ is in $\mathfrak{B}^{1}$. Since that constant term is a function on the one-dimensional $N \backslash G / K \approx A^{+} \approx(0,+\infty)$, as in the previous section we can conclude that this constant term as a function of $t=\log y$ is in the Euclidean Sobolev space $\mathfrak{B}^{1}$ on $\mathbb{R}$. By Sobolev's imbedding on $\mathbb{R}$ [9.7], the constant term is continuous. Since $\psi$ was identically 1 near $\eta=a$, we conclude that $c_{P} \widetilde{E}_{a, s}$ itself is continuous at $\eta=a$, and

$$
A_{s} \cdot a^{s}+B_{s} \cdot a^{1-s}=a^{s} \quad(\text { for all } s)
$$

Let $\operatorname{ch}_{[a, \infty)}$ be the characteristic function of $[a, \infty)$, and

$$
\beta_{a, s}=\operatorname{ch}_{[a, \infty)}(\eta) \cdot\left(A_{s} \eta^{s}+B_{s} \eta^{1-s}-\eta^{s}\right)
$$

and form a pseudo-Eisenstein series

$$
\Phi_{a, s}(g)=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \beta_{a, s}(\gamma g)
$$

The support of $\beta_{a, s}$ is inside the set where $\eta \geq a$, and $a \gg 1$, so for each $g \in G / K$ the series has at most one non-zero summand, so converges for all $s \in \mathbb{C}$.
[11.5.1] Theorem: $A_{s} \cdot E_{s}=\widetilde{E}_{a, s}+\Phi_{a, s}$ and $\widetilde{E}_{a, s}+\Phi_{a, s}=B_{s} \cdot E_{1-s}$. Thus, $E_{s}$ has a meromorphic continuation and $E_{s}-h_{s}$ is a meromorphic $\mathfrak{B}^{1}$-valued function.
Proof: With $\widetilde{\Delta}$ as in [11.1], we have shown that $u=E_{s}-h_{s}$ is the unique solution in $\mathfrak{B}^{1}$ to

$$
\left(\widetilde{\Delta}-\lambda_{s}\right) u=-\left(\Delta-\lambda_{s}\right) h_{s}
$$

Thus, multiplying through by $A_{s}$, it suffices to prove that $\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} \cdot h_{s}$ is in $\mathfrak{B}^{1}$ and satisfies

$$
\left(\widetilde{\Delta}-\lambda_{s}\right)\left(\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} \cdot h_{s}\right)=-\left(\Delta-\lambda_{s}\right)\left(A_{s} \cdot h_{s}\right)
$$

The fact that $\widetilde{E}_{a, s}-h_{s}$ is in $\mathfrak{B}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$ motivates the rearrangement

$$
\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} \cdot h_{s}=\left(\widetilde{E}_{a, s}-h_{s}\right)+\left(\Phi_{a, s}-A_{s} h_{s}+h_{s}\right)
$$

Thus, we must show that the pseudo-Eisenstein series $F=\Phi_{a, s}-A_{s} h_{s}+h_{s}$ is in $\mathfrak{B}^{1}$.

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For integrability, by reduction theory, $\Phi_{a, s}$ is just $\varphi_{s}=\operatorname{ch}_{[a, \infty)}(\eta) \cdot\left(A_{s} \eta^{s}+B_{s} \eta^{1-s}-\eta^{s}\right)$ on $\eta>a$, so on $\eta>a$

$$
F=\Phi_{a, s}-A_{s} h_{s}+h_{s}=\left(A_{s} \eta^{s}+B_{s} \eta^{1-s}-\eta^{s}\right)-A_{s} \eta^{s}+\eta^{s}=B_{s} \eta^{1-s} \quad(\text { for } \eta>a)
$$

For $\operatorname{Re}(s)>1, \eta^{1-s}$ is square-integrable on $\eta>a$, so $F$ is in $L^{2}(\Gamma \backslash G / K)$.
To demonstrate the additional smoothness required for $F$ to be in $\mathfrak{B}^{1}$, from the rewriting of Sobolev norms in [10.7], especially [10.7.5], it suffices to show that the right-translation derivatives $\alpha F$ are in $L^{2}(\Gamma \backslash G)$ for $\alpha \in \mathfrak{g}$. By the left invariance of the right action of $\mathfrak{g}$, it suffices to prove square-integrability, on standard Siegel sets, of the derivatives of the data $\beta_{a, s}-A_{s} \tau \cdot \eta^{s}+\tau \cdot \eta^{s}$ used to form the pseudo-Eisenstein series. This data is smooth everywhere but at $\eta=a$, where it is continuous, since $A_{s} a^{s}+B_{s} a^{s}-a^{s}=0$. Further, it possesses continuous left and right derivatives at $\eta=a$, so is locally in a +1 -index Sobolev space at $\eta=a$. The data is left $N$-invariant and right $K$-invariant, and $A^{+}$normalizes $N$, so we need only consider the differential operator $y \frac{\partial}{\partial y}$ coming from the Lie algebra of $A^{+}$: the derivative of $F$ is discontinuous at $\eta=a$, and as a distribution it is, recalling that $\eta=y^{r}$, so $\eta^{\prime}=r y^{r-1}$, and $\varphi_{s}=\operatorname{ch}_{[a, \infty)} \cdot\left(A_{s} \eta^{s}+B_{s} \eta^{1-s}-\eta^{s}\right)$,

$$
\begin{gathered}
y \frac{\partial}{\partial y} F=y \frac{\partial}{\partial y}\left(\Phi_{a, s}-A_{s} h_{s}+h_{s}\right)=y \frac{\partial}{\partial y}\left(\beta_{a, s}-A_{s} \cdot \tau \cdot \eta^{s}+\tau \cdot \eta^{s}\right) \\
=y \frac{\partial}{\partial y}\left\{\begin{array}{cl}
B_{s} \eta^{1-s} & (\text { for } \eta>a) \\
-A_{s} \cdot \eta^{s}+\eta^{s} & \left(\text { for } a^{\prime} \leq \eta<a\right) \\
-A_{s} \cdot \tau \cdot \eta^{s}+\tau \cdot \eta^{s} & \left(\text { for } a^{\prime \prime} \leq \eta \leq a^{\prime}\right) \\
0 & \left(\text { for } \eta \leq a^{\prime \prime}\right)
\end{array}=\left\{\begin{array}{cl}
B_{s} \cdot(1-s) \cdot \frac{\partial \eta}{\partial y} \cdot \eta^{1-s} & (\text { for } \eta>a) \\
\left(1-A_{s}\right) \cdot s \eta^{s} & \left(\text { for } a^{\prime} \leq \eta<a\right) \\
\left(1-A_{s}\right)\left(\frac{\partial \tau}{\partial y} \cdot \eta^{s}+\tau \cdot \frac{\partial \eta}{\partial y} \cdot s \eta^{s}\right) & \left(\text { for } a^{\prime \prime} \leq \eta \leq a^{\prime}\right) \\
0 & \left(\text { for } \eta \leq a^{\prime}\right)
\end{array}\right.\right.
\end{gathered}
$$

On $a^{\prime \prime} \leq \eta \leq a$, this derivative is bounded, so the truly relevant behavior is in $\eta>a$ : for $\operatorname{Re}(s)>1$ this derivative is square-integrable on standard Siegel sets. Thus, $\Phi_{a, s}-A_{s} h_{s}+h_{s}$ is in $\mathfrak{B}^{1}$, proving that $\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s}$ is in $\mathfrak{B}^{1}$.

To show that $\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s}$ satisfies the expected equation, we justify computing the effect of differential operators on $\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s}$ distributionally, as follows. For $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$, with $\widetilde{\Delta}$ the Friedrichs extension of the restriction of $\Delta$ to $C_{c}^{\infty}(\Gamma \backslash G / K)$ as in [11.1],
$\left\langle\left(\widetilde{\Delta}-\lambda_{s}\right)\left(\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s}\right), f\right\rangle=\left\langle\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s},\left(\Delta-\bar{\lambda}_{s}\right) f\right\rangle=\left\langle\left(\Delta-\lambda_{s}\right)\left(\widetilde{E}_{a, s}+\Phi_{s, f}-A_{s} h_{s}\right), f\right\rangle$
By design, using the invariance of $\Delta$ and the local finiteness of the sum for $\Phi_{s}$, it is legitimate to compute

$$
\begin{gathered}
\left(\Delta-\lambda_{s}\right)\left(\widetilde{E}_{a, s}+\Phi_{a, s}\right)=\left(\Delta-\lambda_{s}\right) \widetilde{E}_{a, s}+\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}\left(\Delta-\lambda_{s}\right) \beta_{a, s} \circ \gamma \\
=-C_{s} \cdot \eta_{a}+C_{s} \cdot \eta_{a}=0 \quad \text { (as distributions) }
\end{gathered}
$$

Thus,

$$
\left(\widetilde{\Delta}-\lambda_{s}\right)\left(\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s}\right)=\left(\Delta-\lambda_{s}\right)\left(\widetilde{E}_{a, s}+\Phi_{a, s}-A_{s} h_{s}\right)=0-A_{s}\left(\Delta-\lambda_{s}\right) h_{s}
$$

as desired, proving $\widetilde{E}_{a, s}+\Phi_{a, s}=A_{s} \cdot E_{s}$ for $\operatorname{Re}(s)>1$. For $\operatorname{Re}(1-s)>1$, the same argument shows that $\widetilde{E}_{a, s}+\Phi_{a, s}=B_{s} \cdot E_{1-s}$. This proves the formulas in the claim. Since not both $A_{s}$ and $B_{s}$ can be identically 0 , we obtain the meromorphic continuation of $E_{s}$.
[11.5.2] Corollary: $A_{s} \cdot E_{s}=B_{s} \cdot E_{1-s}$.
In particular, neither $A_{s}$ nor $B_{s}$ is identically 0 , and with $a(s)=B_{s} / A_{s}, E_{1-s}=a(s) \cdot E_{s}$. The relation $c_{P} E_{s}=\eta^{s}+c_{s} \eta^{1-s}$ gives the meromorphic continuation of $c_{s}$. Since $c_{P} E_{1-s}=\eta^{1-s}+c_{1-s} \eta^{s}$, apparently $c_{s}=a(s)=B_{s} / A_{s}$. Since $1-(1-s)=s$, we obtain $c_{s} \cdot c_{1-s}=1$ :

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[11.5.3] Corollary: $c_{s}$ has a meromorphic continuation, and $c_{s} \cdot c_{1-s}=1$.
On $\operatorname{Im}(s)=0$ and $\operatorname{Re}(s)>1, E_{s}$ and $c_{P} E_{s}$ are real-valued. Thus, the two holomorphic functions $E_{s}$ and $\overline{E_{\bar{s}}}$ agree on $(1,+\infty)$, so agree everywhere. That is, $\overline{E_{s}}=E_{\bar{s}}$. In particular, on $\operatorname{Re}(s)=\frac{1}{2}$, where $\bar{s}=1-s$,

$$
\left|c_{s}\right|^{2}=c_{s} \cdot \overline{c_{s}}=c_{s} \cdot c_{\bar{s}}=c_{s} \cdot c_{1-s}=1 \quad\left(\text { on } \operatorname{Re}(s)=\frac{1}{2}\right)
$$

proving
[11.5.4] Corollary: $\left|c_{s}\right|=1$ on $\operatorname{Re}(s)=\frac{1}{2}$, and $c_{s}$ has no pole on $\operatorname{Re}(s)=\frac{1}{2}$.
Further, we have
[11.5.5] Corollary: $E_{s}$ has no pole on $\operatorname{Re}(s)=\frac{1}{2}$.
Proof: Suppose $E_{s}$ had a pole of order $N>0$ at $s_{o}$ on the critical line $\operatorname{Re}(s)=\frac{1}{2}$. Then $\left(s-s_{o}\right)^{N} \cdot E_{s}$ is holomorphic at $s=s_{o}$, gives a not identically automorphic form, and has vanishing constant term there. From

$$
\wedge^{a}\left(s-s_{o}\right)^{N} E_{s}=\left(s-s_{o}\right)^{N} \wedge^{a} E_{s}
$$

and using the Maass-Selberg relations [1.11] with $s=s_{o}+\varepsilon$ and $r=\bar{s}_{o}+\varepsilon=1-s_{o}+\varepsilon$ with $0<\varepsilon \in \mathbb{R}$, since $\left(s-s_{o}\right) \cdot c_{s} \rightarrow 0$ at $s=s_{o}$, suppressing measure-normalizations,

$$
\begin{gathered}
\left|\left(s-s_{o}\right)^{N} E_{s}\right|^{2}=\varepsilon^{2 N} \cdot\left(\frac{a^{s+r-1}}{s+r-1}+c_{s} \frac{a^{(1-s)+r-1}}{(1-s)+r-1}+c_{r} \frac{a^{s+(1-r)-1}}{s+(1-r)-1}+c_{s} c_{r} \frac{a^{(1-s)+(1-r)-1}}{(1-s)+(1-r)-1}\right) \\
\quad=\varepsilon^{2 N} \cdot\left(\frac{a^{2 \varepsilon}}{2 \varepsilon}+c_{s_{o}+\varepsilon} \frac{a^{1-2 s_{o}-2 \varepsilon}}{1-2 s_{o}-2 \varepsilon}+c_{1-s_{o}+\varepsilon} \frac{a^{2 s_{o}-1+2 \varepsilon}}{2 s_{o}-1+2 \varepsilon}+c_{s_{o}+\varepsilon} c_{1-s_{o}+\varepsilon} \frac{a^{-2 \varepsilon}}{-2 \varepsilon}\right) \longrightarrow 0
\end{gathered}
$$

contradiction. Thus, $E_{s}$ has no pole on the critical line.
Toward proving moderate growth of the meromorphic continuation of $E_{s}$ :
[11.5.6] Claim: $E_{s}$ meromorphically continues as a $C^{\infty}(\Gamma \backslash G / K)$-valued function.
Proof: As earlier, let

$$
\chi_{s}\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=|a|^{2 s}
$$

and put

$$
\varphi_{s}(n m k)=\chi_{s}(m) \quad(\text { for } n \in N, k \in K, \text { and } m \in M)
$$

Up to scalar multiples, $\varphi_{s}$ is the unique function on $G$ that is right $K$-invariant, left $N$-invariant, and $\varphi_{s}(m g)=\chi_{s}(m) \cdot \varphi_{s}(g)$. The function $s \rightarrow \varphi_{s}$ is a holomorphic $C^{o}(G / K)$-valued function on $\mathbb{C}$. For $\psi \in C_{c}^{\infty}(K \backslash G / K)$, the image

$$
\left(\psi \cdot \varphi_{s}\right)(g)=\int_{G} \varphi(h) \varphi_{s}(g h) d h
$$

is again left $N$-invariant, left $M, \chi_{s}$-equivariant, and right $K$-invariant. Thus, $\psi \cdot \varphi_{s}=\mu_{s}(\psi) \cdot \varphi_{s}$ for scalar $\mu_{s}(\psi)$. Since $s \rightarrow \varphi_{s}$ is holomorphic $C^{o}(G / K)$-valued, $s \rightarrow \mu_{s}(\psi)$ is holomorphic $\mathbb{C}$-valued for each such $\psi$. By non-degeneracy [14.1.5], there exists $\psi$ such that the function $s \rightarrow \mu_{s}(\psi)$ is not identically 0 . In the region of convergence $\operatorname{Re}(s)>1$, from $E_{s}=\sum_{\gamma \in(\Gamma \cap P) \backslash \Gamma} \varphi_{s} \circ \gamma$, also $\psi \cdot E_{s}=\mu_{s}(\psi) \cdot E_{s}$. Exactly what we are missing at this point is knowledge of what topological vector space of functions (or distributions) the meromorphically continued Eisenstein series may lie in, so we cannot directly assert much about $\psi \cdot E_{s}$ outside the region of convergence. (Otherwise we could apply the identity principle from complex analysis to the latter identity.) Rather, we approach this a little indirectly, as follows.

Since $\Delta$ commutes with $G, \mathfrak{B}^{1}$ is stable under the action of $\psi \in C_{c}^{\infty}(K \backslash G / K)$. From the meromorphic continuation of $E_{s}-h_{s}$ as $\mathfrak{B}^{1}$-valued function, we have the meromorphic continuation of

$$
\psi \cdot\left(E_{s}-h_{s}\right)=\mu_{s}(\psi) \cdot E_{s}-\psi \cdot h_{s}
$$

as $\mathfrak{B}^{1}$-valued function. In fact, for $F \in \mathfrak{B}^{1}$, by [14.5], $\psi \cdot F$ is in $C^{\infty}(\Gamma \backslash G / K)$. By construction, $h_{s} \in C^{\infty}(\Gamma \backslash G / K)$. Rearranging,

$$
\mu_{s}(\psi) \cdot E_{s}=\psi \cdot\left(E_{s}-h_{s}\right)+\psi \cdot h_{s}
$$

Dividing through by $\mu_{s}(\eta)$ for some $\eta$ with $\mu_{s}(\eta) \neq 0$ exhibits the meromorphically continued $E_{s}$ as a smooth-function-valued function.
[11.5.7] Corollary: $E_{s}$ has a meromorphic continuation as $C^{o}(\Gamma \backslash G / K)$-valued function, so it makes sense to address the issue of its moderate growth.

Finally, we have
[11.5.8] Theorem: Away from poles, the meromorphically continued $E_{s}$ is of moderate growth.
Proof: By [11.5.1] and [11.5.7], (at least) the pointwise values of the meromorphic continuation are given by

$$
E_{s}=A_{s}^{-1} \cdot\left(\widetilde{E}_{a, s}+\Phi_{a, s}\right)
$$

where $\widetilde{E}_{a, s}=h_{s}-\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)^{-1}\left(\Delta-\lambda_{s}\right) h_{s}$ and $\Phi_{a, s}$ is the pseudo-Eisenstein series formed from $\beta_{a, s}=$ $\operatorname{ch}_{[a, \infty)} \cdot\left(A_{s} \eta^{s}+B_{s} \eta^{1-s}-\eta^{s}\right)$. Since $a \gg 1$, in the region $\eta \geq a$ the function $\Phi_{a, s}$ is just $\beta_{a, s}$ itself, which is $A_{s} \eta^{s}+B_{s} \eta^{1-s}-\eta^{s}$, which is of moderate growth in standard Siegel sets. The computation above shows continuity at $\eta=a$. The pseudo-Eisenstein series $h_{s}$ of [11.1] made from $\tau \cdot \eta^{s}$ with smooth cut-off $\tau$ is a locally finite sum, so is smooth, so certainly continuous. For $\eta \geq a$, its value is just $\eta^{s}$, which is of moderate growth for all $s$. Thus, to show that $\widetilde{E}_{a, s}$ is of moderate growth even after meromorphic continuation, it suffices to show that $\left(\widetilde{\Delta}_{a}-\lambda_{s}\right)^{-1}\left(\Delta-\lambda_{s}\right) h_{s}$ is of moderate growth.

Again, the pseudo-Eisenstein series $h_{s}=\Psi_{\tau \cdot \eta^{s}}$ is a locally finite sum, so it is legitimate to compute

$$
\left(\Delta-\lambda_{s}\right) h_{s}=\left(\Delta-\lambda_{s}\right) \Psi_{\tau \cdot \eta^{s}}=\Psi_{\left(\Delta-\lambda_{s}\right)\left(\tau \cdot \eta^{s}\right)}
$$

Since differential operators do not increase support, $f_{s}=\left(\Delta-\lambda_{s}\right)\left(\tau \cdot \eta^{s}\right)$ is smooth and supported in $\left[a^{\prime \prime}, a^{\prime}\right]$. It is visibly a holomorphic $C^{\infty}(0,+\infty)$-valued function of $s \in \mathbb{C}$. Its uniform compact support in $\left[a^{\prime \prime}, a^{\prime}\right]$ implies that $s \rightarrow f_{s}$ is certainly a holomorphic $C_{c}^{\infty}(\Gamma \backslash G / K)$-valued function of $s$, in fact taking values in the Fréchet subspace of functions supported in $\left[a^{\prime \prime}, a^{\prime}\right]$.

Given a uniformly compactly supported holomorphic family $f_{s} \in C_{c}^{\infty}(N \backslash G / K) \approx C_{c}^{\infty}(0,+\infty)$, in light of [11.3.4] we solve equations $\left(\Delta-\lambda_{s}\right) u=f_{s}+c \cdot \eta_{a}$ with $c \in \mathbb{C}$ (depending on $s$ ) for $u$ on $N \backslash G / K \approx(0,+\infty)$ with sufficient decay at $0^{+}$to form a pseudo-Eisenstein series $\Psi_{u}$, giving $\left(\Delta-\lambda_{s}\right) \Psi_{u_{s}}=\Psi_{f_{s}}$. In Iwasawa coordinates, the equation is

$$
y^{2} \frac{\partial^{2}}{\partial y^{2}} u-(r-2) \frac{\partial}{\partial y} u-\lambda_{s} u=f_{s}+c \cdot \delta_{a}
$$

Letting $x=\log y$, with $F_{s}(x)=f_{s}\left(e^{x}\right)$ and $v(x)=u\left(e^{x}\right)$, this becomes

$$
v^{\prime \prime}-(r-1) v^{\prime}-\lambda_{s} v=F+c \cdot \delta_{\log a}
$$

Taking Fourier transform in a normalization that suppresses some factors of $2 \pi$,

$$
(-i \xi)^{2} \widehat{v}-(r-2)(-i \xi) \widehat{v}-\lambda_{s} \widehat{v}=\widehat{F}_{s}+c \cdot a^{-i \xi}
$$

or

$$
\widehat{v}(\xi)=-\frac{\widehat{F}(\xi)+c \cdot a^{-i \xi}}{\xi^{2}-(r-2) i \xi+\lambda_{s}}
$$

Since $F$ is a test function, $\widehat{F}$ is an entire function such that $\widehat{F}\left(x+i y_{o}\right)$ is (uniformly) in the Schwartz space for each fixed $y_{o}$. (We need little more about Paley-Wiener spaces than this idea.) Division by a quadratic

## 11. Meromorphic continuation of Eisenstein series

polynomial produces a function holomorphic in a strip along $\mathbb{R}$ not including either of the two poles. The two poles occur at the zeros of the denominator:

$$
\frac{(r-2) i \pm \sqrt{-(r-2)^{2}-4 \lambda_{s}}}{2}
$$

Fix $\varepsilon>0$. Given a bound $|\operatorname{Re}(s)| \leq B$, for $\operatorname{Im}(s) \gg_{B} 1$, those poles are outside the strip $S=\{z \in \mathbb{C}$ : $|\operatorname{Im}(z)| \leq 1+\varepsilon\}$. Thus, $\widehat{v}$ is holomorphic on an open set containing $S$ and has decay like $1 / \xi^{2}$ on horizontal lines inside that strip. Thus, in the Fourier inversion integral

$$
v(x)=\frac{1}{2 \pi} \int_{\mathbb{R}} e^{i \xi x} \widehat{v}(\xi) d \xi
$$

we can move the contour up to $\mathbb{R}+i(1+\varepsilon)$, giving

$$
v(x)=\frac{1}{2 \pi} \int_{\mathbb{R}} e^{i(\xi+i \varepsilon) x} \widehat{v}(\xi+i \varepsilon) d \xi=e^{-\varepsilon x} \frac{1}{2 \pi} \int_{\mathbb{R}} e^{i \xi x} \widehat{v}(\xi+i \varepsilon) d \xi
$$

Thus, $v(x) \ll e^{-(1+\varepsilon) x}$, giving genuine exponential decrease for $x \rightarrow+\infty$. Similarly, moving the contour down gives exponential decrease $v(x) \ll e^{-(1+\varepsilon)|x|}$ for $x \rightarrow-\infty$. Then $u(y)=u_{s}(y)=v(\log y)$ satisfies $u_{s}(y) \ll y^{1+\varepsilon}$ as $y \rightarrow 0^{+}$, and $u(y) \ll y^{-(1+\varepsilon)}$ as $y \rightarrow+\infty$. Thus, the pseudo-Eisenstein series $\Psi_{u_{s}}$ converges absolutely since the sum for $\Psi_{u_{s}}$ is dominated termwise by the sum for an absolutely convergent Eisenstein series [1.9.1]. Further, as it is termwise dominated by an absolutely convergent Eisenstein series, by [1.9.1] $\Psi_{u_{s}}$ is continuous and of moderate growth.

Having available a choice of the constant $c$ is necessary, since we must adjust $\Psi_{u}$ to have constant term vanishing above height $\eta=a$. Choose it so that $c_{P} \Psi_{u}$ vanishes at $\eta=a$. Since $a \gg 1$, by reduction theory the truncation $\wedge^{a} \Psi_{u}$ has constant term vanishing at and above height $a$. Since $a \gg 1$, this truncation is itself a pseudo-Eisenstein series, and still $\left(\Delta-\lambda_{s}\right) \wedge^{a} \Psi_{u}$ differs from $\Psi_{f}$ only by a multiple of $\eta_{a}$. By [11.3.4], $\left(\widetilde{\Delta}_{a}-\lambda_{s}\right) \wedge^{a} \Psi_{u}=\Psi_{f}$.

Thus, for a given bound $|\operatorname{Re}(s)| \leq B$, there is $C$ sufficiently large so that for $|\operatorname{Im}(s)| \geq C$ we have meromorphic continuation of $E_{s}$ as a (continuous) moderate-growth function.

For $|\operatorname{Im}(s)|<C$, we can express $E_{s}$ as a vector-valued Cauchy integral along a circular path $\gamma$ that lies inside the union $U$ of regions $\operatorname{Re}(s) \geq B, \operatorname{Re}(s) \leq 1-B$, and $|\operatorname{Im}(s)| \geq C$, and does not run through any poles of $E_{s}$. In $\operatorname{Re}(s) \leq 1-B$ the Eisenstein series is (continuous) of moderate growth, via the functional equation. Thus, $E_{s}$ is of moderate growth throughout $U$, and in particular along $\gamma$. Let $Z$ be the collection of poles of $E_{s}$ (as meromorphic $C^{o}(\Gamma \backslash G / K)$-valued function) inside $\gamma$, and $P(z)=\prod_{z_{j} \in Z}\left(z-z_{j}\right)$. For each $g \in G$

$$
P(s) \cdot E_{s}(g)=\frac{1}{2 \pi i} \int_{\gamma} \frac{P(z) \cdot E_{z}(g)}{z-s} d z
$$

In fact, on $\gamma, z \rightarrow\left(s \rightarrow P(z) E_{s} /(z-s)\right.$ is a compactly-supported, continuous, moderate-growth-functionvalued function of $z$, so the vector-valued Cauchy integral

$$
P(s) \cdot E_{s}=\frac{1}{2 \pi i} \int_{\gamma} \frac{P(z) \cdot E_{z}}{z-s} d z
$$

as in [15.2] exists as a Gelfand-Pettis integral [14.1] lying in that same space of functions.

### 11.6 Exotic eigenfunctions: four simple examples

In addition to cuspforms, there must be new, exotic eigenfunctions for the operators $\widetilde{\Delta}_{a}$, which are not eigenfunctions for $\Delta$.
[11.6.1] Claim: Take $a \gg 1$. If $a^{w}+c_{w} a^{1-w}=0$, then the truncation $\wedge^{a} E_{w}$ is an eigenfunction for $\widetilde{\Delta}_{a}$. Conversely, if $\wedge^{a} E_{w}$ is an eigenfunction for $\widetilde{\Delta}_{a}$, then $a^{w}+c_{w} a^{1-w}=0$. In particular, for $a^{w}+c_{w} a^{1-w}=0$, we have $\left(\Delta-\lambda_{w}\right) \wedge^{a} E_{w}=2(1-2 w) a^{w+\frac{1}{r}} \cdot \eta_{a}$. (Proof just below.)

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Since $\widetilde{\Delta}_{a}$ is a non-positive self-adjoint operator, any eigenvalues are non-positive real, giving
[11.6.2] Corollary: If $a^{w}+c_{w} a^{1-w}=0$, then either $\operatorname{Re}(w)=\frac{1}{2}$ or $w \in[0,1]$.
[11.6.3] Remark: An argument-principle discussion shows that there are infinitely-many values $w$ on $\operatorname{Re}(w)=\frac{1}{2}$ such that $a^{w}+c_{w} a^{1-w}=0$.
[11.6.4] Remark: Thus, zeros $w$ of $a^{w}+c_{w} a^{1-w}$ give eigenvalues $\lambda_{w}=w(w-1)$ of $\widetilde{\Delta}_{a}$ for $a \gg 1$. A spectral characterization of the global automorphic Sobolev spaces $\mathfrak{B}^{s}$ will prove a converse, that for $\lambda_{w}<-1 / 4$, the only eigenvalues arise from zeros of $a^{w}+c_{w} a^{1-w}$, and the corresponding exotic eigenfunctions are corresponding trucated Eisenstein series.
Proof: With $a \gg 1$, in a fundamental domain, away from $\eta=a$ we have $\left(\Delta-\lambda_{w}\right) \wedge^{a} E_{w}=0$ locally. In $\eta \gg 1$, the differential operator annihilates all Fourier components of $E_{w}$ but the constant term, and in the lower part of a Siegel set the operator does also annihilate the constant term.

We first do the slightly simpler version of this computation for $S L_{2}(\mathbb{Z})$, in which case $\eta=y$. To compute near $y=a$, let $H$ be the Heaviside function $H(y)=0$ for $y<0$ and $H(y)=1$ for $y>0$. Thus, near $y=a$, as functions of $y$ independent of $x$,

$$
\begin{gathered}
\left(\Delta-\lambda_{w}\right) \wedge^{a} E_{w}=\left(\Delta-\lambda_{w}\right)\left(H(a-y) \cdot\left(y^{w}+c_{w} y^{1-w}\right)\right)=\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-w(w-1)\right)\left(H(a-y) \cdot\left(y^{w}+c_{w} y^{1-w}\right)\right) \\
=y^{2}\left(H^{\prime \prime}(a-y)\left(y^{w}+c_{w} y^{1-w}\right)+2 H^{\prime}(a-y)\left(y^{w}+c_{w} y^{1-w}\right)^{\prime}+H(a-y)\left(y^{w}+c_{w} y^{1-w}\right)^{\prime \prime}\right) \\
-w(w-1) H(a-y)\left(y^{w}+c_{w} y^{1-w}\right) \\
=y^{2}\left(-\delta_{a}^{\prime} \cdot\left(y^{w}+c_{w} y^{1-w}\right)-2 \delta_{a} \cdot\left(w y^{w-1}+(1-w) c_{w} y^{-w}\right)\right)
\end{gathered}
$$

For $a^{w}+c_{w} a^{1-w}=0$, the term with $\delta_{a}^{\prime}$ vanishes, and the rest simplifies to

$$
\left(\Delta-\lambda_{w}\right) \wedge^{a} E_{w}=-2 a \delta_{a} \cdot\left(w a^{w}+(1-w) c_{w} a^{1-w}\right)=-2 \delta_{a} \cdot(2 w-1) a^{w+1}
$$

on functions of $y$ independent of $x$. Thus, this is $2(2 w-1) a^{w+1} \cdot \eta_{a}$. If $a^{w}+c_{w} a^{1-w} \neq 0$, the term with $\delta_{a}^{\prime}$ remains, and is not inside $\mathfrak{B}^{-1}$, so in that case $\wedge^{a} E_{w}$ is not an eigenfunction.

More generally, with height $\eta=y^{r}$ with $r=1,2,3,4, \lambda_{w}=r^{2} \cdot w(w-1)$, and $\Delta=y^{2}\left(\Delta_{x}+\frac{\partial^{2}}{\partial y^{2}}\right)-(r-1) y \frac{\partial}{\partial y}$, the truncated Eisenstein series $\wedge^{a} E_{w}$ is annihilated by $\Delta$ except near (images of) $\eta=a$, at which a messier computation gives

$$
\begin{gathered}
\left(\Delta-\lambda_{w}\right) \wedge^{a} E_{w}=\left(\Delta-\lambda_{w}\right)\left(H(a-y) \cdot\left(y^{r w}+c_{w} y^{r(1-w)}\right)\right) \\
=\left(y^{2} \frac{\partial^{2}}{\partial y^{2}}-r^{2} w(w-1)-(r-1) y \frac{\partial}{\partial y}\right)\left(H(a-y) \cdot\left(y^{r w}+c_{w} y^{r(1-w)}\right)\right) \\
=y^{2}\left(H^{\prime \prime}(a-y)\left(y^{r w}+c_{w} y^{r(1-w)}\right)+2 H^{\prime}(a-y)\left(y^{r w}+c_{w} y^{r(1-w)}\right)^{\prime}+H(a-y)\left(y^{r w}+c_{w} y^{r(1-w)}\right)^{\prime \prime}\right) \\
-(r-1) y\left(H^{\prime}(a-y)\left(y^{r w}+c_{w} y^{r(1-w)}\right)+H(a-y)\left(y^{r w}+c_{w} y^{r(1-w)}\right)^{\prime}\right)-r^{2} w(w-1) H(a-y)\left(y^{r w}+c_{w} y^{r(1-w)}\right) \\
=y^{2}\left(-\delta_{a}^{\prime} \cdot\left(y^{r w}+c_{w} y^{r(1-w)}\right)-2 \delta_{a} \cdot\left(r w y^{r w-1}+r(1-w) c_{w} y^{r(1-w)-1}\right)\right)-(r-1) y\left(-\delta_{a}\left(y^{r w}+c_{w} y^{r(1-w)}\right)\right)
\end{gathered}
$$

At $y^{r}=\eta=a$, for $a^{w}+c_{w} a^{1-w}=0$, the term with $\delta_{a}^{\prime}$ vanishes, as does the $(r-1) y\left(-\delta_{a}\left(y^{r w}+c_{w} y^{r(1-w)}\right)\right)$ term, and the rest simplifies to the indicated expression.

### 11.7 Up to the critical line: $S L_{r}(\mathbb{Z})$

Now take $G=S L_{r}(\mathbb{R}), \Gamma=S L_{r}(\mathbb{Z})$, and $K=S O_{r}(\mathbb{R})$. At various moments, it is convenient to consider $G=Z \backslash G L_{r}$ instead, where $Z$ is the center, but nothing we do depends on the distinction. We only consider Eisenstein series for maximal proper parabolics, and with cuspidal data on the Levi components, as in [3.9] and [3.11]. As in the four simplest cases, it is relatively easy to meromorphically continue these Eisenstein series up to the critical line.

As in chapter 3, there are two qualitatively different types of maximal proper parabolics, namely, the self-associate $P=P^{r, r} \subset G L_{2 r}$, and non-self-associate $P=P^{r_{1}, r_{2}} \subset G L_{r_{1}+r_{2}}$ with $r_{1} \neq r_{2}$, with associate $Q=P^{r_{2}, r_{1}}$.

Fix a maximal proper parabolic $P=P^{r_{1}, r_{2}}$ with Levi decomposition $P=N M$, and fix cuspidal data $f=f_{1} \otimes f_{2}$ on the Levi component $M \approx G L_{r_{1}} \times G L_{r_{2}}$. We assume that $f_{1}$ and $f_{2}$ are cuspforms in a strong sense [: they are eigenfunctions for the corresponding invariant Laplacians, are of rapid decay in Siegel sets (in particular, are bounded on the respective groups), and there exist test functions $\beta_{1}, \beta_{2}$ on the respective groups such that $\beta_{j} \cdot f_{j}=f_{j}$. In particular, all derivatives (whether right invariant under maximal compacts or not) have similarly good decay and smoothness. Recall [3.9] that pseudo-Eisenstein series with cuspidal data $f$ are formed from test functions $\psi \in C_{c}^{\infty}(0,+\infty)$ as follows. Let

$$
\varphi(z n m k)=\varphi_{\psi, f}(z n m k)=\psi\left(\frac{\left|\operatorname{det} m_{1}\right|^{r_{2}}}{\left|\operatorname{det} m_{2}\right|^{r_{1}}}\right) \cdot f_{1}\left(m_{1}\right) \cdot f_{2}\left(m_{2}\right)
$$

and the corresponding pseudo-Eisenstein series is

$$
\Psi_{\varphi}=\Psi_{\psi, f}=\sum_{\gamma \in(P \cap \Gamma) \backslash \Gamma} \varphi_{\psi, f} \circ \gamma
$$

By [3.11.1], because of the cuspidal data, the only non-vanishing constant terms $c_{Q} \Psi_{\varphi}^{P}$ are for $Q=P$, or $Q=P^{r_{2}, r_{1}}$ when $P$ is not self-associate. As in [3.16], these pseudo-Eisenstein series admit spectral decompositions in terms of the (genuine) Eisenstein series $E_{s, f}=E_{s, f}^{P}$ for $P$ with the same cuspidal data $f$ : with

$$
\varphi_{s, f}(n m k)=\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right|^{s} \cdot f_{1}\left(m_{1}\right) \cdot f_{2}\left(m_{2}\right)
$$

with $m=\left(\begin{array}{cc}m_{1} & 0 \\ 0 & m_{2}\end{array}\right) \in M^{P}, n \in N, k \in K$, the corresponding Eisenstein series is

$$
E_{s, f}=\sum_{\gamma \in(P \cap \Gamma) \backslash \Gamma} \varphi_{s, f} \circ \gamma
$$

Again from [3.11.3], $c_{Q} E_{s, f}^{P}=0$ unless $Q=P$ or $Q$ is the associate of $P$. For $f_{1}$ and $f_{2}$ eigenfunctions of the respective Laplacians, by [3.11.11] the function $\varphi_{s, f}^{P}$ is an eigenfunction for the invariant Laplacian on $G$, and $E_{s, f}^{P}$ is an eigenfunction. In particular, letting $\mu_{j}$ be the eigenvalue of $f_{j}$ for the Laplacian on $G L_{r_{j}}$ for $j=1,2$, letting

$$
\lambda_{s, f}=r_{1} r_{2}\left(r_{1}+r_{2}\right)\left(s^{2}-s\right)+\mu_{1}+\mu_{2}
$$

we have

$$
\Delta \cdot E_{s, f}^{P}=\lambda_{s, f} \cdot E_{s, f}^{P}
$$

With this normalization, the eigenvalue is invariant under $s \longrightarrow 1-s$.
Let $\mathfrak{E}(P, f)$ be the space of pseudo-Eisenstein series for $P$ formed with the given cuspidal data $f$. An analogue of [3.11.11] for pseudo-Eisenstein series with cuspidal data:
[11.7.1] Claim: $\mathfrak{E}(P, f)$ is stable under $\Delta$. Explicitly, using coordinate $y>0$ on the ray $(0,+\infty)$, $\Delta \Psi_{\psi, f}=\Psi_{\beta, f}$ with test function $\beta$ given in terms of $\psi$ by

$$
\beta=\left(r_{1} r_{2}\left(r_{1}+r_{2}\right) y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)+\mu_{1}+\mu_{2}\right) \psi
$$

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Proof: This reduces to [3.11.11] via a primitive initial form of the spectral decomposition of $\Psi_{\psi, f}$ in terms of $E_{s, f}$ in the proof of [3.16.1], not requiring any meromorphic continuation of $E_{s, f}$ :

$$
\Psi_{\psi, f}=\frac{1}{2 \pi i} \int_{\sigma-i \infty}^{\sigma+i \infty} \mathcal{M} \psi(s) \cdot E_{s, f} d s \quad(\text { for } \sigma \gg 1)
$$

Applying $\Delta$ multiplies $E_{s, f}$ by $r_{1} r_{2}\left(r_{1}+r_{2}\right) s(s-1)+\lambda_{1}+\lambda_{2}$ by [3.11.11]. At the same time, $y \frac{\partial}{\partial y} y^{s}=s \cdot y^{s}$, so from Mellin inversion

$$
s \cdot \mathcal{M} \psi=\mathcal{M}\left(y \frac{\partial}{\partial y} \psi\right)
$$

This gives the assertion.
Let $\mathfrak{E}^{0}$ be the completion of $\mathfrak{E}(P, f)$ in $L^{2}(\Gamma \backslash G / K)$. Let $T_{f}$ be the restriction of $\Delta$ to $\mathfrak{E}(P, f)$, and $\widetilde{S}_{f}$ its Friedrichs extension. Let $\mathfrak{E}^{1}$ be the completion of $\mathfrak{E}(P, f)$ in the $\mathfrak{B}^{1}$-norm given by

$$
|f|_{\mathfrak{B}^{1}}^{2}=\langle(1-\Delta) f, f\rangle_{L^{2}(\Gamma \backslash G / K)}
$$

and $\mathfrak{E}^{2}$ the completion of $\mathfrak{E}(P, f)$ in the $\mathfrak{B}^{2}$ norm

$$
|f|_{\mathfrak{B}^{2}}^{2}=\left\langle(1-\Delta)^{2} f, f\right\rangle_{L^{2}(\Gamma \backslash G / K)}
$$

The domain of any self-adjoint extension of $T_{f}$ necessarily contains $\mathfrak{E}^{2}$, and the domain of $\widetilde{S}_{f}$ is contained in $\mathfrak{E}^{1}$. More generally, for non-negative integer $k$, let $\mathfrak{E}^{k}$ be the completion of $\mathfrak{E}(P, f)$ in the $\mathfrak{B}^{k}$ norm

$$
|f|_{\mathfrak{B}^{k}}^{2}=\left\langle(1-\Delta)^{k} f, f\right\rangle_{L^{2}(\Gamma \backslash G / K)}
$$

Let $\mathfrak{B}^{\infty}=\bigcap_{k} \mathfrak{B}^{k}=\lim _{k} \mathfrak{B}^{k}$.
[11.7.2] Corollary: For positive integer $k$, the $\mathfrak{E}^{k}$-norm of $\Psi_{\psi, f}$ is the $L^{2}$ norm of $\Psi_{\beta, f}$, where

$$
\beta=\left(1-\left(r_{1} r_{2}\left(r_{1}+r_{2}\right) y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)+\mu_{1}+\mu_{2}\right)\right)^{k} \psi
$$

(Immediate from [11.7.1].)
We grant the general form of the constant term along $P$ (see [3.11.9]): this requires an assumption that in the cuspidal data $f=f_{1} \otimes f_{2}$, both $f_{1}$ and $f_{2}$ are the unique cuspforms on $G L_{r_{1}}$ with their respective eigenvalues (and right invariant under compacts, and left invariant under the respective groups $S L_{r_{i}}(\mathbb{Z})$, as opposed to other subgroups). Then

$$
\left\{\begin{aligned}
c_{P} E_{s, f}^{P} & =\varphi_{s, f}^{P} & & \left(\text { for } r_{1} \neq r_{2}(\text { not self-associate })\right) \\
c_{P} E_{s, f}^{P} & =\varphi_{s, f}^{P}+c_{s, f} \varphi_{1-s, f^{w}}^{P} & & \left(\text { for } r_{1}=r_{2}(\text { self-associate }), \text { meromorphic } c_{s, f}\right) \\
c_{Q} E_{s, f}^{P} & =c_{s, f}^{Q} \cdot \varphi_{1-s, f^{w}}^{Q} & & \left(\text { for } r_{1} \neq r_{2}, Q=P^{r_{2}, r_{1}}, \text { meromorphic } c_{s, f}^{Q}\right)
\end{aligned}\right.
$$

In fact, we do not use the precise nature of $c_{s, f}$.
Let

$$
\alpha(m)=\left|\frac{\left(\operatorname{det} m_{1}\right)^{r_{2}}}{\left(\operatorname{det} m_{2}\right)^{r_{1}}}\right| \quad\left(\text { for } m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right) \in M\right)
$$

For $1 \ll_{b} a^{\prime \prime}<a^{\prime}$, define a real-valued smooth cut-off function by

$$
\tau(y)=\left\{\begin{array}{ll}
1 & \left(\text { for } y>a^{\prime}\right) \\
0 & \left(\text { for } y<a^{\prime \prime}\right)
\end{array} \quad(\text { for } y \in(0,+\infty))\right.
$$

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Form a cuspidal-data pseudo-Eisenstein series $h_{s, f}$ by winding up a smoothly cut-off version of $\varphi_{s, f}$ : with $\psi(y)=\tau(y) \cdot y^{s}$, put

$$
h_{s, f}=\Psi_{\psi, f}
$$

[11.7.3] Lemma: The sum for $h_{s, f}$ is absolutely convergent, uniformly on compacts, for all $s \in \mathbb{C}$. Further, $h_{s, f} \in \mathfrak{B}^{\infty}$.
Proof: By reduction theory [3.3], for $a \gg b$ large enough so that, for $\gamma \in \Gamma$, if $\gamma \mathfrak{S}_{a} \cap \mathfrak{S}_{b} \neq \phi$ then $\gamma \in B \cap \Gamma$ with minimal parabolic $B$. We increase $a^{\prime \prime}$ and $a^{\prime}$ in the definition of $\tau$, if necessary, so that this property holds for them. Of course, $\mathfrak{S}_{a}$ is not $\left(M^{P} \cap \Gamma\right)$-stable, so no (strong-sense) cuspform $f$ on $M$ could be supported on any single copy of $\mathfrak{S}_{a}$. Thus, we need a type of Siegel set adapted to $P$ : with $B$ the standard minimal parabolic, let

$$
\mathfrak{S}_{a}^{P}=\bigcup_{\gamma \in(\Gamma \cap P) /(\Gamma \cap B)} \gamma \mathfrak{S}_{a}
$$

Then $\mathfrak{S}_{a}^{P}$ is $P \cap \Gamma$-stable, and for $\gamma \in \Gamma$, if $\gamma \mathfrak{S}_{a}^{P} \cap \mathfrak{S}_{b}^{P} \neq \phi$ then $\gamma \in P \cap \Gamma$ : Indeed, suppose $\gamma \mathfrak{S}_{a}^{P} \cap \mathfrak{S}_{b}^{P} \neq \phi$. Then $\gamma \gamma_{1} \mathfrak{S}_{a} \cap \gamma_{2} \mathfrak{S}_{b} \neq \phi$ for some $\gamma_{1}, \gamma_{2} \in \Gamma \cap P$. By the choice of $a \gg b$, this implies that $\gamma_{2}^{-1} \gamma \gamma_{1} \in B \cap \Gamma$, or

$$
\gamma \in \gamma_{2}(B \cap \Gamma) \gamma_{1}^{-1}=\gamma_{2} B \gamma_{1}^{-1} \cap \Gamma \subset P \cap \Gamma
$$

Thus, for each $g \in G$ there is at most one non-zero summand in the sum defining $h_{s, f}$. The same is true of its image under $(1-\Delta)^{k}$ for every $k$, so the sum converges in $\mathfrak{B}^{k}$ for every $k$.

Thus, the pseudo-Eisenstein series $h_{s, f}$ is entire as a function-valued function of $s$. Let

$$
\widetilde{E}_{s, f}=h_{s, f}-\left(\widetilde{S}_{f}-\lambda_{s, f}\right)^{-1}\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

[11.7.4] Claim: $\widetilde{E}_{s, f}-h_{s, f}$ is a holomorphic $\mathfrak{E}^{1}$-valued function of $s$ for $\operatorname{Re}(s)>\frac{1}{2}$ and $\operatorname{Im}(s) \neq 0$.
Proof: From Friedrichs' construction [9.2], the resolvent $\left(\widetilde{\Delta}-\lambda_{s, f}\right)^{-1}$ exists as an everywhere-defined, continuous operator for $s \in \mathbb{C}$ for $\lambda_{s, f}$ not a non-positive real number, because of the non-positive-ness of $\Delta$. Further, for $\lambda_{s, f}$ not a non-positive real, this resolvent is a holomorphic operator-valued function. In fact, for such $\lambda_{s, f}$, the resolvent $\left(\widetilde{S}_{f}-\lambda_{s, f}\right)^{-1}$ injects from $L^{2}(\Gamma \backslash G / K)$ to $\mathfrak{E}^{1}$.
[11.7.5] Theorem: With $\lambda_{s, f}$ not non-positive real, $u=\widetilde{E}_{s, f}-h_{s, f}$ is the unique element of the domain of $\widetilde{S}_{f}$ such that

$$
\left(\widetilde{S}_{f}-\lambda_{s, f}\right) u=-\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

Thus, $\widetilde{E}_{s, f}$ is the usual Eisenstein series $E_{s, f}$ of [3.11] for $\operatorname{Re}(s)>1$, and gives an analytic continuation of $E_{s, f}-h_{s, f}$ as $\mathfrak{E}^{1}$-valued function to $\operatorname{Re}(s)>\frac{1}{2}$ with $s \notin\left(\frac{1}{2}, 1\right]$.
Proof: The proof is very similar to that of [11.1.3]. Uniqueness follows from Friedrichs' construction [9.2] and construction of resolvents, because $\widetilde{S}-\lambda_{s}$ is a bijection of its domain to $L^{2}(\Gamma \backslash G / K)$.

On the other hand, for $\operatorname{Re}(s)>\frac{1}{2}$ and $s \notin\left(\frac{1}{2}, 1\right], \widetilde{E}_{s, f}-h_{s, f}$ is in $L^{2}(\Gamma \backslash G / K)$, is smooth, and

$$
\begin{gathered}
\Delta\left(\widetilde{E}_{s, f}-h_{s, f}\right)=\left(\Delta-\lambda_{s, f}\right)\left(\widetilde{E}_{s, f}-h_{s, f}\right)+\lambda_{s, f} \cdot\left(\widetilde{E}_{s, f}-h_{s, f}\right)=\left(\Delta-\lambda_{s, f}\right) h_{s, f}+\lambda_{s, f} \cdot\left(\widetilde{E}_{s, f}-h_{s, f}\right) \\
=\left(\text { element of } \mathfrak{B}^{\infty}\right)+\lambda_{s, f} \cdot\left(\widetilde{E}_{s, f}-h_{s, f}\right)
\end{gathered}
$$

so is in $\mathfrak{B}^{2}$, so certainly in the domain of $\widetilde{\Delta}$. Abbreviating $H_{s, f}=\left(\Delta-\lambda_{s}\right) h_{s, f}$, it is legitimate to compute

$$
\begin{gathered}
\left(\widetilde{S}_{f}-\lambda_{s, f}\right)\left(\widetilde{E}_{s, f}-h_{s, f}\right)=\left(\widetilde{S}_{f}-\lambda_{s, f}\right)\left(\left(h_{s, f}-\left(\widetilde{S}_{f}-\lambda_{s, f}\right)^{-1} H_{s, f}\right)-h_{s, f}\right) \\
=\left(\widetilde{S}_{f}-\lambda_{s, f}\right)\left(-\left(\widetilde{S}_{f}-\lambda_{s, f}\right)^{-1} H_{s, f}\right)=-H_{s, f}
\end{gathered}
$$

Thus, $\widetilde{E}_{s, f}-h_{s, f}$ is a solution. Also, $E_{s, f}-h_{s, f}$ is a solution:

$$
\left(\Delta-\lambda_{s, f}\right)\left(E_{s, f}-h_{s, f}\right)=\left(\Delta-\lambda_{s, f}\right) E_{s, f}-\left(\Delta-\lambda_{s, f}\right) h_{s, f}=0-\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

By uniqueness, we are done.
[11.7.6] Corollary: $E_{s, f}$ has an analytic continuation to $\operatorname{Re}(s)>\frac{1}{2}$ and $s \notin\left(\frac{1}{2}, 1\right]$ as an $h_{s, f}+\mathfrak{E}^{1}$-valued function.
[11.7.7] Corollary: The function $s \rightarrow c_{s, f}$ has a meromorphic continuation to $\Re(s)>\frac{1}{2}$ (off $\left(\frac{1}{2}, 1\right]$ ).

### 11.8 Distributional characterization of pseudo-Laplacians

First, we consider the self-associate maximal proper parabolic $P=P^{r, r} \subset G=S L_{2 r}$, and cuspidal data of the symmetrical form $f=f_{1} \otimes f_{1}$, so that $f^{w}=f$. Further, without loss of generality $\bar{f}=f$. Then the argument is nearly identical to that for the simple examples in [11.3]. However, since strong-sense cuspforms $f_{1}$ are not likely to be compactly supported, the simple local argument for [11.3.2] requires some adaptation. What we do have, from [7.3.19] (and from [7.2.20] and [7.1.20] for simpler situations), is that strong-sense cuspforms are smooth, of rapid decay, that there are test functions $\varphi_{j}$ such that $\varphi_{j} \cdot f_{j}=f_{j}$, and that the $f_{j}$ are eigenfunctions for the Laplacians on the factors of the Levi component. From the theory of the constant term [8.2.5], relations $\varphi_{j} \cdot f_{j}=f_{j}$ imply that all derivatives of such cuspforms (with respect to the universal enveloping algebra of the Lie algebra) are also of rapid decay.

Take symmetrical cuspidal data $f=f_{1} \otimes f_{1}$ on $M=M^{P}$, with $f_{1}$ a cuspform in a strong sense, and with $L^{2}$ norm 1. Put

$$
\mathfrak{E}(P, f)=\left\{\Psi_{\psi, f}^{P}: \psi \in C_{c}^{\infty}(0,+\infty)\right\}
$$

We recall some context from [10.6]. Let $B$ be the standard minimal parabolic, with unipotent radical $N^{B}$ and standard Levi component $M^{B}$. Write Iwasawa decompositions $g=n m_{g} k$ with $n \in N^{B}, m \in M^{B}$. We let $\mathfrak{S}$ be a standard Siegel set stable under the (left) action of $N^{B}$ :

$$
\mathfrak{S}=\mathfrak{S}_{b}=\left\{g \in G:\left|\alpha_{j}\left(m_{g}\right)\right| \geq b, \text { for all simple roots } \alpha_{j}\right\}
$$

Take $0<b \ll 1$ such that $\mathfrak{S}_{b} \rightarrow \Gamma \backslash G$ is a surjection. For $a>b$, let $X_{a}$ be the subset of $\mathfrak{S}$ where $\beta\left(m_{g}\right) \leq a$ for all simple roots $\beta$. The quotient $(\Gamma \cap B) \backslash X_{a}$ is compact, since $\left(N^{B} \cap \Gamma\right) \backslash N$ is compact. For each simple $\operatorname{root} \beta$, let

$$
Y_{a}^{\beta}=\left\{g \in \mathfrak{S}: \beta\left(m_{g}\right) \geq a\right\}
$$

and $Y_{a}=\bigcup_{\beta} Y_{a}^{\beta}$. Thus, $\mathfrak{S}=X_{a} \cup Y_{a}$. Parallel to [10.6], let

$$
\mathfrak{E}(P, f)_{a}=\left\{F \in \mathfrak{E}(P, f): c_{P^{\prime}} F(g)=0, \text { for all } g \in Y_{a}, \text { for all standard parabolics } P^{\prime}\right\}
$$

and

$$
\mathfrak{E}_{a}^{0}=\mathfrak{B}^{0} \text {-closure of } \mathfrak{E}(P, f)_{a} \quad \mathfrak{E}_{a}^{1}=\mathfrak{B}^{1} \text {-closure of } \mathfrak{E}(P, f)_{a} \quad \mathfrak{E}_{a}^{2}=\mathfrak{B}^{2} \text {-closure of } \mathfrak{E}(P, f)_{a}
$$

It suffices to require vanishing of constant terms for maximal proper parabolics $P^{\prime}$. Further, from [3.11.1], since all pseudo-Eisenstein series in $\mathfrak{E}(P, f)$ have cuspidal data, the vanishing condition is automatically satisfied for all parabolics $P^{\prime}$ except $P$.

To be careful, since our unbounded operators should be densely defined, we note
[11.8.1] Lemma: For $a \gg 1, \mathfrak{E}(P, f)_{a}=\mathfrak{E}(P, f) \cap \mathfrak{E}_{a}^{0}$ is dense in $\mathfrak{E}_{a}^{0}$.
Proof: Having restricted our attention to the relatively small space $\mathfrak{E}(P, f)$, with $a \gg 1$, the observation [11.6.2] essentially reduces the issue to a generic, local, one-dimensional issue of smooth cut-offs, much as addressed in the proof of [10.3.1].

Let $S_{a, f}$ be $\Delta$ restricted to $\mathfrak{E}(P, f)_{a}$. Since $\Delta \Psi_{\underset{\sim}{~}, f}^{P}=\Psi_{\beta, f}^{P}$ from [11.6.1], and differential operators do not enlarge supports, $\Delta$ does stabilize $\mathfrak{E}(P, f)_{a}$. Let $\widetilde{S}_{a, f}$ be the Friedrichs extension of $S_{a, f}$ to an unbounded self-adjoint operator on $\mathfrak{E}_{a}^{0}$, with domain contained in $\mathfrak{E}_{a}^{1}$ and containing $\mathfrak{E}_{a}^{2}$.
[11.8.2] Corollary: (of [10.8]) $\widetilde{S}_{a, f}$ has compact resolvent $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1}$ (away from poles).
Proof: As usual, the crucial point is that the inclusion $\mathfrak{E}_{a}^{1} \rightarrow \mathfrak{E}_{a}^{0}$ is a restriction of the inclusion $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}$, the latter shown to be compact in [10.8]. The restrictions of compact operators are compact. The resolvents of the Friedrichs extension are continuous maps $\mathfrak{E}_{a}^{0} \rightarrow \mathfrak{E}_{a}^{1}$ composed with the inclusion $\mathfrak{E}_{a}^{1} \rightarrow \mathfrak{E}_{a}^{0}$. Continuous maps composed with compact maps are compact.

Let $M^{1}$ be the copy of $S L_{r} \times S L_{r}$ inside $M=M^{P}$, and $Z^{M}$ the center of $M$. We take representatives

$$
z_{a}=\left(\begin{array}{cc}
a^{\frac{1}{r_{1}^{2}}} \cdot 1_{r} & 0 \\
0 & 1_{r}
\end{array}\right) \quad\left(\text { for } 0<a \in \mathbb{R}^{\times}\right)
$$

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for the connected component $Z \backslash Z^{M}$ containing $1_{r}$, and let $\eta_{a}$ be the functional on $\mathfrak{E}(P, f)$ defined by

$$
\eta_{a}(F)=\int_{Z\left(\Gamma \cap M_{1}^{P}\right) \backslash M_{1}} c_{P} F\left(m^{\prime} \cdot z_{a}\right) \bar{f}\left(m^{\prime}\right) d m^{\prime} \quad(\text { for } F \in \mathfrak{E}(P, f))
$$

Then $F \in \mathfrak{E}(P, f)_{a}$ if and only if $\eta_{b^{\prime}}(F)=0$ for all $b^{\prime} \geq a$.
As in [3.8.2], pointwise vanishing conditions for constant terms can be rewritten as $L^{2}$ orthogonality to corresponding pseudo-Eisenstein series. With

$$
\Theta=\left\{\Psi_{\psi, f}^{P}: \psi \in C_{c}^{\infty}(0,+\infty) \text { with support inside }[a,+\infty)\right\}
$$

$\mathfrak{E}(P, f)_{a}$ is the intersection of $\mathfrak{E}(P, f)$ with the orthogonal complement to $\Theta$ in $L^{2}(\Gamma \backslash G / K)$.
Let $c$ be the pointwise conjugation map $c: \mathfrak{E}^{0} \rightarrow \mathfrak{E}^{0}$. Let $\mathfrak{E}^{-1}$ be the Hilbert space dual of $\mathfrak{E}^{1}$. Let $j^{*}$ be the adjoint of the inclusion $j: \mathfrak{E}^{1} \rightarrow \mathfrak{E}^{0}$, let $j^{*}$ be its adjoint. Let $\mathfrak{E}_{a}^{-1}$ be the Hilbert space dual of $\mathfrak{E}_{a}^{1}$, let $t: \mathfrak{E}_{a}^{1} \rightarrow \mathfrak{E}^{1}$ be the inclusion, with adjoint $t^{*}: \mathfrak{E}^{-1} \rightarrow \mathfrak{E}_{a}^{-1}$, giving a picture


Let $\Delta_{a}$ be $\Delta$ restricted to $C_{c}^{\infty}(\Gamma \backslash G / K) \cap L_{a}^{2}$, with Friedrichs extension $\widetilde{\Delta}_{a}$. Let $S^{\#}: \mathfrak{E}_{a}^{1} \rightarrow \mathfrak{E}_{a}^{-1}$ be as in [11.2], namely, $S^{\#}(x)(y)=\langle x, y\rangle_{\mathfrak{E}^{1}}$. Recall the re-characterization of Friedrichs extensions in [11.2]: $\widetilde{S}_{a, f} x=y$ for $x \in \mathfrak{E}_{a}^{1}$ and $y \in \mathfrak{E}_{a}^{0}$ if and only if $S^{\#} x=t^{*} \circ\left(j^{*} \circ c\right) y$, and we have
[11.8.3] Corollary: (of [11.2.2]) $\widetilde{S}_{a, f} x=y$ for $x \in \mathfrak{E}_{a}^{1}$ and $y \in \mathfrak{E}_{a}^{0}$ if and only $(\Delta \circ t) x=y+\theta$ for some $\theta$ in the $\mathfrak{B}^{-1}$-closure of $\Theta$.

Thus, as in the simpler cases, the critical fact is
[11.8.4] Claim: For $a \gg_{b} 1$, the intersection of the $(\Delta \circ t) \mathfrak{E}_{a}^{1}$ and the $\mathfrak{E}^{-1}$-closure of $\Theta$ is at most $\mathbb{C} \cdot \eta_{a}$.
Proof: Use Siegel sets $\mathfrak{S}_{b^{\prime}}^{P}$ adapted to $P$, as in the proof of [11.7.3]. Take $b^{\prime}<a$ but still $b^{\prime} \gg_{b} 1$ so that $\mathfrak{S}_{b^{\prime}}^{P}$ has the same features as $a$. The compact abelian group $A=\left(N^{P} \cap \Gamma\right) \backslash N^{P}$ acts on $C_{a^{\prime}}=\left(N^{P} \cap \Gamma\right) \backslash \mathfrak{S}_{a^{\prime}}^{P}$, and $C_{a^{\prime}}$ contains the image $C_{a}$ of $\mathfrak{S}_{a}^{P}$ in the quotient.

On one hand, by the choice of $a>_{b} 1$, for a test function $\psi$ supported in $[a,+\infty)$, on $\mathfrak{S}_{b}^{P}$ the pseudoEisenstein series $\Psi_{\psi, f} \in \Theta$ is just $\varphi_{\psi, f}$. These distributions are $N^{P}$-invariant on $\mathfrak{S}_{a}^{P}$. Taking $\mathfrak{E}^{-1}$ closure does not increase support, and does not harm the $N^{P}$-invariance. Thus, the $\mathfrak{E}^{-1}$-closure of $\Theta$ consists of $A$-invariant distributions supported in $C_{a} \subset C_{b^{\prime}}$.

On the other hand, $A$-invariants in $\mathfrak{E}_{a}^{1}$ are obtained as constant-term integrals, which are averaging integrals over the compact $A$, which exist as Gelfand-Pettis integrals with values at least as distributions. For each small $\varepsilon>0$, there is a sequence $F_{i} \in \mathfrak{E}^{1}$ supported in $C_{a-\varepsilon}$ approaching $\theta \in \mathfrak{E}^{-1}$, with $F_{i}=\Delta u-f$ with $u \in \mathfrak{E}_{a}^{1}$ and $f \in \mathfrak{E}_{a}^{0}$, since $a-\varepsilon \gg 1$, in $C_{b^{\prime}}$

$$
F_{i}=c_{P} \theta=c_{P}(\Delta u-f)=\Delta\left(c_{P} u\right)-c_{P} f
$$

and the intersection of $C_{a-\varepsilon}$ with the supports of $c_{P} u$ and $c_{P} f$ is contained in the complement $C_{a-\varepsilon}-C_{a}$. The differential operator $\Delta$ does not enlarge supports. Thus, the support of $\theta$ is contained in the boundary $\partial C_{a}$.

Thus, for each $\varepsilon>0$, we can approximate in $\mathfrak{E}^{-1}$ such $\theta$ by a sequence $\Psi_{\psi_{i}, f}$ with $\psi_{i} \in C_{c}^{\infty}(0,+\infty)$ supported in $[a-\varepsilon, a+\varepsilon]$. For $\beta \in C_{c}^{\infty}(0,+\infty)$, without loss of generality suppose the support of $\beta$ is similarly restricted. Since $a-\varepsilon \gg 1$, by reduction theory

$$
\int_{\Gamma \backslash G} \Psi_{\psi_{i}, f} \cdot \Psi_{\beta, f}=\int_{C_{a-\varepsilon}} \varphi_{\psi_{i}, f} \cdot \varphi_{\beta, f}
$$

and this integral simplifies to $\int_{a-\varepsilon}^{a+\varepsilon} \psi_{i} \cdot \beta$ with suitable measure on $(0,+\infty)$. Thus, such $\theta$ is specified by a distribution $\theta_{o}$ on $\mathbb{R}$ supported at the point $a$. By the classification of distributions supported at a point [13.14.3], $\theta_{o}$ must be a finite linear combination of Dirac delta $\delta_{a}$ and its derivatives. As in the local computations in the proof of [11.3.2], the condition $\theta \in \mathfrak{E}^{-1}$ requires that $\theta_{o}$ be at worst in $\mathfrak{B}^{-1}(\mathbb{R})$, which then requires that $\theta_{o}$ be a constant multiple of $\delta_{a}$. Thus, the $\mathfrak{E}^{-1}$-limit is a constant multiple of $\eta_{a}$. ///

Thus, the relatively simple characterization of the Friedrichs extension for $a \gg 1$ :
[11.8.5] Corollary: With $a \gg 1, \widetilde{S}_{a, f} x=y$ for $x \in \mathfrak{E}_{a}^{1}$ and $y \in \mathfrak{E}_{a}^{0}$ if and only $(\Delta \circ t) x=y+c \cdot \eta_{a}$ for some constant $c$.

### 11.9 Density lemma for $P^{r, r} \subset S L_{2 r}$

Similar to the description of $E_{s, f}$ as $\widetilde{E}_{s, f}$ above in [11.7], but with $\widetilde{S}_{a, f}$ in place of $\widetilde{S}_{f}$, with the pseudoEisenstein series $h_{s, f}$ formed from the smooth cut-off $\tau \cdot \varphi_{s, f}$ of $\varphi_{s, f}$ as in [11.7], put

$$
\widetilde{E}_{a, s, f}=h_{s, f}-\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1}\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

We already noted in [11.7] that $h_{s, f}$ is an entire, $\mathfrak{E}^{1}$-valued function, for simple reasons, given reduction theory.

For $\lambda_{s, f}$ not a non-positive real, $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1}$ is a bijection of $\mathfrak{E}_{a}^{0}$ to the domain of $\widetilde{S}_{a, f}$, so $u=\widetilde{E}_{a, s, f}-h_{s, f}$ is the unique element of the domain of $\widetilde{S}_{a, f}$ satisfying

$$
\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right) u=-\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

Since $s \rightarrow h_{s, f}$ is entire, the meromorphy of the resolvent $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1}[10.9]$ yields the meromorphy of $\widetilde{E}_{a, s}-h_{s, f}$ as $\mathfrak{E}_{a}^{1}$-valued function, assuming that we have the following lemma, parallel to [11.4.1]:
[11.9.1] Lemma: For $a \gg 1$, $\mathfrak{E}_{a}^{1}=\mathfrak{E}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$. That is, for $a \gg 1, \mathfrak{E}(P, f)_{a}$ is $\mathfrak{E}^{1}$-dense in $\mathfrak{E}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$.
Proof: The proof is also parallel to that of [11.4.1], with a few minor complications. Indeed, after suitable set-up observations, the necessary estimates reduce to essentially one-dimensional estimates as in the proof of [11.4.1].

Since $\mathfrak{E}_{a}^{1}$ is the $\mathfrak{E}^{1}$-closure of $\mathfrak{E}(P, f)_{a}$, the containment is $\mathfrak{E}_{a}^{1} \subset \mathfrak{E}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$ is immediate.
For the opposite containment, given a sequence $\left\{\Psi_{\varphi_{i}, f} \in \mathfrak{E}^{1}\right\}$ of pseudo-Eisenstein series converging to $\Psi \in \mathfrak{E}^{1} \cap L_{a}^{2}(\Gamma \backslash G / K)$ in the $\mathfrak{E}^{1}$-topology, we produce a sequence of pseudo-Eisenstein series in $\mathfrak{E}(P, f)_{a}$ converging to $\Psi$ in the $\mathfrak{E}^{1}$-topology, by smooth cut-offs of the constant terms of the $\Psi_{\varphi_{i}, f}$. Again, by [3.11], as noted in [10.6], all constant terms along parabolics other than $P$ vanish entirely, because of the cuspidal data $f$. As in [11.8], let

$$
\alpha(n m k)=\left|\frac{\operatorname{det} m_{1}}{\operatorname{det} m_{2}}\right|^{r} \quad\left(\text { with } n \in N^{P}, k \in K, m=\left(\begin{array}{cc}
m_{1} & 0 \\
0 & m_{2}
\end{array}\right)\right)
$$

The constant term along $P$ of the limit $\Psi$ of the $\Psi_{\varphi_{i}, f}$ vanishes above $\alpha(g)=a$ by definition, and is in $L_{a}^{2}(\Gamma \backslash G / K)$. We will show that this entails that the part of the constant terms of the $\Psi_{\varphi_{i}, f}$ above $\alpha(g)=a$ must become small. Thus, smoothly cutting off the part of the constant terms above $\alpha(g)=a$ has a vanishingly small effect on the $\Psi_{\varphi_{i}, f}$. More precisely, proceed as follows.

Let $F$ be a smooth real-valued function on $\mathbb{R}$ with $F(t)=0$ for $t<-1,0 \leq F(t) \leq 1$ for $-1 \leq t \leq 0$, and $F(t)=1$ for $t \geq 0$. For $\varepsilon>0$, let $F_{\varepsilon}(t)=F((t-a) / \varepsilon)$. Fix real $b$ with $a>b>1$. Given $\Psi_{\varphi_{i}, f} \rightarrow \Psi \in L_{a}^{2}(\Gamma \backslash G / K)$, the $b$-tail of the $P$-constant term of $\Psi_{\varphi_{i}, f}$ is $\tau_{i}(g)=c_{P} \Psi_{\varphi_{i}, f}(g)$ for $\alpha(g) \geq a^{\prime}$, and $\tau_{i}(g)=0$ for $0<\alpha(g) \leq a^{\prime \prime}$. By design, $\Psi_{\varphi_{i}, f}-\Psi_{F_{\varepsilon} \cdot \tau_{i}, f} \in \mathfrak{E}(P, f)_{a}$ for small $\varepsilon$. We will show that, as $i \rightarrow+\infty$, for $\varepsilon_{i}$ sufficiently small depending on $i$, the $\mathfrak{E}^{1}$-norms of $\Psi_{F_{\varepsilon_{i}} \cdot \tau_{i}, f}$ go to 0 , while still $\Psi_{\varphi_{i}, f}-\Psi_{F_{\varepsilon_{i}} \cdot \tau_{i}, f} \rightarrow f$ in the $\mathfrak{E}^{1}$-norm.

Let $\mathfrak{S}_{a}^{P}$ be Siegel sets adapted to $P$, as in [11.8], and put $C_{a}=(P \cap \Gamma) \backslash \mathfrak{S}_{a}^{P}$. As in [11.8], by reduction theory, for $a \gg 1, C_{a}$ injects to $\Gamma \backslash G / K$. For $\Phi \in C_{c}^{\infty}(\Gamma \backslash G / K)$, let

$$
|\Phi|_{\mathfrak{E}^{1}\left(C_{a}\right)}^{2}=\int_{C_{a}}|\Phi|^{2}-\Delta \Phi \cdot \bar{\Phi}
$$

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We have

$$
|\Phi|_{\mathfrak{E}^{1}\left(C_{a}\right)}^{2} \ll \int_{\Gamma \backslash G / K}|\Phi|^{2}-\Delta \Phi \cdot \bar{\Phi}
$$

For each $a>0$, let $\mathfrak{E}^{1}\left(C_{a}\right)$ be the completion of $C_{c}^{\infty}(\Gamma \backslash G / K)$ with respect to the semi-norm $|\cdot|_{\mathfrak{E}^{1}\left(C_{b}\right)}$ (with collapsing). In contrast to the four simpler examples, $\Gamma \cap M^{P}$ is not finite, so we need a slightly different argument for estimates on tails of constant terms than in [11.4.1].

By reduction theory, there is $0<b \ll 1$ such that $\mathfrak{S}_{b}^{P}$ surjects to $\Gamma \backslash G / K$. Then take $a>_{b} 1$ such that, for $\eta \in C^{\infty}(0,+\infty)$ with support in $[a,+\infty)$, on $\mathfrak{S}_{b}^{P}$

$$
\Psi_{\eta, f}(g)=\varphi_{\eta, f}(g) \quad\left(\text { for } g \in \mathfrak{S}_{b}^{P}\right)
$$

Since $\varphi_{\eta, f}$, for such $\eta$, in $\mathfrak{S}_{b}^{P}$ the $P$-constant term is $c_{P} \Psi_{\eta, f}=\varphi_{\eta, f}$. Thus, for $a \gg_{b} 1$,

$$
\int_{C_{a}}\left|c_{P} \Psi_{\eta, f}\right|^{2}=\int_{C_{b}}\left|\varphi_{\eta, f}\right|^{2}=\int_{C_{a}}\left|\varphi_{\eta, f}\right|^{2}=\left|\Psi_{\eta, f}\right|_{L^{2}(\Gamma \backslash G / K)}^{2}
$$

That is, for such $\eta$, the $P$-constant term has $L^{2}$ norm dominated by (in fact, equal to) that of $\Psi_{\eta, f}$. By [11.7.1],

$$
\Delta \Psi_{\eta, f}=\Psi_{T \eta, f} \quad\left(\text { with } T=2 r^{3} y^{2} \frac{\partial^{2}}{\partial y^{2}}+\mu_{1}+\mu_{2}\right)
$$

and similarly for $\varphi_{\eta, f}$. Thus, we have similar inequalities for the $\mathfrak{E}^{1}\left(C_{a}\right)$-norms for such $\eta$ :

$$
\left|c_{P} \Psi_{\eta, f}\right|_{\mathfrak{E}^{1}\left(C_{a}\right)}^{2}=\left|\varphi_{\eta, f}\right|_{\mathfrak{E}^{1}\left(C_{a}\right)}^{2}=\left|\Psi_{\eta, f}\right|_{\mathbb{E}^{1}\left(C_{a}\right)}^{2} \leq\left|\Psi_{\eta, f}\right|_{\mathfrak{E}^{1}}^{2}
$$

That is, for fixed cuspidal data, for $\eta$ supported in $[a,+\infty)$ with $a>_{b} 1$, the $\mathfrak{E}^{1}$ norm of the constant term of $\Psi_{\eta, f}$ is dominated by that of $\Psi_{\eta, f}$, as desired.

To an extent, we can dodge entirely rewriting the norms as energy norms, instead using the earlier computations [3.11.11] and [11.7.1] which give $\Delta \Psi_{\psi, f}=\Psi_{T \psi, f}$ with

$$
T=2 r^{3} y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)+\mu_{1}+\mu_{2}=2 r^{3} y^{2} \frac{\partial^{2}}{\partial y^{2}}+\mu_{1}+\mu_{2}
$$

to reduce to a one-dimensional computation more directly. As in the very simplest case of $G=S L_{2}(\mathbb{R})$, the coefficient $y^{2}$ on $\frac{\partial^{2}}{\partial y^{2}}$ exactly cancels a denominator in an invariant measure in that coordinate, as we see in the following.

Turning to the main argument: since $a \gg_{b} 1, \Psi_{F_{\varepsilon} \cdot \tau_{i}, f}$ is just $\varphi_{F_{\varepsilon} \cdot \tau_{i}, f}$ on $\mathfrak{S}_{b}$, and the support of $\varphi_{F_{\varepsilon} \cdot \tau_{i}, f}$ is inside the image of the cylinder $C_{a-\varepsilon} \subset C_{b}$, by the triangle inequality,

$$
\begin{gathered}
\left|\Psi_{F_{\varepsilon} \cdot \tau_{i}, f}\right|_{\mathfrak{E}^{1}} \leq\left|\Psi_{F_{\varepsilon} \cdot \tau_{i}, f}\right|_{\mathfrak{E}^{1}\left(C_{b}\right)}=\left|\varphi_{F_{\varepsilon} \cdot \tau_{i}, f}\right|_{\mathfrak{E}^{1}\left(C_{a-\varepsilon}\right)} \\
\leq\left|\varphi_{\left(F_{\varepsilon}-1\right) \cdot \tau_{i}, f}\right|_{\mathfrak{E}^{1}\left(C_{a-\varepsilon}\right)}+\left|\varphi_{\tau_{i}, f}-c_{P} \Psi\right|_{\mathfrak{E}^{1}\left(C_{a-\varepsilon}\right)}+\left|c_{P} \Psi\right|_{\mathfrak{E}^{1}\left(C_{a-\varepsilon}\right)}
\end{gathered}
$$

From above, by design,

$$
\left|\tau_{i}-c_{P} \Psi\right|_{\mathfrak{E}^{1}\left(C_{a-\varepsilon}\right)} \ll\left|c_{P} \Psi_{\varphi_{i}, f}-c_{P} \Psi\right|_{\mathfrak{E}^{1}} \ll\left|\Psi_{\varphi_{i}, f}-\Psi\right|_{\mathfrak{E}^{1}} \longrightarrow 0
$$

so the middle summand goes to 0 . The first and third summands require somewhat more care.
Up to measure normalization constants, integrating away the cuspidal data,

$$
\left|\varphi_{\left(F_{\varepsilon}-1\right) \cdot \tau_{i}, f}\right|_{\mathfrak{E}^{1}\left(C_{a-\varepsilon}\right)}^{2}=\int_{a-\varepsilon}^{\infty}\left|\left(F_{\varepsilon}-1\right) \tau_{i}\right|^{2}-T\left(F_{\varepsilon}-1\right) \tau_{i} \cdot \overline{\left(F_{\varepsilon}-1\right) \tau_{i}} \frac{d y}{y^{2}}
$$

Since $T$ is of the form $A y^{2} \frac{\partial^{2}}{\partial y^{2}}+B$ for real constants $A>0$ and $B \leq 0$, it is a symmetric operator with respect to the measure $d y / y^{2}$, and the previous expression is dominated by

$$
\begin{gathered}
\int_{a-\varepsilon}^{\infty}\left|\left(F_{\varepsilon}-1\right) \tau_{i}\right|^{2} \frac{d y}{y^{2}}+\int_{a-\varepsilon}^{\infty} \frac{\partial}{\partial y}\left(F_{\varepsilon}-1\right) \tau_{i} \cdot \overline{\frac{\partial}{\partial y}\left(F_{\varepsilon}-1\right) \tau_{i}} d y \\
\leq \int_{a-\varepsilon}^{\infty}\left|F_{\varepsilon}-1\right|^{2} \cdot\left(\left|\tau_{i}\right|^{2}+\left|y \tau_{i}^{\prime}\right|^{2}\right) \frac{d y}{y^{2}}+\int_{C_{a-\varepsilon}}\left|F_{\varepsilon}^{\prime}\right|^{2} \cdot\left|\tau_{i}\right|^{2} d y+\int_{C_{a-\varepsilon}} 2\left|F_{\varepsilon}\right| \cdot\left|F_{\varepsilon}^{\prime}\right| \cdot\left|\tau_{i}\right| \cdot\left|\tau_{i}^{\prime}\right| d y
\end{gathered}
$$

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by Leibniz' rule for derivatives. The first summand in the latter expression goes to 0 as $\varepsilon \rightarrow 0^{+}$because $F_{\varepsilon}-1=0$ when $y \geq a$, and $\tau_{i}$ and $\tau_{i}^{\prime}$ are continuous. Thus,

$$
\left|F_{\varepsilon}^{\prime}(y)\right|=\left|\frac{1}{\varepsilon} \cdot F^{\prime}((y-a) / \varepsilon)\right|=\frac{1}{\varepsilon} \cdot\left|F^{\prime}((y-a) / \varepsilon)\right|<_{F} \frac{1}{\varepsilon}
$$

The fundamental theorem of calculus and Cauchy-Schwarz-Bunyakowsky recover an easy instance of a Sobolev inequality:

$$
\left|\tau_{i}(a-v)\right|=\left|0-\int_{0}^{v} \tau_{i}^{\prime}(a-v) d v\right| \leq\left(\int_{0}^{v}\left|\tau_{i}^{\prime}(a-v)\right|^{2} d v\right)^{\frac{1}{2}} \cdot\left(\int_{0}^{v} 1^{2} d v\right)^{\frac{1}{2}} \leq o(1) \cdot \sqrt{v}
$$

with Landau's little-o notation, since $\tau_{i}^{\prime}$ is locally $L^{2}$. Thus,

$$
\begin{gathered}
\int_{C_{a-\varepsilon}}\left|F_{\varepsilon}\right| \cdot\left|F_{\varepsilon}^{\prime}\right| \cdot\left|\tau_{i}\right| \cdot\left|\tau_{i}^{\prime}\right| \leq \frac{1}{\varepsilon} \cdot o(1) \cdot \sqrt{\varepsilon} \cdot \int_{0}^{\varepsilon}\left|\tau_{i}^{\prime}\right| \\
\leq \frac{1}{\varepsilon} \cdot o(1) \cdot \sqrt{\varepsilon} \cdot\left(\int_{0}^{\varepsilon}\left|\tau_{i}^{\prime}\right|^{2}\right)^{\frac{1}{2}} \cdot\left(\int_{0}^{\varepsilon} 1^{2}\right)^{\frac{1}{2}} \ll \tau_{i} \frac{1}{\varepsilon} \cdot o(1) \cdot \sqrt{\varepsilon} \cdot \sqrt{\varepsilon}=o(1)
\end{gathered}
$$

That is, the summand $\int_{C_{a-\varepsilon}}\left|F_{\varepsilon}\right| \cdot\left|F_{\varepsilon}^{\prime}\right|_{\mathfrak{s}} \cdot\left|\tau_{i}\right| \cdot\left|\tau_{i}^{\prime}\right|_{\mathfrak{s}}$ goes to 0 as $\varepsilon \rightarrow 0^{+}$. By the same estimates,

$$
\int_{C_{a-\varepsilon}}\left|F_{\varepsilon}^{\prime}\right|_{\mathfrak{s}}^{2} \cdot\left|\tau_{i}\right|^{2} \ll \frac{1}{\varepsilon^{2}} \int_{0}^{\varepsilon}(o(1) \cdot \sqrt{v})^{2} d v=\frac{1}{\varepsilon^{2}} \cdot o(1) \cdot \frac{\varepsilon^{2}}{2} \longrightarrow 0
$$

Thus, taking the $\varepsilon_{i}$ sufficiently small, the smooth truncations $\Psi_{\varphi_{i}, f}-\Psi_{F_{\varepsilon_{i}} \cdot \tau_{i}, f}$ are in $D \cap L_{a}^{2}(\Gamma \backslash G / K)$, and still converge to $\Psi$ in the $\mathfrak{E}^{1}$-topology.

### 11.10 Beyond the critical line: $P^{r, r} \subset S L_{2 r}$

Returning to the meromorphic continuation, we continue to consider symmetrical cuspidal data, of the symmetrical form $f=f_{1} \otimes f_{1}$. The discussion continues to resemble that for the four simple cases [11.5].

Since $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1} \operatorname{maps}\left(\Delta-\lambda_{s, f}\right) h_{s, f}$ to a function with $P$-constant term vanishing above height $a$, above that height the constant term of $\widetilde{E}_{a, s, f}$ is that of $h_{s, f}$. More generally, evaluate $\widetilde{S}_{a, f}-\lambda_{s, f}$ distributionally by application of $\Delta-\lambda_{s, f}$ : for some constant $C_{s}$, by [11.8.5],
$-\left(\Delta-\lambda_{s, f}\right) h_{s, f}=\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}-h_{s, f}\right)=\left(\Delta-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}-h_{s, f}\right)+C_{s} \cdot \eta_{a} \quad$ (as distributions)
Everything else in the latter equation is meromorphic in $s$, so $C_{s}$ must be, as well. Rearranging,

$$
\left(\Delta-\lambda_{s, f}\right) \widetilde{E}_{a, s, f}=-C_{s} \cdot \eta_{a} \quad \text { (as distributions) }
$$

Since $\Delta$ is $G$-invariant, it commutes with the constant-term map, so the distribution $\left(\Delta-\lambda_{s, f}\right) c_{P} \widetilde{E}_{a, s, f}$ is 0 away from height $a$.

That constant term is of the form $\varphi_{\psi, f}$ for some function $\psi$ on $(0,+\infty)$. By [11.7.1], in coordinates $m^{\prime} \cdot z_{y}$, the distributional differential equation $\left(\Delta-\lambda_{s, f}\right) \varphi_{\psi, f}=0$ has solutions with $\psi$ exactly of the form $A_{s} y^{s} f\left(m^{\prime}\right)+B_{s} y^{1-s} f\left(m^{\prime}\right)$ with constants $A_{s}, B_{s}$, so $c_{P} \widetilde{E}_{a, s}$ must be of this form in $0<y<a$. Since $\widetilde{E}_{a, s, f}$ is meromorphic in $s$, so are $A_{s}, B_{s}$. Thus,

$$
c_{P} \widetilde{E}_{a, s, f}\left(m^{\prime} z_{y}\right)= \begin{cases}y^{s} \cdot f\left(m^{\prime}\right) & (\text { for } y>a) \\ A_{s} y^{s} f\left(m^{\prime}\right)+B_{s} y^{1-s} f\left(m^{\prime}\right) & (\text { for } 0<y<a)\end{cases}
$$

By construction, $h_{s, f}$ is smooth, and $\left(\widetilde{S}_{s, f}-\lambda_{s, f}\right)^{-1} \Psi \in \mathfrak{E}_{a}^{1}$ for all $\Psi \in \mathfrak{E}_{a}^{0}$.

## 11. Meromorphic continuation of Eisenstein series

Next, we claim that the constant term of $\widetilde{E}_{a, s, f}$ is continuous, in particular along the set of values $n m^{\prime} z_{a}$ with $m^{\prime} \in M^{1}, n \in N^{P}$. To this end, for $F \in C_{c}^{\infty}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)$, define a norm via

$$
|F|_{\mathfrak{B}^{1}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)}^{2}=\int_{\left(N^{P} \cap \Gamma\right) \backslash G / K}(1-\Delta) F \cdot \bar{F}
$$

and form the corresponding Hilbert space of functions by completion:

$$
\mathfrak{B}^{1}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)=\mathfrak{B}^{1} \text {-completion of } C_{c}^{\infty}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)
$$

Since the topological space $\left(N^{P} \cap \Gamma\right) \backslash G / K$ is a bit too large, we define a localized version of $\mathfrak{B}^{1}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)$ via seminorms using smooth cut-offs $\beta$ :

$$
\nu_{\beta}(F)=|\beta \cdot F|_{\mathfrak{B}^{1}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)} \quad\left(\text { for } \beta \in C_{c}^{\infty}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)\right)
$$

We can take larger-and-larger smooth cut-offs, so there is a countable co-final subset of these semi-norms so they give a locally convex invariant-metric topology $T$, with completion a Fréchet space:

$$
\mathfrak{B}_{\mathrm{loc}}^{1}=\mathfrak{B}_{\mathrm{loc}}^{1}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)=T \text {-completion of } C_{c}^{\infty}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)
$$

As usual, the compact group $\left(N^{P} \cap \Gamma\right) \backslash N$ acts continuously on $\mathfrak{B}_{\mathrm{loc}}^{1}$, for general reasons. The smoothly cut-off tail $h_{s, f}$ is an entire smooth-function-valued function, so is in $\mathfrak{B}_{\mathrm{loc}}^{1}$. Since $\widetilde{E}_{a, s, f}-h_{s, f}$ is in $\mathfrak{E}_{a}^{1} \subset \mathfrak{B}^{1}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right) \subset \mathfrak{B}_{\text {loc }}^{1}$, certainly $\widetilde{E}_{a, s, f} \in \mathfrak{B}_{\mathrm{loc}}^{1}$. Thus, the constant term $c_{P} \widetilde{E}_{a, s, f}$ exists as a Gelfand-Pettis integral, so is in $\mathfrak{B}_{\text {loc }}^{1}$.

Continuity is a local property, so to prove that the constant term is continuous, it suffices to show that smooth cut-offs are continuous. For sufficiently small support of the smooth cut-offs, we can reduce the local problem to functions on a multi-torus $\mathbb{T}^{n}$. The dimension $n=\operatorname{dim}_{\mathbb{R}}\left(N^{P} \cap \Gamma\right) \backslash G / K$ is too high for local Sobolev imbedding theorems to promise that every element of $\mathfrak{B}_{\mathrm{loc}}^{1}$ is continuous. Indeed, $\mathfrak{B}_{\mathrm{loc}}^{k} \subset C^{o}\left(\left(N^{P} \cap \Gamma\right) \backslash G / K\right)$ for $k>\frac{n}{2}$ [9.5.14]. Fortunately, with strong-sense cuspidal data $f$, the constant term $c_{P} \widetilde{E}_{a, s, f}$ is in a better situation, namely, it is smooth (or even constant) in all but one coordinate. Indeed, it allows separation of variables, being of the form $\varphi_{\psi, f}\left(n m^{\prime} z_{y}\right)=\psi(y) \cdot f\left(m^{\prime}\right)$. Thus, we smoothly truncate in a fashion that preserves the separation of variables, giving a function

$$
F\left(x_{1}, \ldots, x_{n}\right)=F_{1}\left(x_{1}\right) \cdot F_{2}\left(x_{2}, \ldots, x_{n}\right)
$$

on $\mathbb{T}^{n}$, with $F_{2}$ known to be $C^{\infty}$ in $x^{\prime}=x_{2}, \ldots, x_{n}$. We claim that $F \in \mathfrak{B}^{1}\left(\mathbb{T}^{n}\right)$ implies $F_{1} \in \mathfrak{B}^{1}(\mathbb{T})$. Since $F$ is a product, its Fourier coefficients likewise are products: letting $\xi^{\prime}=\left(\xi_{2}, \ldots, \xi_{n}\right)$,

$$
\begin{gathered}
\widehat{F}\left(\xi_{1}, \xi_{2}, \ldots, \xi_{n}\right)=\int_{\mathbb{T}^{n}} F_{1}\left(x_{1}\right) F_{2}\left(x^{\prime}\right) \cdot e^{-2 \pi i\left(x_{1} \xi_{1}+x^{\prime} \cdot \xi^{\prime}\right)} d x_{1} d x^{\prime} \\
=\int_{\mathbb{T}} F_{1}\left(x_{1}\right) e^{-2 \pi i x_{1} \xi_{1}} d x_{1} \cdot \int_{\mathbb{T}^{n-1}} F_{2}\left(x^{\prime}\right) \cdot e^{-2 \pi i x^{\prime} \cdot \xi^{\prime}} d x^{\prime}=\widehat{F}_{1}\left(\xi_{1}\right) \cdot \widehat{F}_{2}\left(\xi^{\prime}\right)
\end{gathered}
$$

The smoothness condition on $F_{2}$ gives rapid decrease of $\widehat{F}_{2}$, that is, for all $N \gg 1$

$$
\left|\widehat{F}_{2}\left(\xi^{\prime}\right)\right| \ll_{N}\left|\xi^{\prime}\right|^{-N}
$$

Then the $\mathfrak{B}^{1}\left(\mathbb{T}^{n}\right)$ condition gives

$$
\begin{aligned}
& \sum_{\xi_{1}}\left|\widehat{F}_{1}\left(\xi_{1}\right)\right|^{2} \cdot\left(1+\left|\xi_{1}\right|^{2}\right) \leq \sum_{\xi \in \mathbb{Z}^{n}}|\widehat{F}(\xi)|^{2} \cdot\left(1+|\xi|^{2}\right)=\sum_{\xi_{1}, \xi^{\prime}}\left|\widehat{F}_{1}\left(\xi_{1}\right)\right|^{2} \cdot\left|\widehat{F}_{2}\left(\xi^{\prime}\right)\right|^{2} \cdot\left(1+|\xi|^{2}\right) \\
& \quad=\sum_{\xi_{1}}\left|\widehat{F}_{1}\left(\xi_{1}\right)\right|^{2} \cdot\left(\sum_{\xi^{\prime}}\left|\widehat{F}_{2}\left(\xi^{\prime}\right)\right|^{2}\right) \cdot\left(1+|\xi|^{2}\right) \ll \sum_{\xi_{1}}\left|\widehat{F}_{1}\left(\xi_{1}\right)\right|^{2} \cdot\left(\sum_{\xi^{\prime}}\left|\xi^{\prime}\right|^{-N}\right) \cdot\left(1+|\xi|^{2}\right)
\end{aligned}
$$

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Taking $N>n-1$ gives convergence of the inner sum, so $F_{1}$ is in the +1 Sobolev space on $\mathbb{T}$. Then Sobolev imbedding [9.5.4] implies that $F_{1}$ is continuous. Thus, $c_{P} \widetilde{E}_{a, s, f}$ is continuous at $\eta=a$, as claimed.

Thus, the values above and below $y=a$ must match:

$$
A_{s} \cdot a^{s} \cdot f\left(m^{\prime}\right)+B_{s} \cdot a^{1-s} \cdot f\left(m^{\prime}\right)=a^{s} \cdot f\left(m^{\prime}\right) \quad(\text { for all } s)
$$

and since $f\left(m^{\prime}\right)$ is not identically 0 ,

$$
A_{s} \cdot a^{s}+B_{s} \cdot a^{1-s}=a^{s} \quad(\text { for all } s)
$$

As in the proof of [11.5] (with somewhat different notation), let $\mathrm{ch}_{[a, \infty)}$ be the characteristic function of $[a, \infty)$, and

$$
\beta_{a, s, f}\left(n m^{\prime} z_{y}\right)=\operatorname{ch}_{[a, \infty)}(y) \cdot\left(A_{s} y^{s}+B_{s} y^{1-s}-y^{s}\right) \cdot f\left(m^{\prime}\right)
$$

and form a pseudo-Eisenstein series

$$
\Phi_{a, s, f}=\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma} \beta_{a, s, f} \circ \gamma
$$

The support of $\beta_{a, s, f}$ is inside $\mathfrak{S}_{a}^{P}$, and $a \gg 1$, so by reduction theory for each $g \in G$ the series has at most one non-zero summand, so converges for all $s \in \mathbb{C}$.
[11.10.1] Theorem: $A_{s} \cdot E_{s, f}=\widetilde{E}_{a, s, f}+\Phi_{a, s, f}$ and $\widetilde{E}_{a, s, f}+\Phi_{a, s, f}=B_{s} \cdot E_{1-s, f}$. Thus, $E_{s, f}$ has a meromorphic continuation and $E_{s, f}-h_{s, f}$ is a meromorphic $\mathfrak{E}^{1}$-valued function.
Proof: With $\widetilde{S}$ as in [11.7], we have shown that $u=E_{s, f}-h_{s, f}$ is the unique solution $u \in \mathfrak{E}^{1}$ to

$$
\left(\widetilde{S}-\lambda_{s, f}\right) u=-\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

Thus, multiplying through by $A_{s}$, it suffices to prove that $\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} \cdot h_{s, f}$ is in $\mathfrak{E}^{1}$ and satisfies

$$
\left(\widetilde{S}-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} \cdot h_{s, f}\right)=-\left(\Delta-\lambda_{s, f}\right)\left(A_{s} \cdot h_{s, f}\right)
$$

That $\widetilde{E}_{a, s, f}-h_{s, f}$ is in $\mathfrak{E}_{a}^{1}$ motivates the rearrangement

$$
\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} \cdot h_{s, f}=\left(\widetilde{E}_{a, s, f}-h_{s, f}\right)+\left(\Phi_{a, s, f}-A_{s} h_{s, f}+h_{s, f}\right)
$$

We claim that the pseudo-Eisenstein series $F=\Phi_{a, s, f}-A_{s} h_{s, f}+h_{s, f}$ is in $\mathfrak{E}^{1}$.
Regarding integrability, by reduction theory, in $\mathfrak{S}_{a}^{P}, \Phi_{a, s, f}$ is just

$$
\beta_{a, s, f}\left(n m^{\prime} z_{y}\right)=\operatorname{ch}_{[a, \infty)}(y) \cdot\left(A_{s} y^{s}+B_{s} y^{1-s}-y^{s}\right) \cdot f\left(m^{\prime}\right)
$$

so in $\mathfrak{S}_{a}^{P}$
$F=\Phi_{a, s, f}-A_{s} h_{s, f}+h_{s, f}=\left(\left(A_{s} y^{s}+B_{s} y^{1-s}-\eta^{s}\right)-A_{s} y^{s}+y^{s}\right) \cdot f\left(m^{\prime}\right)=B_{s} y^{1-s} \cdot f\left(m^{\prime}\right) \quad\left(\right.$ in $\left.\mathfrak{S}_{a}^{P}\right)$
For $\operatorname{Re}(s)>1, y^{1-s} \cdot f\left(m^{\prime}\right)$ is square-integrable on $(P \cap \Gamma) \backslash \mathfrak{S}_{a}^{P}$, so $F$ is in $L^{2}(\Gamma \backslash G / K)$.
To demonstrate the additional smoothness required for $F$ to be in $\mathfrak{E}^{1}$, from the rewriting of Sobolev norms in [10.7], especially [10.7.5], it suffices to show that the right-translation derivatives $x F$ are in $L^{2}(\Gamma \backslash G)$ for $x \in \mathfrak{g}$. By the left invariance of the right action of $\mathfrak{g}$, it suffices to prove square-integrability on the adapted standard Siegel sets $\mathfrak{S}_{t}^{P}$ of the derivatives of the data

$$
\beta_{a, s, f}-A_{s} \tau(y) y^{s} \cdot f\left(m^{\prime}\right)+\tau(y) y^{s} \cdot f\left(m^{\prime}\right)
$$

wound up to form $F=\Phi_{a, s, f}-A_{s} h_{s, f}+h_{s, f}$ This data is smooth everywhere but along $y=a$, where it is at least continuous, since $A_{s} a^{s}+B_{s} a^{s}-a^{s}=0$. Further, it possesses continuous left and right derivatives

## 11. Meromorphic continuation of Eisenstein series

in $y$ at $y=a$, and is smooth in all other directions on $y=a$, so is locally in a +1 -index Sobolev space near $y=a$.

Derivatives in the directions coming from the unipotent radical $\mathfrak{n}^{P}$ give 0 . Since the cuspidal data $f$ is strong-sense, $f$ is smooth, and all derivatives are still of uniform rapid decay in Siegel sets, by [7.3.19] and [7.3.15]. Thus, for such derivatives, the $L^{2}$ estimate above is sufficient.

In coordinates $n m^{\prime} z_{y}$ as above, the remaining differential operator is $y \frac{\partial}{\partial y}$. The derivative of $F$ is discontinuous at $y=a$, although it has left and right limits. As a distribution, it is

$$
\begin{array}{r}
y \frac{\partial}{\partial y} F=y \frac{\partial}{\partial y}\left(\Phi_{s, f}-A_{s} h_{s, f}+h_{s, f}\right)=y \frac{\partial}{\partial y}\left(\varphi_{s, f}-A_{s} \cdot \tau(y) \cdot y^{s} \cdot f\left(m^{\prime}\right)+\tau(y) \cdot y^{s} \cdot f\left(m^{\prime}\right)\right) \\
=y \frac{\partial}{\partial y}\left\{\begin{array}{cl}
B_{s} y^{1-s} \cdot f\left(m^{\prime}\right) & (\text { for } y>a) \\
-A_{s} \cdot y^{s} \cdot f\left(m^{\prime}\right)+y^{s} \cdot f\left(m^{\prime}\right) & \left(\text { for } a^{\prime} \leq y<a\right) \\
-A_{s} \cdot \tau(y) \cdot y^{s} \cdot f\left(m^{\prime}\right)+\tau(y) \cdot y^{s} \cdot f\left(m^{\prime}\right) & \left(\text { for } a^{\prime \prime} \leq y \leq a^{\prime}\right) \\
0 & \left(\text { for } y \leq a^{\prime \prime}\right)
\end{array}\right. \\
=\left\{\begin{array}{cl}
B_{s} \cdot(1-s) \cdot y^{1-s} \cdot f\left(m^{\prime}\right) & (\text { for } y>a) \\
\left(1-A_{s}\right) \cdot s y^{s} \cdot f\left(m^{\prime}\right) & \left(\text { for } a^{\prime} \leq y<a\right) \\
0 & \left(\text { for } a^{\prime \prime} \leq y \leq a^{\prime \prime}\right)
\end{array}\right. \\
\left(1-A_{s}\right)\left(\frac{\partial \tau}{\partial y} \cdot y^{s}+\tau(y) \cdot s y^{s}\right) \cdot f\left(m^{\prime}\right) \\
0
\end{array}
$$

On $a^{\prime} \leq y \leq a$, this derivative is bounded, so the truly relevant behavior is in $y>a$ : for $\operatorname{Re}(s)>1$ this derivative is square-integrable on quotients $(P \cap \Gamma) \backslash \mathfrak{S}_{b^{\prime}}^{P}$. Thus, $\Phi_{s, f}-A_{s} h_{s, f}+h_{s, f}$ is in $\mathfrak{E}^{1}$, proving that $\widetilde{E}_{a, s, f}+\Phi_{s, f}-A_{s} h_{s, f}$ is in $\mathfrak{E}^{1}$.

To show that $\widetilde{E}_{a, s, f}+\Phi_{s, f}-A_{s} h_{s, f}$ satisfies the expected equation, we justify computing the effect of differential operators on $\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} h_{s, f}$ distributionally, as follows. For $F \in C_{c}^{\infty}(\Gamma \backslash G / K)$, with $\widetilde{S}$ the Friedrichs extension of the restriction of $\Delta$ to $C_{c}^{\infty}(\Gamma \backslash G / K)$ as in [11.7],

$$
\begin{gathered}
\left\langle\left(\widetilde{S}-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} h_{s, f}\right), F\right\rangle=\left\langle\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} h_{s, f},\left(\Delta-\bar{\lambda}_{s, f}\right) F\right\rangle \\
=\left\langle\left(\Delta-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} h_{s, f}\right), F\right\rangle
\end{gathered}
$$

By design, using the invariance of $\Delta$ and the local finiteness of the sum for $\Phi_{s, f}$, it is legitimate to compute

$$
\begin{gathered}
\left(\Delta-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}\right)=\left(\Delta-\lambda_{s, f}\right) \widetilde{E}_{a, s, f}+\sum_{\gamma \in \Gamma_{\infty} \backslash \Gamma}\left(\Delta-\lambda_{s, f}\right) \beta_{a, s, f} \circ \gamma \\
=-C_{s} \cdot \eta_{a}+C_{s} \cdot \eta_{a}=0 \quad \text { (as distributions) }
\end{gathered}
$$

Thus,

$$
\left(\widetilde{S}-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} h_{s, f}\right)=\left(\Delta-\lambda_{s, f}\right)\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}-A_{s} h_{s, f}\right)=0-A_{s}\left(\Delta-\lambda_{s, f}\right) h_{s, f}
$$

as desired, proving $\widetilde{E}_{a, s, f}+\Phi_{a, s, f}=A_{s} \cdot E_{s, f}$ for $\operatorname{Re}(s)>1$. For $\operatorname{Re}(1-s)>1$, the same argument shows that $\widetilde{E}_{a, s, f}+\Phi_{a, s, f}=B_{s} \cdot E_{1-s, f}$. This proves the formulas in the claim. Since not both $A_{s}$ and $B_{s}$ can be identically 0 , we obtain the meromorphic continuation of $E_{s, f}$.
[11.10.2] Corollary: $A_{s} \cdot E_{s, f}=B_{s} \cdot E_{1-s, f}$.

In particular, neither $A_{s}$ nor $B_{s}$ is identically 0 , and with $a(s)=B_{s} / A_{s}, E_{1-s, f}=a(s) \cdot E_{s, f}$. The relation $c_{P} E_{s, f}=\left(y^{s}+c_{s, f} y^{1-s} \cdot f\left(m^{\prime}\right)\right.$ gives the meromorphic continuation of $c_{s, f}$. Since $c_{P} E_{1-s, f}=$ $\left(y^{1-s}+c_{1-s, f} y^{s}\right) \cdot f\left(m^{\prime}\right)$, apparently $c_{s, f}=a(s)=B_{s} / A_{s}$. Since $1-(1-s)=s$, we obtain $c_{s, f} \cdot c_{1-s, f}=1$ :
[11.10.3] Corollary: $c_{s, f}$ has a meromorphic continuation, and $c_{s, f} \cdot c_{1-s, f}=1$.
On $\operatorname{Im}(s)=0$ and $\operatorname{Re}(s)>1, E_{s, f}$ and $c_{P} E_{s, f}$ are real-valued. We assume without loss of generality that $f$ is real-valued. Thus, the two holomorphic functions $E_{s, f}$ and $\overline{E_{\bar{s}, f}}$ agree on $(1,+\infty)$, so agree everywhere. That is, $\overline{E_{s, f}}=E_{\bar{s}, f}$. In particular, on $\operatorname{Re}(s)=\frac{1}{2}$, where $\bar{s}=1-s$,

$$
\left|c_{s, f}\right|^{2}=c_{s, f} \cdot \overline{c_{s, f}}=c_{s, f} \cdot c_{\bar{s}, f}=c_{s, f} \cdot c_{1-s, f}=1 \quad\left(\text { on } \operatorname{Re}(s)=\frac{1}{2}\right)
$$

proving
[11.10.4] Corollary: $\left|c_{s, f}\right|=1$ on $\operatorname{Re}(s)=\frac{1}{2}$, and $c_{s, f}$ has no pole on $\operatorname{Re}(s)=\frac{1}{2}$.
Further, we have
[11.10.5] Corollary: $E_{s, f}$ has no pole on $\operatorname{Re}(s)=\frac{1}{2}$.
Proof: Suppose $E_{s, f}$ had a pole of order $N>0$ at $s_{o}$ on the critical line $\operatorname{Re}(s)=\frac{1}{2}$. Then $\left(s-s_{o}\right)^{N} \cdot E_{s, f}$ is holomorphic at $s=s_{o}$, gives a not identically automorphic form, and has vanishing constant term there. From

$$
\wedge^{a}\left(s-s_{o}\right)^{N} E_{s, f}=\left(s-s_{o}\right)^{N} \wedge^{a} E_{s, f}
$$

and using the Maass-Selberg relations [3.14.2] with $s=s_{o}+\varepsilon$ and $r=\bar{s}_{o}+\varepsilon=1-s_{o}+\varepsilon$ with $0<\varepsilon \in \mathbb{R}$, since $\left(s-s_{o}\right) \cdot c_{s} \rightarrow 0$ at $s=s_{o}$, suppressing measure-normalizations,

$$
\begin{gathered}
\left|\left(s-s_{o}\right)^{N} E_{s}\right|^{2}=\varepsilon^{2 N} \cdot\left(\frac{a^{s+r-1}}{s+r-1}+c_{s, f} \frac{a^{(1-s)+r-1}}{(1-s)+r-1}+c_{r, f} \frac{a^{s+(1-r)-1}}{s+(1-r)-1}+c_{s, f} c_{r, f} \frac{a^{(1-s)+(1-r)-1}}{(1-s)+(1-r)-1}\right) \\
\quad=\varepsilon^{2 N} \cdot\left(\frac{a^{2 \varepsilon}}{2 \varepsilon}+c_{s_{o}+\varepsilon, f} \frac{a^{1-2 s_{o}-2 \varepsilon}}{1-2 s_{o}-2 \varepsilon}+c_{1-s_{o}+\varepsilon, f} \frac{a^{2 s_{o}-1+2 \varepsilon}}{2 s_{o}-1+2 \varepsilon}+c_{s_{o}+\varepsilon, f} c_{1-s_{o}+\varepsilon, f} \frac{a^{-2 \varepsilon}}{-2 \varepsilon}\right) \longrightarrow 0
\end{gathered}
$$

contradiction. Thus, $E_{s, f}$ has no pole on the critical line.
Toward proving moderate growth of the meromorphic continuation of $E_{s, f}$ :
[11.10.6] Claim: $E_{s, f}$ meromorphically continues as a $C^{\infty}(\Gamma \backslash G / K)$-valued function.
Proof: As we have assumed throughout, to know the form of the constant term of $E_{s, f}$ with $f=f_{1} \otimes f_{1}$, as in [3.11.9] we need to assume that $f_{1}$ is a $\Delta$-eigenfunction on $S L_{r}(\mathbb{Z}) \backslash S L_{r}(\mathbb{R}) / S O(n, \mathbb{R})$, with eigenvalue $\lambda_{1}$, and up to scalar multiples is the only cuspform there with $\Delta$-eigenvalue $\lambda_{1}$. From the computation in [3.11.11], it follows that $E_{s, f}$ is a $\Delta$-eigenfunction, with eigenvalue $\lambda_{s, f}=2 r^{3} s(s-1)+2 \lambda_{1}$. Thus, in the region $\operatorname{Re}(s)>1$, there is at most a single $s$ making $\lambda_{s, f}$ assume a given value. As above, let

$$
\varphi_{s, f}\left(n m^{\prime} z_{y} k\right)=y^{s} \cdot f\left(m^{\prime}\right) \quad\left(\text { for } n \in N^{P}, k \in K, \text { and } m \in M^{\prime}\right)
$$

The computation in the proof of [3.11.11] also shows that $\varphi_{s, f}$ is a $\Delta$-eigenfunction with eigenvalue $\lambda_{s, f}$. Let $\eta \in C_{c}^{\infty}(K \backslash G / K)$ act on spaces of right $K$-invariant functions on $G$ as usual, by integral operators. From [8.4.1], for every $\eta \in C_{c}^{\infty}(K \backslash G / K)$, there is $\mu_{s, f}(\eta) \in \mathbb{C}$ such that $\eta \cdot \varphi_{s, f}=\mu_{s, f}(\eta) \cdot \varphi_{s, f}$, and there exists $\eta$ giving $\mu_{s, f}(\eta) \neq 0$. In the region of convergence $\operatorname{Re}(s)>1$, from $E_{s, f}=\sum_{\gamma \in(\Gamma \cap P) \backslash \Gamma} \varphi_{s, f} \circ \gamma$, also $\eta \cdot E_{s, f}=\mu_{s, f}(\eta) \cdot E_{s, f}$. Exactly what we are missing at this point is knowledge of what topological vector space of functions (or distributions) the meromorphically continued Eisenstein series may lie in, so we cannot directly assert much about $\eta \cdot E_{s, f}$ outside the region of convergence. (Otherwise we could apply the identity principle from complex analysis to the latter identity.) Rather, we approach this a little indirectly, as follows.

Since $\Delta$ commutes with $G, \mathfrak{E}^{1}$ is stable under the action of $\eta \in C_{c}^{\infty}(K \backslash G / K)$. From the meromorphic continuation of $E_{s, f}-h_{s, f}$ as $\mathfrak{E}^{1}$-valued function, we have the meromorphic continuation of

$$
\eta \cdot\left(E_{s, f}-h_{s, f}\right)=\mu_{s, f}(\eta) \cdot E_{s, f}-\eta \cdot h_{s, f}
$$

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as $\mathfrak{E}^{1}$-valued function. In fact, for $F \in \mathfrak{E}^{1}$, by [14.5], $\eta \cdot F$ is in $C^{\infty}(\Gamma \backslash G / K)$. By construction, $h_{s, f} \in C^{\infty}(\Gamma \backslash G / K)$. Rearranging,

$$
\mu_{s, f}(\eta) \cdot E_{s, f}=\eta \cdot\left(E_{s, f}-h_{s, f}\right)+\eta \cdot h_{s, f}
$$

Dividing through by $\mu_{s, f}(\eta)$ for some $\eta$ with $\mu_{s, f}(\eta) \neq 0$ exhibits the meromorphically continued $E_{s, f}$ as a smooth-function-valued function.
[11.10.7] Corollary: $E_{s}$ has a meromorphic continuation as $C^{o}(\Gamma \backslash G / K)$-valued function, so it makes sense to address the issue of its moderate growth.

Finally, we have
[11.10.8] Theorem: Away from poles, the meromorphically continued $E_{s, f}$ is of moderate growth.
Proof: By [11.10.1] and [11.10.7], (at least) the pointwise values of the meromorphic continuation are given by

$$
E_{s, f}=A_{s}^{-1} \cdot\left(\widetilde{E}_{a, s, f}+\Phi_{a, s, f}\right)
$$

Since $a \gg 1$, in $\mathfrak{S}_{a}^{P}$ the function $\Phi_{a, s, f}$ is just $\varphi_{a, s, f}$ itself, which is $\left(A_{s} y^{s}+B_{s} y^{1-s}-y^{s}\right) \cdot f\left(m^{\prime}\right)$, which is of moderate growth in standard Siegel sets. The computation above shows continuity at $y=a$. The pseudo-Eisenstein series $h_{s, f}$ of [11.7] made from $\tau \cdot y^{s} \cdot f\left(m^{\prime}\right)$ with smooth cut-off $\tau$ is a locally finite sum, so is smooth, so certainly continuous. For $\eta \geq a$, its value is just $\eta^{s}$, which is of moderate growth for all $s$. Thus, to show that $\widetilde{E}_{a, s, f}$ is of moderate growth even after meromorphic continuation, it suffices to show that $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1}\left(\Delta-\lambda_{s, f}\right) h_{s, f}$ is of moderate growth.

Let $T$ be the operator determined in [11.7.1], such that $\Delta \Psi_{\psi, f}=\Psi_{T \psi, f}$, namely, with $\mu_{1}$ the eigenvalue of the cuspform $f_{1}$,

$$
T=2 r^{3} y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)+2 \mu_{1}
$$

Thus,

$$
T-\lambda_{s, f}=2 r^{3} y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)+2 \mu_{1}-2 r^{3} s(s-1)-2 \mu_{1}=2 r^{3} \cdot\left(y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)-s(s-1)\right)
$$

The constant $2 r^{3}$ can be dropped without changing anything.
Again, the pseudo-Eisenstein series $h_{s, f}=\Psi_{\tau \cdot y^{s}, f}$ is a locally finite sum, so it is legitimate to compute

$$
\left(\Delta-\lambda_{s, f}\right) h_{s, f}=\left(\Delta-\lambda_{s, f}\right) \Psi_{\tau \cdot y^{s}, f}=\Psi_{\left(T-\lambda_{s, f}\right)\left(\tau \cdot y^{s}\right), f}
$$

and $H_{s}=\left(T-\lambda_{s, f}\right)\left(\tau \cdot y^{s}\right)$ is smooth and uniformly compactly supported. It suffices to demonstrate solvability of the differential equation $\left(T-\lambda_{s, f}\right) u=H_{s}$ for a function $u$ of sufficient decay at both $0^{+}$ and $+\infty$. Then hope to form a pseudo-Eisenstein series $\Psi_{u, f}$ giving $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right) \Psi_{u, f}=\Psi_{H_{s}, f}$. From the distributional characterization [11.8] of $\widetilde{S}_{a, f}$, this equation is equivalent to

$$
\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right) \Psi_{u, f}=\Psi_{H_{s}, f}+c \cdot \eta_{a} \quad(\text { for some } c \in \mathbb{C})
$$

Thus, in the $y$-coordinate, given $f \in C_{c}^{\infty}(0,+\infty)$, we solve equations $\left(T-\lambda_{s, f}\right) u=H_{s}+c \cdot \delta_{a}$ with $c \in \mathbb{C}$ for $u$ in $C^{\infty}(0,+\infty)$, with behavior at $0^{+}$and $+\infty$ to be adjusted suitably.

From [11.7.1], the differential equation is

$$
\left(y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)-s(s-1)\right) u=H_{s}+c \cdot \delta_{a}
$$

We can divide through by $A$ to suppose without loss of generality that it is 1 , and still the renormalized is $B<0$. Letting $x=\log y$, with $F(x)=H_{s}\left(e^{x}\right)$ and $v(x)=u\left(e^{x}\right)$, this becomes

$$
v^{\prime \prime}-v^{\prime}-s(s-1) v=F+c \cdot \delta_{\log a}
$$

Taking Fourier transform in a normalization that suppresses some factors of $2 \pi$,

$$
(-i \xi)^{2} \widehat{v}-(-i \xi) \widehat{v}-s(s-1) \widehat{v}=\widehat{F}+c \cdot a^{-i \xi}
$$

or

$$
\widehat{v}(\xi)=-\frac{\widehat{F}(\xi)+c \cdot a^{-i \xi}}{\xi^{2}-i \xi+s(s-1)}
$$

Since $F$ is a test function, $\widehat{F}$ is an entire function such that $\widehat{F}\left(x+i y_{o}\right)$ is (uniformly) in the Schwartz space for each fixed $y_{o}$. Division by a quadratic polynomial produces a function holomorphic in a strip along $\mathbb{R}$ not including either of the two poles at the zeros of the denominator:

$$
\frac{i \pm \sqrt{(-i)^{2}-4 s(s-1)}}{2}
$$

Fix $\varepsilon>0$. Given a bound $|\operatorname{Re}(s)| \leq B$, for $\operatorname{Im}(s)>_{B} 1$, those poles are outside the strip $S=\{z \in \mathbb{C}$ : $|\operatorname{Im}(z)| \leq 1+\varepsilon\}$. Thus, $\widehat{v}$ is holomorphic on an open set containing $S$ and has decay like $1 / \xi^{2}$ on horizontal lines inside that strip. Thus, in the Fourier inversion integral

$$
v(x)=\frac{1}{2 \pi} \int_{\mathbb{R}} e^{i \xi x} \widehat{v}(\xi) d \xi
$$

we can move the contour up to $\mathbb{R}+i(1+\varepsilon)$, giving

$$
v(x)=\frac{1}{2 \pi} \int_{\mathbb{R}} e^{i(\xi+i \varepsilon) x} \widehat{v}(\xi+i \varepsilon) d \xi=e^{-\varepsilon x} \frac{1}{2 \pi} \int_{\mathbb{R}} e^{i \xi x} \widehat{v}(\xi+i \varepsilon) d \xi
$$

Thus, $v(x) \ll e^{-(1+\varepsilon) x}$, giving genuine exponential decrease for $x \rightarrow+\infty$. Similarly, moving the contour down gives exponential decrease $v(x) \ll e^{-(1+\varepsilon)|x|}$ for $x \rightarrow-\infty$. Then $u(y)=v(\log y)$ satisfies $u(y) \ll y^{1+\varepsilon}$ as $y \rightarrow 0^{+}$, and $u(y) \ll y^{-(1+\varepsilon)}$ as $y \rightarrow+\infty$. Thus, the pseudo-Eisenstein series $\Psi_{u, f}$ converges absolutely, being dominated termwise by the sum expressing an absolutely convergent Eisenstein series [3.9], [3.11]. Further, being termwise dominated by an absolutely convergent Eisenstein series, $\Psi_{u, f}$ is continuous and of moderate growth.

Having available a choice of the constant $c$ is necessary, to adjust $\Psi_{u, f}$ to have $P$-constant term vanishing above height $y=a$. Choose the constant so that $c_{P} \Psi_{u}$ vanishes at $y=a$. Since $a \gg 1$, by reduction theory the truncation $\wedge^{a} \Psi_{u, f}$ has $P$-constant term vanishing at and above height $a$. Since $a \gg 1$, this truncation is itself a pseudo-Eisenstein series, and still $\left(\Delta-\lambda_{s, f}\right) \wedge^{a} \Psi_{u, f}$ differs from $\Psi_{u, f}$ only by a multiple of $\eta_{a}$. Again by the distributional characterization of $\widetilde{S}_{a, f}$, we have $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right) \wedge^{a} \Psi_{u, f}=\Psi_{H_{s}, f}$.

Thus, for a given bound $|\operatorname{Re}(s)| \leq B$, there is $C$ sufficiently large so that for $|\operatorname{Im}(s)| \geq C$ we have meromorphic continuation of $E_{s, f}$ as a (continuous) moderate-growth function.

For $|\operatorname{Im}(s)|<C$, we can express $E_{s, f}$ as a vector-valued Cauchy integral along a circular path $\gamma$ that lies inside the union $U$ of regions $\operatorname{Re}(s) \geq B, \operatorname{Re}(s) \leq 1-B$, and $|\operatorname{Im}(s)| \geq C$, and does not run through any poles of $E_{s, f}$. In $\operatorname{Re}(s) \leq 1-B$ the Eisenstein series is (continuous) of moderate growth, via the functional equation. Thus, $E_{s, f}$ is of moderate growth throughout $U$, and in particular along $\gamma$. Let $Z$ be the collection of poles of $E_{s, f}$ (as meromorphic $C^{o}(\Gamma \backslash G / K)$-valued function) inside $\gamma$, and $P(z)=\prod_{z_{j} \in Z}\left(z-z_{j}\right)$. For each $g \in G$

$$
P(s) \cdot E_{s, f}(g)=\frac{1}{2 \pi i} \int_{\gamma} \frac{P(z) \cdot E_{z}(g)}{z-s} d z
$$

In fact, on $\gamma, z \rightarrow\left(s \rightarrow P(z) E_{s, f} /(z-s)\right.$ is a compactly-supported, continuous, moderate-growth-functionvalued function of $z$, so the vector-valued Cauchy integral

$$
P(s) \cdot E_{s, f}=\frac{1}{2 \pi i} \int_{\gamma} \frac{P(z) \cdot E_{z, f}}{z-s} d z
$$

as in [15.2] exists as a Gelfand-Pettis integral [14.1] lying in that same space of functions.

### 11.11 Exotic eigenfunctions: $P^{r, r} \subset S L_{2 r}$

Since $\mathfrak{E}(P, f)$ contained no eigenfunctions for $\Delta$ except the finitely-many possible residues of $E_{s, f}^{P}$ in $\operatorname{Re}(s)>\frac{1}{2}$ (see [3.14]), the eigenfunctions for $\widetilde{S}_{a, f}$ cannot be eigenfunctions for $\Delta$, so must be exotic.

Continue to consider symmetrical cuspidal data $f=f_{1} \otimes f_{1}$, so that $f^{w}=f$. On genuine Eisenstein series $E_{s, f}$ the functional $\eta_{a}$ makes sense: unwinding, and using the explicit form of the constant terms, we have absolutely convergent integrals

$$
\begin{gathered}
\eta_{a}\left(E_{s, f}\right)=\int_{Z\left(\Gamma \cap M^{1}\right) \backslash M^{1}} c_{P} E_{s, f}\left(m^{\prime} \cdot z_{a}\right) \bar{f}\left(m^{\prime}\right) d m^{\prime} \\
=\int_{Z\left(\Gamma \cap M^{1}\right) \backslash M^{1}}\left(a^{s} \cdot f\left(m^{\prime}\right)+c_{s, f} a^{1-s} \cdot f^{w}\left(m^{\prime}\right)\right) \cdot \bar{f}\left(m^{\prime}\right) d m^{\prime}=a^{s}+c_{s, f} a^{1-s}
\end{gathered}
$$

[11.11.1] Claim: Take $a \gg 1$. For values of $s$ such that $a^{s}+c_{s, f} a^{1-s}=0$ the truncation $\wedge^{a} E_{s, f}$ is an eigenfunction for $\widetilde{S}_{a, f}$, and $\left(\Delta-\lambda_{s, f}\right) \wedge^{a} E_{s, f}$ is a constant multiple of $\eta_{a}$.
Proof: Let $H$ be the usual Heaviside function on $\mathbb{R}$ : identically 0 on $(-\infty, 0)$ and identically 1 on $(0,+\infty)$. The truncation [3.14] (along $P$ ) of $E_{s, f}^{P}$ is a pseudo-Eisenstein series:

$$
\wedge^{a} E_{s, f}^{P}=\Psi_{H(a-y) \beta, f} \quad\left(\text { with } \beta=c_{P} E_{s, f}^{P}\left(z_{y}\right), \text { with } z_{y}=\left(\begin{array}{cc}
y^{\frac{1}{r}} \cdot 1_{r} & 0 \\
0 & 1_{r}
\end{array}\right)\right)
$$

The identity [11.5.1] shows the effect of applying $\Delta-\lambda_{s, f}$ to pseudo-Eisenstein series: at first for test functions $\psi$,

$$
\left(\Delta-\lambda_{s, f}\right) \Psi_{\psi, f}^{P}=\Psi_{D_{s} \psi, f}
$$

where $D_{s}$ is the differential operator

$$
D_{s}=\left(2 r^{3} y \frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)+\mu_{1}+\mu_{2}\right)-\left(2 r^{3} s(s-1)+\mu_{1}+\mu_{2}\right)=2 r^{3}\left(\frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)-s(s-1)\right)
$$

Extend the identity by continuity. Then

$$
\left(\Delta-\lambda_{s, f}\right) \wedge^{a} E_{s, f}^{P}=\left(\Delta-\lambda_{s, f}\right) \Psi_{H(a-y) \beta, f}=\Psi_{D_{s} H(a-y) \beta, f}
$$

Thus, this computation reduces to an elementary computation on $(0,+\infty)$ :

$$
D_{s}(H(a-y) \beta)=2 r^{3}\left(\frac{\partial}{\partial y}\left(y \frac{\partial}{\partial y}-1\right)-s(s-1)\right)\left(H(a-y) \cdot\left(y^{s}+c_{s, f} y^{1-s}\right)\right)
$$

Apart from the leading coefficient $2 r^{3}$, this is the same expression appearing in the proofs in [11.6] for $S L_{2}(\mathbb{Z})$ and for the other three simple cases. That is, a derivative of a Dirac $\delta$ appears unless $y^{s}+c_{s, f} y^{1-s}=0$, in which case the truncated Eisenstein series is indeed in $\mathfrak{E}^{1}$, and is in the domain of $\widetilde{S}$.
[11.11.2] Corollary: Let $f_{1}$ have eigenvalue $\mu_{1}$ for the Laplacian on $S L_{r}(\mathbb{Z}) \backslash S L_{r}(\mathbb{R}) / S O_{r}(\mathbb{R})$, rewritten as $\mu_{1}=\frac{1}{4}-\tau^{2}$ with $\tau \geq 0$. For $a \gg 1$, if $a^{s}+c_{s, f} a^{1-s}=0$, then either $\operatorname{Re}(s)=\frac{1}{2}$ or $s \in\left[\frac{1}{2}-\tau, \frac{1}{2}+\tau\right]$.
Proof: If $a^{s}+c_{s, f} a^{1-s}=0$, then the corresponding truncated Eisenstein series is an eigenfunction for the the non-positive self-adjoint differential operator $\widetilde{S}_{a, f}$, so the corresponding eigenvalue computed in [3.11.11] must be real and non-positive:

$$
\lambda_{s, f}=2 r^{3}\left(s^{2}-s\right)+\mu_{1}+\mu_{1} \leq 0
$$

Thus, $s(s-1)+\frac{1}{4}-\tau^{2} \leq 0$.

### 11.12 Non-self-associate cases

The general argument for cuspidal-data Eisenstein series for maximal proper parabolics in $S L_{r}$ has the same shape as in the four simple cases, and as for $P^{r, r} \subset S L_{2 r}$ with cuspidal data $f=f_{1} \otimes f_{1}$, treated in [11.8], [11.9], and [11.10]. For self-associate maximal proper parabolic $P=P^{r, r} \subset G=S L_{2 r}$, the case opposite to that already considered is $f_{2} \perp f_{1}$, with strong-sense cuspforms $f_{1}, f_{2}$. We recapitulate the argument in outline for $r_{1} \neq r_{2}$, highlighting only the complications and differences from the previous examples.

For non-self-associate $P=P^{r_{1}, r_{2}}$ with $r_{1} \neq r_{2}$, let $Q=P^{r_{2}, r_{1}}$, and only consider cuspidal data $f=f_{1} \otimes f_{2}$ with strong-sense cuspforms $f_{1}, f_{2}$, and put

$$
\mathfrak{E}(P, Q, f)=\left\{\Psi_{\psi, f}^{P}: \psi \in C_{c}^{\infty}(0,+\infty)\right\}+\left\{\Psi_{\psi, f^{w}}^{Q}: \psi \in C_{c}^{\infty}(0,+\infty)\right\}
$$

Again, recall context from [10.6]. Let $B$ be the standard minimal parabolic, with unipotent radical $N^{B}$ and standard Levi component $M^{B}$. Write Iwasawa decompositions $g=n m_{g} k$ with $n \in N^{B}, m \in M^{B}$. We let $\mathfrak{S}$ be a standard Siegel set stable under the (left) action of $N^{B}$ :

$$
\mathfrak{S}=\mathfrak{S}_{b}=\left\{g \in G:\left|\alpha_{j}\left(m_{g}\right)\right| \geq b, \text { for all simple roots } \alpha_{j}\right\}
$$

Take $0<b \ll 1$ such that $\mathfrak{S}_{b} \rightarrow \Gamma \backslash G$ is a surjection. For $a>b$, let $X_{a}$ be the subset of $\mathfrak{S}$ where $\beta\left(m_{g}\right) \leq a$ for all simple roots $\beta$. The quotient $(\Gamma \cap B) \backslash X_{a}$ is compact, since $\left(N^{B} \cap \Gamma\right) \backslash N$ is compact. For each simple root $\beta$, let

$$
Y_{a}^{\beta}=\left\{g \in \mathfrak{S}: \beta\left(m_{g}\right) \geq a\right\}
$$

and $Y_{a}=\bigcup_{\beta} Y_{a}^{\beta}$. Thus, $\mathfrak{S}=X_{a} \cup Y_{a}$. Parallel to [10.6], let

$$
\mathfrak{E}(P, Q, f)_{a}=\left\{F \in \mathfrak{E}(P, Q, f): c_{P^{\prime}} F(g)=0, \text { for all } g \in Y_{a}, \text { for all standard parabolics } P^{\prime}\right\}
$$

and

$$
\mathfrak{E}_{a}^{0}=\mathfrak{B}^{0} \text {-closure of } \mathfrak{E}(P, Q, f)_{a} \quad \mathfrak{E}_{a}^{1}=\mathfrak{B}^{1} \text {-closure of } \mathfrak{E}(P, Q, f)_{a} \quad \mathfrak{E}_{a}^{2}=\mathfrak{B}^{2} \text {-closure of } \mathfrak{E}(P, Q, f)_{a}
$$

It suffices to require vanishing of constant terms for maximal proper parabolics $P^{\prime}$. Further, from [3.11.1], since all pseudo-Eisenstein series in $\mathfrak{E}(P, Q, f)$ have cuspidal data, the vanishing condition is automatically satisfied for all parabolics $P^{\prime}$ except $P$ (and $Q$, in case $Q \neq P$ ).

As earlier, to be careful, since unbounded operators should be densely defined, we need
[11.12.1] Lemma: For $a \gg 1, \mathfrak{E}(P, Q, f)_{a}=\mathfrak{E}(P, Q, f) \cap \mathfrak{E}_{a}^{0}$ is dense in $\mathfrak{E}_{a}^{0}$.
Proof: On the relatively small space $\mathfrak{E}(P, Q, f)$, with $a \gg 1$, the observation [11.6.2] again reduces the issue to a generic, local, one-dimensional issue of smooth cut-offs, as addressed in the proof of [10.3.1], but now admitting the minor complication that constant terms along $P$ and along $Q$ are related.

Let $S_{a, f}$ be $\Delta$ restricted to $\mathfrak{E}(P, Q, f)_{a}$. Since $\Delta \Psi_{\psi, f}^{P}=\Psi_{\beta, f}^{P}$ from [11.6.1], and similarly for $\Psi_{\psi, f^{w}}^{Q}$, and differential operators do not enlarge supports, $\Delta$ does stabilize $\mathfrak{E}(P, Q, f)_{a}$. Let $\widetilde{S}_{a, f}$ be the Friedrichs extension of $S_{a, f}$ to an unbounded self-adjoint operator on $\mathfrak{E}_{a}^{0}$, with domain contained in $\mathfrak{E}_{a}^{1}$ and containing $\mathfrak{E}_{a}^{2}$.
[11.12.2] Corollary: $\widetilde{S}_{a, f}$ has compact resolvents $\left(\widetilde{S}_{a, f}-\lambda_{s, f}\right)^{-1}$ (away from poles).
Proof: As earlier, the point is that the inclusion $\mathfrak{E}_{a}^{1} \rightarrow \mathfrak{E}_{a}^{0}$ is a restriction of the inclusion $\mathfrak{B}_{a}^{1} \rightarrow L_{a}^{2}$, the latter compact from [10.8]. Restrictions of compact operators are compact. The resolvents of the Friedrichs extension are continuous maps $\mathfrak{E}_{a}^{0} \rightarrow \mathfrak{E}_{a}^{1}$ composed with the inclusion $\mathfrak{E}_{a}^{1} \rightarrow \mathfrak{E}_{a}^{0}$. Continuous maps composed with compact maps are compact.

Let $M^{1}$ be the copy of $S L_{r_{1}} \times S L_{r_{2}}$ inside $M=M^{P}$, and $Z^{M}$ the center of $M$. We take representatives

$$
z_{a}=\left(\begin{array}{cc}
a^{\frac{1}{r_{1} r_{2}}} \cdot 1_{r_{1}} & 0 \\
0 & 1_{r_{2}}
\end{array}\right) \quad\left(\text { for } 0<a \in \mathbb{R}^{\times}\right)
$$

## 11. Meromorphic continuation of Eisenstein series

for the connected component $Z \backslash Z^{M}$ containing $1_{r}$, and let $\eta_{a}$ be the functional on $\mathfrak{E}(P, Q, f)$ defined by

$$
\eta_{a}(F)=\int_{Z\left(\Gamma \cap M_{1}^{P}\right) \backslash M_{1}} c_{P} F\left(m^{\prime} \cdot z_{a}\right) \bar{f}\left(m^{\prime}\right) d m^{\prime} \quad(\text { for } F \in \mathfrak{E}(P, Q, f))
$$

Similarly, both for $P$ self-associate and not, let $M^{1}$ be the copy of $S L_{r_{2}} \times S L_{r_{1}}$ inside $M^{Q}$, let

$$
z_{a}^{\prime}=\left(\begin{array}{cc}
a^{\frac{1}{r_{1} r_{2}}} \cdot 1_{r_{2}} & 0 \\
0 & 1_{r_{1}}
\end{array}\right) \quad\left(\text { for } 0<a \in \mathbb{R}^{\times}\right)
$$

and

$$
\eta_{a}^{w}(F)=\int_{Z\left(\Gamma \cap M^{1}\right) \backslash M_{1}} c_{P} F\left(m^{\prime} \cdot z_{a}\right) \overline{f^{w}}\left(m^{\prime}\right) d m^{\prime} \quad(\text { for } F \in \mathfrak{E}(P, Q, f))
$$

Then $F \in \mathfrak{E}(P, Q, f)_{a}$ if and only if $\eta_{b^{\prime}}(F)=0=\eta_{b^{\prime}}^{w}(F)$ for all $b^{\prime} \geq a$. Certainly $\eta_{a}$ and $\eta_{a}^{w}$ do also depend on the cuspidal data.

On genuine Eisenstein series $E_{s, f}^{P}$ and $E_{s, f^{w}}^{Q}$ the functionals $\eta_{a}$ and $\eta_{a}^{w}$ also make sense: unwinding, and using the explicit form of the constant terms, we have absolutely convergent integrals

$$
\begin{gathered}
\eta_{a}\left(E_{s, f}^{P}\right)=\int_{Z\left(\Gamma \cap M^{1}\right) \backslash M^{1}} c_{P} E_{s, f}\left(m^{\prime} \cdot z_{a}\right) \bar{f}\left(m^{\prime}\right) d m^{\prime} \\
=\int_{Z\left(\Gamma \cap M^{1}\right) \backslash M^{1}}\left(a^{s} \cdot f\left(m^{\prime}\right)\right) \cdot \bar{f}\left(m^{\prime}\right) d m^{\prime}=a^{s} \quad\left(\text { for } r_{1} \neq r_{2}\right)
\end{gathered}
$$

Similarly,

$$
\eta_{a}^{w}\left(E_{s, f}^{P}\right)=c_{s, f}^{Q} a^{1-s} \quad\left(\text { for } r_{1} \neq r_{2}\right)
$$

As earlier in [11.3.4] and [11.8.5], the Friedrichs extension can be usefully recharacterized:
[11.12.3] Claim: $\widetilde{S}_{a, f} x=y$ for $x \in \mathfrak{E}_{a}^{1}$ and $y \in \mathfrak{E}_{a}^{0}$ if and only if $\Delta x=y+A \cdot \eta_{a}+B \cdot \eta_{a}^{w}$ for some constants $A, B$.

As earlier, but now with two different tails to accommodate, form two smooth pseudo-Eisenstein series: let $h_{s, f}$ be the smooth pseudo-Eisenstein series formed from a smooth tail of $c_{P} E_{s, f}^{P}$, and $h_{s, f}^{w}$ a smooth pseudo-Eisenstein series formed from a smooth tail of $c_{Q} E_{s, f}^{P}$. To subtract multiples of $h_{s, f}$ and $h_{s, f}^{w}$ from $E_{s, f}^{P}$ to obtain an element of $\mathfrak{E}_{a}^{1}$, use a linear combination of $h_{s, f}$ and $h_{s, f}^{w}$ whose constant terms along both $P$ and $Q$ are both eventually (that is, sufficiently high up in the corresponding Siegel sets) exactly those of $E_{s, f}^{P}$.

From the computation of constant terms in [3.11.3] and [3.11.5], there is a tight relationship between the constant terms of $E_{s, f}^{P}$ and $E_{s, f w}^{Q}$, which after meromorphic continuation gives the functional equation $E_{1-s, f}=\left(c_{s, f w}^{P}\right)^{-1} \cdot E_{s, f^{w}}^{Q}$ and $c_{1-s, f}^{Q} \cdot c_{s, f w}^{P}=1$.

## 11.A Appendix: distributions supported on submanifolds

The fact that distributions supported at a single point are finite linear combinations of Dirac delta and its derivatives is the simplest special case of the following, which reduces questions about distributions supported on smooth submanifolds to the local situation of Euclidean spaces.
[11.A.1] Theorem: A distribution $u$ on $\mathbb{R}^{m+n} \approx \mathbb{R}^{m} \times \mathbb{R}^{n}$ supported on $\mathbb{R}^{m} \times\{0\}$, is uniquely expressible as a locally finite sum of transverse differentiations followed by restriction and evaluations, namely, a locally finite sum

$$
u=\sum_{\alpha} u_{\alpha} \circ \operatorname{Res}_{\mathbb{R}^{m} \times\{0\}}^{\mathbb{R}^{m} \times \mathbb{R}^{n}} \circ D^{\alpha}
$$

where $\alpha$ is summed over multi-indices $\left(\alpha_{1}, \ldots, \alpha_{n}\right), D^{\alpha}$ is the corresponding differential operator on $\{0\} \times \mathbb{R}^{n}$, and $u_{\alpha}$ are distributions on $\mathbb{R}^{m} \times\{0\}$. Further,

$$
\operatorname{spt} u_{\alpha} \times\{0\} \subset \operatorname{spt} u \quad(\text { for all multi-indices } \alpha)
$$

Proof: For brevity, let

$$
\rho=\operatorname{Res}_{\mathbb{R}^{m} \times\{0\}}^{R^{m} \times R^{n}}: C_{c}^{\infty}\left(\mathbb{R}^{m} \times \mathbb{R}^{n}\right) \longrightarrow C_{c}^{\infty}\left(\mathbb{R}^{m}\right)
$$

be the natural restriction map of test functions on $\mathbb{R}^{m} \times \mathbb{R}^{n}$ to $\mathbb{R}^{m} \times\{0\}$, by

$$
(\rho f)(x)=f(x, 0) \quad\left(\text { for } x \in \mathbb{R}^{m}\right)
$$

The adjoint $\rho^{*}: C_{c}^{\infty}\left(\mathbb{R}^{m}\right)^{*} \rightarrow C_{c}^{\infty}\left(\mathbb{R}^{m+n}\right)^{*}$ is a continuous map of distributions on $\mathbb{R}^{m}$ to distributions on $\mathbb{R}^{m} \times \mathbb{R}^{n}$, defined by

$$
\left(\rho^{*} u\right)(f)=u(\rho(f))
$$

First, if we could apply $u$ to functions of the form $F(x, y)=f(x) \cdot y^{\beta}$, and if $u$ had an expression as a sum as in the statement of the theorem, then

$$
u\left(f(x) \cdot \frac{y^{\alpha}}{\alpha!}\right)=(-1)^{|\beta|} \cdot u_{\beta}(f) \cdot \beta!
$$

since most of the transverse derivatives evaluated at 0 vanish. This is not quite legitimate, since $y^{\alpha}$ is not a test function. However, we can take a test function $\psi$ on $\mathbb{R}^{n}$ that is identically 1 near 0 , and consider $\psi(y) \cdot y^{\alpha}$ instead of $y^{\alpha}$, and reach the same conclusion.

Thus, if there exists such an expression for $u$, it is unique. Further, this computation suggests how to specify the $u_{\alpha}$, namely,

$$
u_{\beta}(f)=u\left(f(x) \otimes \frac{y^{\beta}}{\beta!} \cdot \psi(y) \cdot(-1)^{|\beta|}\right)
$$

This would also show the containment of the supports.
Show that the sum of these $u_{\beta}$ 's does give $u$. Given an open $U$ in $\mathbb{R}^{m+n}$ with compact closure, $u$ on $C_{c}^{\infty} U$ has some finite order $k$. As a slight generalization of the fact that distributions supported on $\{0\}$ are finite linear combinations of Dirac delta and its derivatives, we have
[11.A.2] Lemma: Let $v$ be a distribution of finite order $k$ supported on a compact set $K$. For a test function $\varphi$ whose derivatives up through order $k$ vanish on $K, v(\varphi)=0$.

For any test function $F(x, y)$,

$$
\Phi(x, y)=F(x, y)-\sum_{|\alpha| \leq k}(-1)^{|\alpha|} \frac{y^{\alpha}}{\alpha!} \psi(y)\left(D^{\alpha} F\right)(x, 0)
$$

has all derivatives vanishing to order $k$ on the closure of $U$. Thus, by the lemma, $u(\Phi)=0$, which proves that $u$ is equal to that sum, and also proves the local finiteness.

## 12. Global automorphic Sobolev spaces, Green's functions

1. A simple pre-trace formula
2. Pre-trace formula for compact periods
3. Global automorphic Sobolev spaces $H^{\ell}$
4. Spectral characterization of Sobolev spaces $H^{s}$
5. Continuation of solutions of differential equations
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9. Intertwining operators among principal series

Appendix A: a usual trick with $\Gamma(s)$
The pre-trace formulas below depend on estimates of eigenfunctions of integral operators on automorphic forms, in terms of the eigenvalues for the invariant Laplacian. This is accomplished by two observations. First, the eigenvalues depend only on the isomorphism class of the space generated by the group acting by right translations on the given automorphic form, so that these eigenvalues can be computed on any isomorphic copy of such a space, as representation space for $G$. Second, it happens that there are muchsimpler isomorphic copies of the relevant representations, parametrized essentially by one or more complex numbers, namely, principal series representations. We emphasize the archimedean aspect in this chapter, for which the general result is the subrepresentation theorem of [Casselman 1978/80], [Casselman Miličić 1982], improving the subquotient theorem of [Harish-Chandra 1954]. A simple argument sufficient for the four simple examples follows from older results on asymptotics for solutions of second-order ordinary differential equations, recalled in an appendix (chapter sixteen).

The pre-trace formulas ground the discussion of global automorphic Sobolev spaces. Among other goals, an important one is interpretation and legitimization of term-wise differentiation of $L^{2}$ automorphic spectral expansions, especially by the invariant Laplacian. Of course, Plancherel theorems refer to $L^{2}$ expansions. Significantly, Plancherel theorems do not refer to sup norms of the eigenfunctions (such as cuspforms) entering in a spectral decomposition, nor sup norms of non- $L^{2}$ eigenfunctions (such as Eisenstein series) entering in the $L^{2}$ decomposition. This is already manifest in Plancherel for Fourier inversion on $L^{2}\left(\mathbb{R}^{n}\right)$. Typically, $L^{2}$ expansions do not produce continuous functions, and are not continuously differentiable, so the goal cannot possibly be proving classical differentiability, since it does not hold. Especially with respect to invariant operators such as Casimir operators, and in delicate situations such as automorphic forms, Plancherel theorems most naturally yield corollaries about an extension by continuity of the classical limit-of-difference-quotient notion of differentiation. This $L^{2}$-differentiation is a usefully refined distributional differentiation. Term-wise differentiability of $L^{2}$ spectral expansions in a distributional sense is of course correct, but needlessly very weak, and specifically too weak for many applications, since it is difficult to return from the larger world of distributions to the smaller world of $L^{2}$ functions.

Further, already for Fourier transforms on $\mathbb{R}^{n}$, the apparent integral expressing Fourier inversion is not a superposition of $L^{2}$ functions, since the exponential functions are not in $L^{2}\left(\mathbb{R}^{n}\right)$. Similarly, the spectral decomposition of pseudo-Eisenstein series involves integrals of the corresponding non- $L^{2}$ Eisenstein series.

The global-ness of the automorphic Sobolev spaces first refers to the expression of the norms (on automorphic test functions) as integrals over the whole space $\Gamma \backslash G / K$, rather than as a collection of seminorms given by integrals over smaller sets. Equivalently, the norms have expressions in terms of $L^{2}$ spectral expansions in terms of eigenfunctions for $\Delta$. Global $L^{2}$ Sobolev spaces balance the simplicity of Hilbert space structures with extensions of notions of differentiability, insofar as solving more-or-less elliptic partial differential equations of sufficiently high degree can move back to $L^{2}$ from nearby Sobolev spaces of distributions. That is, in terms of the basic processes of analysis, Sobolev spaces are within finite distance of $L^{2}$.

Among other applications of global automorphic Sobolev spaces, we can immediately write a spectral expansion of an automorphic Green's function.

### 12.1 A simple pre-trace formula

Let $G, \Gamma, K, P, M, N, A^{+}$be as in the four examples from chapter one, with Iwasawa coordinates $x, y[1.3]$ with $x \in \mathbb{R}^{r-1}$ and $0<y \in \mathbb{R}$, with $r=2,3,4,5$ the dimension of $G / K$. From [4.5], [4.6], [4.7], [4.8], the invariant Laplacian is

$$
\Delta=y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\ldots+\frac{\partial^{2}}{\partial x_{r}^{2}}\right)-(r-2) y \frac{\partial}{\partial y}
$$

For complex $s$, let $\lambda_{s}=(r-1)^{2} \cdot s(s-1)$. For cuspforms $F$ in an orthonormal basis, let $s_{F} \in \mathbb{C}$ be such that the $\Delta$-eigenvalue of $F$ is $\lambda_{s_{F}}$, and let $t_{F}$ be the imaginary part of $s_{F}$.
[12.1.1] Theorem: Fix $z_{o}=(x, y) \in G / K$. Then

$$
\sum_{F:\left|t_{F}\right| \leq T}\left|F\left(z_{o}\right)\right|^{2}+\int_{|t| \leq T}\left|E_{\frac{1}{2}+i t}\left(z_{o}\right)\right|^{2} d t \ll_{C} \quad T^{r} \quad(\text { for } T \rightarrow+\infty)
$$

with implied constant uniform for $z_{o}$ in a fixed compact $C$.
Proof: We consider integral operators attached to compactly supported (regular Borel) measures $\eta$ on the group $G$, and their operation on any reasonable representation space $V$ for $G$, for example, Hilbert, Banach, Fréchet, and LF (strict colimits of Fréchet), or, generally, quasi-complete, locally convex spaces. For a continuous action $G \times V \rightarrow V$ of $G$ on such a space $V$, and compactly-supported measure $\eta$, the action is

$$
\eta \cdot v=\int_{G} g \cdot v d \eta(g) \quad(\text { for } v \in V)
$$

as Gelfand-Pettis integral. The further non-trivial fact used in the proof ${ }^{[71]}$ is that the eigenvalues of these integral operators on automorphic forms on $\Gamma \backslash G / K$ depend only on their eigenvalues for the Laplacian. This itself will follow from the fact that for $\Delta f=\lambda_{s} \cdot f$, a suitable topological vector space of functions on $\Gamma \backslash G$ generated by right translates of $f$ is isomorphic to a subquotient of the principal series representation $I_{s}$ (below), a relatively elementary object. The same is true of Eisenstein series $E_{s}$ more immediately, since the Eisenstein series is (the meromorphic continuation of) a wound-up function from $I_{s}$ (below).

Spaces $V, W$ with continuous actions $G \times V \rightarrow V$ and $G \times W \rightarrow W$ are representations of $G$. Continuous $\mathbb{C}$-linear maps $T: V \rightarrow W$ among such spaces, respecting the action of $G: T(g \cdot v)=g \cdot T(v)$ for $v \in V$ and $g \in G$, are $G$-homomorphisms or $G$-intertwinings. The eigenvalues and eigenvectors of integral operators are preserved by $G$-homomorphisms:
[12.1.2] Lemma: Let $T: V \rightarrow W$ be a homomorphism of $G$-spaces. For $\eta$ a compactly-supported measure on $G$, the action of $\eta$ commutes with $T$ :

$$
\eta \cdot T(v)=T(\eta \cdot v) \quad(\text { for all } v \in V)
$$

In particular, for $\eta \cdot v=\lambda \cdot v$ with $v \in V$, the image $T v$ is also an eigenvector, with the same eigenvalue $\lambda$. Proof: By properties of Gelfand-Pettis integrals and the fact that $T$ commutes with the action of $g$,

$$
T(\eta \cdot v)=T \int_{G} g \cdot v d \eta(g)=\int_{G} T(g \cdot v) d \eta(g)=\int_{G} g \cdot T(v) d \eta(g)=\eta \cdot T(v)
$$

as claimed.
For $v$ in a $G$-representation $V$, the subrepresentation generated by $v$ is the topological closure of the span of finite linear combinations of images $g \cdot v$ of $v$ by $g \in G$. Among $K$-invariant vectors in a $G$-representation, eigenvectors for the spherical Hecke-algebra $\mathcal{H}$ occupy a privileged position:
[71] This is a very small instance of a subquotient theorem [Harish-Chandra 1954]. This was strengthened to the subrepresentation theorem [Casselman 1978/80], [Casselman-Miličić 1982].
[12.1.3] Claim: A strong-sense cuspform $f$, or Eisenstein series $E_{s}$, is the unique $K$-invariant vector in the representation it generates under right translation by $G$, up to a constant. More generally, for $v \neq 0$ a $K$ fixed vector in a $G$-representation $V$ with $V$ quasi-complete and locally convex, for $v$ also an $\mathcal{H}$-eigenvector, the subrepresentation generated by $v$ has $K$-fixed vectors exactly $\mathbb{C} \cdot v$.
Proof: Let $\alpha$ be the average-over- $K$ map $v \rightarrow \int_{K} k \cdot v d k$, giving $K$ total mass 1. This Gelfand-Pettis integral maps $V$ to the $K$-fixed vectors $V^{K}$, and is the identity map on $V^{K}$.

First, consider $w=\sum_{i=1}^{n} c_{i} g_{i} \cdot v$ a $K$-fixed vector in the (algebraic) span of images $g_{i} \cdot v$ of $v$ by $g_{i} \in G$. From a basic property [14.1.4] of Gelfand-Pettis integrals, $\varphi_{j} \cdot w \rightarrow w$, for any approximate identity $\left\{\varphi_{j}\right\}$. Since $\alpha \cdot w=w$ and $\alpha \cdot v=v$,

$$
w=\alpha \cdot w=\alpha \cdot \lim _{j} \varphi_{j} \cdot w=\alpha \lim _{j} \varphi_{j} \cdot \sum_{i} c_{i} g_{i} \cdot \alpha \cdot v
$$

By basic properties of Gelfand-Pettis integrals, the operator $\alpha$ commutes with the limit:

$$
w=\lim _{j}\left(\sum_{i} c_{i} \alpha \circ \varphi_{j} \circ g_{i} \circ \alpha\right) \cdot v
$$

The function

$$
\eta_{j}=\sum_{i} c_{i} \alpha \circ \varphi_{j} \circ g_{i} \circ \alpha
$$

is in $C_{c}^{o}(G)$ and is left and right $K$-invariant, so is in the spherical Hecke algebra $\mathcal{H}$. Since $v$ is an $\mathcal{H}$ eigenvector, $\eta_{j} \cdot v=\lambda_{j} \cdot v$ for some scalar $\lambda_{j}$. That is, $\lambda_{j} \cdot v \rightarrow w$, so $w \in \mathbb{C} \cdot v$. Next, for $w=\lim _{j} w_{j}$ a limit of $w_{j}$ of $w_{j}$ in the span of images $g_{i} \cdot v, \alpha w=\alpha \lim w_{j}=\lim \alpha w_{j}$, and by the previous paragraph every $\alpha w_{j}$ is a scalar multiple of $v$, so $w$ must be, as well.

Since Eisenstein series and strong-sense cuspforms are $\mathcal{H}$-eigenfunctions, the conclusion applies to them, as well.

Thus, for any left-and-right $K$-invariant compactly-supported measure $\eta$ the integral operator action

$$
(\eta \cdot f)(x)=\int_{G} f(x y) \cdot d \eta(y)
$$

produces another right $K$-invariant vector in the representation space $V_{f}$ generated by $f$. By the claim, $\eta \cdot f$ is a scalar multiple of $f$. Let $\chi_{f}(\eta)$ denote the eigenvalue:

$$
\eta \cdot f=\chi_{f}(\eta) \cdot f \quad\left(\text { with } \chi_{f}(\eta) \in \mathbb{C}\right)
$$

By the lemma, this is an intrinsic representation-theoretic relation, meaning that the scalar $\chi_{f}(\eta)$ can be computed in any image of $V_{f}$. As demonstrated subsequently, for $\Delta f=(r-1)^{2} \cdot s(s-1)$, the representation generated by $f$ has a common image ${ }^{[72]}$ with an unramified principal series

$$
I_{s}=\left\{\varphi \in C^{\infty}(G): \varphi\left(\left(\begin{array}{ll}
a & * \\
0 & d
\end{array}\right) \cdot g\right)=\left|\frac{a}{d}\right|^{(r-1) s} \cdot \varphi(g)\right\} \quad \quad(\text { with } s \in \mathbb{C})
$$

under right translation action by $G$. The Iwasawa decomposition $G=P \cdot K$ shows that the space of $K$-fixed vectors is one-dimensional. Thus, by the lemma, we can compute eigenvalues of elements of $\mathcal{H}$ on Eisenstein series or strong-sense cuspforms by computing eigenvalues on $I_{s}$.

Choice of integral operators: let $\|g\|$ be the square of the operator norm on $G$ for a standard representation of $G$ on $\mathbb{C}^{2}$ or $\mathbb{C}^{4}$ (depending on cases) by matrix multiplication. In a Cartan decomposition,

$$
\left\|k_{1} \cdot\left(\begin{array}{cc}
e^{\rho / 2} & 0 \\
0 & e^{-\rho / 2}
\end{array}\right) \cdot k_{2}\right\|=e^{r} \quad\left(\text { with } k_{1}, k_{2} \in K, \rho \geq 0\right)
$$

[72] As will be visible, this common image is inside a Whittaker space of smooth functions on $G$ with suitable left equivariance under $N$.

This norm gives a left $G$-invariant metric $d($,$) on G / K$ by

$$
d(g K, h K)=\log \left\|g^{-1} h\right\|=\log \left\|h^{-1} g\right\|
$$

The triangle inequality follows from the submultiplicativity of the norm. Take $\eta$ to be the characteristic function of the left and right $K$-invariant set of group elements of operator norm at most $e^{\delta}$, with small $\delta>0$. That is,

$$
\eta(g)= \begin{cases}1 & \left(\text { for }\|g\| \leq e^{\delta}\right) \\ 0 & \left(\text { for }\|g\|>e^{\delta}\right)\end{cases}
$$

or

$$
\eta\left(k_{1} \cdot\left(\begin{array}{cc}
e^{\rho / 2} & 0 \\
0 & e^{-\rho / 2}
\end{array}\right) \cdot k_{2}\right)=\left\{\begin{array}{ll}
1 & (\text { for } \rho \leq \delta) \\
0 & (\text { for } \rho>\delta)
\end{array} \quad(\text { with } \rho \geq 0)\right.
$$

Upper bound on a kernel: The map $f \rightarrow(\eta \cdot f)(x)$ on automorphic forms $f$ can be expressed as integration of $f$ against a sort of automorphic form $q_{x}$ by winding up the integral, as follows.

$$
\begin{gathered}
(\eta \cdot f)(x)=\int_{G} f(x y) \eta(y) d y=\int_{G} f(y) \eta\left(x^{-1} y\right) d y=\int_{\Gamma \backslash G}\left(\sum_{\gamma \in \Gamma} f(\gamma y) \eta\left(x^{-1} \gamma y\right)\right) d y \\
=\int_{\Gamma \backslash G} f(y) \cdot\left(\sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma y\right)\right) d y
\end{gathered}
$$

Thus, for $x, y \in G$ put

$$
q_{x}(y)=\sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma y\right)
$$

The norm-squared of $q_{x}$, as a function of $y$ alone, is

$$
\left|q_{x}\right|_{L^{2}(\Gamma \backslash G)}^{2}=\int_{\Gamma \backslash G} \sum_{\gamma \in \Gamma} \sum_{\gamma^{\prime} \in \Gamma} \eta\left(x^{-1} \gamma \gamma y\right) \bar{\eta}\left(x^{-1} \gamma^{\prime} y\right) d y=\int_{G} \sum_{\gamma \in \Gamma} \eta\left(x^{-1} \gamma y\right) \bar{\eta}\left(x^{-1} y\right) d y
$$

after unwinding. For both $\eta\left(x^{-1} \gamma y\right)$ and $\eta\left(x^{-1} y\right)$ to be non-zero, the distance from $x$ to both $y$ and $\gamma y$ must be at most $\delta$. By the triangle inequality, the distance from $y$ to $\gamma y$ must be at most $2 \delta$. For $x$ in a fixed compact $C$, this requires that $y$ be in ball of radius $\delta$, and that $\gamma y=y$. Since $K$ is compact and $\Gamma$ is discrete, the isotropy groups of all points in $G / K$ are finite. Thus,

$$
\left|q_{x}\right|_{L^{2}(\Gamma \backslash G)}^{2} \ll \int_{d(x, y) \leq \delta} 1 d y \asymp_{C} \delta^{r} \quad\left(\text { as } \delta \rightarrow 0^{+}\right)
$$

Lower bound on eigenvalues: let $\Delta f=(r-1)^{2} \cdot s_{f}\left(s_{f}-1\right)$, with $s_{f}=\frac{1}{2}+i t_{f}$. A non-trivial lower bound for $\chi_{f}(\eta)$ can be given for $\delta \ll 1 / t_{f}$, as follows. With spherical function

$$
\varphi_{s}^{o}\left(\left(\begin{array}{ll}
a & * \\
0 & d
\end{array}\right) \cdot k\right)=\left|\frac{a}{d}\right|^{(r-1) s}
$$

in the $s^{t h}$ principal series, the corresponding eigenvalue is

$$
\chi_{s}(\eta)=\int_{G} \eta(g) \varphi_{s}^{o}(g) d g=\int_{r \leq \delta} \varphi_{s}^{o}\left(k \cdot\left(\begin{array}{cc}
e^{r / 2} & 0 \\
0 & e^{-r / 2}
\end{array}\right)\right) d g
$$

In fact, a qualitative argument clearly indicates the outcome, although we will also carry out a more explicit computation. For the qualitative argument, we need qualitative metrical properties of the Iwasawa decomposition. Let $g \rightarrow n_{g} a_{g}^{+} k_{g}$ be the Iwasawa decomposition. We claim that $\|g\| \leq \delta$ implies $\left\|n_{g} a_{g}^{+}\right\| \ll \delta$
for small $\delta>0$. This is immediate, since the Jacobian of the map $N \times A^{+} \rightarrow G / K$ near $e \in N A^{+}$is invertible.

But, also, the Iwasawa decomposition is easily computed here for $G=S L_{2}(\mathbb{R})$ and $S L_{2}(\mathbb{C})$, and the integral expressing the eigenvalue can be estimated explicitly: elements of $K$ can be parametrized as

$$
k=\left(\begin{array}{cc}
\bar{\alpha} & \bar{\beta} \\
-\beta & \alpha
\end{array}\right) \quad\left(\text { where }|\alpha|^{2}+|\beta|^{2}=1\right)
$$

and let $a=e^{r / 2}$. Then

$$
k \cdot\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)=\left(\begin{array}{cc}
* & * \\
-a \beta & \alpha / a
\end{array}\right)
$$

Right multiplication by a suitable element $k_{2}$ of $S U(2)$ rotates the bottom row to put the matrix into $N A^{+}$:

$$
k \cdot\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right) \cdot k_{2}=\left(\begin{array}{cc}
* & * \\
0 & \sqrt{(-a|\beta|)^{2}+(|\alpha| / a)^{2}}
\end{array}\right)
$$

Thus,

$$
\chi_{s}(\eta)=\int_{r \leq \delta}\left((-a|\beta|)^{2}+(|\alpha| / a)^{2}\right)^{-s} d g
$$

Rather than compute the integral exactly, make $\delta$ small enough to give a lower bound on the integrand, such as would arise from

$$
\left.\left|\left((-a|\beta|)^{2}+(|\alpha| / a)^{2}\right)^{-s}-1\right|<\frac{1}{2} \quad \text { (for all elements of } K\right)
$$

Since $|\alpha|^{2}+|\beta|^{2}=1$, for small $\rho$,

$$
\left(-e^{\rho / 2}|\beta|\right)^{2}+\left(|\alpha| / e^{\rho / 2}\right)^{2}=e^{\rho}|\beta|^{2}+|\alpha|^{2} / e^{\rho} \asymp(1+\rho)|\beta|^{2}+(1-\rho)|\alpha|^{2} \ll 1+\rho
$$

Thus, for small $0 \leq \rho \leq \delta$,

$$
\left|\left(e^{\rho}|\beta|^{2}+|\alpha|^{2} / e^{\rho}\right)^{-s}-1\right| \ll|s| \cdot \rho
$$

Thus, $0 \leq \rho \leq \delta \ll \frac{1}{|s|}$ suffices to make this less than $\frac{1}{2}$.
From either a qualitative or quantitative approach, we see that with $\eta$ the characteristic function of the $\delta$-ball, we have the lower bound

$$
\left|\chi_{s}(\eta)\right|=\int_{G} \eta(g) \varphi_{s}^{o}(g) d g \gg \int_{\rho \leq \delta} 1=\operatorname{vol}(\delta \text {-ball }) \asymp \delta^{r} \quad(\eta \text { char fcn of } \delta \text {-ball, for }|s| \ll 1 / \delta,)
$$

Taking $\delta$ as large as possible compatible with $\delta \ll 1 /|s|$ gives the lower bound

$$
\chi_{s}(\eta) \gg \delta^{r} \quad(\text { for }|s| \ll 1 / \delta, \eta \text { the characteristic function of } \delta \text {-ball })
$$

Combining the estimates: From the $L^{2}$ automorphic spectral expansion of $q_{x}$, apply Plancherel, dropping the finitely-many terms from residues of Eisenstein series, and dropping normalization constants,

$$
\sum_{F}\left|\left\langle q_{x}, F\right\rangle\right|^{2}+\int_{-\infty}^{+\infty}\left|\left\langle q_{x}, E_{s}\right\rangle\right|^{2} d t \leq\left|q_{x}\right|_{L^{2}(\Gamma \backslash G / K)}^{2} \ll \delta^{r}
$$

Truncating this to Bessel's inequality

$$
\sum_{\left|t_{F}\right| \leq T}\left|\left\langle q_{x}, F\right\rangle\right|^{2}+\int_{-T}^{+T}\left|\left\langle q_{x}, E_{s}\right\rangle\right|^{2} d t \ll \delta^{r}
$$

From the eigenvalue computation above,

$$
\left\langle q_{x}, f\right\rangle=\chi_{s}(\eta) \cdot f
$$

and use the inequality $\chi_{s}(\eta) \gg \delta^{r}$ from above for this restricted parameter range, obtaining

$$
\sum_{\left|t_{F}\right| \leq T}\left(\delta^{r} \cdot|F(x)|\right)^{2}+\int_{-T}^{+T}\left(\delta^{r} \cdot\left|E_{s}(x)\right|\right)^{2} d t \ll \delta^{r}
$$

Multiply through by $T^{2 r} \asymp 1 / \delta^{2 r}$ to obtain the standard estimate or pre-trace formula

$$
\sum_{\left|t_{F}\right| \leq T}|F(x)|^{2}+\int_{-T}^{+T}\left|E_{s}(x)\right|^{2} d t \ll T^{r}
$$

as claimed above. Since the argument succeeds for both $s$ and $1-s$, the ambiguity in determining $s$ from $s(s-1)$ is irrelevant.

### 12.2 Pre-trace formula for compact periods

The argument of the preceding section is a prototype. Now consider somewhat more general $G$, including not only $S L_{2}(\mathbb{R})$, but also $G=S L_{n}(\mathbb{R})$ or $S L_{n}(\mathbb{C})$, with $\Gamma=S L_{n}(\mathbb{Z})$ or $\Gamma=S L_{n}(\mathbb{Z}[i])$, respectively. We will not prove the corresponding subquotient or subrepresentation theorems: see [Harish-Chandra 1954], [Casselman 1978/80], [Casselman-Miličić 1982].

For closed subgroup $H$ of $G$, let $\Theta=H \cap \Gamma$, and suppose that $\Theta \backslash H$ is compact. The $\Theta \backslash H x$-period of $f$ is

$$
\Theta \backslash H x \text {-period of } f=f_{\Theta \backslash H x}=\int_{\Theta \backslash H} f(h x) d h
$$

Similarly, with $\psi$ an automorphic form on $\Theta \backslash H$, the period of $\psi \otimes f$ is

$$
\langle f, \psi\rangle_{\Theta \backslash H x}=\int_{\Theta \backslash H} \psi(h) \cdot f(h x) d h
$$

[12.2.1] Theorem: Using abbreviated notation for the spectral expansions to implicitly include the appropriate integrals of cuspidal-data Eisenstein series, as well as their residues,

$$
\sum_{\operatorname{cfm} F:\left|t_{F}\right| \ll T}\left|F_{\Theta \backslash H x}\right|^{2}+\ldots<_{x, H} T^{\operatorname{dim} X-\operatorname{dim} Y}
$$

and, similarly

$$
\sum_{\operatorname{cfm} F:\left|t_{F}\right| \ll T}|\langle\eta \cdot F, \psi\rangle|^{2}+\ldots \lll x, H, \psi \quad T^{\operatorname{dim} X-\operatorname{dim} Y}
$$

Proof: The usual action of compactly-supported measures $\eta$ on suitable $f$ on $\Gamma \backslash G$ is $(\eta \cdot f)(x)=$ $\int_{G} \eta(g) f(x g) d g$. The $\Theta \backslash H x$-period of $\eta \cdot f$ admits a useful rearrangement

$$
\begin{gathered}
(\eta \cdot f)_{\Theta \backslash H x}=\int_{\Theta \backslash H}(\eta \cdot f)(h x) d h=\int_{\Theta \backslash H} \int_{G} \eta(g) f(h x g) d g d h=\int_{\Theta \backslash H} \int_{G} \eta\left(x^{-1} h^{-1} g\right) f(g) d g d h \\
\quad=\int_{\Theta \backslash H} \int_{\Gamma \backslash G} \sum_{\gamma \in \Gamma} \eta\left(x^{-1} h^{-1} \gamma g\right) f(g) d g d h=\int_{\Gamma \backslash G} f(g)\left(\int_{\Theta \backslash H} \sum_{\gamma \in \Gamma} \eta\left(x^{-1} h^{-1} \gamma g\right) d h\right) d g
\end{gathered}
$$

## 12. Global automorphic Sobolev spaces, Green's functions

Denote the inner sum-and-integral by $q(g)=q_{H, x}(g)$.
Similarly, with $\psi$ an automorphic form on $\Theta \backslash H$, the period of $\psi \otimes f$ rearranges to

$$
\langle\eta \cdot f, \psi\rangle_{\Theta \backslash H x}=\int_{\Gamma \backslash G} f(g)\left(\int_{\Theta \backslash H} \psi(h) \cdot \sum_{\gamma \in \Gamma} \eta\left(x^{-1} h^{-1} \gamma g\right) d h\right) d g
$$

For $\eta$ left-and-right $K$-invariant, for $f$ the spherical vector in a copy of a principal series representation of $G$, necessarily $\eta \cdot f=\lambda_{f}(\eta) \cdot f$ for some constant $\lambda_{f}(\eta)$. Thus, the action of such $\eta$ changes the period by the eigenvalue:

$$
(\eta \cdot f)_{\Theta \backslash H x}=\lambda_{f}(\eta) \cdot f_{\Theta \backslash H x}
$$

An upper bound for the $L^{2}(\Gamma \backslash G)$ norm of $q$, and a lower bound for $\lambda_{f}(\eta)$ contingent on restrictions on the spectral parameter of $f$, yield, by Bessel's inequality, an upper bound for a sum-and-integral of periods $\langle f, \psi\rangle_{\Theta \backslash H x}$ :

Estimating the $L^{2}$ norm,

$$
\begin{gathered}
\int_{\Gamma \backslash G}|q(g)|^{2} d g=\int_{\Gamma \backslash G} \int_{\Theta \backslash H} \int_{\Theta \backslash H} \sum_{\gamma \in \Gamma} \sum_{\gamma_{2} \in \Gamma} \eta\left(x^{-1} h^{-1} \gamma g\right) \bar{\eta}\left(x^{-1} h_{2}^{-1} \gamma_{2} g\right) d h d h_{2} d g \\
=\int_{G} \int_{\Theta \backslash H} \int_{\Theta \backslash H} \sum_{\gamma \in \Gamma} \eta\left(x^{-1} h^{-1} \gamma g\right) \bar{\eta}\left(x^{-1} h_{2}^{-1} g\right) d h d h_{2} d g
\end{gathered}
$$

With $C \subset H$ a large-enough compact to surject to $\Theta \backslash H$,

$$
\int_{\Gamma \backslash G}|q(g)|^{2} d g \leq \int_{G} \int_{C} \int_{C} \sum_{\gamma \in \Gamma}|\eta|\left(x^{-1} h^{-1} \gamma g\right)|\eta|\left(x^{-1} h_{2}^{-1} g\right) d h d h_{2} d g
$$

Let $\eta$ be the characteristic function of a small ball $B_{\varepsilon}$ in $G / K$, of geodesic radius $\varepsilon>0$, for a $G$-invariant metric $d(x, y)=\nu\left(x^{-1} y\right)$ on $G / K$, where $\nu(g)=\log \sup \left(|g|,\left|g^{-1}\right|\right)$, where $|\cdot|$ is operator norm on $G$. The triangle inequality follows from submultiplicativity of operator norm.

Identify $B_{\varepsilon}$ with its pre-image $B_{\varepsilon} \cdot K$ in $G$. The set

$$
\begin{gathered}
\Phi=\Phi_{H, x, \eta}=\left\{\gamma \in \Gamma: \eta\left(x^{-1} h^{-1} \gamma g\right) \eta\left(x^{-1} h_{2}^{-1} g\right) \neq 0 \text { for some } h, h_{2} \in C \text { and } g \in G\right\} \\
=\left\{\gamma \in \Gamma: \gamma \in C x B_{\varepsilon} g^{-1}, g \in C x B_{\varepsilon}\right\} \quad \subset \quad \Gamma \cap C x B_{\varepsilon} \cdot\left(C x B_{\varepsilon}\right)^{-1}=\text { discrete } \cap \text { compact }
\end{gathered}
$$

is finite, and can only shrink as $\varepsilon \rightarrow 0^{+}$.
For each $\gamma \in \Phi$, for each $h \in C, \eta\left(x^{-1} h^{-1} \gamma g\right) \neq 0$ only for $g$ in a ball in $X=G / K$ of radius $\varepsilon$, with volume dominated by $\varepsilon^{\operatorname{dim} X}$. For each $h$ and $g, \eta\left(x^{-1} h_{2}^{-1} g\right) \neq 0$ only for $h_{2} x$ in a ball in $Y=H x K / K$ of radius $\varepsilon$, with volume dominated by $\varepsilon^{\operatorname{dim} Y}$. Thus,

$$
\int_{\Gamma \backslash G}|q(g)|^{2} d g \ll \int_{C} \varepsilon^{\operatorname{dim} X+\operatorname{dim} Y} d h \lll x, H \quad \varepsilon^{\operatorname{dim} X+\operatorname{dim} Y}
$$

By Plancherel for $L^{2}(\Gamma \backslash X)$, with $\eta$ the characteristic function of the $\varepsilon$-ball,

$$
\sum_{\operatorname{cfm} F}\left|\lambda_{F}(\eta)\right|^{2} \cdot\left|F_{\Theta \backslash H x}\right|^{2}+\ldots=\sum_{\operatorname{cfm} F}\left|(\eta \cdot F)_{\Theta \backslash H x}\right|^{2}+\ldots=\left|q_{H, x, \eta}\right|^{2}<_{x, H} \quad \varepsilon^{\operatorname{dim} X+\operatorname{dim} Y}
$$

Similarly,

$$
\sum_{\operatorname{cfm} F}\left|\lambda_{F}(\eta)\right|^{2} \cdot|\langle\eta \cdot F, \psi\rangle|^{2}+\ldots<_{x, H, \psi} \quad \varepsilon^{\operatorname{dim} X+\operatorname{dim} Y}
$$

Next, a bound on the spectral data is determined to give a non-trivial lower bound for $\left|\lambda_{f}(\eta)\right|$. For $f$ the spherical vector in a copy of a principal series representation of $G$, left-and-right $K$-invariant $\eta$ necessarily

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gives $\eta \cdot f=\lambda_{f}(\eta) \cdot f$, since up to scalars $f$ is the unique spherical vector in the irreducible representation it generates.

The eigenvalues $\lambda_{f}(\eta)$ can be computed in the usual model of principal series, as $\eta \cdot \varphi_{s}^{o}=\lambda_{f}(\eta) \cdot \varphi_{s}^{o}$ for $\varphi_{s}^{o}$ the normalized spherical vector for $s \in \mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$, and $\varphi^{o}(1)=1$. Thus,

$$
\lambda_{f}(\eta)=\left(\eta \cdot \varphi_{s}^{o}\right)(1)=\int_{G} \eta(g) \cdot \varphi_{s}^{o}(g) d g=\int_{B_{\varepsilon}} \varphi_{s}^{o}(g) d g
$$

Let $P^{+}$be the connected component of the identity in the minimal parabolic. The Jacobian of the map $P^{+} \times K \rightarrow G$ is non-vanishing at 1 , and $\varphi^{\circ}(1)=1$, so a suitable bound in terms of $\varepsilon$ on the spectral parameter $s \in \mathfrak{a}^{*} \otimes_{\mathbb{R}} \mathbb{C}$ will keep $\varphi_{s}^{o}(g)$ near 1 on $B_{\varepsilon}$. In the example of $S L_{n}(\mathbb{R})$, with

$$
\varphi_{s}^{o}\left(\begin{array}{ccc}
a_{1} & \cdots & * \\
0 & \ddots & \vdots \\
0 & 0 & a_{n}
\end{array}\right)=\left|a_{1}\right|^{s_{1}+\rho_{1}} \cdots\left|a_{n}\right|^{s_{n}+\rho_{n}} \quad \text { (for whatever normalizing constants } \rho_{j} \text { ) }
$$

bounds of the form $\left|s_{j}\right| \ll 1 / \varepsilon$ assure that $\operatorname{Re} \varphi_{s}^{o}(g) \geq \frac{1}{2}$ on $B_{\varepsilon}$, which prevents cancellation in the real part of $\varphi_{s}^{o}(g)$ for $g \in B_{\varepsilon}$, so

$$
\left|\lambda_{f}(\eta)\right|=\left|\int_{B_{\varepsilon}} \varphi_{s}^{o}(g) d g\right| \gg \int_{B_{\varepsilon}} \operatorname{Re} \varphi_{s}^{o}(g) d g \gg \int_{B_{\varepsilon}} \frac{1}{2} d g \gg \varepsilon^{\operatorname{dim} X}
$$

Combining the upper bound on $|q|_{L^{2}}^{2}$ with this lower bound on eigenvalues, letting $T \sim 1 / \varepsilon$,

$$
\left(\varepsilon^{\operatorname{dim} X}\right)^{2} \times\left(\sum_{\operatorname{cfm} F:\left|t_{F}\right| \ll T}\left|F_{\Theta \backslash H x}\right|^{2}+\ldots\right) \lll x, H \quad \varepsilon^{\operatorname{dim} X+\operatorname{dim} Y}
$$

proving the first assertion of the theorem. The proof of the second is essentially identical.

### 12.3 Global automorphic Sobolev spaces $H^{\ell}$

Again, let $G, \Gamma, K$ be any one of the archimedean examples, such as $S L_{2}(\mathbb{R}), S L_{2}(\mathbb{Z}), S O_{2}(\mathbb{R})$ or $S L_{n}(\mathbb{R}), S L_{n}(\mathbb{Z}), S O_{n}(\mathbb{R})$ or $S L_{n}(\mathbb{C}), S L_{n}(\mathbb{Z}[i]), S U_{n}(\mathbb{C})$. Let $\Delta$ be the invariant Laplacian on $\Gamma \backslash G / K$. Functions $f$ in $L^{2}(\Gamma \backslash G / K)$ decompose in an $L^{2}$ sense [1.14], [3.18]. To write the spectral expansion succinctly, let $\Xi$ be a locally compact, Hausdorff, $\sigma$-compact topological space parametrizing cuspforms, Eisenstein series appearing in the spectral decomposition and Plancherel, as well as their residues, with corresponding $\Delta$ eigenfunction $\Phi_{\xi} \in C^{\infty}(\Gamma \backslash G / K)$ for $\xi \in \Xi$, and a positive regular Borel measure $d \xi$ on $\Xi$, to write the expansions of [1.14] and [3.18] uniformly as

$$
f=\int_{\Xi}\left\langle f, \Phi_{\xi}\right\rangle \cdot \Phi_{\xi} d \xi \quad \quad\left(\text { for } f \in C_{c}^{\infty}(\Gamma \backslash G / K)\right)
$$

and Plancherel as

$$
|f|_{L^{2}}^{2}=\int_{\Xi}\left|\left\langle f, \Phi_{\xi}\right\rangle\right|^{2} d \xi \quad\left(\text { for } f \in C_{c}^{\infty}(\Gamma \backslash G / K)\right)
$$

For example, for $S L_{2}(\mathbb{Z})$, and similarly for the other three simplest examples, the explicit spectral expansion

$$
f=\sum_{\operatorname{cfm} F}\langle f, F\rangle \cdot F+\frac{\langle f, 1\rangle \cdot 1}{\langle 1,1\rangle}+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)}\left\langle f, E_{s}\right\rangle \cdot E_{s} d s
$$

for test functions $f$ would parametrize the cuspform components by an infinite discrete set, the constantfunction component by a further point, and the integrals-of-Eisenstein series component by $\mathbb{R}$ or by $[0,+\infty)$. For test functions $f$, the implied integrals $\left\langle f, E_{s}\right\rangle$ against Eisenstein series do converge absolutely.

Although many of the eigenfunctions $\Phi_{\xi}$ are not in $L^{2}(\Gamma \backslash G / K)$, they are all in $C^{\infty}(\Gamma \backslash G / K)$, and we can easily arrange that $\xi \rightarrow \Phi_{\xi}$ is a continuous $C^{\infty}(\Gamma \backslash G / K)$-valued function on $\Xi$. Integration of elements of $C^{\infty}(\Gamma \backslash G / K)$ against a fixed $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$ is a continuous linear functional on $C^{\infty}(\Gamma \backslash G / K)$, so

$$
\left.\xi \longrightarrow \Phi_{\xi} \longrightarrow\left\langle f, \Phi_{\xi}\right\rangle \quad \text { (for fixed } f \in C_{c}^{\infty}(\Gamma \backslash G / K)\right)
$$

is a continuous $\mathbb{C}$-valued function on $\Xi$, and thus has unambiguous pointwise values.
[12.3.1] Remark: For $\Xi_{1} \subset \Xi_{2} \subset$ a sequence of compact subsets of $\Xi$ whose union is $\Xi$, the integrals

$$
\left.\int_{\Xi_{n}}\left\langle f, \Phi_{\xi}\right\rangle \cdot \Phi_{\xi} \quad \quad \text { (for test functions } f\right)
$$

do exist as $C^{\infty}(\Gamma \backslash G / K)$-valued Gelfand-Pettis integrals. However, already the case with ordinary Fourier series on the circle, the fact that finite partial sums are invariably smooth tells little about the nature of the limit. Based on other examples, and on folklore, we imagine that for test functions $f$ the compactly-supported integrals should converge to $f$ in some topology finer than $L^{2}$, but this requires proof, as below.

The implied literal integrals $\left\langle f, \Phi_{\xi}\right\rangle$ against Eisenstein series do not necessarily converge for all $f$ in $L^{2}$, and certainly $L^{2}$ expansions do not reliably converge pointwise. Nevertheless, the Plancherel theorem asserts that the literal integrals $f \rightarrow\left(\xi \rightarrow\left\langle f, \Phi_{\xi}\right\rangle\right)$ on test functions do extend to an isometry

$$
\mathcal{F}: L^{2}(\Gamma \backslash G / K) \longrightarrow L^{2}(\Xi)
$$

Similarly, the spectral synthesis integrals $f=\int_{\Xi} c(\xi) \cdot \Phi_{\xi} d \xi$ do make literal sense (in fact, as $C^{\infty}(\Gamma \backslash G / K)$ valued integrals) for test functions $f$. Then the integrals and pairings for $f \in L^{2}$ are understood as extensions by continuity of the literal integrals. The (extensions of the) integrals $\left\langle f, \Phi_{\xi}\right\rangle$ are the spectral coefficients of $f$.
[12.3.2] Remark: Having said all that, just as is done with the Plancherel extension of the Fourier transform and Fourier inversion on $\mathbb{R}^{n}$, eventually we will write integrals and pairings which do not literally converge, but do exist as extensions by continuity of those integrals and pairings.
[12.3.3] Remark: Notably, Plancherel neither needs nor asserts anything directly about pointwise values of cuspforms or Eisenstein series or residues of Eisenstein series. This is fortunate, since already in the simplest case, various pointwise values of Eisenstein series $E_{s}$ for $S L_{2}(\mathbb{Z})$ are $\zeta_{k}(s) / \zeta(2 s)$ for complex quadratic extensions $k$ of $\mathbb{Q}[2 . \mathrm{C}]$, and sharp pointwise bounds on the critical line presumably include the Lindelöf Hypothesis.
[12.3.4] Claim: The eigenvalues $\lambda_{\xi}$ of $\Delta$ on $\Phi_{\xi}$ are real and non-positive.
Proof: For square-integrable eigenfunctions (such as cuspforms or square-integrable residues of Eisenstein series), these eigenvalues are real because a suitable restriction of $\Delta$ to a dense subspace of $L^{2}$ is symmetric and non-positive, so has a self-adjoint, non-positive Friedrichs extension.

In the four simple cases, the non- $L^{2}$ eigenfunctions entering the spectral expansion and Plancherel are Eisenstein series $E_{s}$ with eigenvalues $s(s-1)$ (up to real constants) and $\operatorname{Re}(s)=\frac{1}{2}$, and residues of $E_{s}$ with $s \in(0,1]$, with eigenvalue $s(s-1)$ (up to real constants). Somewhat more generally, [3.11.11] shows that cuspidal-data Eisenstein series $E_{s, f}^{P}$ for maximal proper parabolics $P$ in $G L_{n}$, with the cuspidal data $f=f_{1} \otimes f_{2}$ eigenfunctions for Casimir operators on the Levi component, are eigenfunctions for Casimir on $G L_{n}$. The explicit formula for eigenvalues shows that they are real and non-positive for $\operatorname{Re}(s)=\mathfrak{H}$, using the corresponding fact for the cuspidal data on the Levi components. The part of the spectral decomposition [3.16.1] and Plancherel [3.17.1], [3.17.4] corresponding to a given maximal proper parabolic uses $\operatorname{Re}(s)=\frac{1}{2}$ and residues (if any) on ( $\left.\frac{1}{2}, 1\right]$.
[12.3.5] Remark: Another proof that these eigenvalues are real would follow from a suitable description of the spectral decomposition of $L^{2}$ and Plancherel in terms of Hilbert integrals of representations of $G$. This would show that all representations appearing must be unitary, allowing various continuations proving the previous claim, using some form of a subquotient theorem.
[12.3.6] Claim: The invariant Laplacian $\Delta$ commutes with pointwise complex conjugation on functions on $G / K$ or on $\Gamma \backslash G / K$.

Proof: Any expression $\Omega=\sum_{i} x_{i} x_{i}^{*}$, with basis $\left\{x_{i}\right\}$ for $\mathfrak{g}$ and dual basis $\left\{x_{i}^{*}\right\}$, expresses Casimir $\Omega$ as a real-linear combination of compositions of operators of the form

$$
\left(x_{i} f\right)(g)=\left.\frac{\partial}{\partial t}\right|_{t=0} f\left(g \cdot e^{t x_{i}}\right)
$$

These visibly commute with pointwise complex conjugation.
[12.3.7] Claim: For $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$, the spectral coefficients of $\Delta f$ are

$$
\left\langle\Delta f, \Phi_{\xi}\right\rangle=\lambda_{\xi} \cdot\left\langle f, \Phi_{\xi}\right\rangle \quad\left(\text { for } f \in C_{c}^{\infty}(X)\right)
$$

Thus, for test functions,

$$
\Delta f=\int_{\Xi}\left\langle\Delta f, \Phi_{\xi}\right\rangle \Phi_{\xi} d \xi=\int_{\Xi} \lambda_{\xi} \cdot\left\langle f, \Phi_{\xi}\right\rangle \Phi_{\xi} d \xi
$$

That is, succinctly,

$$
\Delta f=\mathcal{F}^{-1} \mathcal{F} \Delta f=\mathcal{F}^{-1} \lambda_{\xi} \mathcal{F} f \quad \text { (for test functions } f \text { ) }
$$

Proof: For test functions, integration by parts is legitimate:

$$
\left\langle\Delta f, \Phi_{\xi}\right\rangle=\int_{X} f \Delta \bar{\Phi}_{\xi}=\overline{\lambda_{\xi}} \cdot \int_{X} f \bar{\Phi}_{\xi}=\lambda_{\xi} \cdot \int_{X} f \bar{\Phi}_{\xi} \quad\left(\text { for } f \in C_{c}^{\infty}(X)\right)
$$

as asserted.
[12.3.8] Corollary: The differential operator $\Delta$ differentiates spectral expansions of test functions termwise, in the sense of moving inside the integration and summation over $\Xi$ giving the spectral synthesis. ///

For $0 \leq \ell \in \mathbb{Z}$, the $\ell^{t h}$ Sobolev norm on $C_{c}^{\infty}(\Gamma \backslash G / K)$ is given by

$$
|f|_{H^{\ell}}^{2}=\int_{\Gamma \backslash G / K} \bar{f} \cdot(1-\Delta)^{\ell} f
$$

and

$$
H^{\ell}=H^{\ell}(\Gamma \backslash G / K)=\text { completion of } C_{c}^{\infty}(\Gamma \backslash G / K) \text { with respect to }|\cdot|_{H^{\ell}}
$$

[12.3.9] Claim: $|f|_{H^{k+1}} \geq|f|_{H^{k}}$ for test functions $f$, for $0 \leq k \in \mathbb{Z}$, and there is a canonical continuous injection $H^{k+1}(\Gamma \backslash G / K) \rightarrow H^{k}(\Gamma \backslash G / K)$ with dense image.
Proof: For all test functions $f,\langle-\Delta f, f\rangle \geq 0$. For a polynomial $P$ with non-negative real coefficients, we claim that $P(-\Delta)$ is non-negative on test functions, in the sense that for all test functions $f$

$$
\langle P(-\Delta) f, f\rangle \geq 0
$$

It suffices to prove this for monomials $(-\Delta)^{n}$. For even $n=2 m$,

$$
\left\langle(-\Delta)^{2 m} f, f\right\rangle=\left\langle(-\Delta)^{m} f,(-\Delta)^{m} f\right\rangle \geq 0
$$

For odd $n=2 m+1$,

$$
\left\langle(-\Delta)^{2 m+1} f, f\right\rangle=\left\langle(-\Delta)\left((-\Delta)^{m} f\right),\left((-\Delta)^{m} f\right)\right\rangle \geq 0
$$

For test functions $f$, the desired comparison is

$$
\begin{gathered}
|f|_{H^{k+1}}^{2}=\left\langle(1-\Delta)^{k+1} f, f\right\rangle_{L^{2}}=\left\langle(1+(-\Delta))^{k} f, f\right\rangle_{L^{2}}+\left\langle(1+(-\Delta))^{k}(-\Delta) f, f\right\rangle_{L^{2}} \\
\geq\left\langle(1+(-\Delta))^{k} f, f\right\rangle_{L^{2}}+0=|f|_{H^{k}}^{2}
\end{gathered}
$$

Thus, the identity map $C_{c}^{\infty}(\Gamma \backslash G / K) \rightarrow C_{c}^{\infty}(\Gamma \backslash G / K)$ extends to a continuous injection $H^{k+1} \rightarrow H^{k}$. Since $C_{c}^{\infty}$ is dense in both, necessarily the image is dense.
///
The following result is true by design.
[12.3.10] Claim: The differential operator $\Delta: C_{c}^{\infty}(\Gamma \backslash G / K) \longrightarrow C_{c}^{\infty}(\Gamma \backslash G / K)$ is continuous when the source is given the $H^{\ell+2}$ topology and the target is given the $H^{\ell}$ topology, for $0 \geq \ell \in \mathbb{Z}$.
Proof: Using the latter non-negativity propery of the previous proof,

$$
\begin{gathered}
|\Delta f|_{H^{\ell}}^{2}=\left\langle(1-\Delta)^{\ell}(\Delta f),(\Delta f)\right\rangle=\left\langle(-\Delta)^{2}(1+(-\Delta))^{\ell} f, f\right\rangle \\
\leq\left\langle(-\Delta)^{2}(1+(-\Delta))^{\ell} f, f\right\rangle+\langle(2(-\Delta)+1) f, f\rangle=\left\langle(1+(-\Delta))^{\ell+2} f, f\right\rangle=|f|_{H^{\ell+2}}^{2}
\end{gathered}
$$

as asserted.
[12.3.11] Corollary: $\Delta$ extends by continuity from test functions to a continuous linear map

$$
\Delta: H^{\ell+2}(\Gamma \backslash G / K) \quad \longrightarrow \quad H^{\ell}(\Gamma \backslash G / K) \quad(\text { for } 0 \leq \ell \in \mathbb{Z})
$$

Proof: That is, for test functions $\left\{f_{n}\right\}$ forming a Cauchy sequence in the $H^{\ell+2}$ topology, the continuity in the respective topologies on test functions means that the extension-by-continuity definition

$$
\Delta\left(H^{\ell+2}-\lim _{n} f_{n}\right)=H^{\ell}-\lim _{n} \Delta f_{n}
$$

is well defined and gives a continuous map in those topologies.
[12.3.12] Remark: This extension of $\Delta$ is $L^{2}$-differentiation for non-negative index Sobolev spaces. This extension is a refined version of distributional differentiation. Nevertheless, to examine global automorphic Sobolev spaces, the present discussion of $L^{2}$-differentiation does not directly depend on distributional notions.
[12.3.13] Corollary: For $f$ in $H^{\ell}$ with $\ell \geq 2$,

$$
\mathcal{F}(\Delta f)(\xi)=\lambda_{\xi} \cdot \mathcal{F} f
$$

Proof: Since $\mathcal{F}: L^{2}(\Gamma \backslash G / K) \rightarrow L^{2}(\Xi)$ is an isometric isomorphism obtained by extension by continuity from $\mathcal{F}$ on $C_{c}^{\infty}(\Gamma \backslash G / K)$, the literal integral computation for test functions

$$
(\mathcal{F} \Delta f)(\xi)=\int_{X} \Delta f \Phi_{\xi}=\int_{X} f \Delta \bar{\Phi}_{\xi}=\int_{X} f \lambda_{\xi} \cdot \bar{\Phi}_{\xi}=\lambda_{\xi} \int_{X} f \cdot \bar{\Phi}_{\xi}=\lambda_{\xi} \cdot(\mathcal{F} f)(\xi)
$$

extends by continuity to give the result.
[12.3.14] Corollary: Term-wise differentiation is valid:

$$
\Delta f=\mathcal{F}^{-1} \mathcal{F} \Delta f=\mathcal{F}^{-1} \lambda_{\xi} \mathcal{F} f \quad\left(\text { for } f \in H^{\ell} \text { with } \ell \geq 2\right)
$$

This differentiation is in the extended, non-classical sense.
Negative-index Sobolev spaces are not easily described via differential operators. Instead, characterize negative-index Sobolev spaces as Hilbert-space duals

$$
H^{-\ell}=H^{-\ell}(\Gamma \backslash G / K)=\text { Hilbert-space dual to } H^{\ell} \quad(\text { for } 0 \leq \ell \in \mathbb{Z})
$$

To identify $H^{0}=L^{2}$ with its own dual $\mathbb{C}$-linearly, combine the $\mathbb{C}$-conjugate-linear Riesz-Fréchet map $\Lambda: f \rightarrow\langle-, f\rangle$ with pointwise conjugation $c: f \rightarrow \bar{f}$, so that $\Lambda \circ c: H^{0} \rightarrow H^{0}$ is $\mathbb{C}$-linear. Mapping $\left(H^{0}\right)^{*} \rightarrow H^{-1}$ by the adjoint of the inclusion $H^{1} \rightarrow H^{0}$, and generally $H^{-k} \rightarrow H^{-k-1}$ by the adjoint of the inclusion $H^{k+1} \rightarrow H^{k}$, we have a chain of continuous linear maps

$$
\ldots \longrightarrow H^{2} \xrightarrow{\mathrm{inc}} H^{1} \xrightarrow{\mathrm{inc}} H^{0} \xrightarrow{\Lambda \circ c}\left(H^{0}\right)^{*} \xrightarrow{\mathrm{inc}^{*}} H^{-1} \xrightarrow{\mathrm{inc}^{*}} H^{-2} \xrightarrow{\mathrm{inc}^{*}} \ldots
$$

[12.3.15] Claim: For $0 \leq k \in \mathbb{Z}$, the maps $H^{-k} \rightarrow H^{-k-1}$ adjoint to inclusions $H^{k+1} \rightarrow H^{k}$ are themselves inclusions with dense images. Thus, we have a chain of continuous injections with dense images:

$$
\ldots \subset H^{2} \subset H^{1} \subset H^{0} \approx\left(H^{0}\right)^{*} \subset H^{-1} \subset H^{-2} \subset \ldots
$$

Proof: For $0 \leq k \in \mathbb{Z}$, since $H^{k+1} \rightarrow H^{k}$ has dense image, the adjoint $H^{-k} \rightarrow H^{-k-1}$ is injective. If $H^{-k} \rightarrow H^{-k-1}$ did not have dense image, then its adjoint would not be injective. By the reflexivity of Hilbert spaces, its adjoint is the original $H^{k+1} \rightarrow H^{k}$, which is injective.

In the sequel, we identify $H^{0}$ with its dual via $\Lambda \circ c$.
The continuous $L^{2}$-differentiation $\Delta: H^{2 \ell} \rightarrow H^{2 \ell-2}$ for $2 \ell \geq 2$ on positive-index Sobolev spaces gives an adjoint, still denoted $\Delta$, on negative-index spaces:

$$
\Delta: H^{-2 \ell} \longrightarrow H^{-2 \ell-2} \quad(\text { for } 0 \leq \ell \in \mathbb{Z})
$$

The extension of $\Delta$ to $\Delta: H^{1} \rightarrow H^{-1}$ can be characterized by

$$
\left.((1-\Delta) f)(F)=\langle f, \bar{F}\rangle_{H^{1}} \quad \quad \text { (for } f, F \in H^{1}\right)
$$

The compatibility of this extension with the others is best clarified by the spectral characterizations of the next section. Anticipating that clarification, let

$$
H^{\infty}(\Gamma \backslash G / K)=\bigcap_{k \in \mathbb{Z}} H^{k}(\Gamma \backslash G / K)=\lim _{k \in \mathbb{Z}} H^{k}(\Gamma \backslash G / K)
$$

and

$$
H^{-\infty}(\Gamma \backslash G / K)=\bigcup_{k \in \mathbb{Z}} H^{k}(\Gamma \backslash G / K)=\operatorname{colim}_{k \in \mathbb{Z}} H^{k}(\Gamma \backslash G / K)
$$

Pending the spectral characterization, for the following corollary we can temporarily take

$$
H^{\infty}=\lim _{k} H^{2 k} \quad H^{-\infty}=\operatorname{colim}_{k} H^{2 k}
$$

[12.3.16] Corollary: Both $H^{\infty}(\Gamma \backslash G / K)$ and $H^{-\infty}(\Gamma \backslash G / K)$ are stable under the extension of $\Delta$ (from test functions to global Sobolev spaces), and $\Delta$ gives a continuous linear operator on both.
Proof: For $H^{\infty}=\bigcap_{k} H^{2 k}$, the extended $\Delta$ maps $H^{2 k+2} \rightarrow H^{2 k}$, so $\Delta H^{\infty} \subset H^{2 k}$ for every $k$, and the intersection is $H^{\infty}$. More precisely, by the characterization of (projective) limits, the family of compatible maps

$$
H^{\infty} \xrightarrow{\text { inc }} H^{2 k+2} \xrightarrow{\Delta} H^{2 k}
$$

induces a unique compatible continuous (linear) map $H^{\infty} \rightarrow H^{\infty}$.
Oppositely, the extension-by-adjoint $\Delta$ maps $H^{-2 k} \rightarrow H^{-2 k-2}$, so the image of the ascending union is contained in the ascending union. More precisely, by the characterization of colimits, the compatible family of maps

$$
H^{-2 k} \xrightarrow{\Delta} H^{-2 k-2} \xrightarrow{\mathrm{inc}} H^{-\infty}
$$

gives a unique compatible continuous (linear) map $H^{-\infty} \rightarrow H^{-\infty}$.
[12.3.17] Claim: $H^{-\infty}(\Gamma \backslash G / K)$ is a subset of the space $C_{c}^{\infty}(\Gamma \backslash G / K)$ of distributions on $\Gamma \backslash G / K$.
Proof: First, check that the inclusion $C_{c}^{\infty}(\Gamma \backslash G / K) \subset H^{k}(\Gamma \backslash G / K)$ is continuous for every $k$. Recall [13.9], [6.3] that the space of test functions is the colimit of spaces

$$
C_{E}^{\infty}=\left\{f \in C_{c}^{\infty}: \operatorname{spt} f \subset E\right\}
$$

for compact $E \subset \Gamma \backslash G / K$. By the characterization of colimit, it suffices to prove continuity of each $C_{E}^{\infty} \rightarrow H^{k}$. Among the seminorms defining the topology on $C_{E}^{\infty}$ (see [6.4], [13.9]) are

$$
\nu_{k}(f)=\sup _{x \in E}\left|(1-\Delta)^{k} f(x)\right|
$$

Note that the volume of $\Gamma \backslash G / K$ is finite. Given $f \in C_{E}^{\infty}$, by Cauchy-Schwarz-Bunyakowsky,

$$
\begin{aligned}
|f|_{H^{k}}^{2}= & \int_{\Gamma \backslash G / K}(1-\Delta)^{k} f \cdot \bar{f}=\int_{E}(1-\Delta)^{k} f \cdot \bar{f} \leq\left(\int_{E}|(1-\Delta) f|^{2}\right)^{\frac{1}{2}} \cdot\left(\int_{E}|f|_{L^{2}}\right)^{\frac{1}{2}} \\
& \leq \nu_{k}(f) \cdot \operatorname{meas}(E)^{\frac{1}{2}} \cdot \nu_{0}(f) \cdot \operatorname{meas}(E)^{\frac{1}{2}} \leq \nu_{k}(f) \cdot \nu_{o}(f) \cdot \operatorname{meas}(\Gamma \backslash G / K)
\end{aligned}
$$

Thus, $|f|_{H^{k}}<_{E} \nu_{k}(f)+\nu_{0}(f)$, giving the desired continuity. Thus, $H^{-k}=\left(H^{k}\right)^{*}$ gives continuous linear functionals on $C_{c}^{\infty}$. Thus, the ascending union does so, as well.
[12.3.18] Claim: $H^{-\infty}(\Gamma \backslash G / K)$ is the dual of $H^{\infty}(\Gamma \backslash G / K)$.
Proof: This is an instance of the general fact that every continuous linear functional on a limit of Banach spaces factors through some limitand [13.14.4].

Although we will prove a stronger result in terms of global spectral expansions in the following section, we can reduce to local considerations to prove the following:
[12.3.19] Claim: For $n=\operatorname{dim}_{\mathbb{R}} \Gamma \backslash G / K$, for all $\ell>k+\frac{n}{2}$ and $\ell \in \mathbb{Z}$,

$$
H^{\ell}(\Gamma \backslash G / K) \subset C^{k}(\Gamma \backslash G / K)
$$

Proof: Smoothness is a local property, which allows reduction to a local version of Sobolev spaces. Namely, to show that $f \in H^{\ell}$ is $C^{k}$ at $g_{o}$, it suffices to show that $\eta f$ is in $C^{k}$, for $\eta$ a smooth cut-off function near $g_{o}$. We can take $\eta$ with sufficiently small support so that its compact support $E$ lies inside a small open subset diffeomorphic to a cube of the dimension of $\Gamma \backslash G / K$. Identifying opposite faces of the cube imbeds $E$ into a multi-torus $\mathbb{T}^{n}$. Further, it suffices to show a suitable Sobolev inequality for test functions $f$. Thus, we can apply the literal differential operator $\Delta$ to $f$.

For fixed $\eta$ and $E$, the Laplacian on $\Gamma \backslash G / K$ restricted to $E$ and the Laplacian on $\mathbb{T}^{n}$ are comparable on functions $\eta f$ for test functions $f$, giving constants $0<A_{\ell}, B_{\ell}<\infty$ such that

$$
A_{\ell} \cdot|\eta f|_{H^{k}\left(\mathbb{T}^{n}\right)} \leq|\eta f|_{H^{\ell}(\Gamma \backslash G / K)} \leq B_{k} \cdot|\eta f|_{H^{\ell}\left(\mathbb{T}^{n}\right)}
$$

for all test functions $f$, for $0 \leq \ell \in \mathbb{Z}$. By continuity, the same inequalities hold for $f \in H^{\ell}(\Gamma \backslash G / K)$, for every $k$. This reduces the problem to $H^{\ell}\left(\mathbb{T}^{n}\right) \subset C^{k}\left(\mathbb{T}^{n}\right)$, which we know from [9.5].
[12.3.20] Remark: The quotient $\Gamma \backslash G / K$ can fail to be a smooth manifold at various points, due to the possibility that the isotropy group $G_{x}$ of $x \in G / K$ can have non-trivial intersection with $\Gamma$. However, this is surmountable in various ways. For example, $G_{x}$ is compact, so $G_{x} \cap \Gamma$ is $f i n i t e$, and to examine smoothness of functions at $x$ we can harmlessly shrink $\Gamma$ by finite index to shrink $G_{x} \cap \Gamma$ to act trivially on $\Gamma \backslash G / K$, so that $\Gamma \backslash G / K$ is smooth near $x$. In fact, as earlier, we can identify $C^{\infty}(\Gamma \backslash G / K)$ with the right $K$-fixed vectors $C^{\infty}(\Gamma \backslash G)^{K}$ in $C^{\infty}(\Gamma \backslash G)$.
[12.3.21] Remark: The smooth cut-off device allows elementary local comparison of non-negative integer index Sobolev spaces on $\Gamma \backslash G / K$ with those on $\mathbb{T}^{n}$, but for non-integer index, as in the following section, such a comparison is less elementary.
[12.3.22] Corollary: $C^{\infty}(\Gamma \backslash G / K)^{*} \subset H^{-\infty}(\Gamma \backslash G / K) \subset C_{c}^{\infty}(\Gamma \backslash G / K)^{*}$.

### 12.4 Spectral characterization of Sobolev spaces $H^{s}$

By expressing Sobolev norm and differentiation via spectral transforms $\mathcal{F}$, for $0 \leq \ell \in \mathbb{Z}$ certainly $\mathcal{F} H^{\ell}$ is contained in $V^{\ell}$, where

$$
V^{s}=\left\{\text { measurable } v \text { on } \Xi:\left(1-\lambda_{\xi}\right)^{s / 2} \cdot v \in L^{2}(\Xi)\right\} \quad(\text { for } s \in \mathbb{R})
$$

For any $s \in \mathbb{R}$, give $V^{s}$ the Hilbert-space structure from the expected norm

$$
|v|_{V^{s}}^{2}=\int_{\Xi}\left(1-\lambda_{\xi}\right)^{s}|v|^{2}
$$

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For $s>t$, certainly there is a continuous inclusion $V^{s} \rightarrow V^{t}$ with dense image. The space $V^{-s}$ is naturally the Hilbert space dual $\left(V^{s}\right)^{*}$ of $V^{s}$, with $\mathbb{C}$-bilinear pairing given by integration

$$
\int_{\Xi} v(\xi) w(\xi) d \xi \quad \text { (complex-bilinear, for } v \in V^{s} \text { and } w \in V^{-s} \text { ) }
$$

The asymmetrical extension of the hermitian pairing $V^{0} \times V^{0} \rightarrow \mathbb{C}$ by $v \times w \rightarrow\langle v, w\rangle_{V^{0}}$ to a hermitian pairing on $V^{s} \times V^{-s}$ is

$$
\langle v, w\rangle_{V^{s} \times V^{-s}}=\int_{\Xi} v(\xi) \overline{w(\xi)} d \xi
$$

[12.4.1] Claim: The spectral transform $\mathcal{F}: C_{c}^{\infty}(\Gamma \backslash G / K) \rightarrow L^{2}(\Xi)$ induces a Hilbert space isomorphism $\mathcal{F}: H^{2 \ell} \xrightarrow{\approx} V^{2 \ell}$ for all $0 \leq \ell \in \mathbb{Z}$, and we have commuting rectangles


Proof: Plancherel asserts that $\mathcal{F}: H^{0} \rightarrow V^{0}$ is an isomorphism. By design, $1-\Delta$ gives an isomorphism $H^{2 k} \rightarrow H^{2 k-2}$ for all $k \in \mathbb{Z}$. Even more directly, multiplication by $1-\lambda_{\xi}$ gives an isomorphism $V^{s} \rightarrow V^{s-2}$ for all $s \in \mathbb{R}$. Corollary [12.3.12] shows that $\mathcal{F}$ intertwines $1-\Delta$ and multiplication by $1-\lambda_{\xi}$ on Sobolev spaces with positive index, and dualization gives the same result on negative-index spaces. ///

This allows us to define an isomorphism $\mathcal{F}: H^{-2 k} \rightarrow V^{-2 k}$ for $0>-k \in \mathbb{Z}$ as the adjoint to the isomorphism $\mathcal{F}^{-1}: V^{2 k} \rightarrow H^{2 k}$. Proof of the analogous assertion for odd-index $H^{k}$ and $V^{k}$ is slightly complicated by the fact that the Sobolev space $H^{k}$ does not have an elementary isomorphism to $H^{0}$, so Plancherel does not immediately resolve the issue.
[12.4.2] Claim: The spectral transform $\mathcal{F}: C_{c}^{\infty}(\Gamma \backslash G / K) \rightarrow L^{2}(\Xi)$ induces Hilbert space isomorphisms $\mathcal{F}: H^{2 k+1} \rightarrow V^{2 k+1}$ for all $0 \leq k \in \mathbb{Z}$, and we have commuting rectangles


Proof: In the commuting rectangle

the horizontal maps are injections with dense images, and the left vertical map is an isomorphism, from the previous. Thus, $\mathcal{F}\left(H^{1}\right)$ is dense in $V^{1}$. Since $\mathcal{F}$ intertwines $1-\Delta$ with multiplication by $1-\lambda_{\xi}$, the spectral $\operatorname{map} \mathcal{F}: H^{1} \rightarrow V^{1}$ is an isometry to its image. Since $H^{1}$ is complete, the image is closed. A dense, closed subspace of $V^{1}$ is the whole $V^{1}$.

As in the even-index case, by design, $(1-\Delta): H^{2 k+1} \rightarrow H^{2 k-1}$ is an isomorphism, and multiplication by $1-\lambda_{\xi}$ is an isomorphism $V^{2 k+1} \rightarrow V^{2 k-1}$. The intertwining of these two operators by $\mathcal{F}$, by [12.3.12], gives the commutativity.

Thus, we can define isomorphisms $\mathcal{F}: H^{-1-2 k} \rightarrow V^{-1-2 k}$ as the inverses of the (isomorphism) adjoints $\mathcal{F}^{*}: V^{-1-2 k} \rightarrow H^{-1-2 k}$ to (the isomorphism) $\mathcal{F}: H^{2 k+1} \rightarrow V^{2 k+1}$. Unlike the even-index case, we need
[12.4.3] Lemma: We have a commutative rectangle with all maps isomorphisms:


Proof: Again, the horizontal maps are (isometric) isomorphisms by design, the left vertical map is an (isometric) isomorphism from above, and the right vertical map is the inverse of the adjoint of the left vertical map. Then the commutativity is immediate: composing $\mathcal{F}^{*} \circ\left(1-\lambda_{\xi}\right) \circ \mathcal{F}$ around three sides, for $f, F \in C_{c}^{\infty}(\Gamma \backslash G / K)$, the characterization of adjoints and Plancherel,

$$
\begin{gathered}
\left\langle\left(\mathcal{F}^{*} \circ\left(1-\lambda_{\xi}\right) \circ \mathcal{F}\right) f, F\right\rangle_{H^{-1} \times H^{1}}=\left\langle\left(\left(1-\lambda_{\xi}\right) \circ \mathcal{F}\right) f, \mathcal{F} F\right\rangle_{V^{-1} \times V^{1}}=\langle\mathcal{F}(1-\Delta) f, \mathcal{F} F\rangle_{V^{-1} \times V^{1}} \\
=\langle(1-\Delta) f, F\rangle_{H^{-1} \times H^{1}}=\langle(1-\Delta) f, F\rangle_{H^{-1} \times H^{1}}
\end{gathered}
$$

since the asymettrical pairings are the extensions by continuity of the $L^{2}$ pairing restricted to test functions. Since $H^{-1}=\left(H^{1}\right)^{*}$, this gives the assertion.

Thus, as in the even-index case, we have
[12.4.4] Corollary: For all $k \in \mathbb{Z}$, we have commuting rectangles


That is, for $k \in \mathbb{Z}$, the spaces $V^{k}$ and the multiplication operator $1-\lambda_{\xi}$ are a faithful spectral-side mirror of the spaces $H^{k}$ and operator $1-\Delta$. In fact, on the spectral side, greater flexibility is afforded by the spaces $V^{s}$ for $s \in \mathbb{R}$. Given the compatibility just proven, we can define Sobolev norms

$$
|f|_{H^{s}}=|\mathcal{F} f|_{V^{s}} \quad\left(\text { for } s \in \mathbb{R} \text { and } f \in C_{c}^{\infty}(\Gamma \backslash G / K)\right)
$$

and

$$
H^{s}=H^{s}(\Gamma \backslash G / K)=\text { completion of } C_{c}^{\infty}(\Gamma \backslash G / K) \text { with respect to }|\cdot|_{H^{s}} \quad(\text { for } s \in \mathbb{R})
$$

As an application of the pre-trace formulas of [12.1] and [12.2]:
[12.4.5] Claim: With $n$ the dimension of $\Gamma \backslash G / K$, we have $\int_{\Xi}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \cdot\left(1-\lambda_{\xi}\right)^{-\frac{n}{2}-\varepsilon}<\infty$, with a bound depending uniformly on $z_{o}$ in compacts in $\Gamma \backslash G / K$.
Proof: In current notation, the pre-trace formulas assert that, for $z_{o}$ in a fixed compact $C \subset \Gamma \backslash G / K$,

$$
\int_{\xi:\left|\lambda_{\xi}\right| \leq T^{2}}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \quad<_{C} \quad T^{n} \quad(\text { as } T \rightarrow+\infty)
$$

or

$$
\int_{\xi:\left|\lambda_{\xi}\right| \leq T}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \ll_{C} \quad T^{\frac{n}{2}} \quad(\text { as } T \rightarrow+\infty)
$$

Thus, summing by parts,

$$
\begin{gathered}
\int_{\Xi}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \cdot\left(1-\lambda_{\xi}\right)^{-\frac{n}{2}-\varepsilon}=\sum_{\ell=1}^{\infty} \int_{\xi: \ell-1 \leq\left|\lambda_{\xi}\right|<\ell}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \cdot\left(1-\lambda_{\xi}\right)^{-\frac{n}{2}-\varepsilon} \\
<_{C} \sum_{\ell=1}^{\infty} \int_{\xi:\left|\lambda_{\xi}\right|<\ell}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \cdot\left((1+\ell)^{-\frac{n}{2}-\varepsilon}-(1+(\ell+1))^{-\frac{n}{2}-\varepsilon}\right)<_{C} \sum_{\ell=1}^{\infty} \int_{\xi:\left|\lambda_{\xi}\right|<\ell}\left|\Phi_{\xi}\left(z_{o}\right)\right|^{2} \cdot(1+\ell)^{-\frac{n}{2}-\varepsilon-1} \\
<_{C} \sum_{\ell=1}^{\infty} \ell^{\frac{n}{2}} \cdot(1+\ell)^{-\frac{n}{2}-\varepsilon-1}<\infty
\end{gathered}
$$

as claimed.
Still $n$ is the dimension of $\Gamma \backslash G / K$. Now we can prove
[12.4.6] Theorem: For $s>\frac{n}{2}$, for $f \in C_{c}^{\infty}(\Gamma \backslash G / K)$, and for compact $C \subset \Gamma \backslash G / K$, we have a global Sobolev inequality

$$
\sup _{z \in C}|f(z)|<_{C, s}|f|_{H^{s}}
$$

Thus, $H^{s}(\Gamma \backslash G / K) \subset C^{o}(\Gamma \backslash G / K)$. Further, with $\Xi_{\ell}=\left\{\xi \in \Xi:\left|\lambda_{\xi}\right| \leq \ell\right\}$,

$$
\lim _{\ell} \int_{\Xi_{\ell}}\left\langle f, \Phi_{\xi}\right\rangle \cdot \Phi_{\xi} d \xi=f \quad \quad \text { (in the } C^{o} \text { topology) }
$$

[12.4.7] Remark: That is, for $s>\frac{n}{2}, H^{s}(\Gamma \backslash G / K)$ is an improved version of $C^{o}(\Gamma \backslash G / K)$, in the sense that this $H^{s}$ not only consists of continuous functions, but also the spectral expansion of every $f \in H^{s}$ converges uniformly pointwise to $f$ on compacts. In contrast, already on the circle $\mathbb{T}$, spectral expansions (Fourier series) of continuous functions need not converge pointwise, much less uniformly so.
Proof: For $0 \leq \ell<+\infty, \Xi_{\ell}$ is compact. Every $\Phi_{\xi}$ is smooth, and $\xi \rightarrow \Phi_{\xi}$ is a continuous $C^{\infty}(\Gamma \backslash G / K)$-valued function on $\Xi$. Thus, for a test function $f$,

$$
f_{\ell}=\int_{\Xi_{\ell}}\left\langle f, \Phi_{\xi}\right\rangle \cdot \Phi_{\xi}(z) d \xi
$$

exists as a $C^{\infty}(\Gamma \backslash G / K)$-valued Gelfand-Pettis integral, so is certainly continuous. By the spectral characterization of $H^{s}$, the sequence $\left\{f_{\ell}\right\}$ approaches $f$ in the $H^{s}$ topology.

For $f \in H^{s}$, by Cauchy-Schwarz-Bunyakowsky, for $z_{o} \in C$,

$$
\begin{gathered}
\left|\int_{\Xi} \mathcal{F} f(\xi) \cdot \Phi\left(z_{o}\right) d \xi\right|=\left|\int_{\Xi} \mathcal{F} f(\xi)\left(1-\lambda_{\xi}\right)^{s / 2} \cdot\left(1-\lambda_{\xi}\right)^{-s / 2} \Phi\left(z_{o}\right) d \xi\right| \\
\leq\left(\int_{\Xi}|\mathcal{F} f(\xi)|^{2}\left(1-\lambda_{\xi}\right)^{s} d \xi\right)^{\frac{1}{2}} \cdot\left(\int_{\Xi}\left(1-\lambda_{\xi}\right)^{-s}\left|\Phi\left(z_{o}\right)\right|^{2} d \xi\right)^{\frac{1}{2}}=|f|_{H^{s}} \cdot\left(\int_{\Xi}\left(1-\lambda_{\xi}\right)^{-s}\left|\Phi\left(z_{o}\right)\right|^{2} d \xi\right)^{\frac{1}{2}}
\end{gathered}
$$

The previous claim shows that the latter integral is finite for $s>\frac{n}{2}$, with a bound uniform in $z_{o} \in C$. That is, the sup norm on $C$ of the pointwise function $z \rightarrow \int_{\Xi} \mathcal{F} f(\xi) \cdot \Phi_{\xi}(z) d \xi$ exists, and is dominated by $|f|_{H^{s}}$.

Thus, for test function $f$, the $H^{s}$ convergence of the continuous function $f_{\ell}$ to $f$ implies $C^{o}$ convergence $f_{\ell} \rightarrow f$. That is, the spectral expansion of $f$ converges pointwise to $f$, uniformly on compacts. Extending by continuity, the same result follows for $f \in H^{s}$.

Since test functions are dense in $H^{s}$ and $H^{s}$ convergence implies $C^{o}$ convergence, $H^{s} \subset C^{o}$, and the spectral expansion of every function $f$ in $H^{s}$ converges pointwise to $f$, uniformly on compacts.
[12.4.8] Corollary: The automorphic Dirac $\delta_{z_{o}}$ at $z_{o} \in \Gamma \backslash G / K$ is in $H^{-s}$ for every $s>\frac{n}{2}$, and has spectral expansion

$$
\delta_{z_{o}}=\int_{\Xi} \bar{\Phi}_{\xi}\left(z_{o}\right) \cdot \Phi_{\xi} d \xi \quad \quad\left(\text { convergent in } H^{-s}\right)
$$

[12.4.9] Remark: Unsurprisingly, the indicated integral does not converge pointwise, but there is no claim that it does so.
Proof: Fix $s>\frac{n}{2}, z_{o} \in \Gamma \backslash G / K$, and take $f \in H^{s}$. By the theorem, $H^{s} \subset C^{o}$, so $f \rightarrow f\left(z_{o}\right)=\delta_{z_{o}}(f)$ is a continuous linear functional on $H^{s}$, so is in $H^{-s}$. To determine its spectral coefficients, consider

$$
\delta_{z_{o}}(f)=f\left(z_{o}\right)=\int_{\Xi} \mathcal{F} f(\xi) \cdot \Phi_{\xi}\left(z_{o}\right) d \xi
$$

The claim shows that $\varphi(\xi)=\Phi_{\xi}\left(z_{o}\right)$ is in $V^{-s}$, so this integral is the complex bilinear pairing on $V^{s} \times V^{-s}$ applied to $\mathcal{F} f$ and $\varphi$. By uniqueness of spectral expansions, $\mathcal{F} \delta_{z_{o}}=\varphi$.

### 12.5 Continuation of solutions of differential equations

Given $f \in H^{-\infty}(\Gamma \backslash G / K)$ and $\lambda \in \mathbb{C}$, we want to solve

$$
(\Delta-\lambda) u=f
$$

for $u \in H^{-\infty}(\Gamma \backslash G / K)$, when possible. Here $\Delta$ is the extension of the invariant Laplacian from test functions to $H^{-\infty}$. Applying $\mathcal{F}$ to both sides gives

$$
\mathcal{F} f=\mathcal{F}(\Delta-\lambda) u=\left(\lambda_{\xi}-\lambda\right) \mathcal{F} u
$$

and both $\mathcal{F} f$ and $\mathcal{F} u$ are in weighted $L^{2}$ spaces on $\Xi$. For $\lambda \notin(-\infty, 0]$, the function $\lambda_{\xi}-\lambda$ is bounded away from 0 on $\Xi$, so we can simply divide to obtain

$$
\mathcal{F} u=\frac{\mathcal{F} f}{\lambda_{\xi}-\lambda}
$$

That is,

$$
u=\int_{\Xi} \frac{\mathcal{F} f(\xi)}{\lambda_{\xi}-\lambda} \cdot \Phi_{\xi} d \xi \quad \quad\left(\text { converging in } H^{-\infty}\right)
$$

[12.5.1] Claim: For $\lambda \notin(-\infty, 0]$, the previous solution $u$ to the differential equation $(\Delta-\lambda) u=f$ is the unique solution in $H^{-\infty}$. That is, the corresponding homogeneous equation has no solutions in $H^{-\infty}$.
Proof: The difference $v \in H^{-\infty}$ between the previous solution and any other would be a solution to the homogeneous equation $(\Delta-\lambda) v=0$. Since $v \in H^{-\infty}$, it has a spectral expansion $v=\int_{\Xi} \mathcal{F} v(\xi) \cdot \Phi_{\xi} d \xi$, and the differential equation gives $\left(\lambda_{\xi}-\lambda\right) \mathcal{F} v=0$. Since $\lambda_{\xi}-\lambda \neq 0$ on $\Xi$, this requires that $\mathcal{F} v=0$ almost everywhere, so $v$ is 0 in $H^{-\infty}$.

For $\lambda \in(-\infty, 0]$ there is potential interaction with the eigenvalues $\lambda_{\xi}$ and eigenfunctions $\Phi_{\xi}$. For simplicity, we consider only the simplest example of $S L_{2}(\mathbb{Z})$. The other three simple examples admit nearly identical treatment, with uninteresting minor complications due to normalizations of constants. The spectral expansion is

$$
f=\sum_{\operatorname{cfm} F}\langle f, F\rangle \cdot F+\frac{\langle f, 1\rangle \cdot 1}{\langle 1,1\rangle}+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)}\left\langle f, E_{s}\right\rangle \cdot E_{s} d s \quad \quad \text { (convergent in } H^{-\infty} \text { ) }
$$

where the indicated pairings and integrals are extensions by continuity of the literal pairings and integrals, as above, and $F$ runs through a orthonormal basis for cuspforms, consisting of strong-sense cuspforms.
[12.5.2] Claim: $f \in H^{r}$ if and only if the discrete-spectrum part

$$
\sum_{\operatorname{cfm} F}\langle f, F\rangle \cdot F+\frac{\langle f, 1\rangle \cdot 1}{\langle 1,1\rangle}
$$

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and the continuous-spectrum part

$$
\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)}\left\langle f, E_{s}\right\rangle \cdot E_{s} d s
$$

are both in $H^{r}$, individually.
Proof: Use the spectral characterization.
Use notation $\lambda_{s}=s(s-1)$, and for cuspform eigenfunction $F$ let $s_{F} \in \mathbb{C}$ be such that $\Delta F=\lambda_{s_{F}} \cdot F$. For $\operatorname{Re}(w)>\frac{1}{2}$ and $w \neq 1$, by division, the equation $\left(\Delta-\lambda_{w}\right) u=f$ has solution

$$
u=\sum_{\operatorname{cfm} F} \frac{\langle f, F\rangle \cdot F}{\lambda_{s_{F}}-\lambda_{w}}+\frac{\langle f, 1\rangle \cdot 1}{\left(\lambda_{1}-\lambda_{w}\right) \cdot\langle 1,1\rangle}+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}}{\lambda_{s}-\lambda_{w}} d s
$$

This spectral expansion converges at least in $H^{-\infty}$. For $f \in H^{r}$ with $r \in \mathbb{R}$, the spectral characterization shows that $u \in H^{r+2}$, and the spectral expansion converges in $H^{r+2}$.

To examine possible solutions for all $w \in \mathbb{C}$, it is useful to consider a solution $u=u_{w}$ to $\left(\Delta-\lambda_{w}\right) u=f$ as a holomorphic or meromorphic function-valued function of $w$. By the spectral characterization of $H^{r}$ and $H^{r+2}$, the cuspidal component

$$
\sum_{\operatorname{cfm} F} \frac{\langle f, F\rangle \cdot F}{\lambda_{s_{F}}-\lambda_{w}}
$$

is visibly a meromorphic $H^{r+2}$-valued function of $w \in \mathbb{C}$, with poles at most at $w=s_{F}$ : the decomposition [7.1] of cuspforms shows that the multiplicity of $s_{F}$ is finite, and that the points $s_{F}$ are discrete in $C$. The constant component is similar.

The continuous spectrum component

$$
\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}}{\lambda_{s}-\lambda_{w}} d s
$$

is subtler, and does not generally meromorphically continue as an $H^{r+2}$-valued function, but only in a broader sense, as follows.
[12.5.3] Claim: For $\lambda_{w} \leq-1 / 4$, if $\left(\Delta-\lambda_{w}\right) u=f$ has a solution $u$ in $H^{-\infty}$, then $\left\langle f, E_{w}\right\rangle=0$, in the strong sense that $\left\langle f, E_{s}\right\rangle /\left(\lambda_{s}-\lambda_{w}\right)$ is locally integrable near $s=w$.
Proof: The continuous-spectrum part of the spectral transform of the differential equation gives

$$
\left(\lambda_{s}-\lambda_{w}\right)\left\langle u, E_{s}\right\rangle=\left\langle f, E_{s}\right\rangle \quad \text { (almost everywhere in } s \text { with } \operatorname{Re}(s)=\frac{1}{2} \text { ) }
$$

where the pairings are extensions by continuity of the literal integrals, and are at least locally integrable functions. Since $\lambda_{s}-\lambda_{w}=(s-w)(s-1+w)$, we have the indicated vanishing.
[12.5.4] Theorem: Let $X$ be a quasi-complete, locally convex topological vector space containing both $H^{r+2}$ and Eisenstein series. The function

$$
u_{w}=\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}}{\lambda_{s}-\lambda_{w}} d s \quad \quad \text { (convergent in } H^{r+2} \text { ) }
$$

has a meromorphic continuation as $X$-valued function of $w$, with functional equation

$$
u_{w}=u_{1-w}-\frac{\left\langle f, E_{w}\right\rangle \cdot E_{w}}{2 w-1}
$$

[12.5.5] Remark: Thus, although $u_{1-w}$ is in $H^{r+2}$ for $\operatorname{Re}(w)<\frac{1}{2}$, the extra term is not in $H^{-\infty}$ unless it is 0 , that is, unless $\left\langle f, E_{w}\right\rangle=0$.
[12.5.6] Remark: Despite the seeming symmetry of the spectral integral for $u_{w}$ under $w \rightarrow 1-w$, there is no such symmetry.

Proof: This begins with a natural regularization:

$$
\begin{gathered}
u_{w}=\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s \\
=\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s+\left\langle f, E_{w}\right\rangle E_{w} \frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{d s}{\lambda_{s}-\lambda_{w}} \\
=\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s-\frac{\left\langle f, E_{w}\right\rangle E_{w}}{2(2 w-1)}
\end{gathered}
$$

by residues.
The integral appears to be better behaved near $s=w$, but since it is not necessarily a literal integral the appearance is potentially misleading. With $t=\operatorname{Im}(s)$, rewrite

$$
\int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s=\int_{|t|>T} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s+\int_{|t| \leq T} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s
$$

The first integral is

$$
\int_{|t|>T} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}}{\lambda_{s}-\lambda_{w}} d s-\left\langle f, E_{w}\right\rangle E_{w} \cdot \int_{|t|>T} \frac{d s}{\lambda_{s}-\lambda_{w}} d s
$$

The latter integral on $|t|>T$ is $H^{r+2}$-valued. The extra term is meromorphic in $w$, but takes values in some function space adequate to contain Eisenstein series. Of course, if $\left\langle f, E_{w}\right\rangle=0$, then that extra term disappears.

The second integral

$$
\int_{|t| \leq T} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s
$$

is compactly supported, and is a holomorphic $X$-valued function of two complex variables $s, w$ away from the diagonal $s=w$. By design, there is cancellation on the diagonal, as is visible from a vector-valued power series expansion [15.8]. Thus, the integrand is a holomorphic $X$-valued function of $s, w$.

Let $\operatorname{Hol}(\Omega, X)$ be the space of holomorphic $X$-valued functions on a region $\Omega$, with seminorms

$$
\nu_{\mu}(f)=\sup _{w \in K} \mu(f(w))
$$

where $K \subset \Omega$ is compact, and the topology on $X$ is given by seminorms $\mu$. With this topology, $\operatorname{Hol}(\Omega, X)$ is quasi-complete and locally convex [15.3.2]. By [15.3.3] and [15.3.4], for a complex-analytic $X$-valued function $f(s, w)$ in two variables, on a domain $\Omega_{1} \times \Omega_{2} \subset \mathbb{C}^{2}$, function $\left.s \rightarrow(w \rightarrow f(z, w))\right)$ is a holomorphic $\operatorname{Hol}\left(\Omega_{1}, X\right)$-valued function on $\Omega_{2}$.

Thus, letting $\Omega$ be an appropriate bounded open containing the set where $|t| \leq T$, the integrand in the integral over $|t| \leq T$ is a compactly-supported, continuous, $\operatorname{Hol}(\Omega, X)$-valued function of $s$, and has a Gelfand-Pettis integral in $\operatorname{Hol}(\Omega, X)$. That is, it has a meromorphic continuation as an $X$-valued function of $w$.

To obtain the functional equation of $u_{w}$, from the first part of the proof, at first for $\operatorname{Re}(w)>\frac{1}{2}$,

$$
u_{w}=\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}-\left\langle f, E_{w}\right\rangle \cdot E_{w}}{\lambda_{s}-\lambda_{w}} d s-\frac{\left\langle f, E_{w}\right\rangle E_{w}}{2(2 w-1)}
$$

and then this holds by meromorphic continuation. Now take $\operatorname{Re}(w)<\frac{1}{2}$, and bring the regularizing term back out, producing

$$
\begin{aligned}
u_{w}= & \frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}}{\lambda_{s}-\lambda_{w}} d s-\left\langle f, E_{w}\right\rangle \cdot E_{w} \cdot \frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{d s}{\lambda_{s}-\lambda_{w}}-\frac{\left\langle f, E_{w}\right\rangle E_{w}}{2(2 w-1)} \\
& =\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \frac{\left\langle f, E_{s}\right\rangle \cdot E_{s}}{\lambda_{s}-\lambda_{w}} d s-\frac{\left\langle f, E_{w}\right\rangle E_{w}}{(2 w-1)} \quad\left(\text { for } \operatorname{Re}(w)<\frac{1}{2}\right)
\end{aligned}
$$

That is, the two extra terms do not cancel, but reinforce. For $\operatorname{Re}(w)<\frac{1}{2}, \operatorname{Re}(1-w)>\frac{1}{2}$, so the latter integral is $u_{1-w}$. That is,

$$
u_{w}=u_{1-w}-\frac{\left\langle f, E_{w}\right\rangle E_{w}}{(2 w-1)}
$$

at first for $\operatorname{Re}(w)<\frac{1}{2}$, but then for all $w$ (away from poles), by the identity principle.
[12.5.7] Remark: One simple choice for a topological vector space containing both suitable global automorphic Sobolev spaces and Eisenstein series is as follows. A preliminary choice of topological vector space $\mathfrak{E}$ containing Eisenstein series is needed. An easy choice is the Fréchet space $\mathfrak{E}=C^{\infty}(\Gamma \backslash G / K)$. Others are spaces of moderate-growth continuous functions or moderate-growth smooth functions. The essential point is that $\mathcal{D}=C_{c}^{\infty}(\Gamma \backslash G / K)$ should be dense in $\mathfrak{E}$. Since $\mathcal{D}$ is dense in $H^{r}$, we could hope for a topological vector space $X$ fitting into a pushout diagram ${ }^{[73]}$

meaning that, for every topological vector space $Y$ fitting into a commutative diagram

there is a unique $X \rightarrow Y$ making a commutative diagram


For the usual diagrammatic reasons, there is at most one such $X$, up to unique isomorphism. To prove existence, as with colimits expressed as quotients of coproducts in [13.8], the pushout is a natural quotient of the coproduct:

$$
X=\text { quotient of } H^{r} \oplus \mathfrak{E} \text { by the closure of the anti-diagonal copy } \mathcal{D}^{-\Delta}=\{(\varphi,-\varphi): \varphi \in \mathcal{D}\} \text { of } \mathcal{D}
$$

[^21]
### 12.6 Example: automorphic Green's functions

Continue the situation of the previous section. From [12.4.8], the automorphic Dirac $\delta_{z}$ at $z \in \Gamma \backslash G / K$ is in $H^{-1-\varepsilon}$ for every $\varepsilon>0$. By the previous section, the equation $\left(\Delta-\lambda_{w}\right) u_{w}=\delta_{z}$ has a solution $u_{w}=u_{w, z}$ in $H^{1-\varepsilon}$ for $\operatorname{Re}(w)>\frac{1}{2}$ and $w \neq 1$, with a meromorphic continuation in a topological vector space $X$ large enough to include both $H^{1-\varepsilon}$ and Eisenstein series.

A traditional notation (slightly incompatible with ours) is

$$
G_{w}\left(z, z^{\prime}\right)=u_{w, z}\left(z^{\prime}\right) \quad\left(\text { for } z, z^{\prime} \in \Gamma \backslash G / K\right)
$$

This is often called a Green's function. In traditional (partly heuristic) notation, the differential equation $\left(\Delta-\lambda_{w}\right) u_{w, z}=\delta_{z}$ would be written as

$$
\Delta_{z^{\prime}} G_{w}\left(z, z^{\prime}\right)=\lambda_{w} \cdot \delta_{z}\left(z^{\prime}\right) \quad \text { or, equivalently, } \quad \Delta_{z} G_{w}\left(z, z^{\prime}\right)=\lambda_{w} \cdot \delta_{z^{\prime}}(z)
$$

In contexts where distributions are not acknowledged, the description of $G_{w}\left(z, z^{\prime}\right)$ is considerably more awkward.
[12.6.1] Theorem: Use coordinates $z=x+i y$ on $\mathfrak{H} \approx G / K$. For $\operatorname{Re}(w)>\frac{1}{2}$ and $w \notin\left(\frac{1}{2}, 1\right]$, and for $a \geq \operatorname{Im}(z)$, the solution $u_{w}=u_{w, z}$ in $H^{1-\varepsilon}$ of the equation $\left(\Delta-\lambda_{w}\right) u_{w}=\delta_{z}$ has constant term

$$
c_{P} u_{w}(i a)=\int_{0}^{1} u_{w}(x+i a) d x=a^{1-w} \cdot \frac{E_{w}(z)}{1-2 w}
$$

Proof: From [12.2], since the orbits of $(N \cap \Gamma) \backslash N$ are compact and codimension 1, the distribution

$$
\eta_{a} f=c_{P} f(i a)=\int_{0}^{1} f(x+i a) d x
$$

is in $H^{-\frac{1}{2}-\varepsilon}$ for every $\varepsilon>0$, and, as was argued for $\delta_{z}$ in [12.4.8], has a corresponding spectral expansion

$$
\eta_{a}=\frac{\eta_{a}(1) \cdot 1}{\langle 1,1\rangle}+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \eta_{a} E_{1-s} \cdot E_{s} d s
$$

Since $u_{w} \in H^{1-\varepsilon}$ for every $\varepsilon>0, \eta_{a}$ gives a continuous linear functional on a Sobolev space containing it, and by the extended asymmetrical form of Plancherel,

$$
\begin{aligned}
\eta_{a}\left(u_{w}\right)= & \int_{\Xi} \mathcal{F} \eta_{a} \cdot \mathcal{F} u_{w}=\frac{\eta_{a}(1) \cdot \delta_{z}(1)}{\left(\lambda_{1}-\lambda_{w}\right) \cdot\langle 1,1\rangle}+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)} \eta_{a}\left(E_{1-s}\right) \cdot \frac{\delta_{z}\left(E_{s}\right)}{\lambda_{s}-\lambda_{w}} d s \\
& =\frac{1}{\left(\lambda_{1}-\lambda_{w}\right) \cdot\langle 1,1\rangle}+\frac{1}{4 \pi i} \int_{\left(\frac{1}{2}\right)}\left(a^{1-s}+c_{1-s} a^{s}\right) \cdot \frac{E_{s}(z)}{\lambda_{s}-\lambda_{w}} d s
\end{aligned}
$$

from the computation of the constant term [1.9.4]. Using the functional equation $c_{1-s} E_{s}=E_{1-s}$ and then replacing $s$ by $1-s$, the integral of $c_{1-s} a^{s} \cdot E_{s}(z) /\left(\lambda_{s}-\lambda_{w}\right)$ just produces another copy of the integral of $a^{1-s} E_{s}(z) /\left(\lambda_{s}-\lambda_{w}\right)$, and

$$
\eta_{a}\left(u_{w}\right)=\frac{1}{\left(\lambda_{1}-\lambda_{w}\right) \cdot\langle 1,1\rangle}+\frac{1}{2 \pi i} \int_{\left(\frac{1}{2}\right)} a^{1-s} \cdot \frac{E_{s}(z)}{\lambda_{s}-\lambda_{w}} d s
$$

With $z=x+i y$, from the theory of the constant term [8.1], the Eisenstein series $E_{s}(z)$ is asymptotically dominated by its constant term $y^{s}+c_{s} y^{1-s}$. Thus, for $a \geq y$, by elementary estimates, the contour $\operatorname{Re}(s)=\frac{1}{2}$ can be pushed indefinitely to the right, picking up (negatives of) residues at $s=1$ dues to the pole of $E_{s}$,
and at $s=w$, due to the denominator. The constant-function term exactly cancels the residue at $s=1$. Since $\lambda_{s}-\lambda_{w}=(s-w)(s-1+w)$,

$$
\eta_{a}\left(u_{w}\right)=-\operatorname{Res}_{s=w} a^{1-s} \cdot \frac{E_{s}(z)}{\lambda_{s}-\lambda_{w}}=-a^{1-w} \cdot \frac{E_{w}(z)}{w-1+w}
$$

as asserted.
[12.6.2] Corollary: For $E_{w}(z)=0$, the function $u_{w}=u_{w, z}$ is of rapid decay.
Proof: In the region $\operatorname{Im}\left(z^{\prime}\right)>\operatorname{Im}(z)$, the function $z^{\prime} \rightarrow u_{w, z}\left(z^{\prime}\right)$ is an eigenfunction for $\Delta$. Thus, it is dominated by its constant term, just determined to be $\operatorname{Im}\left(z^{\prime}\right)^{1-w} \cdot E_{w}(z) /(1-2 w)$.
[12.6.3] Remark: Similarly, the constant term of the meromorphic continuation of $u_{w}$ in a topological vector space large enough to include both $H^{1-\varepsilon}$ and Eisenstein series is the meromorphic continuation of $a^{1-w} \cdot E_{w}(z) /(1-2 w)$, assuming that the topology is fine enough so that $\eta_{a}$ is a continuous linear functional on the larger space.
[12.6.4] Remark: As computed more generally in [2.C], finite linear combinations $\theta$ of automorphic Dirac $\delta$ 's applied to Eisenstein series can be arranged to give values of $L$-functions. Thus, solutions $u_{w, \theta}$ to $\left(\Delta-\lambda_{w}\right) u=\theta$ are of rapid decay if and only if $\theta E_{w}=0$.

### 12.7 Whittaker models and a subquotient theorem

Asymptotics of solutions of second-order ordinary differential equations will imply that, for $f$ either an Eisenstein series or a strong-sense cuspform, the representation generated by $f$ has a common $G$ homomorphism image with an unramified principal series $I_{s}$. Specifically $V_{f}$ be the closed subspace of $C^{\infty}(\Gamma \backslash G)$ generated by $f$ under right translations, where $C^{\infty}(\Gamma \backslash G)$ has the Fréchet-space structure given by sups of (Lie algebra) derivatives on compacts.

Let $\psi: N \rightarrow \mathbb{C}^{\times}$be a non-trivial character on $N$, trivial on $\Gamma \cap N$, thus factoring through the (abelian) quotient $(\Gamma \cap N) \backslash N$. For example, one might take

$$
\psi(n)=e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \quad\left(\text { with } x=\left(x_{1}, \ldots, x_{r-1}\right) \in \mathbb{R}^{r-1} \text { and } n=\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\right)
$$

The corresponding Whittaker space is

$$
W_{\psi}=\left\{f \in C^{\infty}(G): f(n g)=\psi_{s}(n) \cdot f(g) \text { for } n \in N \text { and } g \in G\right\}
$$

with Fréchet space structure given by sups of (Lie algebra) derivatives on compacts. The natural $G$ homomorphism $\rho_{\psi}: V_{f} \longrightarrow W_{\psi}$ is given by a Gelfand-Pettis integral:

$$
\rho_{\psi}(F)=\int_{(N \cap \Gamma) \backslash N} \bar{\psi}(n) F(n g) d n \quad\left(\text { for } F \in V_{f}\right)
$$

We will need to have the flexibility to choose $\psi$ for given $f$ so that $\rho_{\psi}(f) \neq 0$.
[12.7.1] Lemma: $\rho_{\psi}: C^{\infty}(\Gamma \backslash G) \rightarrow W_{\psi}$ is continuous.
Proof: Since functions in $W_{\psi}$ are left $\psi$-equivariant, it suffices to show that a compact subset $C$ of $N \backslash G / K$ is covered by a compact subset of $G$. Indeed, the height function $\eta$ assumes a positive inf $\mu$ and finite sup $\sigma$ on $C$. For a sufficiently large compact subset $C_{N}$ of $N$, the compact set

$$
\left\{g=n \cdot\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{y}
\end{array}\right) \cdot K: n \in C_{N}, \mu \leq \eta(g) \leq \sigma\right\} \quad \subset \quad G
$$

surjects to $C$.
[12.7.2] Claim: Let $\Delta f=\lambda \cdot f$ with $\lambda=(r-1)^{2} \cdot s(s-1)$. The image $\rho_{\psi}(f)$ is a constant multiple of

$$
\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{t}
\end{array}\right) \cdot k \quad \longrightarrow \quad e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u(y) \quad(\text { for } k \in K)
$$

where $u$ is the unique (up to scalars) not-rapidly-increasing solution (as $y \rightarrow+\infty$ ) of the differential equation

$$
u^{\prime \prime}-\frac{r-2}{y} \cdot u^{\prime}-\left(4 \pi^{2}+\frac{\lambda}{y^{2}}\right) \cdot u=0
$$

That is, up to scalars, the image $\rho_{\psi}(f)$ in $W_{\psi}$ is uniquely characterized (up to scalar multiples) by satisfaction of that differential equation, and (in fact) rapid decay as $y \rightarrow+\infty$.
Proof: The Casimir operator commutes with $\rho_{\psi}$, and on right $K$-invariant functions is $\Delta$. On $W_{\psi}$, we can separate variables:

$$
\begin{gathered}
\lambda \cdot e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u(y)=\Delta\left(e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u(y)\right. \\
=\left(y^{2}\left(\frac{\partial^{2}}{\partial x_{1}^{2}}+\ldots+\frac{\partial^{2}}{\partial x_{r}^{2}}\right)-(r-2) y \frac{\partial}{\partial y}\right)\left(e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u(y)\right. \\
=y^{2}\left(( 2 \pi i ) ^ { 2 } \cdot \left(e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u(y)+\left(e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u^{\prime \prime}(y)\right)-(r-2) y\left(e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot u^{\prime}(y)\right.\right.\right. \\
=e^{2 \pi i\left(x_{1}+\ldots+x_{r-1}\right)} \cdot\left(-4 \pi^{2} y^{2} \cdot u(y)+y^{2} u^{\prime \prime}(y)-(r-2) y u^{\prime}(y)\right)
\end{gathered}
$$

Thus, the condition is

$$
\lambda \cdot u=-4 \pi^{2} y^{2} \cdot u+y^{2} \cdot u^{\prime \prime}-(r-2) y \cdot u^{\prime}
$$

or

$$
u^{\prime \prime}-\frac{r-2}{y} \cdot u^{\prime}-\left(4 \pi^{2}+\frac{\lambda}{y^{2}}\right) \cdot u=0
$$

The point at $+\infty$ is an irregular singular point of a tractable sort, as in [16.10], [16.B.1]. To see this most clearly, an equation of the form $u^{\prime \prime}+p u^{\prime}+q u=0$ should be rearranged to have no first-derivative term, by the standard procedure. Namely, let $u=v \varphi$ and determine $\varphi$ so that no $v^{\prime}$ term appears in the corresponding differential equation for $v$ :

$$
\begin{gathered}
0=u^{\prime \prime}+p u^{\prime}+q=\left(v^{\prime \prime} \varphi+2 v^{\prime} \varphi^{\prime}+v \varphi^{\prime \prime}\right)+p\left(v^{\prime} \varphi+v \varphi^{\prime}\right)+q(v \varphi) \\
=\varphi \cdot v^{\prime \prime}+\left(2 \varphi^{\prime}+p \varphi\right) \cdot v^{\prime}+\left(\varphi^{\prime \prime}+p \varphi^{\prime}+q \varphi\right) \cdot v
\end{gathered}
$$

Thus, we require $2 \varphi^{\prime}+p \varphi=0$, so $\varphi=e^{-\int p / 2}$, and after dividing through by $\varphi$ the equation is

$$
v^{\prime \prime}+\left(\frac{\varphi^{\prime \prime}}{\varphi}+p \frac{\varphi^{\prime}}{\varphi}+q\right) \cdot v=0
$$

In the case at hand, $p(y)=-(r-2) / y$, so $\varphi(y)=e^{-\int p / 2}=y^{(r-2) / 2}$, and the equation is

$$
v^{\prime \prime}-\left(4 \pi^{2}+\frac{r-2}{2 y^{2}}\right) \cdot v=0
$$

By freezing the coefficients at $y=\infty$, the solutions of the corresponding constant-coefficient differential equation give the correct leading-term asymptotics as $y \rightarrow+\infty$, up to powers of $y$. The frozen equation at $y=\infty$ is $v^{\prime \prime}-4 \pi^{2} \cdot v=0$. The solutions of the frozen equation are linear combinations of $e^{ \pm 2 \pi y}$. From [16.11], these are the leading terms in asymptotics for two linearly independent solutions $v$ of the differential equation. Thus, two linearly independent solutions of the original have asymptotics

$$
u=y^{\frac{r-2}{2}} \cdot e^{ \pm 2 \pi y} \quad(\text { as } y \rightarrow+\infty)
$$

## Garrett: Modern Analysis of Automorphic Forms

Only the scalar multiples of $y^{\frac{r-2}{2}} \cdot e^{-2 \pi y}$ alone, not involving $y^{\frac{r-2}{2}} \cdot e^{+2 \pi y}$, are linear combinations decreasing as $y \rightarrow+\infty$.

By the theory of the constant term [8.1], since by assumption $f$ is a moderate-growth eigenfunction for $\Delta$, the asymptotic behavior of $f$ as $y \rightarrow+\infty$ is dominated (in standard Siegel sets, as height goes to infinity) by its constant term $c_{P} f$, with a rapidly decreasing error term. In particular, applying the Whittaker map $\rho$ to the constant term gives 0 , so $\rho(f)$ is rapidly decreasing. This gives the assertion.

Thus, the image $\rho_{\psi}(f) \in W_{\psi}$ is uniquely determined up to constants, as is $\rho_{\psi}\left(V_{f}\right)$. It is important to note:
[12.7.3] Claim: For given non-constant $f$, there is non-trivial $\psi$ such that $\rho_{\psi}(f) \neq 0$.
Proof: If not, then in Iwasawa coordinates $N \cdot A^{+}$the function $f$ is constant along $N$, and is a function of the $A^{+}$coordinate alone. But apart from constants, there is no such function on $\Gamma \backslash G / K$.

On the other hand, now we will identify the image in $W_{\psi}$ of the corresponding principal series $I_{s}$. We will see that a $G$-homomorphism $I_{s} \rightarrow W_{\psi}$ from a principal series $I_{s}$ also sends the spherical vector $\varphi_{s}^{o} \in I_{s}$ to a function satisfying the same differential equation, and of rapid decay as $y \rightarrow+\infty$. First,
[12.7.4] Claim: On $I_{s}$, the Casimir operator $\Omega$ has eigenvalue $(r-1)^{2} \cdot s(s-1)$.
Proof: This is a computation similar to those in [4.5-4.8]. The computation for $G=S L_{2}(\mathbb{R})$ suffices to illustrate the point. Let

$$
h=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right) \quad X=\left(\begin{array}{cc}
0 & 1 \\
0 & 0
\end{array}\right) \quad Y=\left(\begin{array}{cc}
0 & 0 \\
1 & 0
\end{array}\right) \quad \text { (in the Lie algebra } \mathfrak{g} \text { ) }
$$

For comparison purposes, specify a normalization of the Casimir operator $[4.2] \Omega=\frac{1}{2}\left(\frac{1}{2} h^{2}+X Y+Y X\right) \in U \mathfrak{g}$, so that by the computation in [4.5] in the Iwasawa coordinates $x, y$ on $G / K$,

$$
\left.\Omega=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) \quad \text { (on } G / K\right)
$$

Since $I_{s}$ is defined by a left equivariance condition, it is reasonable to let $\Omega$ act on the left, as the derivative of the left translation action $(g \cdot f)(x)=f\left(g^{-1} x\right)$. In particular, $X$ acts by 0 on $f \in I_{s}$. Thus, $Y X$ acts by 0 . Using the commutation relation,

$$
X Y=Y X+(X Y-Y X)=Y X+[X, Y]=Y X+h
$$

Thus, $X Y$ acts by $h$. Thus, on $I_{s}, \Omega$ acts on the left by $\frac{1}{2}\left(\frac{1}{2} h^{2}+h\right)$. On $I_{s}, h$ acts by

$$
(h \cdot f)(x)=\left.\frac{\partial}{\partial t}\right|_{t=0} f\left(e^{-t} \cdot x\right)=\left.\frac{\partial}{\partial t}\right|_{t=0}\left|\frac{e^{-t}}{e^{t}}\right|^{s} \cdot f(x)=\left.\frac{\partial}{\partial t}\right|_{t=0} e^{-2 t s} \cdot f(x)=-2 s \cdot f(x)
$$

Thus, $\frac{1}{2}\left(\frac{1}{2} h^{2}+h\right)$ acts on $I_{s}$ by

$$
\frac{1}{2}\left(\frac{1}{2}(2 s)^{2}-(2 s)\right)=s(s-1)
$$

as claimed.
At least when the integral converges suitably, the map

$$
(\tau \varphi)(g)=\int_{N} \bar{\psi}(n) \cdot \varphi(w n g) d n \quad\left(\text { with } w=\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)\right)
$$

gives a natural $G$-homomorphism $\tau_{s, \psi}: I_{s} \rightarrow W_{\psi}$. On the spherical vector $\varphi_{s}^{o}$, it is completely determined by its values for $g$ among a set of representatives for $N \backslash G / K$, namely, the Levi component, and by an explicit Iwasawa decomposition [1.3]

$$
\left(\tau_{s, \psi} \varphi_{s}^{o}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{y}
\end{array}\right)=\int_{\mathbb{R}^{r-1}} \bar{\psi}(x) \frac{y^{(r-1) s}}{\left(|x|^{2}+y^{2}\right)^{(r-1) s}} d x
$$

[12.7.5] Claim: The integral for $\tau_{s, \psi}$ on $I_{s}$ converges absolutely for $\operatorname{Re}(s)>\frac{1}{2}$, and produces functions not of rapid growth.
Proof: It suffices to prove convergence for the spherical vector $\varphi_{s}^{o}$, since every other function in $I_{s}$ is dominated by it. Since $|\psi|=1$, letting $\sigma=\operatorname{Re}(s)$,

$$
\left|\int_{\mathbb{R}^{r-1}} \bar{\psi}(x) \frac{y^{(r-1) s}}{\left(|x|^{2}+y^{2}\right)^{(r-1) s}} d x\right| \leq y^{(r-2) \sigma} \int_{\mathbb{R}^{r-1}} \frac{d x}{\left(|x|^{2}+y^{2}\right)^{(r-1) \sigma}} d x=y^{(r-1)(1-\sigma)} \int_{\mathbb{R}^{r-1}} \frac{1}{\left(|x|^{2}+1\right)^{(r-1) \sigma}} d x
$$

by replacing $x \in \mathbb{R}^{r-1}$ by $y \cdot x$. Converting to polar coordinates gives the desired convergence. Further, in that range, the bound is at worst of polynomial growth in $y$, so is not of rapid growth.

The following is necessary for the continuation.
[12.7.6] Theorem: The $G$-map $\tau_{s, \psi}: I_{s} \rightarrow W_{\psi}$ has a meromorphic continuation in $s \in C$, and $\tau_{s, \psi}\left(\varphi_{s}^{o}\right) \neq 0$ except for $s=0,1$. (Proof for $G=S L_{2}(\mathbb{R})$ in the following section.)

Granting the previous theorem, a sufficient subquotient theorem for our purposes follows:
[12.7.7] Theorem: Let $f$ be an Eisenstein series or a strong-sense cuspform, in particular generating an irreducible representation $V_{f}$ of $G$ under right translation on $\Gamma \backslash G$. This entails that $f$ is a $\Delta$-eigenfunction: let $\Delta f=\lambda_{s} \cdot f$ with $\lambda_{s}=(r-1)^{2} \cdot s(s-1)$. Choose additive character $\psi$ on $(N \cap \Gamma) \backslash N$ such that $\rho_{\psi}(f) \neq 0$. Then the image $\rho_{\psi}\left(V_{f}\right) \subset W_{\psi}$ is a subrepresentation of $\tau_{s, \psi}\left(I_{s}\right)$.
Proof: From [12.7.2], the image $W_{\psi}$ contains a unique (up to scalars) right $K$-invariant function $u$ of less than rapid growth with given $\Delta$-eigenvalue $\lambda_{s}$. Since $f$ is at worst of moderate growth, $\rho_{\psi}(f)$ must be a scalar multiple of of that function. Likewise, the image $\tau_{s, \psi}\left(\varphi_{s}^{o}\right)$ of the spherical vector in $I_{s}$ is not of rapid growth, and is non-zero. Thus, the irreducible $\rho_{\psi}\left(V_{f}\right)$ meets $\tau_{s, \psi}\left(I_{s}\right)$ at least in $\mathbb{C} \cdot u$. Thus, $\rho_{\psi}\left(V_{f}\right) \subset \tau_{s, \psi}\left(I_{s}\right)$. ///

### 12.8 Meromorphic continuation of intertwining operators

The analytic continuation of $\tau_{s, \psi}$ to $\operatorname{Re}(s)=\frac{1}{2}$, for real and non-positive $\lambda_{s}=s(s-1)$ possible eigenvalue of $\Delta$ for eigenfunctions in the spectral expansion and Plancherel, is just beyond the range of convergence $\operatorname{Re}(s)>\frac{1}{2}$ of the integral giving $\tau_{s, \psi}$. To prove meromorphic continuation, take $G=S L_{2}(\mathbb{R})$.
[12.8.1] Theorem: The $G$-map $\tau_{s, \psi}: I_{s} \rightarrow W_{\psi}$ has a meromorphic continuation in $s \in C$, and $\tau_{s, \psi}\left(\varphi_{s}^{o}\right) \neq 0$ except for $s=0,1$.
Proof: First, we demonstrate meromorphic continuation of the value of $\tau$ on each of the vectors

$$
\varphi_{s}^{\ell}\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
\sqrt{y} & 0 \\
0 & 1 / \sqrt{y}
\end{array}\right)\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)=y^{s} \cdot e^{2 i \ell \theta}
$$

via Bochner's lemma [3.A]. Using the explicit Iwasawa decomposition [1.3] for $S L_{2}(\mathbb{R})$, after some typical minor rearrangements, up to irrelevant constants the integral for $\tau$ is

$$
\left.\int_{\mathbb{R}} e^{-i x} \frac{1}{(y+i x)^{s+\ell}(y-i x)^{s-\ell}} d x \quad \quad \text { absolutely convergent for } \operatorname{Re}(s)>1\right)
$$

Thus, for $\alpha, \beta \in \mathbb{C}$, on one hand consider

$$
\left.\int_{\mathbb{R}} e^{-i x} \frac{1}{(y+i x)^{\alpha}(y-i x)^{\beta}} d x \quad \text { (absolutely convergent for } \operatorname{Re}(\alpha+\beta)>\frac{1}{2}\right)
$$

On the other hand, the identity $\int_{0}^{\infty} e^{-t(y+i x)} t^{v} d t / t=(i z)^{-v} \cdot \Gamma(v)$ for $y>0$ can be viewed as a computation of a Fourier transform:

$$
\int_{\mathbb{R}} e^{-i t x} \cdot\left\{\begin{array}{cc}
t^{v-1} e^{-t y} & (\text { for } t>0) \\
0 & (\text { for } t<0)
\end{array}\right\} d t=(y+i x)^{-v} \cdot \Gamma(v)
$$

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By Fourier inversion, up to irrelevant constants,

$$
\int_{\mathbb{R}} e^{i x t}(y+i x)^{-v} d x=\left\{\begin{array}{cc}
\frac{1}{\Gamma(v)} e^{-t y} t^{v-1} & (\text { for } t>0) \\
0 & (\text { for } t<0)
\end{array}\right.
$$

Replacing $x$ by $-x$ and $t$ by $-t$ gives the corresponding identity for $(y-i x)^{-v}$ :

$$
\int_{\mathbb{R}} e^{i x t}(y-i x)^{-v} d x=\left\{\begin{array}{cc}
0 & (\text { for } t>0) \\
\frac{1}{\Gamma(v)} e^{-|t| y}|t|^{v-1} & (\text { for } t<0)
\end{array}\right.
$$

The Fourier transform of a product is the convolution, so, up to irrelevant constants, for $\xi>0$,

$$
\begin{aligned}
& \int_{\mathbb{R}} e^{-i x \xi} \frac{1}{(1+i x)^{\alpha}(1-i x)^{\beta}} d x=\frac{1}{\Gamma(\alpha) \cdot \Gamma(\beta)} \int_{\xi-t>0, t<0} e^{-|\xi-t| y}|\xi-t|^{\alpha-1} \cdot e^{-|t| y}|t|^{\beta-1} d t \\
= & \frac{1}{\Gamma(\alpha) \cdot \Gamma(\beta)} \int_{0}^{\infty} e^{-|\xi+t| y}|\xi+t|^{\alpha-1} \cdot e^{-t y} t^{\beta-1} d t=\frac{1}{\Gamma(\alpha) \cdot \Gamma(\beta)} \int_{0}^{\infty} e^{-(\xi+2 t) y}(\xi+t)^{\alpha-1} t^{\beta} \frac{d t}{t}
\end{aligned}
$$

Since $\xi>0$, this is convergent for all $\alpha \in \mathbb{C}$ and for $\operatorname{Re}(\beta)>0$. The convex hull of the union of the regions $\{(\alpha, \beta): \operatorname{Re}(\alpha+\beta)>1\}$ and $\{(\alpha, \beta): \operatorname{Re}(\beta)>0\}$ is all of $\mathbb{C}^{2}$ so Bochner's lemma [3.A] gives the meromorphic continuations of the functions $\tau_{s, \psi}\left(\varphi_{s}^{\ell}\right)$. (The integral expressions show that the vertical growth in $\alpha$ and $\beta$ is mild enough to allow application of Bochner's lemma.)

Via the Iwasawa decomposition $G=P K$, functions $\varphi$ in the smooth principal series $I_{s}$ can be identified with Fourier series on $K \approx S O_{2}(\mathbb{R})$ with rapidly decreasing coefficients. That is,

$$
\varphi\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)=\sum_{\ell \in \mathbb{Z}} c_{\ell} \varphi_{s}^{\ell}\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)=\sum_{\ell \in \mathbb{Z}} c_{\ell} e^{2 i \ell \theta}
$$

with rapidly decreasing $c_{\ell}$. We want to show that the image $\sum_{\ell \in \mathbb{Z}} c_{\ell} \tau_{s, \psi}\left(\varphi_{s}^{\ell}\right)$ is still convergent to a smooth function in $W_{\psi}$. The intertwining operator $\tau_{s, \psi}$ preserves the right $K$-equivariance. Thus, for some constants $C_{s, \ell, \psi}$,

$$
\tau_{s, \psi}\left(\varphi_{s}^{\ell}\right)\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)=C_{s, \ell, \psi} \cdot e^{2 i \ell \theta}
$$

Thus, it suffices to show that the $C_{s, \ell, \psi}$ grow (in $\ell$ ) at most polynomially.
In $\operatorname{Re}(s) \geq \frac{1}{2}+\delta$ for fixed small $\delta>0$, the integral for $\tau_{s, \psi}\left(\varphi_{s}^{\ell}\right)$ converges absolutely, and is uniformly bounded in $\ell \in \mathbb{Z}$ and $\operatorname{Re}(s) \geq \frac{1}{2}+\varepsilon$. In the next section we exhibit an intertwining operator $I_{s} \rightarrow I_{1-s}$ that is an isomorphism for all $s$ with $\left|\operatorname{Re}(s)-\frac{1}{2}\right|<\delta$ with $0<\delta<1$, and that sends $\varphi_{s}^{\ell} \longrightarrow A_{s, \ell} \cdot \varphi_{1-s}^{\ell}$ with polynomial-growth (in $\ell$ ) constants $A_{s, \ell}$. Thus, the analytic continuation demonstrated above extends to smooth vectors in $\frac{1}{4} \leq \operatorname{Re}(s) \leq \frac{1}{2}-\delta$ (for example). By Phragmén-Lindelöf, each individual polynomial growth bound extends to $\frac{1}{4} \leq \operatorname{Re}(s) \leq \frac{3}{4}$, giving the analytic continuation of $\tau_{s, \psi}$ to that region.

### 12.9 Intertwining operators among principal series

To have essentially elementary computations, we consider only $G=S L_{2}(\mathbb{R})$. The standard intertwining operator $T=T_{s}: I_{s} \rightarrow I_{1-s}$ is defined, for $\operatorname{Re}(s)$ sufficiently large, by the integral ${ }^{[74]}$

$$
T_{s} f(g)=\int_{N} f(w n \cdot g) d n
$$

[74] Why this integral? This is an analogue of a finite-group method for writing formulas for intertwining operators from a representation induced from a subgroup $A$ to a representation induced from a subgroup $B$, with intertwining operators roughly corresponding to double cosets $A \backslash G / B$. For finite groups, this goes by the name of Mackey theory, and Bruhat extended the idea to Lie groups and p-adic groups. For non-finite groups, there are issues of convergence and analytic continuation.
with the longest Weyl element

$$
w=\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)
$$

Convergence will be clarified shortly. Since the map is an integration on the left, it does not disturb the right action of $G$. To verify that (assuming convergence) the image really does lie inside $I_{1-s}$, observe that $T_{s} f$ is left $N$-invariant by construction, and that for $m \in M$

$$
\left(T_{s} f\right)(m g)=\int_{N} f(w n \cdot m g) d n=\int_{N} f\left(w m m^{-1} n m \cdot g\right) d n=\chi_{1}(m) \cdot \int_{N} f(w m n \cdot g) d n
$$

by replacing $n$ by $m n m^{-1}$, taking into account the change of measure $d\left(m n m^{-1}\right)=\chi_{1}(m) \cdot d n$ coming from

$$
\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)\left(\begin{array}{cc}
1 & x \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
a & 0 \\
0 & a^{-1}
\end{array}\right)^{-1}=\left(\begin{array}{cc}
1 & a^{2} x \\
0 & 1
\end{array}\right)
$$

This is

$$
\begin{gathered}
\chi_{1}(m) \cdot \int_{N} f\left(w m w^{-1} \cdot w n \cdot g\right) d n=\chi_{1}(m) \cdot \int_{N} f\left(m^{-1} \cdot w n \cdot g\right) d n \\
=\chi_{1}(m) \chi_{s}\left(m^{-1}\right) \cdot \int_{N} f(w n \cdot g) d n=\chi_{1-s}(m) \cdot\left(T_{s} f\right)(g)
\end{gathered}
$$

This verifies that $T_{s}: I_{s} \rightarrow I_{1-s}$.
Parametrize the maximal compact by

$$
K=\left\{\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right): \theta \in \mathbb{R} / 2 \pi i \mathbb{Z}\right\}
$$

and note that the overlap is just $P \cap K= \pm 1$. Thus, a function $f$ in $I_{s}$ is completely determined by its values on $K$, in fact, on $\{ \pm 1\} \backslash K$. Conversely, for fixed $s \in \mathbb{C}$, any smooth function $f_{o}$ on $\{ \pm 1\} \backslash K$ has a unique extension (depending upon $s$ ) to a function $f \in I_{s}$, by $f(p k)=\chi_{s}(p) \cdot f_{o}(k)$. Taking advantage of the simplicity of this situation, we may expand smooth functions on $K$ in Fourier series

$$
f\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)=\sum_{\ell \in \mathbb{Z}} c_{\ell} e^{2 \pi i \ell \theta}
$$

where the Fourier coefficients $c_{n}$ are rapidly decreasing due to the smoothness of $f$.
Initially, we restrict our attention to functions $f \in I_{s}$ which are not merely smooth, but in fact right $K$-finite in the sense that the Fourier expansion of $f$ restricted to $K$ is finite. That is, these are finite sums of functions $\varphi_{s}^{\ell}(p k)=\chi_{s}(p) \cdot \rho_{\ell}(k)$ where

$$
\rho_{\ell}\left(\begin{array}{rr}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)=e^{2 \pi i \ell \theta}
$$

For any function $f$ on $G$ with

$$
f(p k)=f(g) \cdot \rho(k) \quad(\text { for } k \in K)
$$

with $\rho$ among the $\rho_{\ell}$, say that $f$ has (right) $K$-type $\rho$. From the Iwasawa decomposition $G=P K$

$$
\operatorname{dim}_{\mathbb{C}}\left\{f \in I_{s}: f \text { has right } K \text {-type } \rho_{\ell}\right\}=1
$$

The main computation: We will directly compute the effect of the intertwining operator $T_{s}$ on $\varphi_{s}^{\ell}$. Since the left integration over $N$ cannot affect the right $K$-type, $T_{s}$ preserves $K$-types. Since the dimensions of the subspaces of $I_{s}$ and $I_{1-s}$ with given $K$-type $\rho$ are 1 , necessarily $T_{s}$ maps $\varphi_{s}^{\ell}$ to some multiple of $\varphi_{1-s}^{\ell}$. To determine this constant, it suffices to evaluate $\left(T_{s} f\right)(1)$, that is, to evaluate the integral

$$
\left(T_{s} f\right)(1)=\int_{N} f(w n) d n=\int_{\mathbb{R}} f\left(w n_{x}\right) d x \quad\left(\text { with } n_{x}=\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right)\right)
$$

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To evaluate $f\left(w n_{x}\right)$, we give the Iwasawa decomposition $w n_{x}=p k$. One convenient approach is to compute

$$
\left(w n_{x}\right)\left(w n_{x}\right)^{\top}=(p k)(p k)^{\top}=p k k^{-1} p^{\top}=p p^{\top}
$$

since $k$ is orthogonal. Letting

$$
p=\left(\begin{array}{cc}
a & b \\
0 & a^{-1}
\end{array}\right)
$$

and expanding $\left(w n_{x}\right)\left(w n_{x}\right)^{\top}$ gives

$$
\left(\begin{array}{cc}
1 & -x \\
-x & 1+x^{2}
\end{array}\right)=\left(\begin{array}{cc}
a^{2}+b^{2} & b / a \\
b / a & 1 / a^{2}
\end{array}\right)
$$

from which $a^{-2}=1+x^{2}$ and $b / a=-x$, so $a=1 / \sqrt{1+x^{2}}$ and $b=-x / \sqrt{1+x^{2}}$. Then $k=p^{-1} g$, so we find the Iwasawa decomposition

$$
w n_{x}=\left(\begin{array}{cc}
\frac{1}{\sqrt{1+x^{2}}} & \frac{-x}{\sqrt{1+x^{2}}} \\
0 & \sqrt{1+x^{2}}
\end{array}\right) \cdot\left(\begin{array}{cc}
\frac{1}{\sqrt{1+x^{2}}} & \frac{x}{\sqrt{1+x^{2}}} \\
\frac{-x}{\sqrt{1+x^{2}}} & \frac{1}{\sqrt{1+x^{2}}}
\end{array}\right)
$$

Thus, with $K$-type $\rho_{\ell}$,

$$
\begin{gathered}
\varphi_{s}^{\ell}\left(w n_{x}\right)=\varphi_{s}^{\ell}(p k)=\chi_{s}(p) \cdot \rho_{\ell}(k)=\left(\frac{1}{\sqrt{1+x^{2}}}\right)^{2 s} \cdot\left(\frac{1+i x}{\sqrt{1+x^{2}}}\right)^{2 \ell}=\left(1+x^{2}\right)^{-s} \cdot\left(\frac{1+i x}{\sqrt{1+i x} \cdot \sqrt{1-i x}}\right)^{2 \ell} \\
=\left(1+x^{2}\right)^{-s} \cdot\left(\frac{1+i x}{1-i x}\right)^{\ell}=(1+i x)^{-s+\ell}(1-i x)^{-s-\ell}
\end{gathered}
$$

Thus, our intertwining operator when applied to a $f \in I_{s}$ with specified $K$-type $\rho_{\ell}$, evaluated at $1 \in G$ is

$$
\left(T_{s} f\right)(1)=\int_{N} f(w n) d n=\int_{\mathbb{R}}(1+i x)^{-s+\ell}(1-i x)^{-s-\ell} d x
$$

To compute the latter, we use a standard trick employing the gamma function. That is, for complex $z$ in the right half-plane, and for $\operatorname{Re}(s)>0$,

$$
\Gamma(s) \cdot z^{-s}=\int_{0}^{\infty} e^{-t z} t^{-s} \frac{d t}{t}
$$

Thus,

$$
\begin{gathered}
\left(T_{s} \varphi_{s}^{\ell}\right)(1)=\int_{\mathbb{R}}(1+i x)^{-s+\ell}(1-i x)^{-s-\ell} d x \\
=\Gamma(s-\ell)^{-1} \Gamma(s+\ell)^{-1} \cdot \int_{\mathbb{R}} \int_{0}^{\infty} \int_{0}^{\infty} e^{-u(1+i x)} u^{-(s-\ell)} e^{-v(1+i x)} v^{-(s+\ell)} \frac{d u}{u} \frac{d v}{v} d x
\end{gathered}
$$

Changing the order of integration and integrating in $x$ first ${ }^{[75]}$ gives an inner integral

$$
\int_{\mathbb{R}} e^{i x(u-v)} d x=2 \pi \cdot \delta_{u-v}
$$

where $\delta$ is the Dirac delta distribution. Thus, the whole integral becomes
[75] This is not legitimate from an elementary viewpoint. However, it is a compelling heuristic, correctly suggests the true conclusion, and can immediately be justified by Fourier inversion, as is done in the appendix.

$$
\begin{aligned}
& \left(T_{s} \varphi_{s}^{\ell}\right)(1)=\frac{2 \pi}{\Gamma(s-\ell) \Gamma(s+\ell)} \int_{0}^{\infty} e^{-u} u^{-(s-\ell)} e^{-u} u^{-(s+\ell)} u^{-1} u^{-1} d u \\
& \quad=\frac{2 \pi}{\Gamma(s-\ell) \Gamma(s+\ell)} \int_{0}^{\infty} e^{-2 u} u^{-2 s-1} \frac{d u}{u}=\frac{2 \pi 2^{1-2 s} \Gamma(2 s-1)}{\Gamma(s-\ell) \Gamma(s+\ell)}
\end{aligned}
$$

That is, under the intertwining $T_{s}: I_{s} \rightarrow I_{1-s}$, the function $\varphi_{s}^{\ell}$ is mapped to $\varphi_{1-s}^{\ell}$ multiplied by that last constant.

Subrepresentations: For brevity, let

$$
\lambda(s, n)=\frac{2 \pi 2^{1-2 s} \Gamma(2 s-1)}{\Gamma(s-\ell) \Gamma(s+\ell)}
$$

denote the constant computed above. The intertwining operator $T_{s}$ is holomorphic at $s_{o} \in \mathbb{C}$ if for all integers $\ell$ the function $\lambda(s, \ell)$ is holomorphic at $s_{o}$.

The numerator $\Gamma(2 s-1)$ has poles at

$$
\frac{1}{2}, 0,-\frac{1}{2},-1,-\frac{3}{2},-2, \ldots
$$

The half-integer poles are not canceled by the poles of the denominator, so $T_{s}$ has poles at these half-integers. At the non-positive integers, regardless of the value of $\ell$ the poles of the denominator cancel the pole of the numerator. That is,
[12.9.1] Claim: $T_{s}: I_{s} \rightarrow I_{1-s}$ is holomorphic away from

$$
s=\frac{1}{2},-\frac{1}{2},-\frac{3}{2},-\frac{5}{2},-\frac{7}{2}, \ldots
$$

at which it has simple poles.
For $s$ not an integer, the denominator has no poles, so (away from the half-integers at which the numerator has a pole) $\lambda(s, \ell) \neq 0$ for all $K$-types $\rho_{2 n}$. Thus,
[12.9.2] Claim: The intertwining operator $T_{s}: I_{s} \rightarrow I_{1-s}$ has trivial kernel for $s$ not an integer (and away from its poles).

Consider $s=m$ with $0<m \in \mathbb{Z}$. The numerator has no pole at $m$, while the denominator has a pole, yielding a $\lambda(m, \ell)=0$ for all integers

$$
n= \pm m, \pm(m+1), \pm(m+2), \pm(m+3), \ldots
$$

Thus, for $0<m \in \mathbb{Z}, I_{m}$ has a non-trivial infinite-dimensional subrepresentation ${ }^{[76]}$ consisting of these $K$-types in the kernel of $T_{m}: I_{m} \rightarrow I_{1-m}$.

Consider $s=-m$ with $0 \geq-m \in \mathbb{Z}$. The numerator has a pole at $-m$, and the denominator has a double pole for integers

$$
\ell=0, \pm 1, \pm 2, \ldots, \pm m
$$

and a single pole for integers

$$
\ell= \pm(m+1), \pm(m+2), \pm(m+3), \ldots
$$

Thus, $\lambda(-m, \ell)=0$ for the double poles, and the single poles cancel. Thus, for $0 \geq m \in \mathbb{Z}, I_{m}$ has a non-trivial subrepresentation consisting of the finitely-many $K$-types at which the denominator has a double pole. These are (therefore) finite-dimensional representations, the kernels of $T_{-m}: I_{-m} \rightarrow I_{1+m}$.

[^22]Smooth vectors: The explicit computation of the scalar $\lambda(s, \ell)=\left(T_{s} \varphi_{s}^{\ell}\right)(1)$ for $\varphi_{s}^{\ell}$ also shows that $T_{s}$ has an analytic continuation on smooth vectors in $I_{s}$, not merely $K$-finite vectors, as follows. From $\Gamma(s) \cdot s=\Gamma(s+1)$,

$$
\left(T_{s} \varphi_{s}^{\ell}\right)(1)=\frac{2 \pi 2^{1-2 s} \Gamma(2 s-1)}{\Gamma(s-\ell) \Gamma(s+\ell)}=\text { polynomial growth in } \ell
$$

By asymptotics of $\Gamma(s)$, for example from Stirling's formula in [16.1]. Let

$$
f=\sum_{\ell \in \mathbb{Z}} c_{\ell} \cdot \varphi_{s}^{\ell}
$$

be smooth in $I_{s}$. Smoothness is equivalent to the rapid decrease of the Fourier coefficients. Then

$$
T_{s} f=\sum_{\ell \in \mathbb{Z}} \lambda(s, \ell) \cdot c_{\ell} \cdot \varphi_{s}^{\ell}
$$

still has rapidly decreasing coefficients, so is a smooth vector in $I_{1-s}$. That is, away from the poles, the intertwining operator $T_{s}$ when analytically continued is defined on all smooth vectors in $I_{1-s}$, not merely $K$-finite ones.

## 12.A Appendix: a usual trick with $\Gamma(s)$

The property of $\Gamma(s)$ used above is sufficiently important that we review it. The gamma function is given for $\operatorname{Re}(s)>0$ by Euler's integral

$$
\Gamma(s)=\int_{0}^{\infty} e^{-t} t^{s} \frac{d t}{t}
$$

Replacing $t$ by $t y$ with $y>0$

$$
\Gamma(s) \cdot y^{-s}=\int_{0}^{\infty} e^{-t y} t^{s} \frac{d t}{t}
$$

By analytic continuation to the right complex half-plane, for $y>0$ and $x \in \mathbb{R}$

$$
\Gamma(s) \cdot(y+2 \pi i x)^{-s}=\int_{0}^{\infty} e^{-t(y+2 \pi i x)} t^{s} \frac{d t}{t}
$$

Having analytically continued, we may let $y=1$ again, obtaining

$$
\Gamma(s) \cdot(1+2 \pi i x)^{-s}=\int_{0}^{\infty} e^{-t(1+2 \pi i x)} t^{s} \frac{d t}{t}=\int_{0}^{\infty} e^{-2 \pi i x t} e^{-t} t^{s} \frac{d t}{t}
$$

which is the Fourier transform of

$$
\varphi_{s}(t)=\left\{\begin{array}{cc}
e^{-t} t^{s-1} & (t>0) \\
0 & (t<0)
\end{array}\right.
$$

To compute the concrete integral for $\left(T_{s} f\right)(1)$ invoke the Plancherel theorem, that

$$
\int_{\mathbb{R}} f(x) \overline{\varphi(x)} d x=\int_{\mathbb{R}} \hat{f}(x) \overline{\hat{\varphi}(x)} d x
$$

and Fourier inversion. With real $s \gg 0$, replacing $x$ by $2 \pi x$ at the first step, and with real $s$,

$$
\begin{aligned}
& \int_{\mathbb{R}}(1+i x)^{-s+n}(1-i x)^{-s-n} d x=2 \pi \int_{\mathbb{R}}(1+2 \pi i x)^{-s+n}(1-2 \pi i x)^{-s-n} d x \\
&=2 \pi \int_{\mathbb{R}} \hat{\varphi}_{s-n}(x) \overline{\hat{\varphi}_{s+n}(x)} d x=2 \pi \int_{\mathbb{R}} \varphi_{s-n}(x) \overline{\varphi_{s+n}(x)} d x \\
&= \frac{2 \pi}{\Gamma(s-n) \Gamma(s+n)} \int_{0}^{\infty} e^{-u} u^{-(s-n)-1} \cdot e^{-u} u^{-(s+n)-1} d u=\frac{2 \pi \Gamma(2 s-1)}{\Gamma(s-n) \Gamma(s+n)}
\end{aligned}
$$

as computed heuristically earlier. This also exhibits the constant $2 \pi$.

## 13. Examples: topologies on natural function spaces

1. Banach spaces $C^{k}[a, b]$
2. Non-Banach limit $C^{\infty}[a, b]$ of Banach spaces $C^{k}[a, b]$
3. Sufficient notion of topological vector space
4. Unique vectorspace topology on $\mathbb{C}^{n}$
5. Non-Banach limits $C^{k}(\mathbb{R}), C^{\infty}(\mathbb{R})$ of Banach spaces $C^{k}[a, b]$
6. Banach completion $C_{o}^{k}(\mathbb{R})$ of $C_{c}^{k}(\mathbb{R})$
7. Rapid-decay functions, Schwartz functions
8. Non-Fréchet colimit $\mathbb{C}^{\infty}$ of $\mathbb{C}^{n}$, quasi-completeness
9. Non-Fréchet colimit $C_{c}^{\infty}(\mathbb{R})$ of Fréchet spaces
10. LF-spaces of moderate-growth functions
11. Seminorms and locally convex topologies
12. Quasi-completeness theorem
13. Strong operator topology
14. Generalized functions (distributions) on $\mathbb{R}$
15. Tempered distributions and Fourier transforms on $\mathbb{R}$
16. Test functions and Paley-Wiener spaces
17. Schwartz functions and Fourier transforms on $\mathbb{Q}_{p}$

We review natural topological vectorspaces of functions on relatively simple geometric objects, such as $\mathbb{R}$, as opposed to the automorphic examples $\Gamma \backslash G$ and $\Gamma \backslash X$, to separate the geometric and group-theoretic complications from the topological-analytical.

In all cases, we specify a natural topology, in which differentiation and other natural operators are continuous, and so that the space is complete.

Many familiar and useful spaces of continuous or differentiable functions, such as $C^{k}[a, b]$, have natural metric structures, and are complete. Often, the metric $d($,$) comes from a norm |\cdot|$, on the functions, giving Banach spaces.

Other natural function spaces, such as $C^{\infty}[a, b], C^{o}(\mathbb{R})$, are not Banach, but still do have a metric topology and are complete: these are Fréchet spaces, appearing as (projective) limits of Banach spaces, as below. These lack some of the conveniences of Banach spaces, but their expressions as limits of Banach spaces is often sufficient.

Other important spaces, such as compactly-supported continuous functions $C_{c}^{o}(\mathbb{R})$ on $\mathbb{R}$, or compactlysupported smooth functions (test functions) $\mathcal{D}(\mathbb{R})=C_{c}^{\infty}(\mathbb{R})$ on $\mathbb{R}$, are not metrizable so as to be complete. Nevertheless, some are expressible as colimits (sometimes called inductive limits) of Banach or Fréchet spaces, and such descriptions suffice for many applications. An LF-space is a countable ascending union of Fréchet spaces with each Fréchet subspace closed in the next. These are strict colimits or strict inductive limits of Fréchet spaces. These are generally not complete in the strongest sense, but, nevertheless, as demonstrated in [13.12], are quasi-complete, and this suffices for applications.

### 13.1 Banach spaces $C^{k}[a, b]$

We give the vector space $C^{k}[a, b]$ of $k$-times continuously differentiable functions on an interval $[a, b]$ a metric which makes it complete. Mere pointwise limits of continuous functions easily fail to be continuous. First recall the standard
[13.1.1] Claim: The set $C^{o}(K)$ of complex-valued continuous functions on a compact set $K$ is complete with the metric $|f-g|_{C^{o}}$, with the $C^{o}$-norm $|f|_{C^{o}}=\sup _{x \in K}|f(x)|$.
Proof: This is a typical three-epsilon argument. To show that a Cauchy sequence $\left\{f_{i}\right\}$ of continuous functions has a pointwise limit which is a continuous function, first argue that $f_{i}$ has a pointwise limit at every $x \in K$. Given $\varepsilon>0$, choose $N$ large enough such that $\left|f_{i}-f_{j}\right|<\varepsilon$ for all $i, j \geq N$. Then $\left|f_{i}(x)-f_{j}(x)\right|<\varepsilon$ for any $x$ in $K$. Thus, the sequence of values $f_{i}(x)$ is a Cauchy sequence of complex numbers, so has a limit $f(x)$.

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Further, given $\varepsilon^{\prime}>0$ choose $j \geq N$ sufficiently large such that $\left|f_{j}(x)-f(x)\right|<\varepsilon^{\prime}$. For $i \geq N$

$$
\left|f_{i}(x)-f(x)\right| \leq\left|f_{i}(x)-f_{j}(x)\right|+\left|f_{j}(x)-f(x)\right|<\varepsilon+\varepsilon^{\prime}
$$

This is true for every positive $\varepsilon^{\prime}$, so $\left|f_{i}(x)-f(x)\right| \leq \varepsilon$ for every $x$ in $K$. That is, the pointwise limit is approached uniformly in $x \in[a, b]$.

To prove that $f(x)$ is continuous, for $\varepsilon>0$, take $N$ be large enough so that $\left|f_{i}-f_{j}\right|<\varepsilon$ for all $i, j \geq N$. From the previous paragraph $\left|f_{i}(x)-f(x)\right| \leq \varepsilon$ for every $x$ and for $i \geq N$. Fix $i \geq N$ and $x \in K$, and choose a small enough neigborhood $U$ of $x$ such that $\left|f_{i}(x)-f_{i}(y)\right|<\varepsilon$ for any $y$ in $U$. Then

$$
|f(x)-f(y)| \leq\left|f(x)-f_{i}(x)\right|+\left|f_{i}(x)-f_{i}(y)\right|+\left|f(y)-f_{i}(y)\right| \leq \varepsilon+\left|f_{i}(x)-f_{i}(y)\right|+\varepsilon<\varepsilon+\varepsilon+\varepsilon
$$

Thus, the pointwise limit $f$ is continuous at every $x$ in $U$.
Unsurprisingly, but significantly:
[13.1.2] Claim: For $x \in[a, b]$, the evaluation map $f \rightarrow f(x)$ is a continuous linear functional on $C^{o}[a, b]$.
Proof: For $|f-g|_{C^{o}}<\varepsilon$, we have

$$
|f(x)-g(x)| \leq|f-g|_{C^{o}}<\varepsilon
$$

proving the continuity.
As usual, a real-valued or complex-valued function $f$ on a closed interval $[a, b] \subset \mathbb{R}$ is continuously differentiable when it has a derivative which is itself a continuous function. That is, the limit

$$
f^{\prime}(x)=\lim _{h \rightarrow 0} \frac{f(x+h)-f(x)}{h}
$$

exists for all $x \in[a, b]$, and the function $f^{\prime}(x)$ is in $C^{o}[a, b]$. Let $C^{k}[a, b]$ be the collection of $k$-times continuously differentiable functions on $[a, b]$, with the $C^{k}$-norm

$$
|f|_{C^{k}}=\sum_{0 \leq i \leq k} \sup _{x \in[a, b]}\left|f^{(i)}(x)\right|=\sum_{0 \leq i \leq k}\left|f^{(i)}\right|_{\infty}
$$

where $f^{(i)}$ is the $i^{\text {th }}$ derivative of $f$. The associated metric on $C^{k}[a, b]$ is $|f-g|_{C^{k}}$.
Similar to the assertion about evaluation on $C^{o}[a, b]$,
[13.1.3] Claim: For $x \in[a, b]$ and $0 \leq j \leq k$, the evaluation map $f \rightarrow f^{(j)}(x)$ is a continuous linear functional on $C^{k}[a, b]$.
Proof: For $|f-g|_{C^{k}}<\varepsilon$,

$$
\left|f^{(j)}(x)-g^{(j)}(x)\right| \leq|f-g|_{C^{k}}<\varepsilon
$$

proving the continuity.
We see that $C^{k}[a, b]$ is a Banach space:
[13.1.4] Theorem: The normed metric space $C^{k}[a, b]$ is complete.
Proof: For a Cauchy sequence $\left\{f_{i}\right\}$ in $C^{k}[a, b]$, all the pointwise limits $\lim _{i} f_{i}^{(j)}(x)$ of $j$-fold derivatives exist for $0 \leq j \leq k$, and are uniformly continuous. The issue is to show that $\lim _{i} f^{(j)}$ is differentiable, with derivative $\lim _{i} f^{(j+1)}$. It suffices to show that, for a Cauchy sequence $f_{n}$ in $C^{1}[a, b]$, with pointwise limits $f(x)=\lim _{n} f_{n}(x)$ and $g(x)=\lim _{n} f_{n}^{\prime}(x)$ we have $g=f^{\prime}$. By the fundamental theorem of calculus, for any index $i$,

$$
f_{i}(x)-f_{i}(a)=\int_{a}^{x} f_{i}^{\prime}(t) d t
$$

Since the $f_{i}^{\prime}$ uniformly approach $g$, given $\varepsilon>0$ there is $i_{o}$ such that $\left|f_{i}^{\prime}(t)-g(t)\right|<\varepsilon$ for $i \geq i_{o}$ and for all $t$ in the interval, so for such $i$

$$
\left|\int_{a}^{x} f_{i}^{\prime}(t) d t-\int_{a}^{x} g(t) d t\right| \leq \int_{a}^{x}\left|f_{i}^{\prime}(t)-g(t)\right| d t \leq \varepsilon \cdot|x-a| \longrightarrow 0
$$

Thus,

$$
\lim _{i} f_{i}(x)-f_{i}(a)=\lim _{i} \int_{a}^{x} f_{i}^{\prime}(t) d t=\int_{a}^{x} g(t) d t
$$

from which $f^{\prime}=g$.
By design, we have
[13.1.5] Theorem: The map $\frac{d}{d x}: C^{k}[a, b] \rightarrow C^{k-1}[a, b]$ is continuous.
Proof: As usual, for a linear map $T: V \rightarrow W$, by linearity $T v-T v^{\prime}=T\left(v-v^{\prime}\right)$ it suffices to check continuity at 0 . For Banach spaces the homogeneity $|\sigma \cdot v|_{V}=|\alpha| \cdot|v|_{V}$ shows that continuity is equivalent to existence of a constant $B$ such that $|T v|_{W} \leq B \cdot|v|_{V}$ for $v \in V$. Then

$$
\left|\frac{d}{d x} f\right|_{C^{k-1}}=\sum_{0 \leq i \leq k-1} \sup _{x \in[a, b]}\left|\left(\frac{d f}{d x}\right)^{(i)}(x)\right|=\sum_{1 \leq i \leq k} \sup _{x \in[a, b]}\left|f^{(i)}(x)\right| \leq 1 \cdot|f|_{C^{k}}
$$

as desired.

### 13.2 Non-Banach limit $C^{\infty}[a, b]$ of Banach spaces $C^{k}[a, b]$

The space $C^{\infty}[a, b]$ of infinitely differentiable complex-valued functions on a (finite) interval $[a, b]$ in $\mathbb{R}$ is not a Banach space. ${ }^{[77]}$ Nevertheless, the topology is completely determined by its relation to the Banach spaces $C^{k}[a, b]$. That is, there is a unique reasonable topology on $C^{\infty}[a, b]$. After explaining and proving this uniqueness, we also show that this topology is complete metric.

This function space can be presented as

$$
C^{\infty}[a, b]=\bigcap_{k \geq 0} C^{k}[a, b]
$$

and we reasonably require that whatever topology $C^{\infty}[a, b]$ should have, each inclusion $C^{\infty}[a, b] \longrightarrow C^{k}[a, b]$ is continuous.

At the same time, given a family of continuous linear maps $Z \rightarrow C^{k}[a, b]$ from a vector space $Z$ in some reasonable class (specified in the next section), with the compatibility condition of giving commutative diagrams

the image of $Z$ actually lies in the intersection $C^{\infty}[a, b]$. Thus, diagrammatically, for every family of compatible maps $Z \rightarrow C^{k}[a, b]$, there is a unique $Z \rightarrow C^{\infty}[a, b]$ fitting into a commutative diagram


We require that this induced map $Z \rightarrow C^{\infty}[a, b]$ is continuous.
When we know that these conditions are met, we would say that $C^{\infty}[a, b]$ is the (projective) limit of the spaces $C^{k}[a, b]$, written

$$
C^{\infty}[a, b]=\lim _{k} C^{k}[a, b]
$$

[77] It is not essential to prove that there is no reasonable Banach space structure on $C^{\infty}[a, b]$, but this can be readily proven in a suitable context.

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with implicit reference to the inclusions $C^{k+1}[a, b] \rightarrow C^{k}[a, b]$ and $C^{\infty}[a, b] \rightarrow C^{k}[a, b]$.
[13.2.1] Claim: Up to unique isomorphism, there exists at most one topology on $C^{\infty}[a, b]$ such that to every compatible family of continuous linear maps $Z \rightarrow C^{k}[a, b]$ from a topological vector space $Z$ there is a unique continuous linear $Z \rightarrow C^{\infty}[a, b]$ fitting into a commutative diagram as just above.
Proof: Let $X, Y$ be $C^{\infty}[a, b]$ with two topologies fitting into such diagrams, and show $X \approx Y$, and for a unique isomorphism. First, claim that the identity $\operatorname{map}_{\operatorname{id}_{X}}: X \rightarrow X$ is the only map $\varphi: X \rightarrow X$ fitting into a commutative diagram


Indeed, given a compatible family of maps $X \rightarrow C^{k}[a, b]$, there is unique $\varphi$ fitting into


Since the identity map $\operatorname{id}_{X}$ fits, necessarily $\varphi=\operatorname{id}_{X}$. Similarly, given the compatible family of inclusions $Y \rightarrow C^{k}[a, b]$, there is unique $f: Y \rightarrow X$ fitting into


Similarly, given the compatible family of inclusions $X \rightarrow C^{k}[a, b]$, there is unique $g: X \rightarrow Y$ fitting into


Then $f \circ g: X \rightarrow X$ fits into a diagram


Therefore, $f \circ g=\mathrm{id}_{X}$. Similarly, $g \circ f=\mathrm{id}_{Y}$. That is, $f, g$ are mutual inverses, so are isomorphisms of topological vector spaces.
///
Existence of a topology on $C^{\infty}[a, b]$ satisfying the condition above will be proven by identifying $C^{\infty}[a, b]$ as the obvious diagonal closed subspace of the topological product of the limitands $C^{k}[a, b]$ :

$$
C^{\infty}[a, b]=\left\{\left\{f_{k}: f_{k} \in C^{k}[a, b]\right\}: f_{k}=f_{k+1} \text { for all } k\right\}
$$

An arbitrary product of topological spaces $X_{\alpha}$ for $\alpha$ in an index set $A$ is a topological space $X$ with (projections) $p_{\alpha}: X \rightarrow X_{\alpha}$, such that every family $f_{\alpha}: Z \rightarrow X_{\alpha}$ of maps from any other topological
space $Z$ factors through the $p_{\alpha}$ uniquely, in the sense that there is a unique $f: Z \rightarrow X$ such that $f_{\alpha}=p_{\alpha} \circ f$ for all $\alpha$. Pictorially, all triangles commute in the diagram


A similar argument to that for uniqueness of limits proves uniqueness of products up to unique isomorphism. Construction of products is by putting the usual product topology with basis consisting of products $\prod_{\alpha} Y_{\alpha}$ with $Y_{\alpha}=X_{\alpha}$ for all but finitely-many indices, on the Cartesian product of the sets $X_{\alpha}$, whose existence we grant ourselves. Proof that this usual is $a$ product amounts to unwinding the definitions. By uniqueness, in particular, despite the plausibility of the box topology on the product, it cannot function as a product topology since it differs from the standard product topology in general.
[13.2.2] Claim: Giving the diagonal copy of $C^{\infty}[a, b]$ inside $\prod_{k} C^{k}[a, b]$ the subspace topology yields a (projective) limit topology.
Proof: The projection maps $p_{k}: \prod_{j} C^{j}[a, b] \rightarrow C^{k}[a, b]$ from the whole product to the factors $C^{k}[a, b]$ are continuous, so their restrictions to the diagonally imbedded $C^{\infty}[a, b]$ are continuous. Further, letting $i_{k}: C^{k}[a, b] \rightarrow C^{k-1}[a, b]$ be the inclusion, on that diagonal copy of $C^{\infty}[a, b]$ we have $i_{k} \circ p_{k}=p_{k-1}$ as required.

On the other hand, any family of maps $\varphi_{k}: Z \rightarrow C^{k}[a, b]$ induces a map $\widetilde{\varphi}: Z \rightarrow \prod C^{k}[a, b]$ such that $p_{k} \circ \widetilde{\varphi}=\varphi_{k}$, by the property of the product. Compatibility $i_{k} \circ \varphi_{k}=\varphi_{k-1}$ implies that the image of $\widetilde{\varphi}$ is inside the diagonal, that is, inside the copy of $C^{\infty}[a, b]$.

A countable product of metric spaces $X_{k}$ with metrics $d_{k}$ has no canonical single metric, but is metrizable. One of many topologically equivalent metrics is the usual

$$
d\left(\left\{x_{k}\right\},\left\{y_{k}\right\}\right)=\sum_{k=0}^{\infty} 2^{-k} \frac{d_{k}\left(x_{k}, y_{k}\right)}{d_{k}\left(x_{k}, y_{k}\right)+1}
$$

When the metric spaces $X_{k}$ are complete, the product is complete. A closed subspace of a complete metrizable space is complete metrizable, so we have
[13.2.3] Corollary: $C^{\infty}[a, b]$ is complete metrizable.
Abstracting the above, for a (not necessarily countable) family

$$
\ldots \xrightarrow{\varphi_{2}} B_{1} \xrightarrow{\varphi_{1}} B_{o}
$$

of Banach spaces with continuous linear transition maps as indicated, not recessarily requiring the continuous linear maps to be injective (or surjective), a (projective) limit $\lim _{i} B_{i}$ is a topological vector space with continuous linear maps $\lim _{i} B_{i} \rightarrow B_{j}$ such that, for every compatible family of continuous linear maps $Z \rightarrow B_{i}$ there is unique continuous linear $Z \rightarrow \lim _{i} B_{i}$ fitting into


The same uniqueness proof as above shows that there is at most one topological vector space $\lim _{i} B_{i}$. For existence by construction, the earlier argument needs only minor adjustment. The conclusion of complete metrizability would hold when the family is countable.

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Before declaring $C^{\infty}[a, b]$ to be a Fréchet space, we must certify that it is locally convex, in the sense that every point has a local basis of convex opens. Normed spaces are immediately locally convex, because open balls are convex: for $0 \leq t \leq 1$ and $x, y$ in the $\varepsilon$-ball at 0 in a normed space,

$$
|t x+(1-t) y| \leq|t x|+|(1-t) y| \leq t|x|+(1-t)|y|<t \cdot \varepsilon+(1-t) \cdot \varepsilon=\varepsilon
$$

Product topologies of locally convex vectorspaces are locally convex, from the construction of the product. The construction of the limit as the diagonal in the product, with the subspace topology, shows that it is locally convex. In particular, countable limits of Banach spaces are locally convex, hence, are Fréchet. All spaces of practical interest are locally convex for simple reasons, so demonstrating local convexity is rarely interesting.
[13.2.4] Theorem: $\frac{d}{d x}: C^{\infty}[a, b] \rightarrow C^{\infty}[a, b]$ is continuous.
Proof: In fact, the differentiation operator is characterized via the expression of $C^{\infty}[a, b]$ as a limit. We already know that differentiation $d / d x$ gives a continuous map $C^{k}[a, b] \rightarrow C^{k-1}[a, b]$. Differentiation is compatible with the inclusions among the $C^{k}[a, b]$. Thus, we have a commutative diagram


Composing the projections with $d / d x$ gives (dashed) induced maps from $C^{\infty}[a, b]$ to the limitands, inducing a unique (dotted) continuous linear map to the limit, as in


This proves the continuity of differentiation in the limit topology.
In a slightly different vein, we have
[13.2.5] Claim: For all $x \in[a, b]$ and for all non-negative integers $k$, the evaluation map $f \rightarrow f^{(k)}(x)$ is a continuous linear map $C^{\infty}[a, b] \rightarrow \mathbb{C}$.
Proof: The inclusion $C^{\infty}[a, b] \rightarrow C^{k}[a, b]$ is continuous, and the evaluation of the $k^{t h}$ derivative is continuous.

### 13.3 Sufficient notion of topological vector space

To describe a (projective) limit by characterizing its behavior in relation to all topological vectorspaces requires specification of what a topological vectorspace should be.

A topological vector space $V$ (over $\mathbb{C}$ ) is a $\mathbb{C}$-vectorspace $V$ with a topology on $V$ in which points are closed, and so that scalar multiplication

$$
x \times v \longrightarrow x v \quad(\text { for } x \in k \text { and } v \in V)
$$

and vector addition

$$
v \times w \longrightarrow v+w \quad(\text { for } v, w \in V)
$$

are continuous. For subsets $X, Y$ of $V$, let

$$
X+Y=\{x+y: x \in X, y \in Y\}
$$

and

$$
-X=\{-x: x \in X\}
$$

The following trick is elementary, but indispensable. Given an open neighborhood $U$ of 0 in a topological vectorspace $V$, continuity of vector addition yields an open neighborhood $U^{\prime}$ of 0 such that

$$
U^{\prime}+U^{\prime} \subset U
$$

Since $0 \in U^{\prime}$, necessarily $U^{\prime} \subset U$. This can be repeated to give, for any positive integer $n$, an open neighborhood $U_{n}$ of 0 such that

$$
\underbrace{U_{n}+\ldots+U_{n}}_{n} \subset U
$$

In a similar vein, for fixed $v \in V$ the map $V \rightarrow V$ by $x \rightarrow x+v$ is a homeomorphism, being invertible by the obvious $x \rightarrow x-v$. Thus, the open neighborhoods of $v$ are of the form $v+U$ for open neighborhoods $U$ of 0 . In particular, a local basis at 0 gives the topology on a topological vectorspace.
[13.3.1] Lemma: Given a compact subset $K$ of a topological vectorspace $V$ and a closed subset $C$ of $V$ not meeting $K$, there is an open neighborhood $U$ of 0 in $V$ such that

$$
\operatorname{closure}(K+U) \cap(C+U)=\phi
$$

Proof: Since $C$ is closed, for $x \in K$ there is a neighborhood $U_{x}$ of 0 such that the neighborhood $x+U_{x}$ of $x$ does not meet $C$. By continuity of vector addition

$$
V \times V \times V \rightarrow V \quad \text { by } \quad v_{1} \times v_{2} \times v_{3} \rightarrow v_{1}+v_{2}+v_{3}
$$

there is a smaller open neighborhood $N_{x}$ of 0 so that

$$
N_{x}+N_{x}+N_{x} \subset U_{x}
$$

By replacing $N_{x}$ by $N_{x} \cap-N_{x}$, which is still an open neighborhood of 0 , suppose that $N_{x}$ is symmetric in the sense that $N_{x}=-N_{x}$.

Using this symmetry,

$$
\left(x+N_{x}+N_{x}\right) \cap\left(C+N_{x}\right)=\phi
$$

Since $K$ is compact, there are finitely-many $x_{1}, \ldots, x_{n}$ such that

$$
K \subset\left(x_{1}+N_{x_{1}}\right) \cup \ldots \cup\left(x_{n}+N_{x_{n}}\right)
$$

Let $U=\bigcap_{i} N_{x_{i}}$. Since the intersection is finite, $U$ is open. Then

$$
K+U \subset \bigcup_{i=1, \ldots, n}\left(x_{i}+N_{x_{i}}+U\right) \subset \bigcup_{i=1, \ldots, n}\left(x_{i}+N_{x_{i}}+N_{x_{i}}\right)
$$

These sets do not meet $C+U$, by construction, since $U \subset N_{x_{i}}$ for all $i$. Finally, since $C+U$ is a union of opens $y+U$ for $y \in C$, it is open, so even the closure of $K+U$ does not meet $C+U$.

Conveniently, Hausdorff-ness of topological vectorspaces follows from the weaker assumption that points are closed:
[13.3.2] Corollary: A topological vectorspace is Hausdorff.
Proof: Take $K=\{x\}$ and $C=\{y\}$ in the lemma.
[13.3.3] Corollary: The topological closure $\bar{E}$ of a subset $E$ of a topological vectorspace $V$ can be expressed as

$$
\bar{E}=\bigcap_{U}(E+U) \quad(\text { where } U \text { ranges over a local basis at } 0)
$$

Proof: In the lemma, take $K=\{x\}$ and $C=\bar{E}$ for a point $x$ of $V$ not in $C$. Then we obtain an open neighborhood $U$ of 0 so that $x+U$ does not meet $\bar{E}+U$. The latter contains $E+U$, so certainly $x \notin E+U$. That is, for $x$ not in the closure, there is an open $U$ containing 0 so that $x \notin E+U$.

As usual, for two topological vectorspaces $V, W$ over $\mathbb{C}$, a function $f: V \longrightarrow W$ is ( $k$-)linear when $f(\alpha x+\beta y)=\alpha f(x)+\beta f(y)$ for all $\alpha, \beta \in k$ and $x, y \in V$. Almost without exception we care about continuous linear maps, meaning linear maps continuous for the topologies on $V, W$. As expected, the kernel ker $f$ of a linear map is

$$
\operatorname{ker} f=\{v \in V: f(v)=0\}
$$

Being the inverse image of a closed set by a continuous map, the kernel is a closed subspace of $V$.
For a closed subspace $H$ of a topological vectorspace $V$, the quotient $V / H$ is characterized as topological vectorspace with linear quotient $\operatorname{map} q: V \rightarrow V / H$ through which any continuous $f: V \rightarrow W$ with ker $f \supset H$ factors, in the sense that there is a unique continuous linear $\bar{f}: V / H \rightarrow W$ giving a commutative diagram


Uniqueness of the quotient $q: V \rightarrow V / H$, up to unique isomorphism, follows by the usual categorical arguments, as with limits and products above. The existence of the quotient is proven by the usual construction of $V / H$ as the collection of cosets $v+H$, with $q$ given as usual by $q: v \longrightarrow v+H$. We verify that this construction succeeds in the proposition below.

The quotient topology on $V / H$ is the finest topology such that the quotient map $q: V \rightarrow V / H$ is continuous, namely, a subset $E$ of $V / H$ is open if and only if $q^{-1}(E)$ is open.

For non-closed subspaces $H$, the quotient topology on the collection of cosets $\{v+H\}$ would not be Hausdorff. Thus, the proper categorical notion of topological vectorspace quotient, by non-closed subspace, would produce the collection of cosets $v+\bar{H}$ for the closure $\bar{H}$ of $H$.
[13.3.4] Claim: For a closed subspace $W$ of a topological vectorspace $V$, the collection $Q=\{v+W: v \in V\}$ of cosets by $W$ with $\operatorname{map} q(v)=v+W$ is a topological vectorspace and $q$ is a quotient map.
Proof: The algebraic quotient $Q=V / W$ of cosets $v+W$ and $q(v)=v+W$ constructs a vectorspace quotient without any topological hypotheses on $W$. Since $W$ is closed, and since vector addition is a homeomorphism, $v+W$ is closed as well. Thus, its complement $V-(v+W)$ is open, so $q(V-(v+W))$ is open, by definition of the quotient topology. Thus, the complement

$$
q(v)=v+W=q(v+W)=V / W-q(V-(v+W))
$$

of the open set $q(V-(v+W))$ is closed.
Unlike general topological quotient maps,
[13.3.5] Claim: For a closed subspace $H$ of a topological vector space $V$, the quotient map $q: V \rightarrow V / H$ is open, that is, carries open sets to open sets.
Proof: For $U$ open in $V$,

$$
q^{-1}(q(U))=q^{-1}(U+H)=U+H=\bigcup_{h \in H} h+U
$$

This is a union of opens.
[13.3.6] Corollary: For $f: V \rightarrow X$ a linear map with a closed subspace $W$ of $V$ contained in ker $f$, and $\bar{f}$ the induced map $\bar{f}: V / W \rightarrow X$ defined by $\bar{f}(v+W)=f(v), f$ is continuous if and only if $\bar{f}$ is continuous. Proof: Certainly if $\bar{f}$ is continuous then $f=\bar{f} \circ q$ is continuous. The converse follows from the fact that $q$ is open.

This proves that the construction of the quotient by cosets succeeds in producing a quotient: a continuous linear map $f: V \rightarrow X$ factors through any quotient $V / W$ for $W$ a closed subspace contained in the kernel of $f$.

The notions of balanced subset, absorbing subset, directed set, Cauchy net, and completeness are necessary: A subset $E$ of $V$ is balanced when $x E \subset E$ for every $x \in \mathbb{C}$ with $|x| \leq 1$.
[13.3.7] Lemma: Every neighborhood $u$ of 0 in a topological vectorspace $V$ over $\mathbb{C}$ contains a balanced neighborhood $N$ of 0 .
Proof: By continuity of scalar multiplication, there is $\varepsilon>0$ and a neighborhood $U^{\prime}$ of $0 \in V$ so that if $|x|<\varepsilon$ and $v \in U^{\prime}$ then $x v \in U$. Since $\mathbb{C}$ is not discrete, there is $x_{o} \in \mathbb{C}$ with $0<\left|x_{o}\right|<\varepsilon$. Since scalar multiplication by a non-zero element is a homeomorphism, $x_{o} U^{\prime}$ is a neighborhood of 0 and $x_{o} U^{\prime} \subset U$. Put

$$
N=\bigcup_{|y| \leq 1} y x_{o} U^{\prime}
$$

For $|x| \leq 1,|x y| \leq|y| \leq 1$, so

$$
x N=\bigcup_{|y| \leq 1} x\left(y x_{o} U^{\prime}\right) \subset \bigcup_{|y| \leq 1} y x_{o} U^{\prime}=N
$$

producing the desired $N$.
A subset $E$ of vectorspace $V$ over $k$ is absorbing when for every $v \in V$ there is $t_{o} \in R$ so that $v \in \alpha E$ for every $\alpha \in k$ so that $|\alpha| \geq t_{o}$.
[13.3.8] Lemma: Every neighborhood $U$ of 0 in a topological vectorspace is absorbing.
Proof: We may shrink $U$ to assume $U$ is balanced. By continuity of the map $k \rightarrow V$ given by $\alpha \rightarrow \alpha v$, there is $\varepsilon>0$ so that $|\alpha|<\varepsilon$ implies $\alpha v \in U$. By the non-discreteness of $k$, there is non-zero $\alpha \in k$ satisfying any such inequality. Then $v \in \alpha^{-1} U$, as desired.

A poset $S, \leq$ is a partially ordered set. A directed set is a poset $S$ such that, for any two elements $s, t \in S$, there is $z \in S$ so that $z \geq s$ and $z \geq t$.

A net in $V$ is a subset $\left\{x_{s}: s \in S\right\}$ of $V$ indexed by a directed set $S$. A net $\left\{x_{s}: s \in S\right\}$ in a topological vectorspace $V$ is a Cauchy net if, for every neighborhood $U$ of 0 in $V$, there is an index $s_{o}$ so that for $s, t \geq s_{o}$ we have $x_{s}-x_{t} \in U$. A net $\left\{x_{s}: s \in S\right\}$ is convergent if there is $x \in V$ so that, for every neighborhood $U$ of 0 in $V$ there is an index $s_{o}$ so that for $s \geq s_{o}$ we have $x-x_{s} \in U$. Since points are closed, there can be at most one point to which a net converges. Thus, a convergent net is Cauchy. Oppositely, a topological vectorspace is complete if every Cauchy net is convergent.
[13.3.9] Lemma: Let $Y$ be a vector subspace of a topological vector space $X$, complete when given the subspace topology from $X$. Then $Y$ is a closed subset of $X$.
Proof: Let $x \in X$ be in the closure of $Y$. Let $S$ be a local basis of opens at 0 , where we take the partial ordering so that $U \geq U^{\prime}$ if and only if $U \subset U^{\prime}$. For each $U \in S$ choose $y_{U} \in(x+U) \cap Y$. The net $\left\{y_{U}: U \in S\right\}$ converges to $x$, so is Cauchy. It must converge to a point in $Y$, so by uniqueness of limits of nets it must be that $x \in Y$. Thus, $Y$ is closed.

Unfortunately, completeness as above is too strong a condition for general topological vectorspaces, beyond Fréchet spaces. A slightly weaker version of completeness, quasi-completeness or local completeness, does hold for most important natural spaces, as discussed in [13.12].

### 13.4 Unique vectorspace topology on $\mathbb{C}^{n}$

Finite-dimensional topological vectorspaces, and their interactions with other topological vectorspaces, are especially simple:
[13.4.1] Theorem: A finite-dimensional complex vectorspace $V$ has just one topological vectorspace topology, that of the product topology on $\mathbb{C}^{n}$ for $n=\operatorname{dim} V$. A finite-dimensional subspace $V$ of a topological vectorspace $W$ is closed. A $\mathbb{C}$-linear map $X \rightarrow V$ to a finite-dimensional space $V$ is continuous if and only if the kernel is closed.

Proof: The argument is by induction. First treat the one-dimensional situation:
[13.4.2] Claim: For a one-dimensional topological vectorspace $V$ with basis $e$ the map $\mathbb{C} \rightarrow V$ by $x \rightarrow x e$ is a homeomorphism.
Proof: Since scalar multiplication is continuous, we need only show that the map is open. We need only do this at 0 , since translation addresses other points. Given $\varepsilon>0$, by the non-discreteness of $\mathbb{C}$ there is $x_{o}$ in $\mathbb{C}$ so that $0<\left|x_{o}\right|<\varepsilon$. Since $V$ is Hausdorff, there is a neighborhood $U$ of 0 so that $x_{o} e \notin U$. Shrink $U$ so it is balanced. Take $x \in k$ so that $x e \in U$. For $|x| \geq\left|x_{o}\right|,\left|x_{o} x^{-1}\right| \leq 1$, so

$$
x_{o} e=\left(x_{o} x^{-1}\right)(x e) \in U
$$

by balanced-ness of $U$, contradiction. Thus, $x e \in U$ implies that $|x|<\left|x_{o}\right|<\varepsilon$.
[13.4.3] Corollary: For fixed $x_{o} \in \mathbb{C}$, a not-identically-zero $\mathbb{C}$-linear $\mathbb{C}$-valued function $f$ on $V$ is continuous if and only if the affine hyperplane $H=\left\{v \in V: f(v)=x_{o}\right\}$ is closed in $V$.
Proof: Certainly if $f$ is continuous then $H$ is closed. For the converse, consider only the case $x_{o}=0$, since translations (vector additions) are homeomorphisms of $V$ to itself.

For $v_{o}$ with $f\left(v_{o}\right) \neq 0$ and for any other $v \in V$

$$
f\left(v-f(v) f\left(v_{o}\right)^{-1} v_{o}\right)=f(v)-f(v) f\left(v_{o}\right)^{-1} f\left(v_{o}\right)=0
$$

Thus, $V / H$ is one-dimensional. The induced $\mathbb{C}$-linear map $\bar{f}: V / H \rightarrow k$ so that $f=\bar{f} \circ q$, that is, $\bar{f}(v+H)=f(v)$, is a homeomorphism to $\mathbb{C}$, by the previous result, so $f$ is continuous.

For the theorem, for uniqueness of the topology it suffices to prove that for any $\mathbb{C}$-basis $e_{1}, \ldots, e_{n}$ for $V$, the map $\mathbb{C} \times \ldots \times \mathbb{C} \longrightarrow V$ by

$$
\left(x_{1}, \ldots, x_{n}\right) \longrightarrow x_{1} e_{1}+\ldots+x_{n} e_{n}
$$

is a homeomorphism. Prove this by induction on the dimension $n$, that is, on the number of generators for $V$ as a free $\mathbb{C}$-module. The case $n=1$ was treated. Since $\mathbb{C}$ is complete, the lemma above asserting the closed-ness of complete subspaces shows that any one-dimensional subspace is closed. For $n>1$, let $H=\mathbb{C} e_{1}+\ldots+\mathbb{C} e_{n-1}$. By induction, $H$ is closed in $V$, so the quotient $q: V \rightarrow V / H$ is constructed as expected, as the set of cosets $v+H$. The space $V / H$ is a one-dimensional topological vectorspace over $\mathbb{C}$, with basis $q\left(e_{n}\right)$. By induction, $\phi: x q\left(e_{n}\right)=q\left(x e_{n}\right) \longrightarrow x$ is a homeomorphism $V / H \rightarrow \mathbb{C}$.

Likewise, $\mathbb{C} e_{n}$ is a closed subspace and we have the quotient map

$$
q^{\prime}: V \longrightarrow V / \mathbb{C} e_{n}
$$

The image has basis $q^{\prime}\left(e_{1}\right), \ldots, q^{\prime}\left(e_{n-1}\right)$, and by induction

$$
\phi^{\prime}: x_{1} q^{\prime}\left(e_{1}\right)+\ldots+x_{n-1} q^{\prime}\left(e_{n-1}\right) \rightarrow\left(x_{1}, \ldots, x_{n-1}\right)
$$

is a homeomorphism. By the induction hypothesis,

$$
v \longrightarrow(\phi \circ q)(v) \times\left(\phi^{\prime} \circ q^{\prime}\right)(v)
$$

is continuous to $\mathbb{C}^{n-1} \times \mathbb{C} \approx \mathbb{C}^{n}$. On the other hand, by the continuity of scalar multiplication and vector addition,

$$
\mathbb{C}^{n} \longrightarrow V \quad \text { by } \quad x_{1} \times \ldots \times x_{n} \longrightarrow x_{1} e_{1}+\ldots+x_{n} e_{n}
$$

is continuous. These two maps are mutual inverses, certifying the homeomorphism.
Thus, a $n$-dimensional subspace is homeomorphic to $\mathbb{C}^{n}$ with its product topology, so is complete, since a finite product of complete spaces is complete. By the closed-ness of complete subspaces, it is closed.

Continuity of a linear map $f: X \rightarrow \mathbb{C}^{n}$ implies that the kernel $N=\operatorname{ker} f$ is closed. On the other hand, for $N$ closed, the set of cosets $x+N$ constructs a quotient, and is a topological vectorspace of dimension at most $n$. Therefore, the induced map $\bar{f}: X / N \rightarrow V$ is unavoidably continuous. Then $f=\bar{f} \circ q$ is continuous, where $q$ is the quotient map. This completes the induction step.

### 13.5 Non-Banach limits $C^{k}(\mathbb{R}), C^{\infty}(\mathbb{R})$ of Banach spaces $C^{k}[a, b]$

For a non-compact topological space such as $\mathbb{R}$, the space $C^{o}(\mathbb{R})$ of continuous functions is not a Banach space with sup norm, because the sup of the absolute value of a continuous function may be $+\infty$.

But, $C^{o}(\mathbb{R})$ has a Fréchet-space structure: express $\mathbb{R}$ as a countable union of compact subsets $K_{n}=[-n, n]$. Despite the likely non-injectivity of the map $C^{o}(\mathbb{R}) \rightarrow C^{o}\left(K_{i}\right)$, giving $C^{o}(\mathbb{R})$ the (projective) limit topology $\lim _{i} C^{o}\left(K_{i}\right)$ is reasonable: certainly the restriction map $C^{o}(\mathbb{R}) \rightarrow C^{o}\left(K_{i}\right)$ should be continuous, as should all the restrictions $C^{o}\left(K_{i}\right) \rightarrow C^{o}\left(K_{i-1}\right)$, whether or not these are surjective.

The argument in favor of giving $C^{o}(\mathbb{R})$ the limit topology is that a compatible family of maps $f_{i}: Z \rightarrow$ $C^{o}\left(K_{i}\right)$ gives compatible fragments of functions $F$ on $\mathbb{R}$. That is, for $z \in Z$, given $x \in \mathbb{R}$ take $K_{i}$ such that $x$ is in the interior of $K_{i}$. Then for all $j \geq i$ the function $x \rightarrow f_{j}(z)(x)$ is continuous near $x$, and the compatibility assures that all these functions are the same.

That is, the compatibility of these fragments is exactly the assertion that they fit together to make a function $x \rightarrow F_{z}(x)$ on the whole space $X$. Since continuity is a local property, $x \rightarrow F_{z}(x)$ is in $C^{o}(X)$. Further, there is just one way to piece the fragments together. Thus, diagrammatically,


Thus, $C^{o}(X)=\lim _{n} C^{o}\left(K_{n}\right)$ is a Fréchet space. Similarly, $C^{k}(\mathbb{R})=\lim _{n} C^{k}\left(K_{n}\right)$ is a Fréchet space.
[13.5.1] Remark: The question of whether the restriction maps $C^{o}\left(K_{n}\right) \rightarrow C^{o}\left(K_{n-1}\right)$ or $C^{o}(\mathbb{R}) \rightarrow C^{o}\left(K_{n}\right)$ are surjective need not be addressed.

Unsurprisingly, we have
[13.5.2] Theorem: $\frac{d}{d x}: C^{k}(\mathbb{R}) \rightarrow C^{k-1}(\mathbb{R})$ is continuous.
Proof: The argument is structurally similar to the argument for $\frac{d}{d x}: C^{\infty}[a, b] \rightarrow C^{\infty}[a, b]$. The differentiations $\frac{d}{d x}: C^{k}\left(K_{n}\right) \rightarrow C^{k-1}\left(K_{n}\right)$ are a compatible family, fitting into a commutative diagram


Composing the projections with $d / d x$ gives (dashed) induced maps from $C^{k}(\mathbb{R})$ to the limitands, inducing a unique (dotted) continuous linear map to the limit, as in


That is, there is a unique continuous linear map $\frac{d}{d x}: C^{k}(\mathbb{R}) \rightarrow C^{k-1}(\mathbb{R})$ compatible with the differentiations on finite intervals.

Similarly,
[13.5.3] Theorem: $C^{\infty}(\mathbb{R})=\lim _{k} C^{k}(\mathbb{R})$, also $C^{\infty}(\mathbb{R})=\lim _{n} C^{\infty}\left(K_{n}\right)$, and $\frac{d}{d x}: C^{\infty}(\mathbb{R}) \rightarrow C^{\infty}(\mathbb{R})$ is continuous.
Proof: From $C^{\infty}(\mathbb{R})=\lim _{k} C^{k}(\mathbb{R})$ we can obtain the induced map $d / d x$, as follows. Starting with the commutative diagram


Composing the projections with $d / d x$ gives (dashed) induced maps from $C^{k}(\mathbb{R})$ to the limitands, inducing a unique (dotted) continuous linear map to the limit, as in


A novelty is the assertion that (projective) limits commute with each other, so that the limits of $C^{k}\left(K_{n}\right)$ in $k$ and in $n$ can be taken in either order. Generally, in a situation

the maps $\lim _{j}\left(\lim _{i} V_{i j}\right) \rightarrow V_{k \ell}$ induce a map $\lim _{j}\left(\lim _{i} V_{i j}\right) \rightarrow \lim _{\ell} V_{k \ell}$, which induce a unique
$\lim _{j}\left(\lim _{i} V_{i j}\right) \rightarrow \lim _{k}\left(\lim _{\ell} V_{k \ell}\right)$. Similarly, a unique map is induced in the opposite direction, and, for the usual reason, these are mutual inverses.
[13.5.4] Claim: For fixed $x \in \mathbb{R}$ and fixed non-negative integer $k$, the evaluation map $f \rightarrow f^{(k)}(x)$ is continuous.
Proof: Take $n$ large enough so that $x \in[-n, n]$. Evaluation $f \rightarrow f^{(k)}(x)$ was shown in [13.1] to be continuous on $C^{k}[-n, n]$. Composing with the continuous $C^{\infty}(\mathbb{R}) \rightarrow C^{k}(\mathbb{R}) \rightarrow C^{k}[-n, n]$ gives the continuity.

### 13.6 Banach completion $C_{o}^{k}(\mathbb{R})$ of $C_{c}^{k}(\mathbb{R})$

It is reasonable to ask about the completion of the space $C_{c}^{o}(\mathbb{R})$ of compactly-supported continuous functions in the metric given by the sup-norm, and, more generally, about the completion of the space $C_{c}^{k}(\mathbb{R})$ of compactly-supported $k$-times continuously differentiable functions in the metric given by the sum of the sups of the $k$ derivatives.

The spaces $C_{c}^{k}(\mathbb{R})$ are not complete with those norms, because supports can leak out to infinity: for example, in fix any $u$ such that $u(x)=1$ for $|x| \leq 1,0 \leq u(x) \leq 1$ for $1 \leq|x| \leq 2$, and $u(x)=0$ for $|x| \geq 2$. Then

$$
f(x)=\sum_{n=0}^{\infty} \frac{u(x-n)}{n^{2}}
$$

converges in sup-norm, the partial sums have compact support, but the whole does not have compact support.
[13.6.1] Claim: The completion of the space $C_{c}^{o}(\mathbb{R})$ of compactly-supported continuous functions in the metric given by the sup-norm $|f|_{C^{o}}=\sup _{x \in \mathbb{R}}|f(x)|$ is the space $C_{o}^{o}(\mathbb{R})$ of continuous functions $f$ vanishing at infinity, in the sense that, given $\varepsilon>0$, there is a compact interval $K=[-N, N] \subset X$ such that $|f(x)|<\varepsilon$ for $x \notin K$.
[13.6.2] Remark: Since we need to distinguish compactly-supported functions $C_{c}^{o}(\mathbb{R})$ from functions $C_{o}^{o}(\mathbb{R})$ going to 0 at infinity, we cannot use the latter notation for the former, unlike some sources.
Proof: This is almost a tautology. Given $f \in C_{o}^{o}(\mathbb{R})$, given $\varepsilon>0$, let $K=[-N, N] \subset X$ be compact such that $|f(x)|<\varepsilon$ for $x \notin K$. It is easy to make an auxiliary function $\varphi$ that is continuous, compactly-supported, real-valued function such that $\varphi=1$ on $K$ and $0 \leq \varphi \leq 1$ on $X$. Then $f-\varphi \cdot f$ is 0 on $K$, and of absolute value $|\varphi(x) \cdot f(x)| \leq|f(x)|<\varepsilon$ off $K$. That is, $\sup _{\mathbb{R}}|f-\varphi \cdot f|<\varepsilon$, so $C_{c}^{o}(\mathbb{R})$ is dense in $C_{o}^{o}(\mathbb{R})$.

On the other hand, a sequence $f_{i}$ in $C_{c}^{o}(\mathbb{R})$ that is a Cauchy sequence with respect to sup norm gives a Cauchy sequence in each $C^{o}[a, b]$, and converges uniformly pointwise to a continuous function on $[a, b]$ for every $[a, b]$. Let $f$ be the pointwise limit. Given $\varepsilon>0$ take $i_{o}$ such that $\sup _{x}\left|f_{i}(x)-f_{j}(x)\right|<\varepsilon$ for all $i, j \geq i_{o}$. With $K$ the support of $f_{i_{o}}$,

$$
\sup _{x \notin K}|f(x)| \leq \sup _{x \notin K}\left|f(x)-f_{i_{o}}(x)\right|+\sup _{x \notin K}\left|f_{i_{o}}(x)\right|=\sup _{x \notin K}\left|f(x)-f_{i_{o}}(x)\right|+0 \leq \varepsilon<2 \varepsilon
$$

showing that $f$ goes to 0 at infinity.
[13.6.3] Corollary: Continuous functions vanishing at infinity are uniformly continuous.
Proof: For $f \in C_{o}^{o}(\mathbb{R})$, given $\varepsilon>0$, let $g \in C_{c}^{o}(\mathbb{R})$ be such that sup $|f-g|<\varepsilon$. By the uniform continuity of $g$, there is $\delta>0$ such that $|x-y|<\delta$ implies $|g(x)-g(y)|<\varepsilon$, and

$$
|f(x)-f(y)| \leq|f(x)-g(x)|+|f(y)-g(y)|+|g(x)-g(y)|<3 \varepsilon
$$

as desired.
The arguments for $C^{k}(\mathbb{R})$ are completely parallel: the completion of the space $C_{c}^{k}(\mathbb{R})$ of compactly supported $k$-times continuously differentiable functions is the space $C_{c}^{k}(\mathbb{R})$ of $k$-times continuously differentiable functions whose $k$ derivatives go to zero at infinity. Similarly,
[13.6.4] Corollary: The space of $C^{k}$ functions whose $k$ derivatives all vanish at infinity have uniformly continuous derivatives.
[13.6.5] Claim: The $\operatorname{limit}_{\lim } \lim _{o}^{k}(\mathbb{R})$ is the space $C_{o}^{\infty}(\mathbb{R})$ of smooth functions all whose derivatives go to 0 at infinity. All those derivatives are uniformly continuous.
Proof: As with $C^{\infty}[a, b]=\bigcap_{k} C^{k}[a, b]=\lim _{k} C^{k}[a, b]$, by its very definition $C_{o}^{\infty}(\mathbb{R})$ is the intersection of the Banach spaces $C_{o}^{k}(\mathbb{R})$. For any compatible family $Z \rightarrow C_{o}^{k}(\mathbb{R})$, the compatibility implies that the image of $Z$ is in that intersection.
[13.6.6] Corollary: The space $C_{o}^{\infty}(\mathbb{R})$ is a Fréchet space, so is complete.
Proof: As earlier, countable limits of Banach spaces are Fréchet.
[13.6.7] Remark: In contrast, the space of merely bounded continuous functions does not behave so well. Functions such as $f(x)=\sin \left(x^{2}\right)$ are not uniformly continuous. This has the bad side effect that $\sup _{x}|f(x+h)-f(x)|=1$ for all $h \neq 0$, which means that the translation action of $\mathbb{R}$ on that space of functions is not continuous.

### 13.7 Rapid-decay functions, Schwartz functions

A continuous function $f$ on $\mathbb{R}$ is of rapid decay when

$$
\left.\sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{n} \cdot|f(x)|<+\infty \quad \quad \text { (for every } n=1,2, \ldots\right)
$$

With norm $\nu_{n}(f)=\sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{n} \cdot|f(x)|$, let the Banach space $B_{n}$ be the completion of $C_{c}^{o}(\mathbb{R})$ with respect to the metric $\nu_{n}(f-g)$ associated to $\nu_{n}$.
[13.7.1] Lemma: The Banach space $B_{n}$ is isomorphic to $C_{o}^{o}(\mathbb{R})$ by the map $T: f \rightarrow\left(1+x^{2}\right)^{n} \cdot f$. Thus, $B_{n}$ is the space of continuous functions $f$ such that $\left(1+x^{2}\right)^{n} \cdot f(x)$ goes to 0 at infinity.
Proof: By design, $\nu_{n}(f)$ is the sup-norm of $T f$. Thus, the result [13.6] for $C_{o}^{o}(\mathbb{R})$ under sup-norm gives this lemma.
[13.7.2] Remark: Just as we want the completion $C_{o}^{o}(\mathbb{R})$ of $C_{c}^{o}(\mathbb{R})$, rather than the space of all bounded continuous functions, we want $B_{n}$ rather than the space of all continuous functions $f$ with $\sup _{x}\left(1+x^{2}\right) \cdot|f(x)|<\infty$. This distinction disappears in the limit, but it is only via the density of $C_{c}^{o}(\mathbb{R})$ in every $B_{n}$ that it follows that $C_{c}^{o}(\mathbb{R})$ is dense in the space of continuous functions of rapid decay, in the corollary below.
[13.7.3] Claim: The space of continuous functions of rapid decay on $\mathbb{R}$ is the nested intersection, thereby the limit, of the Banach spaces $B_{n}$, so is Fréchet.
Proof: The key issue is to show that rapid-decay $f$ is a $\nu_{n}$-limit of compactly-supported continuous functions for every $n$. For each fixed $n$ the function $f_{n}=\left(1+x^{2}\right)^{n} f$ is continuous and goes to 0 at infinity. From [13.6], $f_{n}$ is the sup-norm limit of compactly supported continuous functions $F_{n j}$. Then $\left(1+x^{2}\right)^{-n} F_{n j} \rightarrow f$ in the topology on $B_{n}$, and $f \in B_{n}$. Thus, the space of rapid-decay functions lies inside the intersection.

On the other hand, a function $f \in \bigcap_{k} B_{k}$ is continuous. For each $n$, since $\left(1+x^{2}\right)^{n+1}|f(x)|$ is continuous and goes to 0 at infinity, it has a finite $\sup \sigma$, and

$$
\sup _{x}\left(1+x^{2}\right)^{n} \cdot|f(x)|=\sup _{x}\left(1+x^{2}\right)^{-1} \cdot\left(1+x^{2}\right)^{n+1}|f(x)| \leq \sup _{x}\left(1+x^{2}\right)^{-1} \cdot \sigma<+\infty
$$

This holds for all $n$, so $f$ is of rapid decay.
[13.7.4] Corollary: The space $C_{c}^{o}(\mathbb{R})$ is dense in the space of continuous functions of rapid decay.
Proof: That every $B_{n}$ is a completion of $C_{c}^{o}(\mathbb{R})$ is essential for this argument.
Use the model of the limit $X=\lim _{n} B_{n}$ as the diagonal in $\prod_{n} B_{n}$, with the product topology restricted to $X$. Let $p_{n}: \prod_{k} B_{k} \rightarrow B_{n}$ be the projection. Thus, given $x \in X$, there is a basis of neighborhood $N$ of $x$ in $X$ of the form $N=X \cap U$ for an open $U$ in the product of the form $U=\prod_{n} U_{n}$ with all but finitely-many $U_{n}=B_{n}$. Thus, for $y \in C_{c}^{o}(\mathbb{R})$ such that $p_{n}(y) \in p_{n}(N)=p_{n}(U)$ for the finitely-many indices such that $U_{n} \neq B_{n}$, we have $y \in N$. That is, approximating $x$ in only finitely-many of the limitands $B_{n}$ suffices to approximate $x$ in the limit. Thus, density in the limitands $B_{n}$ implies density in the limit.
[13.7.5] Remark: The previous argument applies generally, showing that a common subspace dense in all limitands is dense in the limit.

Certainly the operator of multiplication by $1+x^{2}$ preserves $C_{c}^{o}(\mathbb{R})$, and is a continuous map $B_{n} \rightarrow B_{n-1}$. Much as $d / d x$ was treated earlier,
[13.7.6] Claim: Multiplication by $1+x^{2}$ is a continuous map of the space of continuous rapidly-decreasing functions to itself.
Proof: Let $T$ denote the multiplication by $1+x^{2}$, and let $B=\lim _{n} B_{n}$ be the space of rapid-decay continuous functions. From the commutative diagram

composing the projections with $T$ giving (dashed) induced maps from $B$ to the limitands, inducing a unique (dotted) continuous linear map to the limit, as in

giving the continuous multiplication map on the space of rapid-decay continuous functions.
Similarly, adding differentiability conditions, the space of rapidly decreasing $C^{k}$ functions is the space of $k$-times continuously differentiable functions $f$ such that, for every $\ell=0,1,2, \ldots, k$ and for every $n=1,2, \ldots$,

$$
\sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{n} \cdot\left|f^{(\ell)}(x)\right|<+\infty
$$

Let $B_{n}^{k}$ be the completion of $C_{c}^{k}(\mathbb{R})$ with respect to the metric from the norm

$$
\nu_{n}^{k}(f)=\sum_{0 \leq \ell \leq k} \sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{n}\left|f^{(\ell)}(x)\right|
$$

Essentially identical arguments give
[13.7.7] Claim: The space of $C^{k}$ functions of rapid decay on $\mathbb{R}$ is the nested intersection, thereby the limit, of the Banach spaces $B_{n}^{k}$, so is Fréchet.
[13.7.8] Corollary: The space $C_{c}^{k}(\mathbb{R})$ is dense in the space of $C^{k}$ functions of rapid decay.
Identifying $B_{n}^{k}$ as a space of $C^{k}$ functions with additional decay properties at infinity gives the obvious $\operatorname{map} \frac{d}{d x}: B_{n}^{k} \rightarrow B_{n}^{k-1}$.
[13.7.9] Claim: $\frac{d}{d x}: B_{n}^{k} \rightarrow B_{n}^{k-1}$ is continuous.
Proof: Since $B_{n}^{k}$ is the closure of $C_{c}^{k}(\mathbb{R})$, it suffices to check the continuity of $\frac{d}{d x}: C_{c}^{k}(\mathbb{R}) \rightarrow C_{c}^{k-1}(\mathbb{R})$ for the $B_{n}^{k}$ and $B_{n}^{k-1}$ topologies. As usual, that continuity was designed into the situation.

The space of Schwartz functions is

$$
\mathscr{S}(\mathbb{R})=\{\text { smooth functions } f \text { all whose derivatives are of rapid decay }\}
$$

One reasonable topology on $\mathscr{S}(\mathbb{R})$ is as a limit

$$
\mathscr{S}(\mathbb{R})=\bigcap_{k}\left\{C^{k} \text { functions of rapid decay }\right\}=\lim _{k}\left\{C^{k} \text { functions of rapid decay }\right\}
$$

As a countable limit of Fréchet spaces, this makes $\mathscr{S}(\mathbb{R})$ Fréchet.
[13.7.10] Corollary: $\frac{d}{d x}: \mathscr{S}(\mathbb{R}) \rightarrow \mathscr{S}(\mathbb{R})$ is continuous.
Proof: This is structurally the same as before: letting $B_{\infty}^{k}$ be the space of $C^{k}$ functions of rapid decay, from the commutative diagram

composing the projections with $d / d x$ to give (dashed) induced maps from $\mathscr{S}(\mathbb{R})$ to the limitands, inducing a unique (dotted) continuous linear map to the limit:
as desired.


Finally, to induce a canonical continuous map $T: \mathscr{S}(\mathbb{R}) \rightarrow \mathscr{S}(\mathbb{R})$ by multiplication by $1+x^{2}$, examine the behavior of this multiplication map on the auxiliary spaces $B_{n}^{k}$ and its interaction with $\frac{d}{d x}$ :
[13.7.11] Claim: $T: B_{n}^{k} \rightarrow B_{n-1}^{k-1}$ is continuous.
Proof: Of course,

$$
\left|\frac{d}{d x}\left(\left(1+x^{2}\right) \cdot f(x)\right)\right|=\left|2 x \cdot f(x)+\left(1+x^{2}\right) \cdot f^{\prime}(x)\right| \leq 2 \cdot\left(1+x^{2}\right) \cdot|f(x)|+\left(1+x^{2}\right) \cdot\left|f^{\prime}(x)\right|
$$

Thus, $T: C_{c}^{k}(\mathbb{R}) \rightarrow C_{c}^{k-1}(\mathbb{R})$ is continuous with the $B_{n}^{k}$ and $B_{n-1}^{k-1}$ topologies. For general reasons, cofinal limits are isomorphic, so the same argument gives a unique continuous linear map $\mathscr{S}(\mathbb{R})$.

It is worth noting
[13.7.12] Claim: Compactly-supported smooth functions are dense in $\mathscr{S}$.
Proof: At least up to rearranging the order of limit-taking, the description of $\mathscr{S}$ above is as a limit of spaces in each of which compactly-supported smooth functions are dense. Thus, we claim a general result: for a limit $X=\lim _{i} X_{i}$ and compatible maps $f_{i}: V \rightarrow X_{i}$ with dense image, the induced map $f: V \rightarrow X$ has dense image. As in [13.2], the limit is the diagonal

$$
D=\left\{\left\{x_{i}\right\} \in \prod_{i} X_{i}: x_{i} \rightarrow x_{i-1}, \text { for all } i\right\} \subset \prod_{i} X_{i}
$$

with the subspace topology from the product. Suppose we are given a finite collection of neighborhoods $x_{i_{1}} \in U_{i_{1}} \subset X_{i_{1}}, \ldots, x_{i_{n}} \in U_{i_{n}} \subset X_{i_{n}}$, with $x_{i_{j}} \rightarrow x_{i_{k}}$ if $i_{j} \geq i_{k}$. Take $i=\max _{j} i_{j}$, and $U$ a neighborhood of $x_{i}$ such that the image of $U$ is inside every $U_{i_{j}}$, by continuity. Since the image of $V$ is dense in $X_{i}$, there is $v \in V$ such that $f_{i}(v) \in U$. By compatibility, $f_{i_{j}}(v) \in U_{i_{j}}$ for all $j$. Thus, the image of $V$ is dense in the limit.

### 13.8 Non-Fréchet colimit $\mathbb{C}^{\infty}$ of $\mathbb{C}^{n}$, quasi-completeness

Toward topologies in which $C_{c}^{o}(\mathbb{R})$ and $C_{c}^{\infty}(\mathbb{R})$ could be complete, we consider first

$$
\mathbb{C}^{\infty}=\bigcup_{n} \mathbb{C}^{n}
$$

where $i_{n}: \mathbb{C}^{n} \subset C^{n+1}$ by $i_{n}:\left(x_{1}, \ldots, x_{n}\right) \rightarrow\left(x_{1}, \ldots, x_{n}, 0\right)$. We want to topologize $\mathbb{C}^{\infty}$ so that it is complete, in a suitable sense. Above, we saw that finite-dimensional complex vectorspaces have unique vectorspace topologies, so the only question is how to fit them together.

A countable ascending union of complete metric topological vector spaces, each a proper closed subspace of the next, such as $C^{\infty}=\bigcup \mathbb{C}^{n}$, cannot be a complete metric space, because it is exactly presented as a countable union of nowhere-dense closed subsets, contradicting the conclusion of the Baire Category Theorem. The function spaces $C_{c}^{o}(\mathbb{R})$ and $C_{c}^{\infty}(\mathbb{R})$ are also of this type, being the ascending unions of spaces $C_{K}^{o}$ or $C_{K}^{\infty}$, continuous or smooth functions with supports inside compact $K \subset \mathbb{R}$.

Thus, we cannot hope to give such space metric topologies for which they are complete.
Nevertheless, ascending unions are a type of colimit, just as descending intersections are a type of limit. That is, the topology on $\mathbb{C}^{\infty}$ is characterized by a universal property: for every collection of maps $f_{n}: \mathbb{C}^{n} \rightarrow Z$ with the compatibility $i_{n} \circ f_{n}=f_{n+1}$, there is a unique $f: \mathbb{C}^{\infty} \rightarrow Z$ through which all $f_{n}$ 's factor. That is, given a commutative diagram

there is a unique (dotted) map $\mathbb{C}^{\infty} \rightarrow Z$ giving a commutative diagram


To argue that an ascending union $X=\bigcup_{n} X_{n}$ with $X_{1} \subset X_{2} \subset \ldots$ is an example of a colimit, observe that every $x \in X$ lies in some $X_{n}$, so all values $f(x)$ for a map $f: X \rightarrow Z$ are completely determined by the restrictions of $f$ to the limitands $X_{n}$. Thus, on one hand, given a compatible family $f_{n}: X_{n} \rightarrow Z$, there is at most one compatible $f: X \rightarrow Z$. On the other hand, a compatible family $f_{n}: X_{n} \rightarrow Z$ defines a map $X \rightarrow Z$ : given $x \in X$, take $n$ sufficiently large so that $x \in X_{n}$, and define $f(x)=f_{n}(x)$. The compatibility assures that it doesn't matter which sufficiently large $n$ we use.

For the topology of $\mathbb{C}^{\infty}$, the colimit characterization has a possibly-counterintuitive consequence:
[13.8.1] Claim: Every linear map from the space $\mathbb{C}^{\infty}=\operatorname{colim}_{n} \mathbb{C}^{n}$ with the colimit topology to any topological vectorspace is continuous.
Proof: Given arbitrary linear $f: \mathbb{C}^{\infty} \rightarrow Z$, composition with inclusion gives a compatible family of linear $\operatorname{maps} f_{n}: \mathbb{C}^{n} \rightarrow Z$. From [13.4], every linear map from a finite-dimensional space is continuous. The collection $\left\{f_{n}\right\}$ induces a unique continuous map $F: \mathbb{C}^{\infty} \rightarrow Z$ such that $F \circ i_{n}: \mathbb{C}^{n} \rightarrow Z$ is the same as $f \circ i_{n}$. In general, this might not be force $f=F$. However, because $X$ is an ascending union, the values of both $F$ and $f$ are completely determined by their values on the limitands, and these are the same. Thus, $f=F$.

The uniqueness argument for locally convex colimits of locally convex topological vectorspaces, that there is at most one such topology, is identical to the uniqueness argument for limits in [13.2] with arrows reversed.
[13.8.2] Remark: The fact that a colimit of finite-dimensional spaces has a unique canonical topology, from which every linear map from such a colimit is continuous, is often misunderstood and misrepresented as
suggesting that there is no topology on that colimit. Again, there is a unique canonical topology, from which every linear map is continuous.

To prove existence of colimits, just as limits are subobjects of products, colimits are quotients of coproducts, as follows. A locally convex colimit of topological vector spaces $X_{\alpha}$ with transition maps $j_{\beta}^{\alpha}: X_{\alpha} \rightarrow X_{\beta}$ is the quotient of the locally convex coproduct $X$ of the $X_{\alpha}$ by the closure of the subspace $Z$ spanned by vectors

$$
j_{\alpha}\left(x_{\alpha}\right)-\left(j_{\beta} \circ j_{\beta}^{\alpha}\right)\left(x_{\alpha}\right) \quad\left(\text { for all } \alpha<\beta \text { and } x_{\alpha} \in X_{\alpha}\right)
$$

Annihilation of these differences in the quotient forces the desired compatibility relations. Obviously, quotients of locally convex spaces are locally convex.

Locally convex coproducts $X$ of topological vector spaces $X_{\alpha}$ are coproducts (also called direct sums) of the vector spaces $X_{\alpha}$ topologized by the diamond topology, described as follows. ${ }^{[78]}$ For a collection $U_{\alpha}$ of convex neighborhoods of 0 in the $X_{\alpha}$, let

$$
U=\text { convex hull in } X \text { of the union of } j_{\alpha}\left(U_{\alpha}\right) \quad \text { (with } j_{\alpha}: X_{\alpha} \rightarrow X \text { the } \alpha^{t h} \text { canonical map) }
$$

The diamond topology has local basis at 0 consisting of such $U$. Thus, it is locally convex by construction. Closedness of points follows from the corresponding property of the $X_{\alpha}$. Thus, existence of a locally convex coproduct of locally convex spaces is assured by the construction.

A countable colimit of a family $V_{1} \rightarrow V_{2} \rightarrow \ldots$ of topological vectorspaces is a strict colimit, or strict inductive limit, when each $V_{i} \rightarrow V_{i+1}$ is an isomorphism to its image, and each image is closed. A strict colimit of Fréchet spaces is called an LF-space.

Just to be sure:
[13.8.3] Claim: In a colimit indexed by positive integers $V=\operatorname{colim} V_{i}$, if every transition $V_{i} \rightarrow V_{i+1}$ is injective, then every limitand $V_{i}$ injects to the colimit $V$. Further, the colimit is the ascending union of the limitands $V_{i}$, suitably topologized.
Proof: In effect, the argument presents the colimit corresponding to an ascending union more directly, not as a quotient of the coproduct, although it is convenient to already have existence of the colimit. Certainly each $V_{i}$ injects to $W=\bigcup_{n} V_{n}$. We will give $W$ a locally convex topology so that every inclusion $V_{i} \rightarrow W$ is continuous. The universal property of the colimit produces a unique compatible map $V \rightarrow W$, so every $V_{i}$ must inject to $V$ itself.

Since the maps $j_{i}$ of $V_{i}$ to the colimit $V$ are injections, the ascending union $W$ injects to $V$ by $j(w)=j_{i}(w)$ for any index $i$ large enough so that $w \in V_{i}$. The compatibility of the maps among the $V_{i}$ assures that $j$ is welldefined. We claim that $j(W)$ with the subspace topology from $V$, and the inclusions $V_{i} \rightarrow j_{i}\left(V_{i}\right) \subset j(W)$, give a colimit of the $V_{i}$. Indeed for any compatible, family $f_{i}: V_{i} \rightarrow Z$ and induced $f: V \rightarrow Z$, the restriction of $f$ to $j(W)$ gives a map $j(W) \rightarrow Z$ through which the $f_{i}$ factor. Thus, in fact, such a colimit is the ascending union with a suitable topology.

Now we describe a topology on the ascending union $W$ so that all inclusions $V_{i} \rightarrow W$ are continuous. Give $W$ a local basis $\{U\}$ at 0 , by taking arbitrary convex opens $U_{i} \subset V_{i}$ containing 0 , and letting $U$ be the convex hull of $\bigcup_{i} U_{i}$. Every injection $V_{i} \rightarrow W$ is continuous, because the inverse image of such $U \cap V_{i}$ contains $U_{i}$, giving continuity at 0 .

To be sure that points are closed in $W$, given $0 \neq x \in W$, we find a neighborhood of 0 in $W$ not containing $x$. Let $i_{o}$ be the first index such that $x \in V_{i_{o}}$. By Hahn-Banach, there is a continuous linear functional $\lambda_{i_{o}}$ on $V_{i_{o}}$ such that $\lambda_{i_{o}}(x) \neq 0$. Without loss of generality, $\lambda_{i_{o}}(x)=1$ and $\left|\lambda_{i_{o}}\right|=1$. Use Hahn-Banach to extend $\lambda_{i_{o}}$ to a continuous linear functional $\lambda_{i}$ on $V_{i}$ for every $i \geq i_{o}$, with $\left|\lambda_{i}\right| \leq 1$. $\lambda_{i_{o}}$ gives a continuous linear functional on $V_{i}$ for $i<i_{o}$ by composition with the injection $V_{i} \rightarrow V_{i_{o}}$. Then $U_{i}=\left\{y \in V_{i}:\left|\lambda_{i}(y)\right|<1\right\}$ is open in $V_{i}$ and does not contain $x$, for all $i$. The convex hull of the ascending union $\bigcup_{i} U_{i}$ is just $\bigcup_{i} U_{i}$ itself, so does not contain $x$.
[78] The product topology of locally convex topological vector spaces is locally convex, whether in the category of locally convex topological vector spaces or in the larger category of not-necessarily-locally-convex topological vector spaces. However, coproducts behave differently: the locally convex coproduct of uncountably many locally convex spaces is not a coproduct in the larger category of not-necessarily-locally-convex spaces. This already occurs with an uncountable coproduct of lines.

We did not quite prove that this topology is exactly the colimit topology, but we will never need that fact.

Typical colimit topologies are not complete in the strongest possible sense (see below), but are quasicomplete, a property sufficient for all applications. To describe quasi-completeness, we need a notion of boundedness in general topological vectorspaces, not merely metrizable ones. A subset $B$ of a topological vector space $V$ is bounded when, for every open neighborhood $N$ of 0 there is $t_{o}>0$ such that $B \subset t N$ for every $t \geq t_{o}$. A space is quasi-complete when every bounded Cauchy net is convergent.

Nothing new for metric spaces:
[13.8.4] Lemma: Complete metric spaces are quasi-complete. In particular, Cauchy nets converge, and contain cofinal sequences converging to the same limit.
Proof: Let $\left\{s_{i}: i \in I\right\}$ be a Cauchy net in $X$. Given a natural number $n$, let $i_{n} \in I$ be an index such that $d\left(x_{i}, x_{j}\right)<\frac{1}{n}$ for $i, j \geq i_{n}$. Then $\left\{x_{i_{n}}: n=1,2, \ldots\right\}$ is a Cauchy sequence, with limit $x$. Given $\varepsilon>0$, let $j \geq i_{n}$ be also large enough such that $d\left(x, x_{j}\right)<\varepsilon$. Then

$$
d\left(x, x_{i_{n}}\right) \leq d\left(x, x_{j}\right)+d\left(x_{j}, x_{i_{n}}\right)<\varepsilon+\frac{1}{n} \quad(\text { for every } \varepsilon>0)
$$

Thus, $d\left(x, x_{i_{n}}\right) \leq \frac{1}{n}$. The original Cauchy net also converges to $x$ : given $\varepsilon>0$, for $n$ large enough so that $\varepsilon>\frac{1}{n}$,

$$
d\left(x_{i}, x\right) \leq d\left(x_{i}, x_{i_{n}}\right)+d\left(x_{i_{n}}, x\right)<\varepsilon+\varepsilon \quad\left(\text { for } i \geq i_{n}\right)
$$

with the strict inequality coming from $d\left(x_{i_{n}}, x\right)<\varepsilon$.
[13.8.5] Theorem: A bounded subset of an LF-space $X=\operatorname{colim}_{n} X_{n}$ lies in some limitand $X_{n}$. An LF-space is quasi-complete.
Proof: Let $B$ be a bounded subset of $X$. Suppose $B$ does not lie in any $X_{i}$. Then there is a sequence $i_{1}, i_{2}, \ldots$ of positive integers and $x_{i_{\ell}}$ in $X_{i_{\ell}} \cap B$ with $x_{i_{\ell}}$ not lying in $X_{i_{\ell}-1}$. Using $X=\bigcup_{j} X_{i_{\ell}}$, without loss of generality, suppose that $i_{\ell}=\ell$.

By the Hahn-Banach theorem and induction, using the closedness of $X_{i-1}$ in $X_{i}$, there are continuous linear functionals $\lambda_{i}$ on $X_{i}$ 's such that $\lambda_{i}\left(x_{i}\right)=i$ and the restriction of $\lambda_{i}$ to $X_{i-1}$ is $\lambda_{i-1}$, for example. Since $X$ is the colimit of the $X_{i}$, this collection of functionals exactly describes a unique compatible continuous linear functional $\lambda$ on $X$.

But $\lambda(B)$ is bounded since $B$ is bounded and $\lambda$ is continuous, precluding the possibility that $\lambda$ takes on all positive integer values at the points $x_{i}$ of $B$. Thus, it could not have been that $B$ failed to lie inside some single $X_{i}$. The strictness of the colimit implies that $B$ is bounded as a subset of $X_{i}$, proving one direction of the equivalence. The other direction of the equivalence is less interesting.

Thus a bounded Cauchy net lies in some limitand Fréchet space $X_{n}$, so is convergent there, since Fréchet spaces are complete.
[13.8.6] Remark: Strict inductive limits of finite-dimensional spaces do appear as natural function spaces, for example the Schwartz space on $\mathbb{Q}_{p}$, as in [13.17].

### 13.9 Non-Fréchet colimit $C_{c}^{\infty}(\mathbb{R})$ of Fréchet spaces

The space of compactly-supported continuous functions

$$
C_{c}^{o}(\mathbb{R})=\text { compactly-supported continuous functions on } \mathbb{R}
$$

is an ascending union of the subspaces

$$
C_{[-n, n]}^{o}=\left\{f \in C^{o}(\mathbb{R}): \operatorname{spt} f \subset[-n, n]\right\}
$$

Each space $C_{[-n, n]}^{o}$ is a Banach space, being a closed subspace of the Banach space $C^{o}[-n, n]$, further requiring vanishing of the functions on the boundary of $[-n, n]$. A closed subspace of a Banach space is a Banach space. Thus, $C_{c}^{o}(\mathbb{R})$ is an LF-space, and is quasi-complete.

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Similarly,

$$
C_{c}^{k}(\mathbb{R})=\text { compactly-supported } C^{k} \text { functions on } \mathbb{R}
$$

is an ascending union of the subspaces

$$
C_{[-n, n]}^{k}=\left\{f \in C^{k}(\mathbb{R}): \operatorname{spt} f \subset[-n, n]\right\}
$$

Each space $C_{[-n, n]}^{k}$ is a Banach space, being a closed subspace of the Banach space $C^{k}[-n, n]$, further requiring vanishing of the functions and derivatives on the boundary of $[-n, n]$. A closed subspace of a Banach space is a Banach space. Thus, $C_{c}^{k}(\mathbb{R})$ is an LF-space, and is quasi-complete.

The space of test functions is

$$
\mathcal{D}(\mathbb{R})=C_{c}^{\infty}(\mathbb{R})=\text { compactly-supported } C^{\infty} \text { functions on } \mathbb{R}
$$

is an ascending union of the subspaces

$$
\mathcal{D}_{[-n, n]}=C_{[-n, n]}^{\infty}=\left\{f \in C^{\infty}(\mathbb{R}): \operatorname{spt} f \subset[-n, n]\right\}
$$

Each space $\mathcal{D}_{[-n, n]}$ is a Fréchet space, being a closed subspace of the Fréchet space $C^{\infty}[-n, n]$, by further requiring vanishing of the functions and derivatives on the boundary of $[-n, n]$. A closed subspace of a Fréchet space is a Fréchet space. Thus, $\mathcal{D}(\mathbb{R})=C_{c}^{\infty}(\mathbb{R})$ is an LF-space, and is quasi-complete.

The operator $\frac{d}{d x}: C^{k}[-n, n] \rightarrow C^{k-1}[-n, n]$ is continuous, and preserves the vanishing conditions at the endpoints, so restricts to a continuous map $\frac{d}{d x}: C_{[-n, n]}^{k} \rightarrow C_{[-n, n]}^{k-1}$ on the Banach sub-spaces of functions vanishing suitably at the endpoints. Composing with the inclusions $C_{[-n, n]}^{k-1} \rightarrow C_{c}^{k-1}(\mathbb{R})$ gives a compatible family of continuous maps $\frac{d}{d x}: C_{[-n, n]}^{k} \rightarrow C_{c}^{k-1}(\mathbb{R})$. This induces a unique continuous map on the colimit: $\frac{d}{d x}: C_{c}^{k}(\mathbb{R}) \rightarrow C_{c}^{k-1}(\mathbb{R})$.

Similarly, $\frac{d}{d x}: C^{\infty}[-n, n] \rightarrow C^{\infty}[-n, n]$ is continuous, and preserves the vanishing conditions at the endpoints, so restricts to a continuous map $\frac{d}{d x}: \mathcal{D}_{[-n, n]} \rightarrow \mathcal{D}_{[-n, n]}$ on the Frechet sub-spaces of functions vanishing to all orders at the endpoints. Composing with the inclusions $\mathcal{D}_{[-n, n]} \rightarrow \mathcal{D}(\mathbb{R})$ gives a compatible family of continuous maps $\frac{d}{d x}: \mathcal{D}_{[-n, n]} \rightarrow \mathcal{D}(\mathbb{R})$. This induces a unique continuous map on the colimit: $\frac{d}{d x}: \mathcal{D}(R) \rightarrow \mathcal{D}(\mathbb{R})$. Diagrammatically,


That is, $\frac{d}{d x}$ is continuous in the LF-space topology on test functions $\mathcal{D}(\mathbb{R})=C_{c}^{\infty}(\mathbb{R})$.
[13.9.1] Claim: For fixed $x \in \mathbb{R}$ and non-negative integer $k$, the evaluation map $f \rightarrow f^{(k)}(x)$ on $\mathcal{D}(\mathbb{R})=C_{c}^{\infty}(\mathbb{R})$ is continuous.
Proof: This evaluation map is continuous on $C^{\infty}[-n, n]$ for every large-enough $n$ so that $x \in[-n, n]$, so is continuous on the closed subspace $\mathcal{D}_{[-n, n]}$ of $C^{\infty}[-n, n]$. The inclusions among these spaces are extend-by- 0 , so the evaluation map is the 0 map on $\mathcal{D}_{[-n, n]}$ if $|x| \geq n$. These maps to $\mathbb{C}$ fit together into a compatible family, so extend uniquely to a continuous linear map of the colimit $\mathcal{D}(\mathbb{R})$ to $\mathbb{C}$.
[13.9.2] Claim: For $F \in C^{\infty}(\mathbb{R})$, the map $f \rightarrow F \cdot f$ is a continuous map of $\mathcal{D}(\mathbb{R})$ to itself.
Proof: By the colimit characterization, it suffices to show that such a map is continuous on $C_{[-n, n]}^{\infty}$, or on the larger Fréchet space $C^{\infty}[-n, n]$ without vanishing conditions on the boundary. This is the limit of $C^{k}[-n, n]$, so it suffices to show that $f \rightarrow F \cdot f$ is a continuous map $C^{k}[-n, n] \rightarrow C^{k}[-n, n]$ for every $k$. The sum of sups of derivatives is

$$
\sum_{0 \leq i \leq k} \sup _{|x| \leq n}\left|\left(\frac{d}{d x}\right)^{i}(F f)(x)\right| \leq 2^{k}\left(\sum_{0 \leq i \leq k} \sup _{|x| \leq n}\left|F^{(i)}(x)\right|\right) \cdot\left(\sum_{0 \leq i \leq k} \sup _{|x| \leq n}\left|f^{(i)}(x)\right|\right)
$$

Although $F$ and its derivatives need not be bounded, this estimate only uses their boundedness on $[-n, n]$. This is a bad estimate, but sufficient for continuity.
[13.9.3] Claim: The inclusion $\mathcal{D}(\mathbb{R}) \rightarrow \mathscr{S}(\mathbb{R})$ is continuous, and the image is dense.
Proof: At least after changing order of limits, $\mathscr{S}(\mathbb{R})$ is described as a limit of spaces in which $\mathcal{D}(\mathbb{R})$ is dense, so $\mathcal{D}(\mathbb{R})$ is dense in that limit.

The slightly more serious issue is that $\mathcal{D}(\mathbb{R})$ with its LF-space topology maps continuously to $\mathscr{S}(\mathbb{R})$. Since $\mathcal{D}(\mathbb{R})$ is a colimit, we need only check that the limitands (compatibly) map continuously. On a limitand $C_{[-n, n]}^{\infty}$, the norms

$$
\nu_{N, k}(f)=\sup _{x}\left(1+x^{2}\right)^{N} \cdot\left|f^{(k)}(x)\right|
$$

differ from the norms $\sup _{x}\left|f^{(k)}(x)\right|$ defining the topology on $C_{[-n, n]}^{\infty}$ merely by constants, namely, the sups of $\left(1+x^{2}\right)^{N}$ on $[-n, n]$. Thus, we have the desired continuity on the limitands.

### 13.10 LF-spaces of moderate-growth functions

The space $C_{\text {mod }}^{o}(\mathbb{R})$ of continuous functions of moderate growth on $\mathbb{R}$ is

$$
C_{\mathrm{mod}}^{o}(\mathbb{R})=\left\{f \in C^{o}(\mathbb{R}): \sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{-N} \cdot|f(x)|<+\infty \text { for some } N\right\}
$$

Literally, it is an ascending union

$$
C_{\bmod }^{o}(\mathbb{R})=\bigcup_{N}\left\{f \in C^{o}(\mathbb{R}): \sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{-N} \cdot|f(x)|<+\infty\right\}
$$

However, it is ill-advised to use the individual spaces

$$
B_{N}=\left\{f \in C^{o}(\mathbb{R}): \sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{-N} \cdot|f(x)|<+\infty\right\}
$$

with norms $\nu_{N}(f)=\sup _{x \in \mathbb{R}}\left(1+x^{2}\right)^{-N} \cdot|f(x)|$ because $C_{c}^{o}(\mathbb{R})$ is not dense in these spaces $B_{N}$. Indeed, in the simple case $N=0$, the norm $\nu_{0}$ is the sup-norm, and the sup-norm closure of $C_{c}^{o}(\mathbb{R})$ is continuous functions going to 0 at infinity, which excludes many bounded continuous functions.

In particular, there are many bounded continuous functions $f$ which are not uniformly continuous, and the translation action of $\mathbb{R}$ on such functions cannot be continuous: no matter how small $\delta>0$ is, $\sup _{x \in \mathbb{R}}|f(x+\delta)-f(x)|$ may be large. For example, $f(x)=\sin \left(x^{2}\right)$ has this feature.

This difficulty does not mean that the characterization of the whole set of moderate-growth functions is incorrect, nor that the norms $\nu_{N}$ are inappropriate, but only that the Banach spaces $B_{N}$ are too large, and that the topology of the whole should not be the strict colimit of the Banach spaces $B_{N}$. Instead, take the smaller

$$
V_{N}=\text { completion of } C_{c}^{o}(\mathbb{R}) \text { with respect to } \nu_{N}
$$

As in the case of completion of $C_{c}^{o}(\mathbb{R})$ with respect to sup-norm $\nu_{0}$,
[13.10.1] Claim: $V_{N}=\left\{\right.$ continuous $f$ such that $\left(1+x^{2}\right)^{-N} f(x)$ goes to 0 at infinity $\}$.
///
Of course, if $\left(1+x^{2}\right)^{-N} f(x)$ is merely bounded, then $\left(1+x^{2}\right)^{-(N+1)} f(x)$ then goes to 0 at infinity. Thus, as sets, $B_{N} \subset V_{N+1}$, but this inclusion cannot be continuous, since $C_{c}^{o}(\mathbb{R})$ is dense in $V_{N+1}$, but not in $B_{N}$. That is, there is a non-trivial effect on the topology in setting

$$
C_{\mathrm{mod}}^{o}=\operatorname{colim}_{N} V_{N}
$$

instead of the colimit of the too-large spaces $B_{N}$.

### 13.11 Seminorms and locally convex topologies

So far, the vectorspace topologies have been described as Banach spaces, limits of Banach spaces, and colimits of limits of Banach spaces. By design, these descriptions facilitate proof of (quasi-) completeness. Weaker topologies are not usually described in this fashion. For example, for a topological vectorspace $V$, with (continuous) dual

$$
V^{*}=\{\text { continuous linear maps } V \rightarrow \mathbb{C}\}
$$

the weak dual topology ${ }^{[79]}$ on $V^{*}$ has a local sub-basis at 0 consisting of sets

$$
U=U_{v, \varepsilon}=\left\{\lambda \in V^{*}:|\lambda(v)|<\varepsilon\right\} \quad \text { (for fixed } v \in V \text { and } \varepsilon>0 \text { ) }
$$

Unless $V$ is finite-dimensional, this topology on $V^{*}$ is much coarser than a Banach, Fréchet, or LF-topology. The map $\lambda \rightarrow|\lambda(v)|$ is a natural example of a seminorm. It is not a norm, because $\lambda(v)=0$ can easily happen.

Seminorms are a general device to describe topologies on vectorspaces. These topologies are invariably locally convex, in the sense of having a local basis at 0 consisting of convex sets.

Description of a vectorspace topology by seminorms does not generally give direct information about completeness. Nevertheless, we can prove quasi-completeness for an important class of examples, just below.

A seminorm $\nu$ on a complex vectorspace $V$ is a real-valued function on $V$ so that $\nu(x) \geq 0$ for all $x \in V$ (non-negativity), $\nu(\alpha x)=|\alpha| \cdot \nu(x)$ for all $\alpha \in \mathbb{C}$ and $x \in V$ (homogeneity), and $\nu(x+y) \leq \nu(x)+\nu(y)$ for all $x, y \in V$ (triangle inequality). This differs from the notion of norm only in the significant point that we allow $\nu(x)=0$ for $x \neq 0$.

To compensate for the possibility that an individual seminorm can be 0 on a particular non-zero vector, since we want Hausdorff topologies, we mostly care about separating families $\left\{\nu_{i}: i \in I\right\}$ of semi-norms: for every $0 \neq x \in V$ there is $\nu_{i}$ so that $\nu_{i}(x) \neq 0$.
[13.11.1] Claim: The collection $\Phi$ of all finite intersections of sets

$$
U_{i, \varepsilon}=\left\{x \in V: \nu_{i}(x)<\varepsilon\right\} \quad(\text { for } \varepsilon>0 \text { and } i \in I)
$$

is a local basis at 0 for a locally convex topology on $V$.
Proof: As expected, we intend to define a topological vector space topology on $V$ by saying a set $U$ is open if and only if for every $x \in U$ there is some $N \in \Phi$ so that $x+N \subset U$ This would be the induced topology associated to the family of seminorms.

That we have a topology does not use the hypothesis that the family of seminorms is separating, although points will not be closed without the separating property. Arbitrary unions of sets containing sets of the form $x+N$ containing each point $x$ have the same property. The empty set and the whole space $V$ are visibly in the collection. The least trivial issue is to check that finite intersections of such sets are again of the same form. Looking at each point $x$ in a given finite intersection, this amounts to checking that finite intersections of sets in $\Phi$ are again in $\Phi$. But $\Phi$ is defined to be the collection of all finite intersections of sets $U_{i, \varepsilon}$, so this succeeds: we have closure under finite intersections, and a topology on $V$.

To verify that this topology makes $V$ a topological vectorspace is to verify the continuity of vector addition and continuity of scalar multiplication, and closed-ness of points. None of these verifications is difficult:

The separating property implies that for each $x \in V$ the intersection of all the sets $x+N$ with $N \in \Phi$ is just $x$. Given $y \in V$, for each $x \neq y$ let $U_{x}$ be an open set containing $x$ but not $y$. Then

$$
U=\bigcup_{x \neq y} U_{x}
$$

is open and has complement $\{y\}$, so the singleton $\{y\}$ is closed.
For continuity of vector addition, it suffices to prove that, given $N \in \Phi$ and given $x, y \in V$ there are $U, U^{\prime} \in \Phi$ so that

$$
(x+U)+\left(y+U^{\prime}\right) \subset x+y+N
$$

[79] The weak dual topology is traditionally called the weak-*-topology, but replacing * by dual is more explanatory.

The triangle inequality implies that for a fixed index $i$ and for $\varepsilon_{1}, \varepsilon_{2}>0$

$$
U_{i, \varepsilon_{1}}+U_{i, \varepsilon_{2}} \subset U_{i, \varepsilon_{1}+\varepsilon_{2}}
$$

Then

$$
\left(x+U_{i, \varepsilon_{1}}\right)+\left(y+U_{i, \varepsilon_{2}}\right) \subset(x+y)+U_{i, \varepsilon_{1}+\varepsilon_{2}}
$$

Thus, given

$$
N=U_{i_{1}, \varepsilon_{1}} \cap \ldots \cap U_{i_{n}, \varepsilon_{n}}
$$

take

$$
U=U^{\prime}=U_{i_{1}, \varepsilon_{1} / 2} \cap \ldots \cap U_{i_{n}, \varepsilon_{n} / 2}
$$

proving continuity of vector addition.
For continuity of scalar multiplication, prove that for given $\alpha \in k, x \in V$, and $N \in \Phi$ there are $\delta>0$ and $U \in \Phi$ so that

$$
\left(\alpha+B_{\delta}\right) \cdot(x+U) \subset \alpha x+N \quad\left(\text { with } B_{\delta}=\{\beta \in k:|\alpha-\beta|<\delta\}\right)
$$

Since $N$ is an intersection of the sub-basis sets $U_{i, \varepsilon}$, it suffices to consider the case that $N$ is such a set. Given $\alpha$ and $x$, for $\left|\alpha^{\prime}-\alpha\right|<\delta$ and for $x-x^{\prime} \in U_{i, \delta}$,

$$
\begin{gathered}
\nu_{i}\left(\alpha x-\alpha^{\prime} x^{\prime}\right)=\nu_{i}\left(\left(\alpha-\alpha^{\prime}\right) x+\left(\alpha^{\prime}\left(x-x^{\prime}\right)\right) \leq \nu_{i}\left(\left(\alpha-\alpha^{\prime}\right) x\right)+\nu_{i}\left(\alpha^{\prime}\left(x-x^{\prime}\right)\right)\right. \\
=\left|\alpha-\alpha^{\prime}\right| \cdot \nu_{i}(x)+\left|\alpha^{\prime}\right| \cdot \nu_{i}\left(x-x^{\prime}\right) \leq\left|\alpha-\alpha^{\prime}\right| \cdot \nu_{i}(x)+(|\alpha|+\delta) \cdot \nu_{i}\left(x-x^{\prime}\right) \\
\leq \delta \cdot\left(\nu_{i}(x)+|\alpha|+\delta\right)
\end{gathered}
$$

Thus, for the joint continuity, take $\delta>0$ small enough so that

$$
\delta \cdot\left(\delta+\nu_{i}(x)+|\alpha|\right)<\varepsilon
$$

Taking finite intersections presents no further difficulty, taking the corresponding finite intersections of the sets $B_{\delta}$ and $U_{i, \delta}$, finishing the demonstration that separating families of seminorms give a structure of topological vectorspace.

Last, check that finite intersections of the sets $U_{i, \varepsilon}$ are convex. Since intersections of convex sets are convex, it suffices to check that the sets $U_{i, \varepsilon}$ themselves are convex, which follows from the homogeneity and the triangle inequality: with $0 \leq t \leq 1$ and $x, y \in U_{i, \varepsilon}$,

$$
\nu_{i}(t x+(1-t) y) \leq \nu_{i}(t x)+\nu_{i}((1-t) y)=t \nu_{i}(x)+(1-t) \nu_{i}(y) \leq t \varepsilon+(1-t) \varepsilon=\varepsilon
$$

Thus, the set $U_{i, \varepsilon}$ is convex.
The converse, that every locally convex topology is given by a family of seminorms, is more difficult:
Let $U$ be a convex open set containing 0 in a topological vectorspace $V$. Every open neighborhood of 0 contains a balanced neighborhood of 0 , so shrink $U$ if necessary so it is balanced, that is, $\alpha v \in U$ for $v \in U$ and $|\alpha| \leq 1$. The Minkowski functional $\nu_{U}$ associated to $U$ is

$$
\nu_{U}(v)=\inf \{t \geq 0: v \in t U\}
$$

[13.11.2] Claim: The Minkowski functional $\nu_{U}$ associated to a balanced convex open neighborhood $U$ of 0 in a topological vectorspace $V$ is a seminorm on $V$, and is continuous in the topology on $V$.
Proof: The argument is as expected:
By continuity of scalar multiplication, every neighborhood $U$ of 0 is absorbing, in the sense that every $v \in V$ lies inside $t U$ for large enough $|t|$. Thus, the set over which we take the infimum to define the Minkowski functional is non-empty, so the infimum exists.

Let $\alpha$ be a scalar, and let $\alpha=s \mu$ with $s=|\alpha|$ and $|\mu|=1$. The balanced-ness of $U$ implies the balanced-ness of $t U$ for any $t \geq 0$, so for $v \in t U$ also

$$
\alpha v \in \alpha t U=s \mu t U=s t U
$$

From this,

$$
\{t \geq 0: \alpha v \in \alpha U\}=|\alpha| \cdot\{t \geq 0: \alpha v \in t U\}
$$

from which follows the homogeneity property required of a seminorm:

$$
\nu_{U}(\alpha v)=|\alpha| \cdot \nu_{U}(v) \quad(\text { for scalar } \alpha)
$$

For the triangle inequality use the convexity. For $v, w \in V$ and $s, t>0$ such that $v \in s U$ and $w \in t U$,

$$
v+w \in s U+t U=\left\{s u+t u^{\prime}: u, u^{\prime} \in U\right\}
$$

By convexity,

$$
s u+t u^{\prime}=(s+t) \cdot\left(\frac{s}{s+t} \cdot u+\frac{t}{s+t} \cdot u^{\prime}\right) \in(s+t) \cdot U
$$

Thus,

$$
\nu_{U}(v+w)=\inf \{r \geq 0: v+w \in r U\} \leq \inf \{r \geq 0: v \in r U\}+\inf \{r \geq 0: w \in r U\}=\nu_{U}(v)+\nu_{U}(w)
$$

Thus, the Minkowski functional $\nu_{U}$ attached to balanced, convex $U$ is a continuous seminorm.
[13.11.3] Theorem: The topology of a locally convex topological vectorspace $V$ is given by the collection of seminorms obtained as Minkowski functionals $\nu_{U}$ associated to a local basis at 0 consisting of convex, balanced opens.
Proof: The proof is straightforward, once we decide to tolerate an extravagantly large collection of seminorms. With or without local convexity, every neighborhood of 0 contains a balanced neighborhood of 0 . Thus, a locally convex topological vectorspace has a local basis $X$ at 0 of balanced convex open sets.

We claim that every open $U \in X$ can be recovered from the corresponding seminorm $\nu_{U}$ by

$$
U=\left\{v \in V: \nu_{U}(v)<1\right\}
$$

Indeed, for $v \in U$, the continuity of scalar multiplication gives $\delta>0$ and a neighborhood $N$ of $v$ such that $z \cdot v-1 \cdot v \in U$ for $|1-z|<\delta$. Thus, $v \in(1+\delta)^{-1} \cdot U$, so

$$
\nu_{U}(v)=\inf \{t \geq 0: v \in t \cdot U\} \leq \frac{1}{1+\delta}<1
$$

On the other hand, for $\nu_{U}(v)<1$, there is $t<1$ such that $v \in t U \subset U$, since $U$ is convex and contains 0 . Thus, the seminorm topology is at least as fine as the original.

Oppositely, the same argument shows that every seminorm local basis open

$$
\left\{v \in V: \nu_{U}(v)<t\right\}
$$

is simply $t U$. Thus, the original topology is at least as fine as the seminorm topology.
The comparison of descriptions of topologies is straightforward, as follows. For a seminorm $\nu$ on a topological vectorspace $V$, we can form a Banach space completing with respect to the pseudo-metric $\nu(x-y)$. In particular, unlike completions with respect to genuine metrics, there can be collapsing, so that the natural map of $V$ to this completion need not be an injection.
[13.11.4] Claim: Let $V$ be a topological vectorspace with topology given by a (separating) family of seminorms $S=\{\nu\}$. Order the set of finite subsets of $S$ by inclusion, and

$$
\nu_{F}=\sum_{\nu \in F} \nu \quad(\text { for finite subset } F \text { of } S)
$$

Then $V$ with its seminorm topology is a dense subspace of the limit $\lim _{F \in \Phi} V_{F}$ of the Banach-space completions $V_{F}$ with respect to $\nu_{F}$.
Proof: As earlier, the seminorm topology is literally the subspace topology on the diagonal copy of $V$ in the product of the $V_{F}$.

Of course, the poset of finite subsets of $S$ is more complicated than the poset of positive integers, so such a limit can be large. Certainly $V$ has a natural map to every $V_{F}$. Indeed, by definition of the seminorm topology, the open sets in $V$ are exactly the inverse images in $V$ of open sets in the various $V_{F}$.

For $F \subset F^{\prime}$, since $\nu_{F^{\prime}} \geq \nu_{F}$, there is a natural continuous linear map $V_{F^{\prime}} \rightarrow V_{F}$. The maps $V \rightarrow V_{F}$ are compatible, in the sense that the composite $V \rightarrow V_{F^{\prime}} \rightarrow V_{F}$ is the same as $V \rightarrow V_{F}$, for $F \subset F^{\prime}$. This induces a unique continuous linear map of $V$ to the limit of the $V_{F}$.

As in [13.2], the limit is the diagonal

$$
D=\left\{\left\{v_{F}\right\} \in \prod_{F} V_{F}: v_{F^{\prime}} \rightarrow v_{F}, \text { for all } F^{\prime} \supset F\right\} \subset \prod_{F} V_{F}
$$

with the subspace topology. Repeating part of an earlier argument, given a finite collection of finite subsets $F_{1}, \ldots, F_{n}$ of $S$, for $\left\{v_{F}\right\} \in D$, take neighborhoods $U_{i} \subset V_{F_{i}}$ containing $v_{F_{i}}$. Let $\Phi=\bigcup_{i} F_{i}$. The compatibility implies that there is $v_{\Phi} \in V_{\Phi}$ such that $v_{\Phi} \rightarrow v_{F_{i}}$ for all $i$. Also, there is a sufficiently small neighborhood $U$ of $v_{\Phi}$ such that its image in every $V_{F_{i}}$ is inside the neighborhood $U_{i}$ of $v_{F_{i}}$. Since the image of $V$ is dense in $V_{\Phi}$, take $v \in V$ with image inside $U$. Then the image of $v$ is inside $U_{i}$ for all $i$. Thus, the image of $V$ is dense in the limit.

Although it turns out that we only care about locally convex topological vectorspaces, there do exist complete-metric topological vectorspaces which fail to be locally convex. This underscores the need to explicitly specify that a Fréchet space should be locally convex. The usual example of a not-locally-convex complete-metric space is the sequence space

$$
\ell^{p}=\left\{x=\left(x_{1}, x_{2}, \ldots\right): \sum_{i}\left|x_{i}\right|^{p}<\infty\right\}
$$

for $0<p<1$ with metric

$$
d(x, y)=\sum_{i}\left|x_{i}-y_{i}\right|^{p} \quad \text { (note: no } p^{t h} \text { root, unlike the } p \geq 1 \text { case) }
$$

This example's interest is mostly as a counterexample to a naive presumption that local convexity is automatic.

### 13.12 Quasi-completeness theorem

We have already seen that LF-spaces such as the space of test functions $\mathcal{D}(\mathbb{R})=C_{c}^{\infty}(\mathbb{R})$, although not complete metrizable, are quasi-complete. It is fortunate that most important topological vector spaces are quasi-complete.

At the end of this section, we show that the fullest notion of completeness easily fails to hold, even for quasi-complete spaces.

It is clear that closed subspaces of quasi-complete spaces are quasi-complete. Products and finite sums of quasi-complete spaces are quasi-complete.

Let $\operatorname{Hom}(X, Y)$ be the space of continuous linear functions from a topological vectorspace $X$ to another topological vectorspace $Y$. Give $\operatorname{Hom}(X, Y)$ the topology by seminorms $p_{x, U}$ where $x \in X$ and $U$ is a convex, balanced neighborhood of 0 in $Y$, defined by

$$
\left.p_{x, U}(T)=\inf \{t>0: T x \in t U\} \quad \text { (for } T \in \operatorname{Hom}(X, Y)\right)
$$

For $Y=\mathbb{C}$, this gives the weak dual topology on $X^{*}$.
[13.12.1] Theorem: For $X$ a Fréchet space or LF-space, and $Y$ quasi-complete, the space $\operatorname{Hom}(X, Y)$, with the topology induced by the seminorms $p_{x, U}$, is quasi-complete.

Proof: Some preparation is required. A set $E$ of continuous linear maps from one topological vectorspace $X$ to another topological vectorspace $Y$ is equicontinuous when, for every neighborhood $U$ of 0 in $Y$, there is a neighborhood $N$ of 0 in $X$ so that $T(N) \subset U$ for every $T \in E$.
[13.12.2] Claim: Let $V$ be a strict colimit of a well-ordered countable collection of locally convex closed subspaces $V_{i}$. Let $Y$ be a locally convex topological vectorspace. Let $E$ be a set of continuous linear maps from $V$ to $Y$. Then $E$ is equicontinuous if and only if for each index $i$ the collection of continuous linear maps $\left\{\left.T\right|_{V_{i}}: T \in E\right\}$ is equicontinuous.
Proof: Given a neighborhood $U$ of 0 in $Y$, shrink $U$ if necessary so that $U$ is convex and balanced. For each index $i$, let $N_{i}$ be a convex, balanced neighborhood of 0 in $V_{i}$ so that $T N_{i} \subset U$ for all $T \in E$. Let $N$ be the convex hull of the union of the $N_{i}$. By the convexity of $N$, still $T N \subset U$ for all $T \in E$. By the construction of the diamond topology, $N$ is an open neighborhood of 0 in the coproduct, hence in the colimit, which is a quotient of the coproduct. This gives the equicontinuity of $E$. The other direction of the implication is easy.
[13.12.3] Claim: (Banach-Steinhaus/uniform boundedness theorem) Let $X$ be a Fréchet space or LF-space and $Y$ a locally convex topological vector space. A set $E$ of linear maps $X \rightarrow Y$, such that every set $E x=\{T x: T \in E\}$ is bounded in $Y$, is equicontinuous.
Proof: First consider $X$ Fréchet. Given a neighborhood $U$ of 0 in $Y$, let $A=\bigcap_{T \in E} T^{-1} \bar{U}$. By assumption, $\bigcup_{n} n A=X$. By the Baire category theorem, the complete metric space $X$ is not a countable union of nowhere dense subsets, so at least one of the closed sets $n A$ has non-empty interior. Since (non-zero) scalar multiplication is a homeomorphism, $A$ itself has non-empty interior, containing some $x+N$ for a neighborhood $N$ of 0 and $x \in A$. For every $T \in E$,

$$
T N \subset T\{a-x: a \in A\} \subset\left\{u_{1}-u_{2}: u_{1}, u_{2} \in \bar{U}\right\}=\bar{U}-\bar{U}
$$

By continuity of addition and scalar multiplication in $Y$, given an open neighborhood $U_{o}$ of 0 , there is $U$ such that $\bar{U}-\bar{U} \subset U_{o}$. Thus, $T N \subset U_{o}$ for every $T \in E$, and $E$ is equicontinuous.

For $X=\bigcup_{i} X_{i}$ an LF-space, this argument already shows that $E$ restricted to each $X_{i}$ is equicontinuous. As in the previous claim, this gives equicontinuity on the strict colimit.

For the proof of the theorem on quasi-completeness, let $E=\left\{T_{i}: i \in I\right\}$ be a bounded Cauchy net in $\operatorname{Hom}(X, Y)$, where $I$ is a directed set. Of course, attempt to define the limit of the net by $T x=\lim _{i} T_{i} x$. For $x \in X$ the evaluation map $S \rightarrow S x$ from $\operatorname{Hom}(X, Y)$ to $Y$ is continuous. In fact, the topology on $\operatorname{Hom}(X, Y)$ is the coarsest with this property. Therefore, by the quasi-completeness of $Y$, for each fixed $x \in X$ the net $T_{i} x$ in $Y$ is bounded and Cauchy, so converges to an element of $Y$ suggestively denoted $T x$.

To prove linearity of $T$, fix $x_{1}, x_{2}$ in $X, a, b \in \mathbb{C}$ and fix a neighborhood $U_{o}$ of 0 in $Y$. Since $T$ is in the closure of $E$, for any open neighborhood $N$ of 0 in $\operatorname{Hom}(X, Y)$, there exists

$$
T_{i} \in E \cap(T+N)
$$

In particular, for any neighborhood $U$ of 0 in $Y$, take

$$
N=\left\{S \in \operatorname{Hom}(X, Y): S\left(a x_{1}+b x_{2}\right) \in U, S\left(x_{1}\right) \in U, S\left(x_{2}\right) \in U\right\}
$$

Then

$$
\begin{gathered}
T\left(a x_{1}+b x_{2}\right)-a T\left(x_{1}\right)-b T\left(x_{2}\right) \\
=\left(T\left(a x_{1}+b x_{2}\right)-a T\left(x_{1}\right)-b T\left(x_{2}\right)\right)-\left(T_{i}\left(a x_{1}+b x_{2}\right)-a T_{i}\left(x_{1}\right)-b T_{i}\left(x_{2}\right)\right)
\end{gathered}
$$

since $T_{i}$ is linear. The latter expression is

$$
\begin{gathered}
T\left(a x_{1}+b x_{2}\right)-\left(a x_{1}+b x_{2}\right)+a\left(T\left(x_{1}\right)-T_{i}\left(x_{1}\right)+b\left(T\left(x_{2}\right)-T_{i}\left(x_{2}\right)\right.\right. \\
\in U+a U+b U
\end{gathered}
$$

By choosing $U$ small enough so that

$$
U+a U+b U \subset U_{o}
$$

we find that

$$
T\left(a x_{1}+b x_{2}\right)-a T\left(x_{1}\right)-b T\left(x_{2}\right) \in U_{o}
$$

Since this is true for every neighborhood $U_{o}$ of 0 in $Y$,

$$
T\left(a x_{1}+b x_{2}\right)-a T\left(x_{1}\right)-b T\left(x_{2}\right)=0
$$

which proves linearity.
Continuity of the limit operator $T$ exactly requires equicontinuity of $E=\left\{T_{i} x: i \in I\right\}$. Indeed, for each $x \in X,\left\{T_{i} x: i \in I\right\}$ is bounded in $Y$, so by Banach-Steinhaus, $\left\{T_{i}: i \in I\right\}$ is equicontinuous.

Fix a neighborhood $U$ of 0 in $Y$. Invoking the equicontinuity of $E$, let $N$ be a small enough neighborhood of 0 in $X$ so that $T(N) \subset U$ for all $T \in E$. Let $x \in N$. Choose an index $i$ sufficiently large so that $T x-T_{i} x \in U$, via the definition of the topology on $\operatorname{Hom}(X, Y)$. Then

$$
T x \in U+T_{i} x \subset U+U
$$

The usual rewriting, replacing $U$ by $U^{\prime}$ such that $U^{\prime}+U^{\prime} \subset U$, shows that $T$ is continuous.
Finally, we demonstrate that weak duals of reasonable topological vector spaces, such as infinitedimensional Hilbert, Banach, or Fréchet spaces, are definitely not complete in the strongest sense. That is, in these weak duals there are Cauchy nets which do not converge.
[13.12.4] Theorem: The weak dual of a locally-convex topological vector space $V$ is complete if and only if every linear functional on $V$ is continuous.
Proof: A vectorspace $V$ can be (re-) topologized as the colimit $V_{\text {init }}$ of all its finite-dimensional subspaces. Although the poset of finite-dimensional subspaces is much larger than the poset of positive integers, the earlier argument still applies: this colimit really is the ascending union with a suitable topology.
[13.12.5] Claim: For a locally-convex topological vector space $V$ the identity map $V_{\text {init }} \rightarrow V$ is continuous. That is, $V_{\text {init }}$ is the finest locally convex topological vector space topology on $V$.
Proof: Finite-dimensional topological vector spaces have unique topologies [13.4]. Thus, for any finitedimensional vector subspace $X$ of $V$ the inclusion $X \rightarrow V$ is continuous with that unique topology on $X$. These inclusions form a compatible family of maps to $V$, so by the characterization of colimit there is a unique continuous map $V_{\text {init }} \rightarrow V$. This map is the identity on every finite-dimensional subspace, so is the identity on the underlying set $V$.
[13.12.6] Claim: Every linear functional $\lambda: V_{\text {init }} \rightarrow \mathbb{C}$ is continuous.
Proof: The restrictions of a given linear function $\lambda$ on $V$ to finite-dimensional subspaces are compatible with the inclusions among finite-dimensional subspaces. Every linear functional on a finite-dimensional space is continuous, so the characterizing property of the colimit implies that $\lambda$ is continuous on $V_{\text {init }}$.
[13.12.7] Claim: The weak dual $V^{*}$ of a locally-convex topological vector space $V$ injects continuously to the limit of the finite-dimensional Banach spaces

$$
\left.V_{\Phi}^{*}=\text { completion of } V^{*} \text { under seminorm } p_{\Phi}(\lambda)=\sum_{v \in \Phi}|\lambda(v)| \quad \quad \text { finite } \Phi \subset V\right)
$$

and the weak dual topology is the subspace topology.
Proof: The weak dual topology on the continuous dual $V^{*}$ of a topological vector space $V$ is given by the seminorms

$$
p_{v}(\lambda)=|\lambda(v)| \quad\left(\text { for } \lambda \in V^{*} \text { and } v \in V\right)
$$

The corresponding local basis is finite intersections

$$
\left\{\lambda \in V^{*}:|\lambda(v)|<\varepsilon, \text { for all } v \in \Phi\right\} \quad \text { (for arbitrary finite sets } \Phi \subset V \text { ) }
$$

These sets contain, and are contained in, sets of the form

$$
\left\{\lambda \in V^{*}: \sum_{v \in \Phi}|\lambda(v)|<\varepsilon\right\} \quad \text { (for arbitrary finite sets } \Phi \subset V \text { ) }
$$

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Therefore, the weak dual topology on $V^{*}$ is also given by semi-norms

$$
\left.p_{\Phi}(\lambda)=\sum_{v \in \Phi}|\lambda(v)| \quad \text { (finite } \Phi \subset V\right)
$$

These have the convenient feature that they form a projective family, indexed by (reversed) inclusion. Let $V^{*}(\Phi)$ be $V^{*}$ with the $p_{\Phi}$-topology: this is not Hausdorff, so continuous linear maps $V^{*} \rightarrow V^{*}(\Phi)$ descend to maps $V^{*} \rightarrow V_{\Phi}^{*}$ to the completion $V_{\Phi}^{*}$ of $V^{*}$ with respect to the pseudo-metric attached to $p_{\Phi}$. The quotient map $V^{*}(\Phi) \rightarrow V_{\Phi}^{*}$ typically has a large kernel, since

$$
\operatorname{dim}_{\mathbb{C}} V_{\Phi}^{*}=\operatorname{card} \Phi \quad(\text { for finite } \Phi \subset V)
$$

The maps $V^{*} \rightarrow V_{\Phi}^{*}$ are compatible with respect to (reverse) inclusion $\Phi \supset Y$, so $V^{*}$ has a natural induced map to the $\lim _{\Phi} V_{\Phi}^{*}$. Since $V$ separates points in $V^{*}, V^{*}$ injects to the limit. The weak topology on $V^{*}$ is exactly the subspace topology from that limit.
[13.12.8] Claim: The weak dual $V_{\text {init }}^{*}$ of $V_{\text {init }}$ is the limit of the finite-dimensional Banach spaces

$$
\left.V_{\Phi}^{*}=\text { completion of } V_{\text {init }}^{*} \text { under seminorm } p_{\Phi}(\lambda)=\sum_{v \in \Phi}|\lambda(v)| \quad \quad \text { (finite } \Phi \subset V\right)
$$

Proof: The previous proposition shows that $V_{\mathrm{init}}^{*}$ injects to the limit, and that the subspace topology from the limit is the weak dual topology. On the other hand, the limit consists of linear functionals on $V$, without regard to topology or continuity. Since all linear functionals are continuous on $V_{\text {init }}$, the limit is naturally a subspace of $V_{\text {init }}^{*}$.

Returning to the proof of the theorem, $\lim _{\Phi} V_{\Phi}^{*}$ is a closed subspace of the corresponding product, so is complete in the fullest sense. Any other locally convex topologization $V_{\tau}$ of $V$ has weak dual $\left(V_{\tau}\right)^{*} \subset\left(V_{\mathrm{init}}\right)^{*}$ with the subspace topology, and the image is dense in $\left(V_{\text {init }}\right)^{*}$. Thus, unless $\left(V_{\tau}\right)^{*}=\left(V_{\text {init }}\right)^{*}$, the weak dual $V_{\tau}^{*}$ is not complete.

### 13.13 Strong operator topology

For $X$ and $Y$ Hilbert spaces, the topology on $\operatorname{Hom}(X, Y)$ given by seminorms

$$
\left.p_{x, U}(T)=\inf \{t>0: T x \in t U\} \quad \text { (for } T \in \operatorname{Hom}(X, Y)\right)
$$

where $x \in X$ and $U$ is a convex, balanced neighborhood of 0 in $Y$, is the strong operator topology. Indeed, every neighborhood of 0 in $Y$ contains an open ball, so this topology can also be given by seminorms

$$
q_{x}(T)=|T x|_{Y} \quad(\text { for } T \in \operatorname{Hom}(X, Y))
$$

where $x \in X$. The strong operator topology is weaker than the uniform topology given by the operator norm $|T|=\sup _{|x| \leq 1}|T x|_{Y}$.

The uniform operator-norm topology makes the space of operators a Banach space, certainly simpler than the strong operator topology, but the uniform topology is too strong for many purposes.

For example, group actions on Hilbert spaces are rarely continuous for the uniform topology: letting $\mathbb{R}$ act on $L^{2}(\mathbb{R})$ by $T_{x} f(y)=f(x+y)$, no matter how small $\delta>0$ is, there is an $L^{2}$ function $f$ with $|f|_{L^{2}}=1$ such that $\left|T_{\delta} f-f\right|_{L^{2}}=\sqrt{2}$.

Despite the strong operator topology being less elementary than the uniform topology, the theorem of the previous section [13.12] shows that $\operatorname{Hom}(X, Y)$ with the strong operator topology is quasi-complete.

### 13.14 Generalized functions (distributions) on $\mathbb{R}$

The most immediate definition of the space of distributions or generalized functions on $\mathbb{R}$ is as the dual $\mathcal{D}^{*}=\mathcal{D}(\mathbb{R})^{*}=C_{c}^{\infty}(\mathbb{R})^{*}$ to the space $\mathcal{D}=\mathcal{D}(\mathbb{R})$ of test functions, with the weak dual topology by seminorms $\nu_{f}(u)=|u(f)|$ for test functions $f$ and distributions $u$.

Similarly, the tempered distributions are the weak dual $\mathscr{S}^{*}=\mathscr{S}(\mathbb{R})^{*}$, and the compactly-supported distributions are the weak dual $\mathcal{E}^{*}=\mathcal{E}(\mathbb{R})^{*}$, in this context writing $\mathcal{E}(\mathbb{R})=C^{\infty}(\mathbb{R})$. Naming $\mathcal{E}^{*}$ compactlysupported will be justified below.

By dualizing, the continuous containments $\mathcal{D} \subset \mathscr{S} \subset \mathcal{E}$ give continuous maps $\mathcal{E}^{*} \rightarrow \mathscr{S}^{*} \rightarrow \mathcal{D}^{*}$. When we know that $\mathcal{D}$ is dense in $\mathscr{S}$ and in $\mathcal{E}$, it will follow that these are injections. The most straightforward argument for density uses Gelfand-Pettis integrals, as in [14.5]. Thus, for the moment, we cannot claim that $\mathcal{E}^{*}$ and $\mathscr{S}^{*}$ are distributions, but only that they naturally map to distributions.
[13.12] shows that $\mathcal{D}^{*}, \mathscr{S}^{*}$, and $\mathcal{E}^{*}$ are quasi-complete, despite not being complete in the strongest possible sense.

The description of the space of distributions as the weak dual to the space of test functions falls far short of explaining its utility. There is a natural imbedding $\mathcal{D}(\mathbb{R}) \rightarrow \mathcal{D}(\mathbb{R})^{*}$ of test functions into distributions, by

$$
f \rightarrow u_{f} \quad \text { by } \quad u_{f}(g)=\int_{\mathbb{R}} f(x) g(x) d x \quad(\text { for } f, g \in \mathcal{D}(\mathbb{R}))
$$

That is, via this imbedding we consider distributions to be generalized functions. Indeed, [14.5] shows that test functions $\mathcal{D}(\mathbb{R})$ are dense in $\mathcal{D}(\mathbb{R})^{*}$.

The simplest example of a distribution not obtained by integration against a test function on $\mathbb{R}$ is the Dirac delta, the evaluation map $\delta(f)=f(0)$. From [13.9] and other earlier results, this is continuous for the LF-space topology on test functions.

This imbedding, and integration by parts, explain how to define $\frac{d}{d x}$ on distributions in a form compatible with the imbedding $\mathcal{D} \subset \mathcal{D}^{*}$ : noting the sign, due to integration by parts,

$$
\left(\frac{d}{d x} u\right)(f)=-u\left(\frac{d}{d x} f\right) \quad\left(\text { for } u \in \mathcal{D}^{*} \text { and } f \in \mathcal{D}\right)
$$

[13.14.1] Claim: $\frac{d}{d x}: \mathcal{D}^{*} \rightarrow \mathcal{D}^{*}$ is continuous.
Proof: By the nature of the weak dual topology, it suffices to show that for each $f \in \mathcal{D}$ and $\varepsilon>0$ there are $g \in \mathcal{D}$ and $\delta>0$ such that $|u(g)|<\delta$ implies $\left|\left(\frac{d}{d x} u\right)(f)\right|<\varepsilon$. Taking $g=\frac{d}{d x} f$ and $\delta=\varepsilon$ succeeds. ///

Multiplications by $F \in C^{\infty}(\mathbb{R})$ give continuous maps $\mathcal{D}$ to itself, compatible with the imbedding $\mathcal{D} \rightarrow \mathcal{D}^{*}$ in the sense that

$$
\int_{\mathbb{R}}(F \cdot u)(x) f(x) d x=\int_{\mathbb{R}} u(x)(F \cdot f)(x) d x \quad \quad\left(\text { for } F \in C^{\infty}(\mathbb{R}) \text { and } u, f \in \mathcal{D}(\mathbb{R})\right)
$$

Extend this to a map $\mathcal{D}^{*} \rightarrow \mathcal{D}^{*}$ by

$$
(F \cdot u)(f)=u(F \cdot f) \quad\left(\text { for } F \in C^{\infty}, u \in \mathcal{D}^{*}, \text { and } f \in \mathcal{D}\right)
$$

[13.14.2] Claim: Multiplication operators $\mathcal{D}^{*} \rightarrow \mathcal{D}^{*}$ by $F \in C^{\infty}$ are continuous.
Proof: By the nature of the weak dual topology, it suffices to show that for each $f \in \mathcal{D}$ and $\varepsilon>0$ there are $g \in \mathcal{D}$ and $\delta>0$ such that $|u(g)|<\delta$ implies $\mid F \cdot u)(f) \mid<\varepsilon$. Taking $g=F \cdot f$ and $\delta=\varepsilon$ succeeds. $\quad / / /$

Since $\mathscr{S}$ is mapped to itself by Fourier transform [13.13], this gives a way to define Fourier transform on $\mathscr{S}^{*}$, by a duality extending the Plancherel theorem:

$$
\widehat{u}(f)=u(\widehat{f}) \quad\left(\text { for } f \in \mathscr{S} \text { and } u \in S c h w^{*}\right)
$$

Recall that the support of a function is the closure of the set on which it is non-zero, slightly complicating the notion of support for a distribution $u$ : support of $u$ is the complement of the union of all open sets $U$ such that $u(f)=0$ for all test functions $f$ with support inside $U$.
[13.14.3] Theorem: A distribution with support $\{0\}$ is a finite linear combination of Dirac's $\delta$ and its derivatives.
Proof: Since $\mathcal{D}$ is a colimit of $\mathcal{D}_{K}$ over $K=[-n, n]$, it suffices to classify $u$ in $\mathcal{D}_{K}^{*}$ with support $\{0\}$. We claim that a continuous linear functional on $\mathcal{D}_{K}=\lim _{k} C_{K}^{k}$ factors through some limitand

$$
C_{K}^{k}=\left\{f \in C^{k}(K): f^{(i)} \text { vanishes on } \partial K \text { for } 0 \leq i \leq k\right\}
$$

This is a special case of
[13.14.4] Claim: Let $X=\lim _{n} B_{n}$ be a limit of Banach spaces, with the image of $X$ dense in each $B_{n}$. A continuous linear map $T: \lim _{n} B_{n} \rightarrow Z$ from a, to a normed space $Z$ factors through some limitand $B_{n}$. For $Z=\mathbb{C}$, the same conclusion holds without the density assumption.
Proof: Let $X=\lim _{i} B_{i}$ with projections $p_{i}: X \rightarrow B_{i}$. Each $B_{i}$ is the closure of the image of $X$. By the continuity of $T$ at 0 , there is an open neighborhood $U$ of 0 in $X$ such that $T U$ is inside the open unit ball at 0 in $Z$. By the description of the limit topology as the product topology restricted to the diagonal, there are finitely-many indices $i_{1}, \ldots, i_{n}$ and open neighborhoods $V_{i_{t}}$ of 0 in $B_{i_{t}}$ such that

$$
\bigcap_{t=1}^{n} p_{i_{t}}^{-1}\left(p_{i_{t}} X \cap V_{i_{t}}\right) \subset U
$$

We can make a smaller open in $X$ by a condition involving a single limitand, as follows. Let $j$ be any index with $j \geq i_{t}$ for all $t$, and

$$
N=\bigcap_{t=1}^{n} p_{i_{t}, j}^{-1}\left(p_{i_{t}, j} B_{j} \cap V_{i_{t}}\right) \subset B_{j}
$$

By the compatibility $p_{i_{t}}^{-1}=p_{j}^{-1} \circ p_{i_{t}, j}^{-1}$, we have $p_{i_{t}, j} N \subset V_{i_{t}}$ for $i_{1}, \ldots, i_{n}$, and $p_{j}^{-1}\left(p_{j} X \cap N\right) \subset U$. By the linearity of $T$, for any $\varepsilon>0$,

$$
T\left(\varepsilon \cdot p_{j}^{-1}\left(p_{j} X \cap N\right)\right)=\varepsilon \cdot T\left(p_{j}^{-1}\left(p_{j} X \cap N\right)\right) \subset \varepsilon \text {-ball in } Z
$$

We claim that $T$ factors through $p_{j} X$ with the subspace topology from $B_{j}$. One potential issue in general is that $p_{j}: X \rightarrow B_{j}$ can have a non-trivial kernel, and we must check that ker $p_{j} \subset \operatorname{ker} T$. By the linearity of $T$,

$$
T\left(\frac{1}{n} \cdot p_{j}^{-1}\left(p_{j} \cap N\right)\right) \subset \frac{1}{n} \text {-ball in } Z
$$

so

$$
T\left(\bigcap_{n} \frac{1}{n} \cdot p_{j}^{-1}\left(p_{j} X \cap N\right)\right) \subset \frac{1}{m} \text {-ball in } Z \quad(\text { for all } m)
$$

and then

$$
T\left(\bigcap_{n} \frac{1}{n} \cdot p_{j}^{-1}\left(p_{j} \cap N\right)\right) \subset \bigcap_{m} \frac{1}{m} \text {-ball in } Z=\{0\}
$$

Thus,

$$
\bigcap_{n} p_{j}^{-1}\left(p_{j} X \cap \frac{1}{n} \cdot N\right)=\bigcap_{n} \frac{1}{n} \cdot p^{-1}\left(p_{j} X \cap N\right) \subset \operatorname{ker} T
$$

Thus, for $x \in X$ with $p_{j} x=0$, certainly $p_{j} x \in \frac{1}{n} N$ for all $n=1,2, \ldots$, and

$$
x \in \bigcap_{n} p_{j}^{-1}\left(p_{j} X \cap \frac{1}{n} N\right) \subset \operatorname{ker} T
$$

This proves the subordinate claim that $T$ factors through $p_{j}: X \rightarrow B_{j}$ via a (not necessarily continuous) linear map $T^{\prime}: p_{j} X \rightarrow Z$. The continuity follows from continuity at 0 , which is

$$
T\left(\varepsilon \cdot p_{j}^{-1}\left(p_{j} X \cap N\right)\right)=\varepsilon \cdot T\left(p_{j}^{-1}\left(p_{j} X \cap N\right)\right) \subset \varepsilon \text {-ball in } Z
$$

Then $T^{\prime}: p_{j} X \rightarrow Z$ extends to a map $B_{j} \rightarrow Z$ by continuity: given $\varepsilon>0$, take symmetric convex neighborhood $U$ of 0 in $B_{j}$ such that $\left|T^{\prime} y\right|_{Z}<\varepsilon$ for $y \in p_{j} X \cap U$. Let $y_{i}$ be a Cauchy net in $p_{j} X$ approaching $b \in B_{j}$. For $y_{i}$ and $y_{j}$ inside $b+\frac{1}{2} U,\left|T^{\prime} y_{i}-T^{\prime} y_{j}\right|=\left|T^{\prime}\left(y_{i}-y_{j}\right)\right|<\varepsilon$, since $y_{i}-y_{j} \in \frac{1}{2} \cdot 2 U=U$. Then unambiguously define $T^{\prime} b$ to be the $Z$-limit of the $T^{\prime} y_{i}$. The closure of $p_{j} X$ in $B_{j}$ is $B_{j}$, giving the desired map.

When $u$ is a functional, that is, a map to $\mathbb{C}$, we can extend it by Hahn-Banach.
Returning to the proof of the theorem: thus, there is $k \geq 0$ such that $u$ factors through a limitand $C_{K}^{k}$. In particular, $u$ is continuous for the $C^{k}$ topology on $\mathcal{D}_{K}$.

We need an auxiliary gadget. Fix a test function $\psi$ identically 1 on a neighborhood of 0 , bounded between 0 and 1 , and (necessarily) identically 0 outside some (larger) neighborhood of 0 . For $\varepsilon>0$ let

$$
\psi_{\varepsilon}(x)=\psi\left(\varepsilon^{-1} x\right)
$$

Since the support of $u$ is just $\{0\}$, for all $\varepsilon>0$ and for all $f \in \mathcal{D}\left(\mathbb{R}^{n}\right)$ the support of $f-\psi_{\varepsilon} \cdot f$ does not include 0 , so

$$
u\left(\psi_{\varepsilon} \cdot f\right)=u(f)
$$

Thus, for implied constant depending on $k$ and $K$, but not on $f$,

$$
\left|\psi_{\varepsilon} f\right|_{k}=\sup _{x \in K} \sum_{0 \leq i \leq k}\left|\left(\psi_{\varepsilon} f\right)^{(i)}(x)\right| \ll \sum_{i \leq k} \sum_{0 \leq j \leq i} \sup _{x} \varepsilon^{-j}\left|\psi^{(j)}\left(\varepsilon^{-1} x\right) f^{(i-j)}(x)\right|
$$

For test function $f$ vanishing to order $k$ at 0 , that is, $f^{(i)}(0)=0$ for all $0 \leq i \leq k$, on a fixed neighborhood of 0 , by a Taylor-Maclaurin expansion, $|f(x)| \ll|x|^{k+1}$, and, generally, for $i^{t h}$ derivatives with $0 \leq i \leq k$, $\left|f^{(i)}(x)\right| \ll|x|^{k+1-i}$. By design, all derivatives $\psi^{\prime}, \psi^{\prime \prime}, \ldots$ are identically 0 in a neighborhood of 0 , so, for suitable implied constants independent of $\varepsilon$,

$$
\begin{gathered}
\left|\psi_{\varepsilon} f\right|_{k} \ll \sum_{0 \leq i \leq k} \sum_{0 \leq j \leq i} \varepsilon^{-j} \cdot\left|\psi^{(j)}\left(\varepsilon^{-1} x\right) f^{(i-j)}(x)\right| \ll \sum_{0 \leq i \leq k} \sum_{j=0} \varepsilon^{-j} \cdot 1 \cdot \varepsilon^{k+1-i} \\
=\sum_{0 \leq i \leq k} \varepsilon^{k+1-i} \ll \varepsilon^{k+1-k}=\varepsilon
\end{gathered}
$$

Thus, for sufficiently small $\varepsilon>0$, for smooth $f$ vanishing to order $k$ at $0,|u(f)|=\left|u\left(\psi_{\varepsilon} f\right)\right| \ll \varepsilon$, and $u(f)=0$. That is,

$$
\operatorname{ker} u \supset \bigcap_{0 \leq i \leq k} \operatorname{ker} \delta^{(i)}
$$

The conclusion, that $u$ is a linear combination of the distributions $\delta, \delta^{\prime}, \delta^{(2)}, \ldots, \delta^{(k)}$, follows from
[13.14.5] Claim: A linear functional $\lambda \in V^{*}$ vanishing on the intersection $\bigcap_{i} \operatorname{ker} \lambda_{i}$ of kernels of a finite collection $\lambda_{1}, \ldots, \lambda_{n} \in V^{*}$ is a linear combination of the $\lambda_{i}$.
Proof: The linear map

$$
q: V \longrightarrow \mathbb{C}^{n} \quad \text { by } \quad v \longrightarrow\left(\lambda_{1} v, \ldots, \lambda_{n} v\right)
$$

is continuous since each $\lambda_{i}$ is continuous, and $\lambda$ factors through $q$, as $\lambda=L \circ q$ for some linear functional $L$ on $\mathbb{C}^{n}$. We know all the linear functionals on $\mathbb{C}^{n}$, namely, $L$ is of the form

$$
L\left(z_{1}, \ldots, z_{n}\right)=c_{1} z_{1}+\ldots+c_{n} z_{n} \quad\left(\text { for some } c_{i} \in \mathbb{C}\right)
$$

Thus,

$$
\lambda(v)=(L \circ q)(v)=L\left(\lambda_{1} v, \ldots, \lambda_{n} v\right)=c_{1} \lambda_{1}(v)+\ldots+c_{n} \lambda_{n}(v)
$$

expressing $\lambda$ as a linear combination of the $\lambda_{i}$.
The order of a distribution $u: \mathcal{D} \rightarrow \mathbb{C}$ is the integer $k$, if such exists, such that $u$ is continuous when $\mathcal{D}$ is given the weaker topology from $\operatorname{colim}_{K} C_{K}^{k}$. Not every distribution has finite order, but there is a useful technical application of the previous discussion:

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[13.14.6] Corollary: A distribution $u \in \mathcal{D}^{*}$ with compact support has finite order.
Proof: Let $\psi$ be a test function that is identically 1 on an open containing the support of $u$. Then

$$
u(f)=u((1-\psi) \cdot f)+u(\psi \cdot f)=0+u(\psi \cdot f)
$$

since $(1-\psi) \cdot f$ is a test function with support not meeting the support of $u$. With $K=\operatorname{spt} \psi$, this suggests that $u$ factors through a subspace of $\mathcal{D}_{K}$ via $f \rightarrow \psi \cdot f \rightarrow u(\psi \cdot f)$, but there is the issue of continuity. Distinguishing things a little more carefully, the compatibility embodied in the commutative diagram

gives

$$
u(f)=u(\psi \cdot f)=u(i(\psi f))=u_{K}(\psi f)
$$

The map $u_{K}$ is continuous, as is the multiplication $f \rightarrow \psi f$. The map $u_{K}$ is from the limit $\mathcal{D}_{K}$ of Banach spaces $C_{K}^{k}$ to the normed space $\mathbb{C}$, so factors through some limitand $C_{K}^{k}$, by [13.4.4]. As in proof that multiplication is continuous in the $C^{\infty}$ topology, by Leibniz' rule, the $C^{k}$ norm of $\psi f$ is

$$
\begin{aligned}
|\psi f|_{k} & =\sum_{0 \leq i \leq k} \sup _{x \in K}\left|(\psi f)^{(i)}(x)\right| \ll \sum_{i \leq k} \sum_{0 \leq j \leq i} \sup _{x}\left|\psi^{(j)}(x) f^{(i-j)}(x)\right| \\
& \ll \sum_{0 \leq i \leq k} \sup _{x \in K}\left|f^{(i)}(x)\right| \cdot \sum_{j \leq k} \sup _{x}\left|\psi^{(j)}(x)\right|=|f|_{C^{k}} \cdot|\psi|_{C^{k}}
\end{aligned}
$$

Since $\psi$ is fixed, this gives continuity in $f$ in the $C^{k}$ topology.
[13.14.7] Claim: In the inclusion $\mathcal{E}^{*} \subset \mathscr{S}^{*} \subset \mathcal{D}^{*}$, the image of $\mathcal{E}^{*}$ really is the collection of distributions with compact support.
Proof: On one hand the previous shows that $u \in \mathcal{D}^{*}$ with compact support can be composed as $u(f)=$ $u_{K}(\psi f)$ for suitable $\psi \in \mathcal{D}$. The map $f \rightarrow \psi \cdot f$ is also continuous as a map $\mathcal{E} \rightarrow \mathcal{D}$, so the same expression $f \rightarrow \psi f \rightarrow u_{K}(\psi f)$ extends $u \in \mathcal{D}^{*}$ to a continuous linear functional on $\mathcal{E}$.

On the other hand, let $u \in \mathcal{E}^{*}$. Composition of $u$ with $\mathcal{D} \rightarrow \mathcal{E}$ gives an element of $\mathcal{D}^{*}$, which we must check has compact support. From [13.5], $\mathcal{E}$ is a limit of the Banach spaces $C^{k}(K)$ with $K=[-n, n]$, without claiming that the image of $\mathcal{E}$ is necessarily dense in any of these. By [13.4.4], $u$ factors through some limitand $C^{k}(K)$. The map $\mathcal{D} \rightarrow \mathcal{E}$ is compatible with the restriction maps $\operatorname{Res}_{K}: \mathcal{D} \rightarrow C^{k}(K)$ : the diagram

commutes. For $f \in \mathcal{D}$ with support disjoint from $K, \operatorname{Res}_{K}(f)=0$, and $u(f)=0$. This proves that the support of the (induced) distribution is contained in $K$, so is compact.

### 13.15 Tempered distributions and Fourier transforms on $\mathbb{R}$

One normalization of the Fourier transform integral is

$$
\widehat{f}(\xi)=\mathcal{F} f(\xi)=\int_{\mathbb{R}} \bar{\psi}_{\xi}(x) f(x) d x \quad \quad\left(\text { with } \psi_{\xi}(x)=e^{2 \pi i \xi x}\right)
$$

This converges nicely for $f$ in the space $\mathscr{S}(\mathbb{R})$ of Schwartz functions.
[13.15.1] Theorem: Fourier transform is a topological isomorphism of $\mathscr{S}(\mathbb{R})$ to itself, with Fourier inversion $\operatorname{map} \varphi \rightarrow \check{\varphi}$ given by

$$
\check{\varphi}(x)=\int_{\mathbb{R}} \psi_{\xi}(x) \widehat{f}(\xi) d \xi
$$

Proof: Using the idea [14.3] that Schwartz functions extend to smooth functions on a suitable one-point compactification of $\mathbb{R}$ vanishing to infinite order at the point at infinity, Gelfand-Pettis integrals justify moving a differentiation under the integral,

$$
\begin{gathered}
\frac{d}{d \xi} \widehat{f}(\xi)=\frac{d}{d \xi} \int_{\mathbb{R}} \bar{\psi}_{\xi}(x) f(x) d x=\int_{\mathbb{R}} \frac{\partial}{\partial \xi} \bar{\psi}_{\xi}(x) f(x) d x \\
=\int_{\mathbb{R}}(-2 \pi i x) \bar{\psi}_{\xi}(x) f(x) d x=(-2 \pi i) \int_{\mathbb{R}} \bar{\psi}_{\xi}(x) x f(x) d x=(-2 \pi i) \widehat{x f}(\xi)
\end{gathered}
$$

Similarly, with an integration by parts,

$$
-2 \pi i \xi \cdot \widehat{f}(\xi)=\int_{\mathbb{R}} \frac{\partial}{\partial x} \bar{\psi}_{\xi}(x) \cdot f(x) d x=-\mathcal{F} \frac{d f}{d x}(\xi)
$$

Thus, $\mathcal{F}$ maps $\mathscr{S}(\mathbb{R})$ to itself.
The natural idea to prove Fourier inversion for $\mathscr{S}(\mathbb{R})$, that unfortunately begs the question, is the obvious:

$$
\int_{\mathbb{R}} \psi_{\xi}(x) \widehat{f}(\xi) d \xi=\int_{\mathbb{R}} \psi_{\xi}(x)\left(\int_{\mathbb{R}} \bar{\psi}_{\xi}(t) f(t) d t\right) d \xi=\int_{\mathbb{R}} f(t)\left(\int_{\mathbb{R}} \psi_{\xi}(x-t) d t\right) d t
$$

If we could justify asserting that the inner integral is $\delta_{x}(t)$, which it is, then Fourier inversion follows. However, Fourier inversion for $\mathscr{S}(\mathbb{R})$ is used to make sense of that inner integral in the first place.

Despite that issue, a dummy convergence factor will legitimize the idea. For example, let $g(x)=e^{-\pi x^{2}}$ be the usual Gaussian. Various computations show that it is its own Fourier transform. For $\varepsilon>0$, as $\varepsilon \rightarrow 0^{+}$, the dilated Gaussian $g_{\varepsilon}(x)=g(\varepsilon \cdot x)$ approaches 1 uniformly on compacts. Thus,

$$
\int_{\mathbb{R}} \psi_{\xi}(x) \widehat{f}(\xi) d \xi=\int_{\mathbb{R}} \lim _{\varepsilon \rightarrow 0^{+}} g(\varepsilon \xi) \psi_{\xi}(x) \widehat{f}(\xi) d \xi=\lim _{\varepsilon \rightarrow 0^{+}} \int_{\mathbb{R}} g(\varepsilon \xi) \psi_{\xi}(x) \widehat{f}(\xi) d \xi
$$

by monotone convergence or more elementary reasons. Then the iterated integral is legitimately rearranged:

$$
\int_{\mathbb{R}} g(\varepsilon \xi) \psi_{\xi}(x) \widehat{f}(\xi) d \xi=\int_{\mathbb{R}} \int_{\mathbb{R}} g(\varepsilon \xi) \psi_{\xi}(x) \bar{\psi}_{\xi}(t) f(t) d t d \xi=\int_{\mathbb{R}} \int_{\mathbb{R}} g(\varepsilon \xi) \psi_{\xi}(x-t) f(t) d \xi d t
$$

By changing variables in the definition of Fourier transform, $\widehat{g_{\varepsilon}}=\frac{1}{\varepsilon} g_{1 / \varepsilon}$. Thus,

$$
\int_{\mathbb{R}} \psi_{\xi}(x) \widehat{f}(\xi) d \xi=\int_{\mathbb{R}} \frac{1}{\varepsilon} g\left(\frac{x-t}{\varepsilon}\right) f(t) d t=\int_{\mathbb{R}} \frac{1}{\varepsilon} g\left(\frac{t}{\varepsilon}\right) \cdot f(x+t) d t
$$

The sequence of function $g_{1 / \varepsilon} / \varepsilon$ is not an approximate identity in the strictest sense, since the supports are the entire line. Nevertheless, the integral of each is 1 , and as $\varepsilon \rightarrow 0^{+}$, the mass is concentrated on smaller and smaller neighborhoods of $0 \in \mathbb{R}$. Thus, for $f \in \mathscr{S}(\mathbb{R})$,

$$
\lim _{\varepsilon \rightarrow 0^{+}} \int_{\mathbb{R}} \frac{1}{\varepsilon} g\left(\frac{t}{\varepsilon}\right) \cdot f(x+t) d t=f(x)
$$

This proves Fourier inversion. In particular, this proves that Fourier transform bijects the Schwartz space to itself.

With Fourier inversion in hand, we can prove the Plancherel identity for Schwartz functions:
[13.15.2] Corollary: For $f, g \in \mathscr{S}$, the Fourier transform is an isometry in the $L^{2}(\mathbb{R})$ topology, that is, $\langle\widehat{f}, \widehat{g}\rangle=\langle f, g\rangle$.
Proof: There is an immediate preliminary identity:

$$
\int_{\mathbb{R}} \widehat{f}(\xi) h(\xi) d \xi=\int_{\mathbb{R}} \int_{\mathbb{R}} e^{-2 \pi i \xi x} f(x) h(\xi) d \xi d x=\int_{\mathbb{R}} \int_{\mathbb{R}} e^{-2 \pi i \xi x} f(x) h(\xi) d x d \xi=\int_{\mathbb{R}} f(x) \widehat{h}(x) d x
$$

To get from this identity to Plancherel requires, given $g \in \mathscr{S}$, existence of $h \in \mathscr{S}$ such that $\widehat{h}=\bar{g}$, with complex conjugation. By Fourier inversion on Schwartz functions, $h=(\bar{g})^{\llcorner }$succeeds.
[13.15.3] Corollary: Fourier transform extends by continuity to an isometry $L^{2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$.
Proof: Schwartz functions are dense in in $L^{2}(\mathbb{R})$.
[13.15.4] Corollary: Fourier transform extends to give a bijection of the space tempered distributions $\mathscr{S}^{*}$ to itself, by

$$
\widehat{u}(\varphi)=u(\widehat{\varphi}) \quad(\text { for all } \varphi \in \mathscr{S})
$$

Proof: Fourier transform is a topological isomorphism of $\mathscr{S}$ to itself.

### 13.16 Test functions and Paley-Wiener spaces

Of course, the original [Paley-Wiener 1934] referred to $L^{2}$ functions, not distributions. The distributional aspect is from [Schwartz 1952]. An interesting point is that rate-of-growth of the Fourier transforms in the imaginary part determines the support of the inverse Fourier transforms.

The class $P W$ of entire functions appearing in the following theorem is the Paley-Wiener space in one complex variable. The assertion is that, in contrast to the fact that Fourier transform maps the Schwartz space to itself, on test functions the Fourier transform has less symmetrical behavior, bijecting to the PaleyWiener space.
[13.16.1] Theorem: A test function $f$ supported on $[-r, r] \subset \mathbb{R}$ has Fourier transform $\widehat{f}$ extending to an entire function on $\mathbb{C}$, with

$$
|\widehat{f}(z)| \ll N N(1+|z|)^{-N} e^{r \cdot|y|} \quad(\text { for } z=x+i y \in \mathbb{C}, \text { for every } N)
$$

Conversely, an entire function satisfying such an estimate has (inverse) Fourier transform which is a test function supported in $[-r, r]$.
Proof: First, the integral for $\widehat{f}(z)$ is the integral of the compactly-supported, continuous, entire-functionvalued function,

$$
\xi \longrightarrow\left(z \rightarrow f(\xi) \cdot e^{-i \xi z}\right)
$$

where the space of entire functions is given the sups-on-compacts semi-norms $\sup _{z \in K}|f(z)|$. Since $\mathbb{C}$ can be covered by countably-many compacts, this topology is metrizable. Cauchy's integral formula proves completeness, so this space is Fréchet. Thus, the Gelfand-Pettis integral exists, and is entire. Multiplication by $z$ is converted to differentiation inside the integral,

$$
(-i z)^{N} \cdot \widehat{f}(z)=\int_{|\xi| \leq r} \frac{\partial^{N}}{\partial \xi^{N}} e^{-i z \cdot \xi} \cdot f(\xi) d \xi=(-1)^{N} \int_{|\xi| \leq r} e^{-i z \cdot \xi} \cdot \frac{\partial^{N}}{\partial \xi^{N}} f(\xi) d \xi
$$

by integration by parts. Differentiation does not enlarge support, so

$$
\begin{gathered}
|\widehat{f}(z)|<_{N}(1+|z|)^{-N} \cdot\left|\int_{|\xi| \leq r} e^{-i z \cdot \xi} f^{(N)}(\xi) d \xi\right| \leq(1+|z|)^{-N} \cdot e^{r \cdot|y|} \cdot\left|\int_{|\xi| \leq r} e^{-i x \cdot \xi} f^{(N)}(\xi) d \xi\right| \\
\leq(1+|z|)^{-N} \cdot e^{r \cdot|y|} \cdot \int_{|\xi| \leq r}\left|f^{(N)}(\xi)\right| d \xi<_{f, N}(1+|z|)^{-N} \cdot e^{r \cdot|y|}
\end{gathered}
$$

Conversely, for an entire function $F$ with the indicated growth and decay property, we show that

$$
\varphi(\xi)=\int_{\mathbb{R}} e^{i x \xi} F(x) d x
$$

is a test function with support inside $[-r, r]$. The assumptions on $F$ do not directly include any assertion that $F$ is Schwartz, so we cannot directly conclude that $\varphi$ is smooth. Nevertheless, a similar obvious computation would give

$$
\int_{\mathbb{R}}(i x)^{N} \cdot e^{i x \xi} F(x) d x=\int_{\mathbb{R}} \frac{\partial^{N}}{\partial \xi^{N}} e^{i x \xi} F(x) d x=\frac{\partial^{N}}{\partial \xi^{N}} \int_{\mathbb{R}} e^{i x \xi} F(x) d x
$$

Moving the differentiation outside the integral is necessary, justified via Gelfand-Pettis integrals by a compactification device, as in [14.3], as follows. Since $F$ strongly vanishes at $\infty$, the integrand extends continuously to the stereographic-projection one-point compactification of $\mathbb{R}$, giving a compactly-supported smooth-function-valued function on this compactification. The measure on the compactification can be adjusted to be finite, taking advantage of the rapid decay of $F$ :

$$
\varphi(\xi)=\int_{\mathbb{R}} e^{i x \xi} F(x) d x=\int_{\mathbb{R}} e^{i x \xi} F(x)\left(1+x^{2}\right)^{N} \frac{d x}{\left(1+x^{2}\right)^{N}}
$$

Thus, the Gelfand-Pettis integral exists, and $\varphi$ is smooth. Thus, in fact, the justification proves that such an integral of smooth functions is smooth without necessarily producing a formula for derivatives.

To see that $\varphi$ is supported inside $[-r, r]$, observe that, taking $y$ of the same sign as $\xi$,

$$
\left|F(x+i y) \cdot e^{i \xi(x+i y)}\right|<_{N}(1+|z|)^{-N} \cdot e^{(r-|\xi|) \cdot|y|}
$$

Thus,

$$
|\varphi(\xi)|<_{N} \int_{\mathbb{R}}(1+|z|)^{-N} \cdot e^{(r-|\xi|) \cdot|y|} d x \leq e^{(r-|\xi|) \cdot|y|} \cdot \int_{\mathbb{R}} \frac{d x}{(1+|x|)^{-N}}
$$

For $|\xi|>r$, letting $|y| \rightarrow+\infty$ shows that $\varphi(\xi)=0$.
[13.16.2] Corollary: We can topologize $P W$ by requiring that the linear bijection $\mathcal{D} \rightarrow P W$ be a topological vector space isomorphism.
[13.16.3] Remark: The latter topology on $P W$ is finer than the sups-on-compacts topology on all entire functions, since the latter cannot detect growth properties.
[13.16.4] Corollary: Fourier transform can be defined on all distributions $u \in \mathcal{D}^{*}$ by $\widehat{u}(\varphi)=u(\widehat{\varphi})$ for $\varphi \in P W$, giving an isomorphism $\mathcal{D}^{*} \rightarrow P W^{*}$ to the dual of the Paley-Wiener space.

For example, the exponential $t \rightarrow e^{i z \cdot t}$ with $z \in \mathbb{C}$ but $z \notin \mathbb{R}$ is not a tempered distribution, but is a distribution, and its Fourier transform is the Dirac delta $\delta_{z} \in P W^{\prime}$.

Compactly-supported distributions have a similar characterization:
[13.16.5] Theorem: The Fourier transform $\widehat{u}$ of a distribution $u$ supported in $[-r, r]$, of order $N$, is (integration against) the function $x \rightarrow u\left(\xi \rightarrow e^{-i x \xi}\right)$, which is smooth, and extends to an entire function satisfying

$$
|\widehat{u}(z)| \ll(1+|z|)^{N} \cdot e^{r \cdot|y|}
$$

Conversely, an entire function meeting such a bound is the Fourier transform of a distribution of order $N$ supported inside $[-r, r]$.

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Proof: The Fourier transform $\widehat{u}$ is the tempered distribution defined for Schwartz functions $\varphi$ by

$$
\widehat{u}(\varphi)=u(\widehat{\varphi})=u\left(\xi \rightarrow \int_{\mathbb{R}} e^{-i x \xi} \varphi(x) d x\right)=\int_{\mathbb{R}} u\left(\xi \rightarrow e^{-i x \xi}\right) \varphi(x) d x
$$

since $x \rightarrow\left(\xi \rightarrow e^{-i x \xi} \varphi(\xi)\right.$ extends to a continuous smooth-function-valued function on the one-point compactification of $\mathbb{R}$, and Gelfand-Pettis applies. Thus, as expected, $\widehat{u}$ is integration against $x \rightarrow u(\xi \rightarrow$ $\left.e^{-i x \xi}\right)$.

The smooth-function-valued function $z \rightarrow\left(\xi \rightarrow e^{-i z \xi}\right)$ is holomorphic in $z$. Compactly-supported distributions constitute the dual of $C^{\infty}(\mathbb{R})$. Application of $u$ gives a holomorphic scalar-valued function $z \rightarrow u\left(\xi \rightarrow e^{-i z \xi}\right)$.

Let $\nu_{N}$ be the $N^{t h}$-derivative seminorm on $C^{\infty}[-r, r]$, so

$$
|u(\varphi)|<_{\varepsilon} \nu_{N}(\varphi)
$$

Then

$$
|\widehat{u}(z)|=\left|u\left(\xi \rightarrow e^{-i z \xi}\right)\right| \lll \varepsilon \nu_{N}\left(\xi \rightarrow e^{-i z \xi}\right) \ll \sup _{[-r, r]}\left|(1+|z|)^{N} e^{-i z \xi}\right| \leq(1+|z|)^{N} e^{r \cdot|y|}
$$

Conversely, let $F$ be an entire function with $|F(z)| \ll(1+|z|)^{N} e^{r \cdot|y|}$. Certainly $F$ is a tempered distribution, so $F=\widehat{u}$ for a tempered distribution. We show that $u$ is of order at most $N$ and has support in $[-r, r]$.

With $\eta$ supported on $[-1,1]$ with $\eta \geq 0$ and $\int \eta=1$, make an approximate identity $\eta_{\varepsilon}(x)=\eta(x / \varepsilon) / \varepsilon$ for $\varepsilon \rightarrow 0^{+}$. By the easy half of Paley-Wiener for test functions, $\widehat{\eta}_{\varepsilon}$ is entire and satisfies

$$
\left|\widehat{\eta}_{\varepsilon}(z)\right|<_{\varepsilon, N}(1+|z|)^{-N} \cdot e^{\varepsilon \cdot|y|} \quad(\text { for all } N)
$$

Note that $\widehat{\eta}_{\varepsilon}(x)=\widehat{\eta}(\varepsilon \cdot x)$ goes to 1 as tempered distribution
By the more difficult half of Paley-Wiener for test functions, $F \cdot \widehat{\eta}_{\varepsilon}$ is $\widehat{\varphi}_{\varepsilon}$ for some test function $\varphi_{\varepsilon}$ supported in $[-(r+\varepsilon), r+\varepsilon]$. Note that $F \cdot \widehat{\eta}_{\varepsilon} \rightarrow F$.

For Schwartz function $g$ with the support of $\widehat{g}$ not meeting $[-r, r], \widehat{g} \cdot \varphi_{\varepsilon}$ for sufficiently small $\varepsilon>0$. Since $F \cdot \widehat{\eta}_{\varepsilon}$ is a Cauchy net as tempered distributions,

$$
u(\widehat{g})=\widehat{u}(g)=\int F \cdot g=\int \lim _{\varepsilon}\left(F \cdot \widehat{\eta}_{\varepsilon}\right) g=\lim _{\varepsilon} \int\left(F \cdot \widehat{\eta}_{\varepsilon}\right) g=\lim _{\varepsilon} \int \widehat{\varphi}_{\varepsilon} g=\lim _{\varepsilon} \int \varphi_{\varepsilon} \widehat{g}=0
$$

This shows that the support of $u$ is inside $[-r, r]$.

### 13.17 Schwartz functions and Fourier transforms on $\mathbb{Q}_{p}$

For simplicity, we only look at Fourier analysis on $\mathbb{Q}_{p}$, rather than on general $p$-adic fields. The same ideas apply to the general case, with minor modifications.

Fix a prime $p$, let $\mathbb{Q}_{p}$ be the $p$-adic field and $\mathbb{Z}_{p}$ the $p$-adic integers. Give $\mathbb{Q}_{p}$ the additive Haar measure that gives $\mathbb{Z}_{p}$ total measure 1 . This determines the measure of every set $x+p^{n} \mathbb{Z}_{p}$ with $n \geq 0$, by translationinvariance, and the fact that $\mathbb{Z}_{p}$ is a disjoint union of such translates, as $x$ ranges over $\mathbb{Z}_{p} / p^{n} \mathbb{Z}_{p} \approx \mathbb{Z} / p^{n} \mathbb{Z}$. The standard choice of additive character, trivial on $\mathbb{Z}_{p}$, is $\psi_{1}(x)=e^{-2 \pi i x^{\prime}}$, where $x^{\prime} \in \mathbb{Z}\left[\frac{1}{p}\right]$ is such that $x-x^{\prime} \in \mathbb{Z}_{p}$. Parametrize additive characters by $\psi_{\xi}(x)=\psi_{1}(\xi \cdot x)$.

Unsurprisingly, the Fourier transform on $\mathbb{C}$-valued $L^{1}$ functions on $\mathbb{Q}_{p}$ is

$$
\mathscr{F} f(\xi)=\widehat{f}(\xi)=\int_{\mathbb{Q}_{p}} \bar{\psi}_{\xi}(x) f(x) d x
$$

The space of Schwartz functions $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ on $\mathbb{Q}_{p}$ should be mapped to itself homeomorphically under Fourier transform, should consist of very simple functions, and should be dense in $L^{2}\left(\mathbb{Q}_{p}\right)$. We will show that the following choice succeeds: take

$$
\mathscr{S}\left(\mathbb{Q}_{p}\right)=\{\text { compactly-supported, locally constant, } \mathbb{C} \text {-valued functions }\}
$$

where $f$ being locally constant means that every $x \in \mathbb{Q}_{p}$ has a neighborhood $U$ such that $f\left(x^{\prime}\right)=f(x)$ for $x^{\prime} \in U$.
[13.17.1] Remark: The local constancy turns out to be the appropriate $p$-adic notion of smoothness. Unlike the archimedean case, $p$-adic Schwartz functions are compactly supported. That is, in the $p$-adic case, test functions and Schwartz functions are the same classes of functions.
[13.17.2] Claim: $f \in \mathscr{S}\left(\mathbb{Q}_{p}\right)$ is uniformly locally constant: there is a (compact, open) subgroup $U=p^{n} \mathbb{Z}_{p}$ such that $f(x+u)=f(x)$ for all $x \in \mathbb{Q}_{p}$, and for all $u \in U$.
Proof: Since $\bigcup_{m \geq 0} p^{-m} \mathbb{Z}_{p}=\mathbb{Q}_{p}$, the support of a given $f \in \mathscr{S}\left(\mathbb{Q}_{p}\right)$ is contained in some $p^{-m} \mathbb{Z}_{p}$. For each $x \in p^{-m} \mathbb{Z}_{p}$, there is a neighborhood $x+p^{n_{x}} \mathbb{Z}_{p}$ on which $f$ is constant. By compactness of $p^{-m} \mathbb{Z}_{p}$, there are finitely-many points $x_{1}, \ldots, x_{\ell}$ so that the corresponding neighborhoods cover $p^{-m} \mathbb{Z}_{p}$. Let $n=\max _{1 \leq i \leq \ell} n_{x_{i}}$ and $U=p^{n} \mathbb{Z}_{p}$. A given $x \in p^{-m} \mathbb{Z}_{p}$ lies in $x_{j}+p^{n_{x_{j}}} \mathbb{Z}_{p}$ for some $j$, and

$$
x+U \subset x_{j}+p^{n_{x_{j}}} \mathbb{Z}_{p}+U \subset x_{j}+p^{n_{x_{j}}} \mathbb{Z}_{p}+p^{n_{x_{j}}} \mathbb{Z}_{p}+U=x_{j}+p^{n_{x_{j}}} \mathbb{Z}_{p}
$$

since every $p^{n} \mathbb{Z}_{p}$ is closed under addition. Thus, $f$ is locally constant on $x+U$.
[13.17.3] Corollary: $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ is a strict colimit of the finite-dimensional subspaces

$$
V_{m, n}=\left\{f \in \mathscr{S}\left(\mathbb{Q}_{p}\right): \operatorname{spt} f \subset p^{-m} \mathbb{Z}_{p}, f(x+u)=f(x) \text { for all } x, \text { for all } u \in p^{n} \mathbb{Z}_{p}\right\}
$$

In particular, $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ consists of finite linear combinations of characteristic functions of sets $x_{o}+p^{n} \mathbb{Z}_{p}$.
Proof: The lemma asserts that $\mathscr{S}\left(\mathbb{Q}_{p}\right)=\bigcup_{m, n} V_{m, n}$. Since $p^{-m} \mathbb{Z}_{p}$ is the disjoint union of $p^{m+n}$ distinct cosets $x_{o}+p^{n} \mathbb{Z}_{p}$, the subspace $V_{m, n}$ is the collection of linear combinations of characteristic functions of these sets.

Thus, the Schwartz space $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ is not Fréchet, but is the simplest type of LF-space, namely, a strict colimit of finite-dimensional spaces (and finite-dimensional spaces have unique topologies [13.4]) like $\mathbb{C}^{\infty}$ in [13.8].

The following holds for Schwartz functions by direct computation, and then will follow for $L^{2}$ functions by denseness of $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ in $L^{2}\left(\mathbb{Q}_{p}\right)$ and extending by $L^{2}$-continuity.
[13.17.4] Theorem: For Schwartz functions, Fourier inversion holds:

$$
f(x)=\int_{k_{v}} \psi_{\xi}(x) \widehat{f}(\xi) d \xi \quad\left(\text { for } f \in \mathscr{S}\left(\mathbb{Q}_{p}\right)\right)
$$

and Plancherel's theorem holds:

$$
\int_{\mathbb{Q}_{p}}|f|^{2}=\int_{\mathbb{Q}_{p}}|\widehat{f}|^{2} \quad\left(\text { for } f \in \mathscr{S}\left(\mathbb{Q}_{p}\right)\right)
$$

Proof: For Schwartz functions, we prove more, by giving sample computations of Fourier transforms which are useful. In particular, we observe simply-described functions on $\mathbb{Q}_{p}$ whose Fourier transforms are of a similar nature. For example, certain natural functions in $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ are their own Fourier transform, analogous to the Gaussian in the archimedean case.
[13.17.5] Claim: The characteristic function $f$ of $\mathbb{Z}_{p}$ is its own Fourier transform.
Proof: Computing directly,

$$
\widehat{f}(\xi)=\int_{\mathbb{Q}_{p}} \bar{\psi}_{\xi}(x) f(x) d x=\int_{\mathbb{Z}_{p}} \bar{\psi}_{1}(\xi \cdot x) d x=\int_{\mathbb{Z}_{p}} \psi_{1}(-\xi \cdot x) d x
$$

Recall a form of the cancellation lemma: (a tiny case of Schur orthogonality)
[13.17.6] Lemma: Let $\psi: K \rightarrow \mathbb{C}^{\times}$be a continuous group homomorphism on a compact group $K$. Then

$$
\int_{K} \psi(x) d x=\left\{\begin{array}{cl}
\operatorname{meas}(K) & (\text { for } \psi=1) \\
0 & (\text { for } \psi \neq 1)
\end{array}\right.
$$

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Proof: (of Lemma) Yes, of course, the measure is a Haar measure on $K$. Since $K$ is compact, it is unimodular. For $\psi$ trivial, of course the integral is the total measure of $K$. For $\psi$ non-trivial, there is $y \in K$ such that $\psi(y) \neq 1$. Using the invariance of the measure, change variables by replacing $x$ by $x y$ :

$$
\int_{K} \psi(x) d x=\int_{K} \psi(x y) d(x y)=\int_{K} \psi(x) \psi(y) d x=\psi(y) \int_{K} \psi(x) d x
$$

Since $\psi(y) \neq 1$, the integral is 0 .
Apply the lemma to the integrals computing the Fourier transform of the characteristic function $f$ of $\mathbb{Z}_{p}$. Since $\mathbb{Z}_{p}$ has measure 1,

$$
\widehat{f}(\xi)=\int_{\mathbb{Z}_{p}} \psi_{1}(-\xi \cdot x) d x= \begin{cases}1 & \left(\psi_{1}(-\xi x)=1 \text { for } x \in \mathbb{Z}_{p}\right) \\ 0 & \text { (otherwise) }\end{cases}
$$

On one hand, for $\xi \in \mathbb{Z}_{p}$, certainly $\psi_{1}(\xi x)=1$ for $x \in \mathbb{Z}_{p}$. On the other hand, for $\xi \notin \mathbb{Z}_{p}$, there is $x \in \mathbb{Z}_{p}$ such that, for example, $\xi \cdot x=1 / p$. Then

$$
\psi_{1}(-\xi \cdot x)=\psi_{1}\left(\frac{-1}{p}\right)=e^{+2 \pi i \cdot \frac{1}{p}} \neq 1
$$

Thus, $\psi_{\xi}$ is not trivial on $\mathbb{Z}_{p}$, so the integral is 0 . Thus, the characteristic function of $\mathbb{Z}_{p}$ is its own Fourier transform.
[13.17.7] Claim: The Fourier transform of the characteristic function of $p^{k} \mathbb{Z}_{p}$ is $p^{-k}$ times the characteristic function of $p^{-k} \mathbb{Z}_{p}$.
Proof: Let $f$ be the characteristic function of $p^{k} \mathbb{Z}_{p}$, so

$$
\begin{gathered}
\widehat{f}(\xi)=\int_{\mathbb{Q}_{p}} \bar{\psi}_{\xi}(x) f(x) d x=\int_{p^{k} \mathbb{Z}_{p}} \bar{\psi}_{1}(\xi \cdot x) d x \\
=\left|p^{k}\right|_{p} \cdot \int_{\mathbb{Z}_{p}} \psi_{1}\left(-\xi \cdot x / p^{k}\right) d x=p^{-k} \cdot \int_{\mathbb{Z}_{p}} \psi_{1}\left(-\xi \cdot x / p^{k}\right) d x
\end{gathered}
$$

This reduces to the previous computation: by cancellation, for $\xi / p^{k} \notin \mathbb{Z}_{p}$ the character $x \rightarrow \psi_{1}\left(-\xi x / p^{k}\right)$ is non-trivial, so the integral is 0 . Otherwise, the integral is 1 .
[13.17.8] Claim: The Fourier transform of the characteristic function of $\mathbb{Z}_{p}+y$ is $\psi_{y}$ times the characteristic function of $\mathbb{Z}_{p}$.
Proof: Let $f$ be the characteristic function of $\mathbb{Z}_{p}+y$, so

$$
\begin{gathered}
\widehat{f}(\xi)=\int_{\mathbb{Q}_{p}} \bar{\psi}_{\xi}(x) f(x) d x=\int_{\mathbb{Z}_{p}+y} \bar{\psi}_{1}(\xi \cdot x) d x=\int_{\mathbb{Z}_{p}} \psi_{1}(-\xi \cdot(x+y)) d x \\
=\psi_{1}(-\xi \cdot y) d x \int_{\mathbb{Z}_{p}} \psi_{1}(-\xi \cdot x) d x=\psi_{1}(-\xi \cdot y) \cdot f(\xi)
\end{gathered}
$$

by the previous computation.
Combining the two computations above,

$$
\mathscr{F}\left(\operatorname{char} \mathrm{fcn} p^{k} \mathbb{Z}_{p}+y\right)=\psi_{y} \cdot p^{-k} \cdot\left(\operatorname{char} \operatorname{fcn} p^{-k} \mathbb{Z}_{p}\right)
$$

Conveniently, products $\psi_{y} \cdot\left(\operatorname{char}\right.$ fcn $\left.p^{-k} \mathbb{Z}_{p}\right)$ are in the same class of functions, since $\psi_{y}$ has a kernel which is an open (and compact) neighborhood of 0 , so we this class of functions is mapped to itself under Fourier transform.

We have essentially proven Fourier inversion, in the computations above, as follows. Let $f^{o}$ be the characteristic function of $\mathbb{Z}_{p}$. We computed $\widehat{f^{o}}=f$. Let $\delta_{t}$ be the dilation operator $\delta_{t} f(x)=f(t \cdot x)$ for $t \in \mathbb{Q}_{p}^{\times}$. We computed, by changing variables in the integral defining the Fourier transform, that

$$
\mathscr{F}\left(\delta_{t} f\right)=\frac{1}{|t|_{p}} \cdot \delta_{1 / t}(\mathscr{F} f)
$$

Let $\tau_{y}$ be the translation operator $\tau_{y} f(x)=f(x+y)$. By changing variables,

$$
\mathscr{F}\left(\tau_{y} f\right)=\psi_{y} \cdot(\mathscr{F} f)
$$

It is convenient to also compute that

$$
\mathscr{F}\left(\psi_{y} \cdot f\right)(\xi)=\int_{\mathbb{Q}_{p}} \bar{\psi}_{\xi}(x) \cdot \psi_{y}(x) f(x) d x=\int_{\mathbb{Q}_{p}} \bar{\psi}_{\xi-y}(x) f(x) d x=\widehat{f}(\xi-y)=\tau_{-y}(\mathscr{F} f)
$$

Let $\mathscr{F}^{*}$ be the integral for Fourier inversion, namely,

$$
\mathscr{F}^{*} f(x)=\int_{\mathbb{Q}_{p}} \psi_{\xi}(x) f(\xi) d \xi
$$

Similar computations give

$$
\mathscr{F}^{*}\left(\delta_{t} f\right)=\frac{1}{|t|_{p}} \delta_{1 / t}\left(\mathscr{F}^{*} f\right) \quad \mathscr{F}^{*}\left(\tau_{y} f\right)=\psi_{-y}\left(\mathscr{F}^{*} f\right)
$$

and

$$
\mathscr{F}^{*}\left(\psi_{y} f\right)=\tau_{y}\left(\mathscr{F}^{*} f\right)
$$

Since every element of $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ is a linear combination of images of $f^{o}$ under dilation and translation, it suffices to give a sort of inductive proof of Fourier inversion:

$$
\begin{gathered}
\mathscr{F}^{*} \mathscr{F}\left(\tau_{y} f\right)==\mathscr{F}^{*} \psi_{y} \mathscr{F} f=\tau_{y} \mathscr{F}^{*} \mathscr{F} f \\
\mathscr{F}^{*} \mathscr{F}\left(\delta_{t} f\right)=\mathscr{F}^{*} \frac{1}{|t|_{p}} \delta_{1 / t} \mathscr{F} f=\frac{1}{|t|_{p}} \frac{1}{|1 / t|_{p}} \delta_{t} \mathscr{F}^{*} \mathscr{F} f=\delta_{t} \mathscr{F}^{*} \mathscr{F} f
\end{gathered}
$$

Similarly for multiplication by $\psi_{y}$. Since $\mathscr{F} * \mathscr{F} f^{o}=\mathscr{F}^{*} f^{o}=f^{o}$, we have Fourier inversion on $\mathscr{S}\left(\mathbb{Q}_{p}\right)$.
The surjectivity of $\mathscr{F}: \mathscr{S}\left(\mathbb{Q}_{p}\right) \rightarrow \mathscr{S}\left(\mathbb{Q}_{p}\right)$ is made explicit in the computations above. Then we have the Plancherel theorem on $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ :

$$
\begin{gathered}
\int_{\mathbb{Q}_{p}} f \cdot \bar{g}=\int_{\mathbb{Q}_{p}} f \cdot \overline{\mathscr{F}-1 \widehat{g}}=\int_{\mathbb{Q}_{p}} \int_{\mathbb{Q}_{p}} f(x) \cdot \psi_{1}(-\xi x) \cdot \overline{\bar{g}}(\xi) d \xi d x \\
\quad=\int_{\mathbb{Q}_{p}}\left(\int_{\mathbb{Q}_{p}} f(x) \cdot \psi_{1}(-\xi x) d x\right) \cdot \bar{g}(\xi) d \xi=\int_{\mathbb{Q}_{p}} \hat{f} \cdot \bar{g}
\end{gathered}
$$

This proves the theorem for Schwartz functions.
Similarly, and as for $C_{c}^{o}(\mathbb{R})$, the space $C_{c}^{o}(\mathbb{Q})$ of compactly-supported, continuous, $\mathbb{C}$-valued functions on $\mathbb{Q}_{p}$ is an LF-space, the strict colimit of the spaces of continuous functions supported on $p^{-n} \mathbb{Z}_{p}$. Much as in [6.2] and [13.9], we have
[13.17.9] Claim: The translation action $\mathbb{Q}_{p} \times C_{c}^{o}(\mathbb{Q}) \longrightarrow C_{c}^{o}(\mathbb{Q})$ by $(x \cdot f)(y)=f(y+x)$ is (jointly) continuous. Proof: Since $\mathbb{Q}_{p}$ itself is the colimit of $p^{-m} \mathbb{Z}_{p}$ (as additive topological group), and $C_{c}^{o}\left(\mathbb{Q}_{p}\right)$ is a colimit, it suffices to show that $p^{-m} \mathbb{Z}_{p} \times C^{o}\left(p^{-n} \mathbb{Z}_{p}\right) \longrightarrow C_{c}^{o}\left(\mathbb{Q}_{p}\right)$ is continuous for all $m, n$. Indeed, $p^{-m} \mathbb{Z}_{p} \times C^{o}\left(p^{-n} \mathbb{Z}_{p}\right)$ maps to $C^{o}\left(p^{-\max (m, n)} \mathbb{Z}_{p}\right)$, and sup-norms are preserved.

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[13.17.10] Claim: $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ is dense in $C_{c}^{o}\left(\mathbb{Q}_{p}\right)$.
Proof: This is a simple $p$-adic analogue of the smoothing of distributions [14.5], and of Gårding 's theorem [14.6], asserting that smooth vectors are dense in a representation, also following from the basic result [14.1.4] about approximate identities and Gelfand-Pettis integrals. Namely, let $\varphi_{n}=p^{n} \cdot \chi_{n}$, where $\chi_{n}$ is the characteristic function of $p^{n} \mathbb{Z}_{p}$. These are continuous, compactly-supported functions, and form an approximate identity in $C_{c}^{o}\left(\mathbb{Q}_{p}\right)$ in the sense that they are non-negative, their integrals are all 1 , and their supports shrink to $\{0\}$. By the previous claim, $\mathbb{Q}_{p}$ acts continuously on $C_{c}^{o}\left(\mathbb{Q}_{p}\right)$, giving integral operators

$$
\left(\varphi_{n} \cdot f\right)(x)=\int_{\mathbb{Q}_{p}} \varphi_{n}(y) f(x+y) d y
$$

on $f \in C_{c}^{o}\left(\mathbb{Q}_{p}\right)$. By [14.1.4], $\varphi_{n} \cdot f \longrightarrow f$ in $C_{c}^{o}\left(\mathbb{Q}_{p}\right)$.
Analogous to the archimedean discussion in the proof of smoothing of distributions theorem [14.5], we check that each $\varphi_{n} \cdot f$ is locally constant and compactly supported, so is in $\mathscr{S}\left(\mathbb{Q}_{p}\right)$. The compact support is clear, since the support of $\varphi_{n} \cdot f$ is contained in $\operatorname{spt}\left(\varphi_{n}\right)+\operatorname{spt}(f)$, which is compact, being the image of the compact $\operatorname{spt}\left(\varphi_{n}\right) \times \operatorname{spt}(f) \subset \mathbb{Q}_{p} \times \mathbb{Q}_{p}$ under the continuous map $x \times y \rightarrow x+y$. For local constancy, for $u \in p^{n} \mathbb{Z}_{p}$,
$\left(\varphi_{n} \cdot f\right)(x+u)=\int_{\mathbb{Q}_{p}} \varphi_{n}(y) f(x+u+y) d y=p^{n} \int_{p^{n} \mathbb{Z}_{p}} f(x+u+y) d y=p^{n} \int_{p^{n} \mathbb{Z}_{p}} f(x+y) d y=\left(\varphi_{n} \cdot f\right)(x)$
by changing variables $y \rightarrow y-u$, since $p^{n} \mathbb{Z}_{p}$ is a group.
[13.17.11] Corollary: $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ is dense in $L^{1}\left(\mathbb{Q}_{p}\right)$ and in $L^{2}\left(\mathbb{Q}_{p}\right)$.
Proof: $C_{c}^{o}\left(\mathbb{Q}_{p}\right)$ is dense in both $L^{1}\left(\mathbb{Q}_{p}\right)$ and $L^{2}\left(\mathbb{Q}_{p}\right)$, essentially by Urysohn's Lemma [9.E.2], as in [6.1], [6.2], so $\mathscr{S}\left(\mathbb{Q}_{p}\right)$ is dense in both, by the previous.
[13.17.12] Corollary: $\mathscr{F}$ extends to $L^{2}\left(\mathbb{Q}_{p}\right)$ by continuity, giving the Fourier-Plancherel transform $\mathscr{F}: L^{2}\left(\mathbb{Q}_{p}\right) \longrightarrow L^{2}\left(\mathbb{Q}_{p}\right)$, no longer defined literally by the integrals, but still satisfying Fourier inversion and Plancherel theorem.

## 14. Vector-valued integrals

## 14. Vector-valued integrals

1. Characterization and basic results
2. Differentiation of parametrized integrals
3. Fourier transforms
4. Uniqueness of invariant distributions
5. Smoothing of distributions
6. Density of smooth vectors
7. Quasi-completeness and convex hulls of compacts
8. Existence proof

Appendix A: Hahn-Banach theorems

Quasi-complete, locally convex topological vector spaces $V$ have the useful property that continuous compactly-supported $V$-valued functions have integrals with respect to finite, regular Borel measures. Rather than being constructed as limits, these vector-valued integrals are characterized. Uniqueness follows from the Hahn-Banach theorem, and existence follows from a construction.

An immediate application is justification of differentiation with respect to a parameter inside an integral, under mild, easily understood hypotheses, a special case of the general assertion that Gelfand-Pettis integrals commute with continuous operators, as in the first section. A subtler application is to passage of compactly-supported distributions inside the integrals expressing Fourier inversion, as in [14.3]. Uniqueness of group-invariant measures, distributions, and other functionals is another corollary. Other applications are to holomorphic vector-valued functions, to holomorphically parametrized families of generalized functions (distributions), as in chapter 14. Many distributions which are not classical functions appear naturally as residues or analytic continuations of meromorphic families of classical functions.

### 14.1 Characterization and basic results

For a topological vectorspace $V$ over $\mathbb{C}$ and for $f$ a continuous $V$-valued function on a topological space $X$ with a regular Borel measure, a Gelfand-Pettis integral of $f$ is a vector $I_{f} \in V$ so that

$$
\lambda\left(I_{f}\right)=\int_{X} \lambda \circ f \quad\left(\text { for all } \lambda \in V^{*}\right)
$$

If it exists and is unique, this vector $I_{f}$ is reasonably denoted

$$
I_{f}=\int_{X} f
$$

In contrast to construction of integrals as limits, this characterization surely should apply to any reasonable notion of integral, without asking how the property comes to be. Since the property of allowing continuous linear functionals to pass inside the integral is an irreducible minimum, the Gelfand-Pettis integral is sometimes called a weak integral.

We only consider locally convex vectorspaces, so uniqueness of the integral is immediate, since $V^{*}$ separates points on $V$, by Hahn-Banach. Similarly, for such $V$, linearity of $f \rightarrow I_{f}$ follows by Hahn-Banach. The issue is existence. [80] We only consider $V$-valued functions that are continuous on compact measure spaces with regular Borel measures. Under these assumptions, all the $\mathbb{C}$-valued integrals

$$
f \longrightarrow f \circ \lambda \longrightarrow \int_{X} \lambda \circ f \quad\left(\text { for } \lambda \in V^{*}\right)
$$

[80] We want the integral to be in $V$ itself, rather than in a larger space containing $V$, such as a double dual $V^{* *}$, for example, to make existence trivial, but then leaving technical issues. Some discussions of vector-valued integration do allow integrals to exist in larger spaces, but this only delays certain issues, rather than resolving them directly.

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exist for elementary reasons, being integrals of compactly-supported $\mathbb{C}$-valued continuous functions on compact sets with respect to a regular Borel measure.

For existence of Gelfand-Pettis integrals of compactly-supported, continuous $V$-valued functions, the literal requirement on $V$ turns out to be that the closure of the convex hull of a compact set is compact. We show below that local convexity and quasi-completeness suffice. For the following, a probability measure is a positive, regular, Borel measure with total measure 1.
[14.1.1] Theorem: Let $X$ be a compact Hausdorff topological space with a probability measure. Let $V$ be a quasi-complete, locally convex vectorspace. Then continuous $V$-valued functions $f$ on $X$ have Gelfand-Pettis integrals. The basic estimate holds:

$$
\int_{X} f \in(\text { closure of convex hull of } f(X))
$$

substituting for the estimate of a $\mathbb{C}$-valued integral by the integral of its absolute value. (Proof in [14.8].)
[14.1.2] Corollary: In the situation of the theorem, but when the total measure of $X$ is finite but not necessarily 1 , the basic estimate becomes

$$
\int_{X} f \in(\text { closure of convex hull of } f(X)) \cdot \int_{X} 1
$$

(Replace the measure by a constant multiple.)
[14.1.3] Corollary: For a continuous linear map of locally convex, quasi-complete topological vectorspaces $T: V \rightarrow W$, and $f$ a continuous, compactly-supported $V$-valued function on a finite, regular, positive Borel measure space $X$. Then

$$
T\left(\int_{X} f\right)=\int_{X} T \circ f
$$

Proof: To verify that the left-hand side of the asserted equality is a Gelfand-Pettis integral of $T \circ f$, show that

$$
\mu(\text { left-hand side })=\int_{X} \mu \circ(T \circ f) \quad\left(\text { for all } \mu \in W^{*}\right)
$$

Starting with the left-hand side,

$$
\begin{aligned}
\mu\left(T\left(\int_{X} f\right)\right) & =(\mu \circ T)\left(\int_{X} f\right) \quad \text { (associativity) } \\
& =\int_{X}(\mu \circ T) \circ f \quad\left(\mu \circ T \in V^{*} \text { and } \int_{X} f\right. \text { is a weak integral) } \\
& =\int_{X} \mu \circ(T \circ f) \quad \text { (associativity) }
\end{aligned}
$$

proving that $T\left(\int_{X} f\right)$ is a weak integral of $T \circ f$.
A representation of $G$ on a locally convex, quasi-complete topological vectorspace $V$ is a continuous map $G \times V \rightarrow V$ that is linear in $V$, has the associativity $(g h) \cdot v=g \cdot(h \cdot v)$ for $g, h \in G$, and $1_{G} \cdot v=v$ for all $v \in V$.

For any $v \in V$ and $\varphi \in C_{c}^{o}(G)$, we have a $V$-valued Gelfand-Pettis integral

$$
\varphi \cdot f=\int_{G} \varphi(g) T_{g} f d g \in \text { closure of convex hull of }\{\varphi(g) f: g \in G\} \subset V
$$

For present purposes, a continuous approximate identity on a topological group $G$ is a sequence $\left\{\varphi_{i}\right\}$ of non-negative, continuous, real-valued functions such that $\int_{G} \varphi_{i}=1$ for all $i$, and such that the supports shrink to $\{1\}$, in the sense that for every neighborhood $N$ of 1 in $G$, there is an index $i_{o}$ so that the support of $\varphi_{i}$ is inside $N$ for all $i \geq i_{o}$. From Urysohn's Lemma [9.E.2], there always exists a continuous approximate identity. Let $T_{g} f(y)=f(y g)$ be right translation. With right translation-invariant measure $d g$ on $G$, since

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$\int_{G} \varphi_{i}(g) d g=1$ and $\varphi_{i}$ is non-negative, $\varphi_{i}(g) d g$ is a probability measure (total mass 1 ) on the (compact) support of $\varphi_{i}$.
[14.1.4] Corollary: Given a representation of $G$ on a quasi-complete, locally convex topological vector space $V$, for every approximate identity $\left\{\varphi_{i}\right\}$ on $G, \varphi_{i} \cdot v \rightarrow v$ for every $v \in V$.
Proof: By continuity, given a neighborhood $N$ of 0 in $V$, we have $\varphi_{i} \cdot f \in f+N$ for all sufficiently large $i$. That is, $\varphi_{i} \cdot f \rightarrow f$.
[14.1.5] Corollary: Given a representation of $G$ on a quasi-complete, locally convex topological vector space $V$, the action of $C_{c}^{o}(G)$ on $V$ is non-degenerate, in the sense that, for every $0 \neq v \in V$, there exists $\varphi \in C_{c}^{o}(G)$ such that $\varphi \cdot v \neq 0$.
Proof: For every approximate identity $\left\{\varphi_{i}\right\}$, for every $v \in V, \varphi_{i} \cdot v \rightarrow v$. Thus, for all sufficiently large $i$, $\varphi_{i} \cdot v \neq 0$ for $v \neq 0$.

A $G$-subrepresentation $W \subset V$ of a representation of $G$ on a quasi-complete, locally convext topological vector space $V$ is a (topologically) closed $G$-stable vector subspace of $V$. Similarly, a $C_{c}^{o}(G)$-subrepresentation $W \subset V$ of a $G$-representation is a (topologically) closed $C_{c}^{o}(G)$-stable vector subspace of $V$.
[14.1.6] Corollary: A $C_{c}^{o}(G)$-subrepresentation of a $G$-representation on a quasi-complete, locally convex topological vector space $V$ is a $G$-representation.
Proof: Let $\left\{\varphi_{i}\right\}$ be an approximate identity, fix $w$ in the subrepresentation $W$, and take $g \in G$. We will show that $g \cdot w \in W$. On one hand, from [14.1.4] above, $\varphi_{i} \cdot(g \cdot w) \rightarrow g \cdot w$. On the other hand,

$$
\varphi_{i} \cdot(g \cdot w)=\int_{G} \varphi_{i}(h) h \cdot(g \cdot v) d h=\int_{G} \varphi_{i}\left(h g^{-1}\right) h \cdot v d h
$$

by changing variables. The function $h \rightarrow \varphi_{i}\left(h g^{-1}\right)$ is still in $W$, by assumption, so $\varphi_{i} \cdot(g \cdot w)$ is a sequence of vectors in $W$. Since $W$ is closed, and the sequence converges to $g \cdot w$, necessarily $g \cdot w \in W$. ///

A representation of $G$ on $W$ is (topologically) $G$-irreducible when there is no proper $G$-subrepresentation. A representation of $G$ on $W$ is (topologically) $C_{c}^{o}(G)$-irreducible when there is no proper $C_{c}^{o}(G)$ subrepresentation.
[14.1.7] Corollary: For a representation of $G$ on a quasi-complete, locally convext topological vector space $V$, every irreducible $C_{c}^{o}(G)$-subrepresentation is an irreducible $G$-subrepresentation.
Proof: From the previous corollary, every $C_{c}^{o}(G)$-subrepresentation $W \subset V$ is a $G$-subrepresentation. If $W$ had a proper $G$-subrepresentation $W^{\prime}$, then $W^{\prime}$ would be a proper $C_{c}^{o}(G)$-subrepresentation, as well. ///

### 14.2 Differentiation of parametrized integrals

Differentiation under the integral is an immediate corollary, in many useful situations.
[14.2.1] Claim: A $\mathbb{C}$-valued $C^{k}$ function $F$ on $[a, b] \times[c, d]$ gives a continuous $C^{k}[c, d]$-valued function $f(x)=F(x,-)$ of $x \in[a, b]$.
Proof: For each $0 \leq i \leq k$, the function $(x, y) \rightarrow \frac{\partial^{i}}{\partial y^{i}} F(x, y)$ is continuous as a function of two variables. For each $\varepsilon>0$ and each $x_{o} \in[a, b]$, we want $\delta>0$ such that

$$
\left|x-x_{o}\right|<\delta \quad \Longrightarrow \quad \sup _{y}\left|\frac{\partial^{i}}{\partial y^{i}} F(x, y)-\frac{\partial^{i}}{\partial y^{i}} F\left(x_{o}, y\right)\right|<\varepsilon
$$

The continuous function $(x, y) \rightarrow \frac{\partial^{k}}{\partial y^{i}} F(x, y)$ is uniformly continuous on the compact $[a, b] \times[c, d]$, so there is $\delta>0$ such that

$$
\left.\left|\frac{\partial^{k}}{\partial y^{i}} F\left(x_{1}, y_{1}\right)-\frac{\partial^{i}}{\partial y^{i}} F\left(x_{2}, y_{2}\right)\right|<\varepsilon \quad \text { (for all }\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right) \text { with }\left|x_{1}-x_{2}\right|<\delta \text { and }\left|y_{1}-y_{2}\right|<\delta\right)
$$

In particular, this holds for all $y_{1}=y_{2}$, and $x_{1}=x$, and $x_{2}=x_{o}$.

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[14.2.2] Corollary: For a $\mathbb{C}$-valued $C^{k}$ function $F$ on $[a, b] \times[c, d]$,

$$
\frac{\partial}{\partial y} \int_{a}^{b} F(x, y) d x=\int_{a}^{b} F \frac{\partial}{\partial y}(x, y) d x
$$

Proof: The function-valued function $x \rightarrow(y \rightarrow F(x, y))$ is a continuous, $C^{k}[c, d]$-valued function, and $\frac{\partial}{\partial y}$ is a continuous linear map $C^{k}[c, d] \rightarrow C^{k-1}[c, d]$, so the Gelfand-Pettis property allows interchange of the operator and the integral.

### 14.3 Fourier transforms

Certainly an integral expressing Fourier inversion [13.15]

$$
f(x)=\frac{1}{2 \pi} \int_{\mathbb{R}} e^{i \xi x} \widehat{f}(\xi) d \xi
$$

for Schwartz function $f$ cannot converge as a Schwartz-function-valued integral, because $x \rightarrow e^{i \xi x}$ is in $C^{\infty}(\mathbb{R})$, but not Schwartz. Multiplying by $\widehat{f}$ does not affect decay in $x$, so does not alter the situation. Examination of the situation is complicated by the fact that the integrand is not compactly supported, but we can follow Schwartz' device of suitably compactifying $\mathbb{R}^{n}$ to a sphere $S^{n}$, and then invoke the GelfandPettis property for compactly-supported functions. Then we will see that the integral does converge as a $C^{\infty}(\mathbb{R})$-valued Gelfand-Pettis integral. First,
[14.3.1] Claim: For any $\Phi \in C^{\infty}\left(\mathbb{R}^{2}\right)$, the $C^{\infty}(\mathbb{R})$-valued function $\xi \rightarrow \Phi(-, \xi)$ that is, $\xi \rightarrow(x \rightarrow \Phi(x, \xi))$ is a continuous, $C^{\infty}(\mathbb{R})$-valued function on $\mathbb{R}$. (Similarly, it is a smooth $C^{\infty}(\mathbb{R})$-valued function, but we do not need this.)
Proof: The function $(x, \xi) \rightarrow \Phi(x, \xi)$ is $C^{\infty}$ as a function of two variables. In particular, $(x, \xi) \rightarrow \frac{\partial^{k}}{\partial x^{k}} \Phi(x, \xi)$ is continuous as a function of two variables. For each $k$, compact $C \subset \mathbb{R}, \varepsilon>0$ and each $\xi_{o} \in \mathbb{R}$, we want $\delta>0$ such that

$$
\left|\xi-\xi_{o}\right|<\delta \Longrightarrow \sup _{x \in C}\left|\frac{\partial^{k}}{\partial x^{k}} \Phi(x, \xi)-\frac{\partial^{k}}{\partial x^{k}} \Phi\left(x, \xi_{o}\right)\right|<\varepsilon
$$

Let $I$ be the interval $\left[\xi_{1}, \xi_{o}+1\right]$. The continuous function $(x, \xi) \rightarrow \frac{\partial^{k}}{\partial x^{k}} \Phi(x, \xi)$ is uniformly continuous on the compact $C \times I$ that is, there is $\delta>0$ such that

$$
\left.\left|\frac{\partial^{k}}{\partial x^{k}} \Phi\left(x_{1}, \xi_{1}\right)-\frac{\partial^{k}}{\partial x^{k}} \Phi\left(x_{2}, \xi_{2}\right)\right|<\varepsilon \quad \text { (for all }\left(x_{1}, \xi_{1}\right),\left(x_{2}, \xi_{2}\right) \in I \text { with }\left|x_{1}-x_{2}\right|<\delta \text { and }\left|\xi_{1}-\xi_{2}\right|<\delta\right)
$$

In particular, this holds for all $x_{1}=x_{2}$, and $\xi=\xi_{o} \in I$, and $\xi_{2}=\xi_{o}$, giving the desired continuity. This previous applies to $\Phi(x, \xi)=e^{i \xi x}$. Since $\xi \rightarrow F(\xi)$ is a continuous $\mathbb{C}$-valued function, the product $x \rightarrow F(\xi) \cdot e^{i \xi x}$ is a continuous $C^{\infty}(\mathbb{R})$-valued function of $\xi \in \mathbb{R}$.

Compactify $\mathbb{R}$ to the circle $\mathbb{T} \subset \mathbb{R}^{2}$ via by stereographic projection

$$
\sigma: x \longrightarrow\left(\frac{x}{\sqrt{1+x^{2}}}, \frac{1}{\sqrt{1+x^{2}}}\right)
$$

and adding the point $\infty=(0,1)$.
[14.3.2] Claim: $\xi \rightarrow F(\xi) \cdot \psi_{\xi}$ extends (by $0 \in C^{\infty}(\mathbb{R})$ ) to a continuous, $C^{\infty}(\mathbb{R})$-valued function on the compactification $\mathbb{T}$ of $\mathbb{R}$.
Proof: We must check continuity in $\xi$ near $\infty$. That is, for each $k$, compact $C \subset \mathbb{R}$, and $\varepsilon>0$, we want (large) $B$ such that

$$
|\xi|>B \Longrightarrow \sup _{x \in C}\left|F(\xi) \cdot \frac{\partial^{k}}{\partial x^{k}} e^{i \xi x}-0\right|<\varepsilon
$$

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The exponential function is easy to estimate: for example, with $M$ a bound so that $\left|\left(1+\xi^{2}\right)^{k} \cdot F(\xi)\right| \leq M$,

$$
\sup _{x \in C}\left|F(\xi) \cdot \frac{\partial^{k}}{\partial x^{k}} e^{i \xi x}\right|=\left|F(\xi) \cdot(i \xi)^{k}\right| \cdot 1 \leq \frac{M \cdot|\xi|^{k}}{\left(1+\xi^{2}\right)^{k}}
$$

Take $B$ large enough so that $M \cdot B^{k} /\left(1+B^{2}\right)^{k}<\varepsilon$. For any continuous linear functional, $\xi \rightarrow \lambda \circ\left(\psi_{\xi} \cdot F(\xi)\right)$ is a continuous scalar-valued function on the compact set $\mathbb{T}$, so is bounded. The same is true of any $\xi \rightarrow \lambda \circ\left(\psi_{\xi} \cdot\left(1+\xi^{2}\right)^{N} F(\xi)\right)$, so $\xi \rightarrow \lambda \circ\left(\psi_{\xi} \cdot F(\xi)\right)$ is rapid decreasing. Adjust the measure on $\mathbb{R}$ to give total measure 1:

$$
\int_{\mathbb{R}} \psi_{\xi} \cdot F(\xi) d \xi=\int_{\mathbb{R}} \psi_{\xi} \cdot \pi(1+\xi)^{2} F(\xi) \frac{d \xi}{\pi\left(1+\xi^{2}\right)}
$$

The function $\pi(1+\xi)^{2} F(\xi)$ is still continuous and of rapid decay. Being continuous and compactly supported on a measure space with total measure 1 , with values in a quasi-complete, locally convex topological vector space, $\xi \rightarrow \pi\left(1+\xi^{2}\right) F(\xi) \cdot \psi_{\xi}$ has a Gelfand-Pettis integral $J$ with respect to the measure $d \xi / \pi\left(1+\xi^{2}\right)$, lying inside the closed convex hull of the image. That is,

$$
\lambda(J)=\int_{\mathbb{R}} \lambda\left(\psi_{\xi}\right) \cdot \pi(1+\xi)^{2} F(\xi) \frac{d \xi}{\pi\left(1+\xi^{2}\right)}
$$

for every continuous linear functional $\lambda$. In the the latter scalar-valued integral the adjustment factors cancel:

$$
\int_{\mathbb{R}} \lambda\left(\psi_{\xi} \cdot \pi(1+\xi)^{2} F(\xi)\right) \frac{d \xi}{\pi\left(1+\xi^{2}\right)}=\int_{\mathbb{R}} \lambda\left(\psi_{\xi}\right) \cdot \pi(1+\xi)^{2} F(\xi) \frac{d \xi}{\pi\left(1+\xi^{2}\right)}=\int_{\mathbb{R}} \lambda\left(\psi_{\xi}\right) \cdot F(\xi) d \xi
$$

That is, $\lambda(J)=\int_{\mathbb{R}} \lambda\left(\psi_{\xi}\right) \cdot F(\xi) d \xi$, and the Gelfand-Pettis integral $J$ of the mutually adjusted function and measure is the Gelfand-Pettis integral of the original.
[14.3.3] Corollary: For rapidly decreasing $F \in C^{o}(\mathbb{R})$, for any continuous linear $T: C^{\infty} \rightarrow V$ for another topological vector space $V$,

$$
T\left(\int_{\mathbb{R}} \psi_{\xi} \cdot F(\xi) d \xi\right)=\int_{\mathbb{R}} T\left(\psi_{\xi} \cdot F(\xi)\right) d \xi=\int_{\mathbb{R}} T\left(\psi_{\xi}\right) \cdot F(\xi) d \xi
$$

as $V$-valued Gelfand-Pettis integral.
[14.3.4] Corollary: For rapidly decreasing $F \in C^{o}(\mathbb{R})$, for any continuous, for any compactly-supported distribution $u$,

$$
u\left(\int_{\mathbb{R}} \psi_{\xi} \cdot F(\xi) d \xi\right)=\int_{\mathbb{R}} u\left(\psi_{\xi} \cdot F(\xi)\right) d \xi=\int_{\mathbb{R}} u\left(\psi_{\xi}\right) \cdot F(\xi) d \xi
$$

with absolutely convergent integral.
[14.3.5] Corollary: For rapidly decreasing $F \in C^{o}(\mathbb{R})$, the Fourier transform is a $C^{\infty}$ function on $\mathbb{R}$, and its derivative is computed by the expected expression

$$
\frac{\partial}{\partial x}\left(\int_{\mathbb{R}} \psi_{\xi} \cdot F(\xi) d \xi\right)=\int_{\mathbb{R}} \frac{\partial \psi_{\xi}}{\partial x} \cdot F(\xi) d \xi=i \int_{\mathbb{R}} \psi_{\xi} \cdot \xi F(\xi) d \xi
$$

since $\partial / \partial x$ is a continuous map of $C^{\infty}(\mathbb{R})$ to itself.

### 14.4 Uniqueness of invariant distributions

We prove uniqueness of invariant functionals on suitable function spaces $V$ on topological spaces $X$ on which a topological group acts transitively. This includes uniqueness of invariant (Haar) measures, and uniqueness of invariant distributions, as special cases.

A translation-invariant function $f$ on the real line, that is, a function with $f(x+y)=f(x)$ for all $x, y \in \mathbb{R}$, is constant, by a point-wise argument:

$$
f(x)=\left(T_{x} f\right)(0)=f(0)
$$

where $T_{x} f(y)=f(x+y)$ is translation. The same conclusion holds for translation-invariant distributions, but we cannot argue in terms of point-wise values.

Let $G$ be a topological group, ${ }^{[81]}$ with right translation-invariant measure $d g$, meaning that

$$
\int_{G} f(g \cdot h) d g=\int_{G} f(g) d g \quad(\text { for all } h \in G)
$$

We assume only existence of a right translation-invariant measure. The theorem proves uniqueness.
For present purposes, a continuous approximate identity on a topological group $G$ is a sequence $\left\{\varphi_{i}\right\}$ of non-negative, continuous, real-valued functions such that $\int_{G} \varphi_{i}=1$ for all $i$, and such that the supports shrink to $\{1\}$, in the sense that for every neighborhood $N$ of 1 in $G$, there is an index $i_{o}$ so that the support of $\varphi_{i}$ is inside $N$ for all $i \geq i_{o}$. From Urysohn's Lemma [9.E.2], there always exists a continuous approximate identity. Not all classes of functions contain an approximate identity in this strict sense: (real-) analytic functions on a non-compact group, such as $\mathbb{R}$ cannot be compactly supported, so a compromise notion would be needed. The following theorem refers to the strict sense that supports shrink to $\{1\}$ :
[14.4.1] Theorem: Let $V \subset C_{c}^{o}(G)$ be a quasi-complete, locally convex topological vector space of complexvalued functions on $G$ stable under left and right translations, so that $G \times V \rightarrow V$ is continuous, and containing an approximate identity $\left\{\varphi_{i}\right\}$. Then there is a unique right $G$-invariant element of the dual space $V^{*}$ (up to constant multiples), and it is

$$
f \rightarrow \int_{G} f(g) d g \quad \text { (with right translation-invariant measure) }
$$

Proof: Let $T_{g} f(y)=f(y g)$ be right translation. With right translation-invariant measure $d g$ on $G$, since $\int_{G} \varphi_{i}(g) d g=1$ and $\varphi_{i}$ is non-negative, $\varphi_{i}(g) d g$ is a probability measure (total mass 1 ) on the (compact) support of $\varphi_{i}$. Thus, for any $f \in V$, we have a $V$-valued Gelfand-Pettis integral

$$
T_{\varphi_{i}} f=\int_{G} \varphi_{i}(g) T_{g} f d g \in \text { closure of convex hull of }\left\{\varphi_{i}(g) f: g \in G\right\} \subset V
$$

By continuity, given a neighborhood $N$ of 0 in $V$, we have $T_{\varphi_{i}} f \in f+N$ for all sufficiently large $i$. That is, $T_{\varphi_{i}} f \rightarrow f$. For a right-invariant (continuous) functional $u \in V^{*}$,

$$
u(f)=\lim _{i} u\left(g \rightarrow \int_{G} \varphi_{i}(h) f(g h) d h\right)
$$

This is

$$
u\left(g \rightarrow \int_{G} f(h g) \varphi_{i}\left(h^{-1}\right) d h\right)=u\left(g \rightarrow \int_{G} f(h) \varphi_{i}\left(g h^{-1}\right) d h\right)
$$

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by replacing $h$ by $h g^{-1}$. By properties of Gelfand-Pettis integrals, and since $f$ is guaranteed to be a compactly-supported continuous function, we can move the functional $u$ inside the integral: the above becomes

$$
\int_{G} f(h) u\left(g \rightarrow \varphi_{i}\left(g h^{-1}\right)\right) d h
$$

Using the right $G$-invariance of $u$ the evaluation of $u$ with right translation by $h^{-1}$ gives

$$
\int_{G} f(h) u\left(g \rightarrow \varphi_{i}(g)\right) d h=u\left(\varphi_{i}\right) \cdot \int_{G} f(h) d h
$$

By assumption the latter expressions approach $u(f)$ as $i \rightarrow \infty$. For $f$ so that the latter integral is non-zero, we see that the limit of the $u\left(\varphi_{i}\right)$ exists, and then we conclude that $u(f)$ is a constant multiple of the indicated integral with right Haar measure.

### 14.5 Smoothing of distributions

Every locally integrable ${ }^{[82]}$ function $f$ on $\mathbb{R}^{n}$, for example, gives a distribution $u_{f}$ by integrating against it:

$$
u_{f}(\varphi)=\int_{\mathbb{R}^{n}} \varphi \cdot f \quad\left(\text { for } \varphi \in \mathcal{D}\left(\mathbb{R}^{n}\right)\right)
$$

Conversely, we prove here that the distributions $u_{f}$ from $f \in \mathcal{D}=C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ are dense in the whole space $\mathcal{D}^{*}$ of distributions, with the weak dual topology. Further, a sequence of such smooth functions approaching a given distribution can be expressed in terms of smoothing or mollifying $u$.

Let $g \rightarrow T_{g}$ be the regular representation of $\mathbb{R}^{n}$ on test functions $f \in \mathcal{D}=\mathcal{D}\left(\mathbb{R}^{n}\right)$ by $\left(T_{g} f\right)(x)=f(x+g)$, for $x, g \in \mathbb{R}^{n}$. The map, $x \times f \rightarrow T_{x} f$ gives a continuous map $\mathbb{R}^{n} \times \mathcal{D} \longrightarrow \mathcal{D}$. The corresponding adjoint action of $\mathbb{R}^{n}$ on distributions $u$ is

$$
\left(T_{g}^{*} u\right)(f)=u\left(T_{g}^{-1} f\right)
$$

For the usual reasons, this gives a continuous map $x \times u \longrightarrow x \cdot u=T_{x}^{*} u$ with the weak dual topology: for $f \in \mathcal{D}$, let $\nu_{f}$ be the semi-norm $\nu_{f}(u)=|u(f)|$ on $\mathcal{D}^{*}$, and then

$$
\begin{gathered}
\nu\left(T_{g}^{*} u-T_{h}^{*} v\right)=\left|u\left(T_{g}^{-1} f\right)-v\left(T_{h}^{-1} f\right)\right| \leq\left|u\left(T_{g}^{-1} f\right)-u\left(T_{h}^{-1} f\right)\right|+\left|u\left(T_{h}^{-1} f\right)-v\left(T_{h}^{-1} f\right)\right| \\
=\left|u\left(T_{g}^{-1} f-T_{h}^{-1} f\right)\right|+\nu_{T_{h}^{-1} f}(u-v)
\end{gathered}
$$

For $g$ close to $h$, since the translation action of $\mathbb{R}^{n}$ on $\mathcal{D}$ is continuous and $u$ is a continuous functional, $\left|u\left(T_{g}^{-1} f\right)-u\left(T_{h}^{-1} f\right)\right|$ is small. And for $u$ close to $v$ in the weak dual topology, the second term is small. This proves the continuity.

As earlier and throughout, the action of a function $\varphi \in C_{c}^{o}\left(\mathbb{R}^{n}\right)$ on distributions $u$ is by integrating the group action

$$
T_{\varphi}^{*} u=\int_{\mathbb{R}^{n}} \varphi(x) T_{x}^{*} u d x \in \mathcal{D}^{*}
$$

Suppressing the $T^{*}$, this is

$$
\varphi \cdot u=\int_{\mathbb{R}^{n}} \varphi(x) x \cdot u d x \in \mathcal{D}^{*}
$$

A smooth approximate identity on $\mathbb{R}^{n}$ is a sequence $\left\{\psi_{i}\right\} \subset \mathcal{D}$ which are non-negative, real-valued, have $\int_{\mathbb{R}^{n}} \psi_{i}=1$, and supports shrink to $\{0\} \subset \mathbb{R}^{n}$.
[14.5.1] Theorem: For a smooth approximate identity $\left\{\psi_{i}\right\}$ and distribution $u$, the distributions $T_{\psi_{i}}^{*} u$ go to $u$ in the weak dual topology on $\mathcal{D}^{*}$, and are (integration against) the functions $x \rightarrow u\left(T_{x}^{-1} \psi_{i}\right)$, which are smooth functions.
[82] Again, locally integrable means that $|f|$ is in $L^{1}(K)$ for every compact $K$. This makes best sense for positive regular Borel measures, so that the measures of compact sets are finite.

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Proof: That $T_{\psi_{i}}^{*} u \rightarrow u$ as distributions is an instance of a general property of such Gelfand-Pettis integrals from [14.1.4]. To prove that every $T_{f} u$ for $f \in \mathcal{D}$ is (integration against) a continuous or smooth function, we first guess what that continuous function is, by determining its point-wise values. Indeed, if $u=u_{\varphi}$ were known to be integration against a continuous function $\varphi$, then with an approximate identity $\left\{\psi_{i}\right\}$

$$
\lim _{i} u_{\varphi}\left(\psi_{i}\right)=\lim _{i} \int_{\mathbb{R}^{n}} \varphi(x) \psi_{i}(x) d x=\varphi(0)
$$

Thus, we anticipate determining values of the alleged continuous function $f \cdot u$ by computing

$$
\text { alleged value }(f \cdot u)(0)=\lim _{i}(f \cdot u)\left(\psi_{i}\right)
$$

For a continuous function $F$ on $\mathbb{R}^{n}$, let $F^{\vee}(x)=F(-x)$. For for $f$ and $\psi$ in $\mathcal{D}$, since Gelfand-Pettis integrals commute with continuous linear maps,

$$
\begin{gathered}
\left(T_{f}^{*} u\right)(\psi)=\left(\int_{\mathbb{R}^{n}} f(x) T_{x}^{*} u d x\right)(\psi)=\int_{\mathbb{R}^{n}} f(x)\left(T_{x}^{*} u\right)(\psi) d x \\
=\int_{\mathbb{R}^{n}} f(x) u\left(T_{x}^{-1} \cdot \psi\right) d x=u\left(\int_{\mathbb{R}^{n}} f(x)\left(T_{x}^{-1} \cdot \psi\right) d x\right)=u\left(\int_{\mathbb{R}^{n}} f(-x)\left(T_{x} \psi\right) d x\right)=u\left(T_{f^{\vee}} \psi\right)
\end{gathered}
$$

The function $T_{f} \vee \psi$ admits a rewriting that reverses the roles of $f$ and $\psi$, namely

$$
\begin{gathered}
\left(T_{f^{\vee}} \psi\right)(y)=\int_{\mathbb{R}^{n}} f(-x) \psi(y+x) d x=\int_{\mathbb{R}^{n}} f(y-x) \psi(x) d x \\
=\int_{\mathbb{R}^{n}} f(y+x) \psi(-x) d x=\int_{\mathbb{R}^{n}} f(y+x) \psi^{\vee}(x) d x=\left(T_{\psi^{\vee}} f\right)(y)
\end{gathered}
$$

Thus,

$$
\left(T_{f}^{*} \cdot u\right)(\psi)=u\left(T_{f} \vee \psi\right)=u\left(T_{\psi} \vee f\right)=\left(T_{\psi}^{*} u\right)(f)
$$

We already know that $T_{\psi_{i}}^{*} u \rightarrow u$ for an approximate identity $\psi_{i}$, so the limit exists, and has an understandable value:

$$
\left(T_{f}^{*} u\right)\left(\psi_{i}\right)=\left(T_{\psi_{i}}^{*} u\right)(f) \longrightarrow u(f)=\text { supposed value of } f \cdot u \text { at } 0
$$

Thus, we would guess that $T_{f}^{*} u$ should be a function with value $u(f)$ at 0 . More generally, for the distribution $u_{\varphi}$ given by integration against $\varphi$, we have

$$
\left(T_{z}^{*} u_{\varphi}\right)\left(\psi_{i}\right)=u_{\varphi}\left(T_{z}^{-1} \psi_{i}\right)=\int_{\mathbb{R}^{n}} \varphi(x) \psi_{i}(x-z) d x=\int_{\mathbb{R}^{n}} \varphi(x+z) \psi_{i}(x) d x \rightarrow \varphi(z)
$$

The analogous computation suggests the values of the function $T_{f}^{*} u$ at $z$. First, a more elaborate version of the identity reverses the roles of test functions $f$ and $\varphi$, namely

$$
\begin{aligned}
& \left(T_{f^{\vee}} T_{z}^{-1} \psi\right)(y)=\int_{\mathbb{R}^{n}} f(-x) \psi(y+x-z) d x=\int_{\mathbb{R}^{n}} f(y-x-z) \psi(x) d x \\
= & \int_{\mathbb{R}^{n}} f(y+x-z) \psi(-x) d x=\int_{\mathbb{R}^{n}}\left(T_{z}^{-1} f\right)(y+x) \psi^{\vee}(x) d x=\left(T_{\psi} \vee T_{z}^{-1} f\right)(y)
\end{aligned}
$$

The same sort of computation gives

$$
\begin{aligned}
& \left.\left(T_{y}^{*}\left(T_{f}^{*} u\right)\right)\left(\psi_{i}\right)=\left(T_{f}^{*} u\right)\right)\left(T_{y}^{-1} \psi_{i}\right)=u\left(T_{f} \vee T_{y}^{-1} \psi_{i}\right)=u\left(T_{\psi_{i}^{\vee}} T_{y}^{-1} f\right) \\
= & \left(T_{y}^{*}\left(T_{\psi_{i}}^{*} u\right)\right)(f) \rightarrow\left(T_{y}^{*} u\right)(f)=u\left(T_{y}^{-1} f\right)=\text { supposed value of } f \cdot u \text { at } y
\end{aligned}
$$

Since $\mathbb{R}^{n} \times \mathcal{D} \rightarrow \mathcal{D}$ is continuous, and $u$ is continuous, the composition

$$
y \times f \longrightarrow T_{y}^{-1} f \longrightarrow u\left(T_{y}^{-1} f\right)
$$

is indeed continuous as a function of $y \in \mathbb{R}^{n}$.
Now we check that the distribution $f \cdot u$ is truly given by integration against the continuous function

$$
\varphi(y)=u\left(T_{y}^{-1} f\right)
$$

that apparently gives the pointwise values of $T_{f}^{*} u$. Letting $h \in \mathcal{D}$,

$$
\int_{\mathbb{R}^{n}} \varphi(x) h(x) d x=\int_{\mathbb{R}^{n}} u\left(T_{x}^{-1} f\right) h(x) d x=\left(\int_{\mathbb{R}^{n}} h(x) x \cdot u d x\right)(f)=\left(T_{h}^{*} u\right)(f)
$$

We already computed directly that

$$
\left(T_{h}^{*} u\right)(f)=u\left(T_{h} \vee\right)=u\left(T_{f} \vee h\right)=\left(T_{f}^{*} u\right)(h)
$$

which shows that integration against the continuous function $\varphi(y)=u\left(T_{y}^{-1} f\right)$ gives the distribution $T_{f}^{*} u$.
Smoothness of $y \rightarrow u\left(T_{y}^{-1} f\right)$ would follow from the assertion that $y \rightarrow T_{y}^{-1} f$ is a smooth, $\mathcal{D}$-valued function, since $u$ is a continuous linear functional on $\mathcal{D}$. The latter assertion is existence of limits

$$
\lim _{t \rightarrow 0} \frac{T_{y+t X}^{-1} f-T_{y}^{-1} f}{t} \quad\left(\text { for } X \in \mathbb{R}^{n} \text { and } y \in \mathbb{R}^{n}\right)
$$

in $\mathcal{D}$ for each $X \in \mathbb{R}^{n}$, and iterates thereof. It suffices to consider $y=0$. By design, differentiation (such as this directional derivative in the $X$ direction) is a continuous map of $\mathcal{D}$ to itself [13.9]. This gives the smoothness of $y \rightarrow u\left(T_{y}^{-1} f\right)$.
[14.5.2] Remark: That is, given the idea that $f \cdot u$ has been smoothed, determination of it as a classical function is straightforward. The proof that $T_{\psi_{i}}^{*} u \rightarrow u$ did not use the specifics of the situation: the same argument applies to representations of Lie groups.

### 14.6 Density of smooth vectors

Let $G$ be a Lie group, so that the notion of $C^{\infty}$ function on $G$ makes sense. A representation of $G$ on a locally convex, quasi-complete topological vectorspace $V$ is a continuous map $G \times V \rightarrow V$ that is linear in $V$, and has the associativity $(g h) \cdot v=g \cdot(h \cdot v)$ for $g, h \in G$. The subspace $V^{\infty}$ of smooth vectors is

$$
V^{\infty}=\left\{v \in V: g \rightarrow g \cdot v \text { is a } C^{\infty} V \text {-valued function on } G\right\}
$$

It suffices to consider derivatives associated to the Lie algebra $\mathfrak{g}$ of $G$ :

$$
(x \cdot f)(g)=\left.\frac{\partial}{\partial t}\right|_{t=0}\left(\left(g e^{t x}\right) \cdot v\right) \quad(\text { for } x \in \mathfrak{g})
$$

where $x \rightarrow e^{x}$ is the exponential map $\mathfrak{g} \rightarrow G$.
Note that in the representation of $\mathbb{R}^{n}$ on distributions $\mathcal{D}^{*}$ every distribution is a smooth vector, since every distribution is infinitely differentiable as a distribution. Thus, smooth vectors are not necessarily smooth functions. Nevertheless, as in the previous section, distributions are approximable by smooth functions. For general representations $G \times V \rightarrow V$, the following is the appropriate corollary of [14.1.4]:
[14.6.1] Theorem: (Gårding) For quasi-complete, locally convex $V$ with a continuous action of a real Lie group $G, V^{\infty}$ is dense in $V$.
Proof: Let $\left\{\psi_{i}\right\}$ be an approximate identity in $\mathcal{D}(G)$. On one hand, by [14.1.4], for each $v \in V, T_{\psi_{i}} v \rightarrow v$. On the other hand, we claim that $T_{\psi_{i}} v$ is a smooth vector in $V$. That is, for any $\psi \in \mathcal{D}$, we claim that $g \rightarrow T_{g}\left(T_{\psi} v\right)$ is a smooth function of $g \in G$. By the weak-to-strong result [15.1.1] it suffices to show that, for all $\lambda \in V^{*}, g \rightarrow \lambda\left(T_{g} T_{\psi} v\right)$ is a smooth scalar-valued function. By properties of Gelfand-Pettis integrals,

$$
\lambda\left(T_{g} T_{\psi} v\right)=\lambda T_{g} \int_{G} \psi(h) T_{h} v d h=\int_{G} \psi(h) \lambda\left(T_{g} T_{h} v\right) d h=\int_{G} \psi\left(g^{-1} h\right) \lambda\left(T_{h} v\right) d h
$$

To show differentiability near a given $g_{o}$, without loss of generality we can multiply by a smooth, compactlysupported cut-off function $\eta(g)$ which is identically 1 near $g_{o}$. Then $h \rightarrow\left(g \rightarrow \eta(g) \psi\left(g^{-1} h\right)\right)$ is a smooth, compactly-supported function on $G \times G$, and the integrand $h \rightarrow\left(g \rightarrow \eta(g) \psi\left(g^{-1} h\right) \lambda\left(T_{h} v\right)\right)$ is a continuous, compactly-supported, $\mathcal{D}(G)$-valued function on $G$. Thus, it admits a Gelfand-Pettis integral $g \rightarrow \int_{G} \psi\left(g^{-1} h\right) \lambda\left(T_{h} v\right) d h$ that is a smooth function. This holds for every $\lambda \in V^{*}$, so by [15.1.1] shows that $g \rightarrow T_{g} T_{\psi} v$ is a smooth $V$-valued function on $G$.

### 14.7 Quasi-completeness and convex hulls of compacts

A subset $E$ of a complete metric space $X$ is totally bounded if, for every $\varepsilon>0$ there is a covering of $E$ by finitely-many open balls of radius $\varepsilon$. The property of total boundedness in a metric space is generally stronger than mere boundedness. It is immediate that any subset of a totally bounded set is totally bounded. Recall:
[14.7.1] Proposition: A subset of a complete metric space has compact closure if and only if it is totally bounded.
Proof: Certainly if a set has compact closure then it admits a finite covering by open balls of arbitrarily small (positive) radius. On the other hand, suppose that a set $E$ is totally bounded in a complete metric space $X$. To show that $E$ has compact closure it suffices to show that any sequence $\left\{x_{i}\right\}$ in $E$ has a Cauchy subsequence.

Choose such a subsequence as follows. Cover $E$ by finitely-many open balls of radius 1 . In at least one of these balls there are infinitely-many elements from the sequence. Pick such a ball $B_{1}$, and let $i_{1}$ be the smallest index so that $x_{i_{1}}$ lies in this ball.

The set $E \cap B_{1}$ is still totally bounded, and contains infinitely-many elements from the sequence. Cover it by finitely-many open balls of radius $1 / 2$, and choose a ball $B_{2}$ with infinitely-many elements of the sequence lying in $E \cap B_{1} \cap B_{2}$. Choose the index $i_{2}$ to be the smallest one so that both $i_{2}>i_{1}$ and so that $x_{i_{2}}$ lies inside $E \cap B_{1} \cap B_{2}$.

Inductively, suppose that indices $i_{1}<\ldots<i_{n}$ have been chosen, and balls $B_{i}$ of radius $1 / i$, so that

$$
x_{i} \in E \cap B_{1} \cap B_{2} \cap \ldots \cap B_{i}
$$

Cover $E \cap B_{1} \cap \ldots \cap B_{n}$ by finitely-many balls of radius $1 /(n+1)$ and choose one, call it $B_{n+1}$, containing infinitely-many elements of the sequence. Let $i_{n+1}$ be the first index so that $i_{n+1}>i_{n}$ and so that

$$
x_{n+1} \in E \cap B_{1} \cap \ldots \cap B_{n+1}
$$

For $m<n$ we have $d\left(x_{i_{m}}, x_{i_{n}}\right) \leq \frac{1}{m}$ so this subsequence is Cauchy.
In a not-necessarily metric topological vectorspace $V$, a subset $E$ is totally bounded if, for every neighborhood $U$ of 0 there is a finite subset $F$ of $V$ so that $E \subset F+U$, where

$$
F+U=\bigcup_{v \in F} v+U=\{v+u: v \in F, u \in U\}
$$

[14.7.2] Proposition: A totally bounded subset $E$ of a locally convex topological vectorspace $V$ has totally bounded convex hull.
Proof: First, recall that the convex hull of a finite set $F=\left\{x_{1}, \ldots, x_{n}\right\}$ in a topological vectorspace is compact, since it is the continuous image of the compact set $\left\{\left(c_{1}, \ldots, c_{n}\right) \in \mathbb{R}^{n}: \sum_{i} c_{i}=1,0 \leq c_{i} \leq 1\right.$, for all $\left.i\right\} \subset \mathbb{R}^{n}$ under $\left(c_{1}, \ldots, c_{n}\right) \rightarrow \sum_{i} c_{i} x_{i}$.

Given a neighborhood $U$ of 0 in $V$, let $U_{1}$ be a convex neighborhood of 0 so that $U_{1}+U_{1} \subset U$. For some finite subset $F$ we have $E \subset F+U_{1}$, by total boundedness. The convex hull $K$ of $F$ is compact. Then $E \subset K+U_{1}$, and the latter is convex. Therefore, the convex hull $H$ of $E$ lies inside $K+U_{1}$. Since $K$ is compact, it is totally bounded, so can be covered by a finite union $\Phi+U_{1}$ of translates of $U_{1}$. Thus, since $U_{1}+U_{1} \subset U, H \subset\left(\Phi+U_{1}\right)+U_{1} \subset \Phi+U$. Thus, $H$ lies inside this finite union of translates of $U$. This holds for any open $U$ containing 0 , so $H$ is totally bounded.

## 14. Vector-valued integrals

[14.7.3] Corollary: In a Fréchet space, the closure of the convex hull of a compact set is compact.
Proof: A compact set in a Fréchet space (or in any complete metric space) is totally bounded, as recalled above. By the previous, the convex hull of a totally bounded set in a Fréchet space is totally bounded. Thus, this convex hull has compact closure, since totally bounded sets in complete metric spaces have compact closure.

The general case reduces to the case of Fréchet spaces.
[14.7.4] Proposition: In a quasi-complete, locally convex topological vectorspace $X$, the closure of the convex hull of a compact set is compact.
Proof: Since $X$ is locally convex, its topology is given by a collection of seminorms $v$. For each seminorm $v$, let $X_{v}$ be the completion of the quotient $X /\{x \in X: v(x)=0\}$ with respect to the metric that $v$ induces on the latter quotient. Thus, $X_{v}$ is a Banach space. Consider $Z=\prod_{v} X_{v}$ with product topology, with the natural injection $j: X \rightarrow Z$, and with projection $p_{v}$ to the $v^{t h}$ factor. By construction, and by definition of the topology given by the seminorms, $j$ is a (linear) homeomorphism to its image. That is, $X$ is homeomorphic to the subset $j X$ of $Z$, given the subspace topology.

Let $K \subset X$ be compact, with convex hull $H$, and $C$ the closure of $H$. The continuous image $p_{v} j K$ of compact $K$ is compact. Since $X_{v}$ is Fréchet, the convex hull $H_{v}$ of $p_{v} j K$ has compact closure $C_{v}$. The convex hull $j H$ of $j K$ is contained in the product $\prod_{v} H_{v}$ of the convex hulls $H_{v}$ of the projections $p_{v} j K$. By Tychonoff's theorem, the product $\prod_{v} C_{v}$ is compact.

Since $j C$ is contained in the compact set $\prod_{v} C_{v}$, to prove that the closure $j C$ of $j H$ in $j X$ is compact, it suffices to prove that $j C$ is closed in $Z$. Since $j C$ is a subset of the compact set $\prod_{v} C_{v}$, it is totally bounded and so is certainly bounded (in $Z$, hence in $X \approx j X$ ). By the quasi-completeness, a Cauchy net in $j C$ is necessarily bounded and converges to a point in $j C$. Since any point in the closure of $j C$ in $Z$ has a Cauchy net in $j C$ converging to it, $j C$ is closed in $Z$.

### 14.8 Existence proof

To simplify, divide by a constant to make $X$ have total measure 1 . The closure $H$ of the convex hull of $f(X)$ in $V$ is compact by hypothesis. We will show that there is an integral of $f$ inside $H$.

For finite $L \subset V^{*}$, let

$$
V_{L}=\left\{v \in V: \lambda v=\int_{X} \lambda \circ f, \forall \lambda \in L\right\} \quad \text { and } \quad I_{L}=H \cap V_{L}
$$

Since $H$ is compact and $V_{L}$ is closed, $I_{L}$ is compact. Certainly $I_{L} \cap I_{L^{\prime}}=I_{L \cup L^{\prime}}$ for two finite subsets $L, L^{\prime}$ of $V^{*}$. If all the $I_{L}$ are non-empty, then the intersection of all these compact sets $I_{L}$ is non-empty, by the finite intersection property, giving existence.

To prove that each $I_{L}$ is non-empty for finite subsets $L$ of $V^{*}$, choose an ordering $\lambda_{1}, \ldots, \lambda_{n}$ of the elements of $L$. Make a continuous linear mapping $\Lambda=\Lambda_{L}$ from $V$ to $\mathbb{R}^{n}$ by $\Lambda(v)=\left(\lambda_{1} v, \ldots, \lambda_{n} v\right)$. Since this map is continuous, the image $\Lambda(f(X))$ is compact in $\mathbb{R}^{n}$.

For a finite set $L$ of functionals, the integral $y=y_{L}=\int_{X} \Lambda f(x) d x$ is readily defined by component-wise integration. Take $y$ in the convex hull of $\Lambda(f(X))$. Since $\Lambda_{L}$ is linear, $y=\Lambda_{L} v$ for some $v$ in the convex hull of $f(X)$. Then

$$
\Lambda_{L} v=y=\left(\ldots, \int \lambda_{i} f(x) d x, \ldots\right)
$$

Thus, $v \in I_{L}$ as desired. It remains to show that $y$ lies in the convex hull of $\Lambda_{L}(f(x))$.
Suppose not. From the lemma below, in a finite-dimensional space the convex hull of a compact set is still compact, without taking closure. By the finite-dimensional case of the Hahn-Banach theorem, there would be a linear functional $\eta$ on $\mathbb{R}^{n}$ so that $\eta y>\eta z$ for all $z$ in this convex hull. That is, letting $y=\left(y_{1}, \ldots, y_{n}\right)$, there would be real $c_{1}, \ldots, c_{n}$ so that for all $\left(z_{1}, \ldots, z_{n}\right)$ in the convex hull $\sum_{i} c_{i} z_{i}<\sum c_{i} y_{i}$. In particular, for all $x \in X$

$$
\sum_{i} c_{i} \lambda_{i}(f(x))<\sum_{i} c_{i} y_{i}
$$

Integration of both sides of this over $X$ preserves ordering, giving the impossible $\sum_{i} c_{i} y_{i}<\sum_{i} c_{i} y_{i}$. Thus, $y$ does lie in this convex hull.
[14.8.1] Lemma: The convex hull of a compact set $K$ in $\mathbb{R}^{n}$ is compact.
Proof: First claim that, for $E \subset \mathbb{R}^{n}$ and for any $x$ a point in the convex hull of $E$, there are $n+1$ points $x_{0}, x_{1}, \ldots, x_{n}$ in $E$ of which $x$ is a convex combination.

By induction, it suffices to consider a convex combination $v=c_{1} v_{1}+\ldots+c_{N} v_{N}$ of vectors $v_{i}$ with $N>n+1$ and show that $v$ is actually a convex combination of $N-1$ of the $v_{i}$. Further, without loss of generality that all the coefficients $c_{i}$ are non-zero. Define a linear map

$$
L: \mathbb{R}^{N} \longrightarrow \mathbb{R}^{n} \times \mathbb{R} \quad \text { by } \quad L\left(x_{1}, \ldots, x_{N}\right) \longrightarrow\left(\sum_{i} x_{i} v_{i}, \sum_{i} x_{i}\right)
$$

By dimension-counting, since $N>n+1$ the kernel of $L$ is non-trivial. Let $\left(x_{1}, \ldots, x_{N}\right)$ be a non-zero vector in the kernel. Since $c_{i}>0$ for every index, and since there are only finitely-many indices altogether, there is a constant $c$ so that $\left|c x_{i}\right| \leq c_{i}$ for every index $i$, and so that $c x_{i_{o}}=c_{i_{o}}$ for at least one index $i_{o}$. Then

$$
v=v-0=\sum_{i} c_{i} v_{i}-c \cdot \sum x_{i} v_{i}=\sum_{i}\left(c_{i}-c x_{i}\right) v_{i}
$$

Since $\sum_{i} x_{i}=0$ this is still a convex combination, and since $c x_{i_{o}}=c_{i_{o}}$ at least one coefficient has become zero. This is the induction proving the claim.

By this claim, a point $v$ in the convex hull of $K$ is a convex combination $c_{o} v_{o}+\ldots+c_{n} v_{n}$ of $n+1$ points $v_{o}, \ldots, v_{n}$ of $K$. Let $\sigma$ be the compact set $\left(c_{o}, \ldots, c_{n}\right)$ with $0 \leq c_{i} \leq 1$ and $\sum_{i} c_{i}=1$. The convex hull of $K$ is the image of the compact set $\sigma \times K^{n+1}$ under the continuous map

$$
L:\left(c_{o}, \ldots, c_{n}\right) \times\left(v_{o}, v_{1}, \ldots, v_{n}\right) \longrightarrow \sum_{i} c_{i} v_{i}
$$

so is compact. This proves the lemma, finishing the proof of the theorem.

## 14.A Appendix: Hahn-Banach theorems

For a locally convex vectorspace $V$, functionals $\lambda \in V^{*}$ separate points, and convex sets can be separated by linear functionals. Continuous linear functionals on arbitrary subspaces have continuous extensions to the whole space. In contrast, in general, linear maps from subspaces $W$ to not-finite-dimensional topological vectorspaces need not extend to $V$. Indeed, if the identity map $T: W \rightarrow W$ extended to $T^{\prime}: V \rightarrow W$, then ker $T^{\prime}$ would be a complementary subspace, which need not exist even for closed subspaces $W$.

Let $k$ be either $\mathbb{R}$ or $\mathbb{C}$, and let $V$ be a $k$-vectorspace, without any assumptions about topologies for the moment. A $k$-linear $k$-valued function on $V$ is a linear functional. A linear functional $\lambda$ on $V$ is bounded when there is a neighborhood $U$ of 0 in $V$ and constant $c$ so that $|\lambda x| \leq c$ for $x \in U$, where $|\mid$ is the usual absolute value on $k$. The following proposition is the general analogue of the corresponding assertion for Banach spaces, in which boundedness has a different sense.
[14.A.1] Proposition: The following conditions on a linear functional $\lambda$ on a topological vectorspace $V$ over $k$ are equivalent: (i) $\lambda$ is continuous, (ii) $\lambda$ is continuous at 0 , (iii) $\lambda$ is bounded.
Proof: The first assertion certainly implies the second. Assume the second. Then, given $\varepsilon>0$, there is a neighborhood $U$ of 0 so that $|\lambda|$ is bounded by $\varepsilon$ on $U$. This proves boundedness. Finally, suppose that $|\lambda(x)| \leq c$ on a neighborhood $U$ of 0 . Then given $x \in V$ and given $\varepsilon>0$, we claim that for $y \in x+\frac{\varepsilon}{2 c} U$ we have $|\lambda(x)-\lambda(y)|<\varepsilon$. Indeed, letting $x-y=\frac{\varepsilon}{2 c} u$ with $u \in U$, we have

$$
|\lambda(x)-\lambda(y)|=\frac{\varepsilon}{2 c}|\lambda(u)| \leq \frac{\varepsilon}{2 c} \cdot c=\frac{\varepsilon}{2}<\varepsilon
$$

This proves the proposition.

## 14. Vector-valued integrals

The immediate goal is to extend a linear functional while preserving a comparison to another function (denoted $p$ below). For this, we need not suppose that the vectorspaces involved are topological vectorspaces. Let $V$ be a real vectorspace, without any assumption about topologies. Let $p: V \rightarrow \mathbb{R}$ be a non-negative real-valued function on $V$ so that

$$
\begin{gathered}
p(t v)=t \cdot p(v) \quad(\text { for } t \geq 0) \quad(\text { positive-homogeneity) } \\
p(v+w) \leq p(v)+p(w) \quad(\text { triangle inequality })
\end{gathered}
$$

Lacking a description of $p(t v)$ for $t<0, p$ is not quite a semi-norm.
[14.A.2] Theorem: Let $\lambda$ be a real-linear function on a real vector subspace $W$ of $V$, so that $\lambda(w) \leq p(w)$ for all $w \in W$. There is an extension of $\lambda$ to a real-linear function $\Lambda$ on all of $V$, so that $-p(-v) \leq \Lambda(v) \leq p(v)$ for all $v \in V$.
Proof: The key issue is extending the functional one step. That is, for $v_{o} \in V$, attempt to extend $\lambda^{\prime}$ of $\lambda$ to $W+\mathbb{R} v_{o}$ by $\lambda^{\prime}\left(w+t v_{o}\right)=\lambda(w)+c t$ and examine the resulting conditions on $c$.

For all $w, w^{\prime} \in W$

$$
\begin{gathered}
\lambda(w)-p\left(w-v_{o}\right)=\lambda\left(w+w^{\prime}\right)-\lambda\left(w^{\prime}\right)-p\left(w-v_{o}\right) \\
\leq p\left(w+w^{\prime}\right)-\lambda\left(w^{\prime}\right)-p\left(w-v_{o}\right)=p\left(w-v_{o}+w^{\prime}+v_{o}\right)-\lambda\left(w^{\prime}\right)-p\left(w-v_{o}\right) \\
\leq p\left(w-v_{o}\right)+p\left(w^{\prime}+v_{o}\right)-\lambda\left(w^{\prime}\right)-p\left(w-v_{o}\right)=p\left(w^{\prime}+v_{o}\right)-\lambda\left(w^{\prime}\right)
\end{gathered}
$$

That is,

$$
\lambda(w)-p\left(w-v_{o}\right) \leq p\left(w^{\prime}+v_{o}\right)-\lambda\left(w^{\prime}\right) \quad\left(\text { for all } w, w^{\prime} \in W\right)
$$

Let $\sigma$ be the sup of all the left-hand sides as $w$ ranges over $W$. Since the right-hand side is finite, this sup is finite. With $\mu$ the inf of the right-hand side as $w^{\prime}$ ranges over $W$,

$$
\lambda(w)-p\left(w-v_{o}\right) \leq \sigma \leq \mu \leq p\left(w^{\prime}+v_{o}\right)-\lambda\left(w^{\prime}\right)
$$

Take any real number $c$ so that $\sigma \leq c \leq \mu$ and define $\lambda^{\prime}\left(w+t v_{o}\right)=\lambda(w)+t c$.
To compare with $p$ is easy: in the inequality $\lambda(w)-p\left(w-v_{o}\right) \leq \sigma$ replace $w$ by $w / t$ with $t>0$, multiply by $t$ and invoke the positive-homogeneity to obtain $\lambda(w)-p\left(w-t v_{o}\right) \leq t \sigma$ from which

$$
\lambda^{\prime}\left(w-t v_{o}\right)=\lambda(w)-t c \leq \lambda(w)-t \sigma \leq p\left(w-t v_{o}\right)
$$

Likewise, from $\mu \leq p\left(w+v_{o}\right)-\lambda(w)$ a similar trick produces

$$
\lambda^{\prime}\left(w+t v_{o}\right)=\lambda(w)+t c \leq \lambda(w)+t \mu \leq p\left(w+t v_{o}\right)
$$

for $t>0$, the other half of the desired inequality. Thus, for all $v \in W+R v_{o}$ we have $\lambda^{\prime}(v) \leq p(v)$. Using the linearity of $\lambda^{\prime}, \lambda^{\prime}(v)=-\lambda^{\prime}(-v) \geq-p(-v)$ giving the bottom half of the comparison of $\lambda^{\prime}$ and $p$.

Extend to a functional on the whole space dominated by $p$ by transfinite induction, as follows. Let $\mathcal{X}$ be the collection of all pairs $(X, \mu)$, where $X$ is a subspace of $V$ (containing $W$ ), and where $\mu$ is real-linear real-valued function on $X$ so that $\mu$ restricted to $W$ is $\lambda$, and so that $-p(-x) \leq \mu(x) \leq p(x)$ for all $x \in X$. Order these by writing $(X, \mu) \leq(Y, \nu)$ when $X \subset Y$ and $\left.\nu\right|_{X}=\mu$. By the Hausdorff Maximality Principle, there is a maximal totally ordered subset $\mathcal{Y}$ of $\mathcal{X}$. Let

$$
V^{\prime}=\bigcup_{(X, \mu) \in \mathcal{Y}} X
$$

be the ascending union of all the subspaces in $\mathcal{Y}$. Define a linear functional $\lambda^{\prime}$ on this union as follows: for $v \in V^{\prime}$, take any $X$ so that $(X, \mu) \in \mathcal{Y}$ and $v \in X$ and define $\lambda^{\prime}(v)=\mu(v)$. The total ordering on $\mathcal{Y}$ makes the choice of $(X, \mu)$ not affect the definition of $\lambda^{\prime}$. If $V^{\prime}$ were not the whole space $V$ the first part of the proof would create an extension to a properly larger subspace, contradicting the maximality.
[14.A.3] Theorem: For a non-empty convex open subset $X$ of a locally convex topological vectorspace $V$, and a non-empty convex set $Y$ in $V$ with $X \cap Y=\phi$, there is a continuous real-linear real-valued functional $\lambda$ on $V$ and a constant $c$ so that $\lambda(x)<c \leq \lambda(y)$ for all $x \in X$ and $y \in Y$.
Proof: Fix $x_{o} \in X$ and $y_{o} \in Y$. Since $X$ is open, $X-x_{o}$ is open, and thus

$$
U=\left(X-x_{o}\right)-\left(Y-y_{o}\right)=\left\{\left(x-x_{o}\right)-\left(y-y_{o}\right): x \in X, \quad y \in Y\right\}
$$

is open. Further, since $x_{o} \in X$ and $y_{o} \in Y, U$ contains 0 . Since $X, Y$ are convex, $U$ is convex. The Minkowski functional $p=p_{U}$ attached to $U$ is $p(v)=\inf \{t>0: v \in t U\}$. The convexity assures that this function $p$ has the positive-homogeneity and triangle-inequality properties of the auxiliary functional $p$ above.

Let $z_{o}=-x_{o}+y_{o}$. Since $X \cap Y=\phi, z_{o} \notin U$, so $p\left(z_{o}\right) \geq 1$. Define a linear functional $\lambda$ on $\mathbb{R} z_{o}$ by $\lambda\left(t z_{o}\right)=t$. Check that $\lambda$ is dominated by $p$ in the sense of the previous section:

$$
\lambda\left(t z_{o}\right)=t \leq t \cdot p\left(z_{o}\right)=p\left(t z_{o}\right) \quad(\text { for } t \geq 0)
$$

while

$$
\lambda\left(t z_{o}\right)=t<0 \leq p\left(t z_{o}\right) \quad(\text { for } t<0)
$$

Thus, $\lambda\left(t z_{o}\right) \leq p\left(t z_{o}\right)$ for all real $t$, and $\lambda$ extends to a real-linear real-valued functional $\Lambda$ on $V$, still so that $-p(-v) \leq \Lambda(v) \leq p(v)$ for all $v \in V$. From the definition of $p,|\Lambda| \leq 1$ on $U$. Thus, on $\frac{\varepsilon}{2} U$ we have $|\Lambda|<\varepsilon$. That is, the linear functional $\Lambda$ is bounded, so is continuous at 0 , so is continuous on $V$.

For arbitrary $x \in X$ and $y \in Y$,

$$
\Lambda x-\Lambda y+1=\Lambda\left(x-y+z_{o}\right) \leq p\left(x-y+z_{o}\right)<1
$$

since $x-y+z_{o} \in U$. Thus, $\Lambda x-\Lambda y<0$ for all such $x, y$. Therefore, $\Lambda(X)$ and $\Lambda(Y)$ are disjoint convex subsets of $\mathbb{R}$. Since $\Lambda$ is not the zero functional, it is surjective to $\mathbb{R}$, and so is an open map. Thus, $\Lambda(X)$ is open, and $\Lambda(X)<\sup \Lambda(X) \leq \Lambda(Y)$ as desired.

The analogous results for complex scalars are corollaries of the real-scalar cases, as follows. Let $V$ be a complex vectorspace. Given a complex-linear complex-valued functional $\lambda$ on $V$, let its real part be

$$
u(v)=\operatorname{Re} \lambda(v)=\frac{\lambda(v)+\overline{\lambda(v)}}{2}
$$

where the overbar denotes complex conjugation. On the other hand, given a real-linear real-valued functional $u$ on $V$, its complexification is $C u(x)=u(x)-i u(i x)$ where $i=\sqrt{-1}$.
[14.A.4] Lemma: For a real-linear functional $u$ on the complex vectorspace $V$, the complexification $C u$ is a complex-linear functional so that $\operatorname{Re}(C u)=u$ and for a complex-linear functional $\lambda C(\operatorname{Re} \lambda)=\lambda$. (Straightforward computation).
[14.A.5] Corollary: Let $p$ be a seminorm on the complex vectorspace $V$. Let $\lambda$ be a complex-linear function on a complex vector subspace $W$ of $V$, so that $|\lambda(w)| \leq p(w)$ for all $w \in W$. Then there is an extension of $\lambda$ to a complex-linear function $\Lambda$ on all of $V$, so that $|\Lambda(v)| \leq p(v)$ for all $v \in V$.
Proof: Certainly if $|\lambda| \leq p$ then $|\operatorname{Re} \lambda| \leq p$. By the theorem for real-linear functionals, there is an extension $u$ of $\operatorname{Re} \lambda$ to a real-linear functional so that still $|u| \leq p$. Let $\Lambda=C u$. In light of the lemma, it remains to show that $|\Lambda| \leq p$. To this end, given $v \in V$, let $\mu$ be a complex number of absolute value 1 so that $|\Lambda(v)|=\mu \Lambda(v)$. Then

$$
|\Lambda(v)|=\mu \Lambda(v)=\Lambda(\mu v)=\operatorname{Re} \Lambda(\mu v) \leq p(\mu v)=p(v)
$$

using the seminorm property of $p$. Thus, the complex-linear functional made by complexifying the real-linear extension of the real part of $\lambda$ satisfies the desired bound.
[14.A.6] Corollary: Let $X$ be a non-empty convex open subset of a locally convex topological vectorspace $V$, and let $Y$ be an arbitrary non-empty convex set in $V$ so that $X \cap Y=\phi$. Then there is a continuous complex-linear complex-valued functional $\lambda$ on $V$ and a constant $c$ so that

$$
\operatorname{Re} \lambda(x)<c \leq \operatorname{Re} \lambda(y) \quad(\text { for all } x \in X \text { and } y \in Y)
$$

## 14. Vector-valued integrals

Proof: Invoke the real-linear version of the theorem to make a real-linear functional $u$ so that $u(x)<c \leq u(y)$ for all $x \in X$ and $y \in Y$. By the lemma, $u$ is the real part of its own complexification.
[14.A.7] Corollary: Let $V$ be a locally convex topological vectorspace. Let $K$ and $C$ be disjoint sets, where $K$ is a compact convex non-empty subset of $V$, and $C$ is a closed convex subset of $V$. Then there is a continuous linear functional $\lambda$ on $V$ and there are real constants $c_{1}<c_{2}$ so that

$$
\operatorname{Re} \lambda(x) \leq c_{1}<c_{2} \leq \operatorname{Re} \lambda(y) \quad(\text { for all } x \in K \text { and } y \in C)
$$

Proof: Take a small-enough convex neighborhood $U$ of 0 in $V$ so that $(K+U) \cap C=\phi$. Apply the separation theorem to $X=K+U$ and $Y=C$. The constant $c_{2}$ can be taken to be $c_{2}=\sup \operatorname{Re} \lambda(K+U)$. Since $\operatorname{Re} \lambda(K)$ is a compact subset of $\operatorname{Re} \lambda(K+U)$, its sup $c_{1}$ is strictly less than $c_{2}$.
[14.A.8] Corollary: Let $V$ be a locally convex topological vectorspace, $W$ a subspace, and $v_{o} \in V$. Let $\bar{W}$ be the topological closure of $W$. Then $v_{o} \notin \bar{W}$ if and only if there is a continuous linear functional $\lambda$ on $V$ so that $\lambda(W)=0$ while $\lambda(v)=1$.
Proof: On one hand, if $v_{o}$ lies in the closure of $W$, then any continuous function which is 0 on $W$ must be 0 on $v_{o}$, as well. On the other hand, suppose that $v_{o}$ does not lie in the closure of $W$. Then apply the previous corollary with $K=\left\{v_{o}\right\}$ and $C=\bar{W}$. We find that

$$
\operatorname{Re} \lambda\left(\left\{v_{o}\right\}\right) \cap \operatorname{Re} \lambda(\bar{W})=\phi
$$

Since $\operatorname{Re} \lambda(\bar{W})$ is a vector subspace of the real line, and is not the whole real line, it is just $\{0\}$, and $\operatorname{Re} \lambda\left(v_{o}\right) \neq 0$. Divide $\lambda$ by the constant $\operatorname{Re} \lambda\left(v_{o}\right)$ to obtain a continuous linear functional zero on $W$ but 1 on $v_{o}$.
[14.A.9] Corollary: Let $V$ be a locally convex topological (real) vectorspace. Let $\lambda$ be a continuous linear functional on a subspace $W$ of $V$. Then there is a continuous linear functional $\Lambda$ on $V$ extending $\lambda$.
Proof: Without loss of generality, take $\lambda \neq 0$. Let $W_{o}$ be the kernel of $\lambda$ (on $W$ ), and pick $w_{1} \in W$ so that $\lambda w_{1}=1$. Evidently $w_{1}$ is not in the closure of $W_{o}$, so there is $\Lambda$ on the whole space $V$ so that $\left.\Lambda\right|_{W_{o}}=0$ and $\Lambda w_{1}=1$. It is easy to check that this $\Lambda$ is an extension of $\lambda$.
[14.A.10] Corollary: Let $V$ be a locally convex topological vectorspace. Given two distinct vectors $x \neq y$ in $V$, there is a continuous linear functional $\lambda$ on $V$ so that $\lambda(x) \neq \lambda(y)$
Proof: The set $\{x\}$ is compact convex non-empty, and the set $\{y\}$ is closed convex non-empty, so we can apply a corollary just above.

## 15. Differentiable vector-valued functions

1. Weak-to-strong differentiability
2. Holomorphic vector-valued functions
3. Holomorphic $\operatorname{Hol}(\Omega, V)$-valued functions
4. Banach-Alaoglu: compactness of polars
5. Variant Banach-Steinhaus/uniform boundedness
6. Weak boundedness implies (strong) boundedness
7. Proof that weak $C^{1}$ implies strong $C^{o}$
8. Proof that weak holomorphy implies continuity

Appendix A: vector-valued power series
Appendix B: two forms of the Baire category theorem
Appendix C: Hartogs' theorem on joint analyticity

### 15.1 Weak-to-strong differentiability

A $V$-valued function $f:[a, c] \rightarrow V$ on an interval $[a, c] \subset \mathbb{R}$ is differentiable if for every $x_{o} \in[a, c]$

$$
f^{\prime}\left(x_{o}\right)=\lim _{x \rightarrow x_{o}}\left(x-x_{o}\right)^{-1}\left(f(x)-f\left(x_{o}\right)\right)
$$

exists. The function $f$ is continuously differentiable when it is differentiable and $f^{\prime}$ is continuous. A $k$-times continuously differentiable function is $C^{k}$, and a continuous function is $C^{o}$.

A $V$-valued function $f$ is weakly $C^{k}$ when for every $\lambda \in V^{*}$ the scalar-valued function $\lambda \circ f$ is $C^{k}$. This sense of weak differentiability of a function $f$ does not refer to distributional derivatives, but to differentiability of every scalar-valued function $\lambda \circ f$ where $\lambda \in V^{*}$ for $V$-valued $f$.
[15.1.1] Theorem: For quasi-complete, locally convex $V$, a weakly $C^{k} V$-valued function $f$ on an interval $[a, c]$ is strongly $C^{k-1}$.
Proof: This is a corollary of [15.7.1] below. To have $f$ be (strongly) differentiable at fixed $b \in[a, c]$ is to have (strong) continuity at $b$ of

$$
g(x)=\frac{f(x)-f(b)}{x-b} \quad(\text { for } x \neq b)
$$

Weak $C^{2}$-ness of $f$ implies that every $\lambda \circ g$ extends to a $C^{1}$ scalar-valued function on $[a, c]$. We need to get from this to a (strongly) continuous extension of $g$ to the whole interval.

The (strong) continuity of $f^{\prime}$ will follow from consideration of the function of two variables (initially for $x \neq y$ )

$$
g(x, y)=\frac{f(x)-f(y)}{x-y}
$$

The weak $C^{2}$-ness of $f$ assures that $g$ extends to a weakly $C^{1}$ function on $[a, c] \times[a, c]$. In particular, the function $x \rightarrow g(x, x)$ of (the extended) $g$ is weakly $C^{1}$, and $x \rightarrow g(x, x)$ is $f^{\prime}(x)$, so $f^{\prime}$ is weakly $C^{1}$. By [15.7.1], $f^{\prime}$ is (strongly) $C^{o}$. Suppose that we already know that $f$ is $C^{\ell}$, for $\ell<k-1$. As the $\ell^{t h}$ derivative $g=f^{(\ell)}$ of $f$ is weakly $C^{2}$, it is (strongly) $C^{1}$ by the first part of the argument. That is, $f$ is $C^{\ell+1}$. ///

### 15.2 Holomorphic vector-valued functions

Let $V$ be a quasi-complete, locally convex topological vector space. A $V$-valued function $f$ on a nonempty open set $\Omega \subset \mathbb{C}$ is (strongly) complex-differentiable when $\lim _{z \rightarrow z_{o}}\left(f(z)-f\left(z_{o}\right)\right) /\left(z-z_{o}\right)$ exists (in $V)$ for all $z_{o} \in \Omega$, where $z \rightarrow z_{o}$ specificially means for complex $z$ approaching $z_{o}$. The function $f$ is weakly holomorphic when the $\mathbb{C}$-valued functions $\lambda \circ f$ are holomorphic for all $\lambda$ in $V^{*}$. The useful version of vectorvalued meromorphy of $f$ at $z_{o}$ is that $\left(z-z_{o}\right)^{n} \cdot f(z)$ extends to a vector-valued holomorphic function at $z_{o}$ for some $n$. After some preparation, we will prove
[15.2.1] Theorem: Weakly holomorphic $V$-valued functions $f$ are continuous. (Proof in [15.8.1])
[15.2.2] Corollary: Weakly holomorphic $V$-valued functions are (strongly) holomorphic. The Cauchy integral formula applies:

$$
f(z)=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w-z} d w \quad \text { (as Gelfand-Pettis } V \text {-valued integral) }
$$

Proof: Since $f(z)$ is continuous, the integral

$$
I(z)=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w-z} d w
$$

exists as a Gelfand-Pettis integral [14.1]. Thus, for any $\lambda \in V^{*}$

$$
\lambda(I(z))=\frac{1}{2 \pi i} \int_{\gamma} \frac{(\lambda \circ f)(w)}{w-z} d w=(\lambda \circ f)(z)
$$

by the holomorphy of $\lambda \circ f$. By Hahn-Banach, linear functionals separate points, so $I(z)=f(z)$, giving the Cauchy integral formula for $f$ itself.

To prove (strong) complex-differentiability of $f$ at $z_{o}$, take $z_{o}=0$ and use $f(0)=0$, for convenience. There is a disk $|z|<3 r$ such that for every $\lambda \in V^{*}$

$$
F_{\lambda}(z)=\frac{(\lambda \circ f)(z)}{z} \quad(\text { on } 0<|z|<r)
$$

extends to a holomorphic function on $|z|<r$. Continuity of $f$ assures existence of

$$
\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w} \frac{d w}{w-z}
$$

By Cauchy theory for $\mathbb{C}$-valued functions, and Gelfand-Pettis,

$$
\lambda\left(\frac{f(z)}{z}\right)=F_{\lambda}(z)=\frac{1}{2 \pi i} \int_{\gamma} \frac{(\lambda \circ f)(w)}{w} \frac{d w}{w-z}=\lambda\left(\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w} \frac{d w}{w-z}\right)
$$

Since functionals separate points,

$$
\frac{f(z)}{z}=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w} \frac{d w}{w-z}
$$

From

$$
\frac{1}{w(w-z)}=\frac{1}{w^{2}}+\frac{z}{w^{2}(w-z)}
$$

we have

$$
\frac{f(z)}{z}=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w^{2}} d w+z \cdot \frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w^{2}(w-z)} d w
$$

Using the continuity of $f$, given a convex balanced neighborhood $U$ of 0 in $V$, the compact set $K=\{f(w):|w|=2 r\}$ is contained in some multiple $t_{o} U$ of $U$. Thus, for $|z|<r$,

$$
\frac{f(z)}{z}-\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w^{2}} d w \quad \in \quad|z| \cdot \frac{1}{(2 r)^{2} r} \cdot t_{o} U
$$

so $\lim _{z \rightarrow 0} f(z) / z$ exists. Since $f(0)=0$,

$$
\lim _{z \rightarrow z_{o}} \frac{f(z)-f\left(z_{o}\right)}{z-z_{o}}=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w) d w}{\left(w-z_{o}\right)^{2}}
$$

giving the complex differentiability of $f$.
[15.2.3] Corollary: The usual Cauchy-theory integral formulas apply. In particular, weakly holomorphic $f$ is (strongly) infinitely differentiable, in fact expressible as a convergent power series with coefficients given by Cauchy's formulas:

$$
f(z)=\sum_{n \geq 0} c_{n}\left(z-z_{o}\right)^{n} \quad \text { with } \quad c_{n}=\frac{f^{(n)}\left(z_{o}\right)}{n!}=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{\left(w-z_{o}\right)^{n+1}} d w
$$

for $\gamma$ a path with winding number +1 around $z_{o}$.
Proof: Without loss of generality, treat $z_{o}=0$, and $|z|<\rho|w|$ with $\rho<1$, and $|w|=r$. The expansion

$$
\frac{1}{w-z}=\frac{1}{w} \frac{1}{1-\frac{z}{w}}=\frac{1}{w}\left(1+\frac{z}{w}+\left(\frac{z}{w}\right)^{2}+\ldots+\left(\frac{z}{w}\right)^{N}+\frac{(z / w)^{N+1}}{1-\frac{z}{w}}\right)
$$

combined with an integration around $\gamma$ against $f(w)$, and the basic Cauchy integral formula, give

$$
f(z)=\sum_{n=0}^{N} \frac{1}{2 \pi i} \int_{\gamma} \frac{f(w) d w}{w^{n+1}} \cdot z^{n}+\frac{1}{2 \pi i} \int_{\gamma} \frac{1}{w^{N+1}} \frac{f(w) d w}{w-z} \cdot z^{N+1}
$$

Much as in the previous proof, given a convex balanced neighborhood $U$ of 0 in $V$, the compact set $K=\{f(w):|w|=r\}$ is contained in some multiple $t_{o} U$ of $U$, and

$$
\frac{1}{2 \pi i} \int_{\gamma} \frac{1}{w^{N+1}} \frac{f(w) d w}{w-z} \cdot z^{N+1} \in \frac{1}{r^{N+1}} \cdot t_{o} U \cdot \frac{1}{r(1-\rho)} \cdot(\rho r)^{N+1}=U \frac{t_{o}}{r(1-\rho)} \rho^{N+1}
$$

Since $0<\rho<1, \rho^{N+1} / r(1-\rho)<1$ for sufficiently large $N$, so the leftover term is inside given $U$.
Appendix [15.A] discusses the differentiability of power series with coefficients in topological vector spaces.
The next section collects some important corollaries of the main result, prior to preparation for the proof that weak holomorphy implies continuity,

### 15.3 Holomorphic $\operatorname{Hol}(\Omega, V)$-valued functions

The vector-valued versions of Cauchy's formulas have useful corollaries. First, recall some aspects of the classical scalar-valued case.

For open $\phi \neq \Omega \subset \mathbb{C}$, give the space $\operatorname{Hol}(\Omega)$ of holomorphic functions on $\Omega$ the topology given by the seminorms $\mu_{K}(f)=\sup _{z \in K}|f(z)|$ for compacts $K \subset \Omega$.
[15.3.1] Claim: $\operatorname{Hol}(\Omega)$ is a Fréchet space.
Proof: Let $\left\{f_{n}\right\}$ be a Cauchy sequence in that topology. As in [13.5], the pointwise limit $f(z)=\lim _{n} f_{n}(z)$ is at least continuous. Then, for a small circle $\gamma$ inside $\Omega$ and enclosing $z$,

$$
f(z)=\lim _{n} f_{n}(z)=\lim _{n} \frac{1}{2 \pi i} \int_{\gamma} \frac{f_{n}(w)}{w-z} d w
$$

Since $\gamma$ is compact and the limit is uniformly approached on compacts, this gives

$$
f(z)=\frac{1}{2 \pi i} \int_{\gamma} \lim _{n} \frac{f_{n}(w)}{w-z} d w=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w-z} d w
$$

Direct estimates (simpler than in the previous section) show that the latter integral is complex-differentiable in $w$.

Let $V$ be quasi-complete, locally convex, with topology given by seminorms $\{\nu\}$. The space $\operatorname{Hol}(\Omega, V)$ of holomorphic $V$-valued functions on $\Omega$ has the natural topology given by seminorms

$$
\mu_{\nu, K}(f)=\sup _{z \in K} \nu(f(z)) \quad(\text { compacts } K \subset \Omega, \text { seminorms } \nu \text { on } V)
$$

This topology is obviously the analogue of the sups-on-compacts seminorms on scalar-valued holomorphic functions, and there is the analogous corollary of the vector-valued Cauchy formulas:
[15.3.2] Corollary: $\operatorname{Hol}(\Omega, V)$ is locally convex, quasi-complete.
Proof: Let $\left\{f_{n}\right\}$ be a bounded Cauchy net. Just as in the scalar case, the pointwise limits $\lim _{n} f_{n}(z)$ exist. The same three-epsilon argument as for scalar-valued functions will show that the pointwise limit exists and is continuous, as follows. First, using compact $K=\{z\}$, the value $\mu_{\{z\}, \nu}(f)$ is just $\nu(f(z))$. Thus, by quasicompleteness of $V$, for each fixed $z$ the bounded Cauchy net $f_{n}(z)$ converges to a value $f(z)$. Given $\varepsilon>0$ and $z_{o} \in \Omega$, let $K$ be a compact neighborhood of $z_{o}$, and take $N$ sufficiently large so that $\nu\left(f_{m}(z)-f_{n}\left(z^{\prime}\right)\right)<\varepsilon$ for all $z, z^{\prime} \in K$ and all $m, n \geq N$. Then

$$
\mu_{K, \nu}\left(f(z)-f\left(z_{o}\right)\right) \leq \mu_{K, \nu}\left(f(z)-f_{n}(z)\right)+\mu_{K, \nu}\left(f_{n}(z)-f_{n}\left(z_{o}\right)\right)+\mu_{K, \nu}\left(f_{n}\left(z_{o}\right)-f\left(z_{o}\right)\right) \leq 3 \varepsilon
$$

proving the continuity of the pointwise limit. Then, as in the previous scalar-valued argument, the vectorvalued Cauchy formula gives, for a small circle $\gamma$ inside $\Omega$ and enclosing $z$,

$$
f(z)=\lim _{n} f_{n}(z)=\lim _{n} \frac{1}{2 \pi i} \int_{\gamma} \frac{f_{n}(w)}{w-z} d w
$$

with Gelfand-Pettis integrals. Since $\gamma$ is compact and the limit is uniformly approached on compacts, this gives

$$
f(z)=\frac{1}{2 \pi i} \int_{\gamma} \lim _{n} \frac{f_{n}(w)}{w-z} d w=\frac{1}{2 \pi i} \int_{\gamma} \frac{f(w)}{w-z} d w
$$

Again, the differentiability of latter integral is directly verifiable, and $f$ is holomorphic.
It is occasionally useful to iterate the previous ideas: A $V$-valued function $f(z, w)$ on a non-empty open subset $\Omega \subset \mathbb{C}^{2}$ is complex analytic when it is locally expressible as a convergent power series in $z$ and $w$, with coefficients in $V$. The two-variable version of the discussion of convergence of power series with coefficients in $V$ in appendix [15.A] succeeds without incident in the two-variable case. [83]
[15.3.3] Corollary: Let $f(z, w)$ be complex-analytic $\mathbb{C}$-valued in two variables, on a domain $\Omega_{1} \times \Omega_{2} \subset \mathbb{C}^{2}$. Then the function $w \longrightarrow(z \rightarrow f(z, w))$ is a holomorphic $\operatorname{Hol}\left(\Omega_{1}\right)$-valued function on $\Omega_{2}$.
Proof: The issue is the uniformity in $z$ in compacts $K$ of the limit

$$
\lim _{h \rightarrow 0} \frac{f(z, w+h)-f(z, w)}{h}
$$

Using the scalar-valued Cauchy integral, for a small circle $\gamma$ about $w$, letting $f_{2}$ be the partial derivative of $f$ with respect to its second argument,

$$
\left.\left.\begin{array}{c}
\frac{f(z, w+h)-}{} \quad f(z, w) \\
h \\
=\frac{1}{2 \pi i} \int_{2}(z, w)=\frac{1}{2 \pi i} \int_{\gamma} f(z, \zeta)\left(\frac{\frac{1}{\zeta-(w+h)}-\frac{1}{\zeta-w}}{h}-\frac{1}{(\zeta-w)^{2}}\right) d \zeta \\
(\zeta-(w+h))(\zeta-h) \\
(\zeta-w)^{2}
\end{array}\right) d \zeta\right)
$$

The two-variable analytic function $z, \zeta \rightarrow f(z, \zeta)$ is certainly continuous as a function of two variables, so is uniformly continuous on compacts $K \times \gamma$. Thus, the limit as $h \rightarrow 0$ is approached uniformly.

[^24]
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Application of the vector-valued form of Cauchy's integrals gives the same result for $f(z, w)$ taking values in a quasi-complete, locally convex $V$ :
[15.3.4] Corollary: Let $V$ be quasi-complete, locally convex. Let $f(z, w)$ be complex-analytic $V$-valued in two variables, on a domain $\Omega_{1} \times \Omega_{2} \subset \mathbb{C}^{2}$. Then the function $w \longrightarrow(z \rightarrow f(z, w))$ is a holomorphic $\operatorname{Hol}\left(\Omega_{1}, V\right)$-valued function on $\Omega_{2}$.

### 15.4 Banach-Alaoglu: compactness of polars

The polar $U^{o}$ of an open neighborhood $U$ of 0 in a topological vector space $V$ is

$$
U^{o}=\left\{\lambda \in V^{*}:|\lambda u| \leq 1, \text { for all } u \in U\right\}
$$

[15.4.1] Theorem: (Banach-Alaoglu) In the weak dual topology on $V^{*}$ the polar $U^{o}$ of an open neighborhood $U$ of 0 in $V$ is compact.
Proof: For every $v$ in $V$ there is real $t_{v}$ sufficiently large such that $v \in t_{v} \cdot U$, and $|\lambda v| \leq t_{v}$ for $\lambda \in U^{o}$. Tychonoff gives compactness of the product

$$
P=\prod_{v \in V}\left\{z \in \mathbb{C}:|z| \leq t_{v}\right\} \subset \prod_{v \in V} \mathbb{C}
$$

Map $V^{*}$ to $\prod_{v \in V} C$ by $j(\lambda)=\{\lambda(v): v \in V\}$. By design, $j\left(U^{o}\right) \subset P$. To prove the compactness of $U^{o}$ it suffices to show that the weak dual topology on $U^{o}$ is identical to the subspace topology on $j\left(U^{o}\right)$ inherited from $P$, and that $j\left(U^{o}\right)$ is closed in $P$.

The sub-basis sets

$$
\left\{\lambda \in V^{*}:\left|\lambda v-\lambda_{o} v\right|<\delta\right\} \quad(\text { for } v \in V \text { and } \delta>0)
$$

for $V^{*}$ are mapped by $j$ to the sub-basis sets

$$
\left\{p \in P:\left|p_{v}-\lambda_{o} v\right|<\delta\right\} \quad(\text { for } v \in V \text { and } \delta>0)
$$

for the product topology on $P$. That is, $j$ maps $U^{o}$ with the weak star-topology homeomorphically to $j\left(U^{o}\right)$.
To show that $j\left(U^{o}\right)$ is closed in $P$, consider $L$ in the closure of $U^{o}$ in $P$. Given $x, y \in V, a, b \in \mathbb{C}$, the sets

$$
\left\{p \in P:\left|(p-L)_{x}\right|<\delta\right\} \quad\left\{p \in P:\left|(p-L)_{y}\right|<\delta\right\} \quad\left\{p \in P:\left|(p-L)_{a x+b y}\right|<\delta\right\}
$$

are open in $P$ and contain $L$, so meet $j\left(U^{o}\right)$. Let $\lambda \in j\left(U^{o}\right)$ lie in the intersection of these three sets and $j\left(U^{o}\right)$. Then

$$
\begin{gathered}
\left|a L_{x}+b L_{y}-L_{a x+b y}\right| \leq|a| \cdot\left|L_{x}-\lambda x\right|+|b| \cdot\left|L_{y}-\lambda y\right|+\left|L_{a x+b y}-\lambda(a x+b y)\right|+|a \lambda x+b \lambda y-\lambda(a x+b y)| \\
\leq|a| \cdot \delta+|b| \cdot \delta+\delta+0 \quad(\text { for every } \delta>0)
\end{gathered}
$$

so $L$ is linear. Given $\varepsilon>0$, for $N$ be a neighborhood of 0 in $V$ such that $x-y \in N$ implies $\lambda x-\lambda y \in N$,

$$
\left|L_{x}-L_{y}\right|=\left|L_{x}-\lambda x\right|+\left|L_{y}-\lambda y\right|+|\lambda x-\lambda y| \delta+\delta+\varepsilon
$$

Thus, $L$ is continuous. Also, $\left|L_{x}-\lambda x\right|<\delta$ for all $x \in U$ and all $\delta>0$, so $L \in j\left(U^{o}\right)$, and $j\left(U^{o}\right)$ is closed, giving compactness.

### 15.5 Variant Banach-Steinhaus/uniform boundedness

This variant of the Banach-Steinhaus (uniform boundedness) theorem is used with Banach-Alaoglu to show that weak boundedness implies boundedness in a locally convex space, the starting point for weak-tostrong principles. It uses the version of Baire category for locally compact Hausdorff spaces, rather than complete metric spaces.
[15.5.1] Theorem: (Variant Banach-Steinhaus) Let $K$ be a compact convex set in a topological vectorspace $X$, and $\mathscr{T}$ a set of continuous linear maps $X \rightarrow Y$ from $X$ to another topological vectorspace $Y$. Suppose that for every individual $x \in K$ the collection of images $\mathscr{T} x=\{T x: T \in \mathscr{T}\}$ is bounded in $Y$. Then $B=\bigcup_{x \in K} \mathscr{T} x$ is bounded in $Y$.
Proof: Let $U, V$ be balanced neighborhoods of 0 in $Y$ so that $\bar{U}+\bar{U} \subset V$, and let

$$
E=\bigcap_{T \in \mathscr{T}} T^{-1}(\bar{U})
$$

By the boundedness of $\mathscr{T} x$, there is a positive integer $n$ such that $\mathscr{T} x \subset n U$, and then $x \in n E$. For every $x \in K$ there is such $n$, so

$$
K=\bigcup_{n}(K \cap n E)
$$

Since $E$ is closed, the version of the Baire category theorem for locally compact Hausdorff spaces implies that at least one set $K \cap n E$ has non-empty interior in $K$. For such $n$, let $x_{o}$ be an interior point of $K \cap n E$. Pick a balanced neighborhood $W$ of 0 in $X$ such that

$$
K \cap\left(x_{o}+W\right) \subset n E
$$

Since $K$ is compact, it is bounded, so $K-x_{o}$ is bounded, and $K \subset x_{o}+t W$ for large enough positive real $t$. Since $K$ is convex, $\left(1-t^{-1}\right) x+t^{-1} x \in K$ for any $x \in K$ and $t \geq 1$. At the same time,

$$
z-x_{o}=t^{-1}\left(x-x_{o}\right) \in W \quad(\text { for large enough } t)
$$

by the boundedness of $K$, so $z \in x_{o}+W$. Thus, $z \in K \cap\left(x_{o}+V\right) \subset n E$. From the definition of $E, T E \subset \bar{U}$, so $T(n E)=n T(E) \subset n \bar{U}$. And $x=t z-(t-1) x_{o}$ yields

$$
T x \in \operatorname{tn} \bar{U}-(t-1) n \bar{U} \subset \operatorname{tn}(\bar{U}+\bar{U})
$$

by the balanced-ness of $U$. Since $\bar{U}+\bar{U} \subset V$, we have $B \subset \operatorname{tn} V$. Since $V$ was arbitrary, this proves the boundedness of $B$.

### 15.6 Weak boundedness implies (strong) boundedness

[15.6.1] Theorem: Let $V$ be a locally convex topological vectorspace. A subset $E$ of $V$ is bounded if and only if it is weakly bounded.
Proof: For the proof, we need the notion of second polar $N^{o o}$ of an open neighborhood $N$ of 0 in a topological vector space $V$ :

$$
N^{o o}=\left\{v \in V:|\lambda v| \leq 1 \text { for all } \lambda \in N^{o}\right\}
$$

where $N^{o}$ is the polar of $N$. Conveniently,
[15.6.2] Claim: (On second polars) For $V$ a locally convex topological vectorspace and $N$ a convex, balanced neighborhood of 0 , the second polar $N^{o o}$ of $N$ is the closure $\bar{N}$ of $N$.
Proof: Certainly $N$ is contained in $N^{o o}$, and in fact $\bar{N}$ is contained in $N^{o o}$ since $N^{o o}$ is closed. By the local convexity of $V$, Hahn-Banach implies that for $v \in V$ but $v \notin \bar{N}$ there is $\lambda \in V^{*}$ such that $\lambda v>1$ and $\left|\lambda v^{\prime}\right| \leq 1$ for all $v^{\prime} \in \bar{N}$. Thus, $\lambda$ is in $N^{o}$, and every element $v \in N^{o o}$ is in $\bar{N}$, so $N^{o o}=\bar{N}$.

Returning to the proof of the theorem: clearly boundedness implies weak boundedness. On the other hand, take $E$ weakly bounded, and $U$ be a neighborhood of 0 in $V$ in the original topology. By local convexity, there is a convex (and balanced) neighborhood $N$ of 0 such that the closure $\bar{N}$ is contained in $U$.

By the weak boundedness of $E$, for each $\lambda \in V^{*}$ there is a bound $b_{\lambda}$ such that $|\lambda x| \leq b_{\lambda}$ for $x \in E$. By Banach-Alaoglu the polar $N^{o}$ of $N$ is compact in $V^{*}$. The functions $\lambda \rightarrow \lambda x$ are continuous, so by variant Banach-Steinhaus there is a uniform constant $b<\infty$ such that $|\lambda x| \leq b$ for $x \in E$ and $\lambda \in N^{o}$. Thus, $b^{-1} x$ is in the second polar $N^{o o}$ of $N$, shown by the previous proposition to be the closure $\bar{N}$ of $N$. That is, $b^{-1} x \in \bar{N}$. By the balanced-ness of $N, E \subset t \bar{N} \subset t U$ for any $t>b$, so $E$ is bounded.

### 15.7 Proof that weak $C^{1}$ implies strong $C^{o}$

The claim below, needed to complete the proof of [15.1.1] that weak $C^{k}$ implies (strong) $C^{k-1}$, is an application of the fact that weak boundedness implies boundedness.
[15.7.1] Claim: Let $V$ be a quasi-complete locally convex topological vector space. Fix real numbers $a \leq b \leq c$. Let $g$ be a $V$-valued function defined on $[a, b) \cup(b, c]$. Suppose that for $\lambda \in V^{*}$ the scalar-valued function $\lambda \circ g$ extends to a $C^{1}$ function $F_{\lambda}$ on the whole interval $[a, c]$. Then $g(b)$ can be chosen such that the extended $g(x)$ is (strongly) continuous on $[a, c]$.
Proof: For each $\lambda \in V^{*}$, let $F_{\lambda}$ be the extension of $\lambda \circ g$ to a $C^{1}$ function on $[a, c]$. The differentiability of $F_{\lambda}$ implies that for each $\lambda$ the function

$$
\Phi_{\lambda}(x, y)=\frac{F_{\lambda}(x)-F_{\lambda}(y)}{x-y} \quad(\text { for } x \neq y)
$$

has a continuous extension $\tilde{\Phi}_{\lambda}$ to the compact set $[a, c] \times[a, c]$. The image $C_{\lambda}$ of $[a, c] \times[a, c]$ under this continuous map is compact in $\mathbb{R}$, so bounded. Thus, the subset

$$
\left\{\frac{\lambda f(x)-\lambda f(y)}{x-y}: x \neq y\right\} \subset C_{\lambda}
$$

is bounded in $\mathbb{R}$. That is,

$$
E=\left\{\frac{g(x)-g(y)}{x-y}: x \neq y\right\} \subset V
$$

is weakly bounded. Because weakly bounded implies (strongly) bounded, $E$ is (strongly) bounded. That is, for a balanced, convex neighborhood $N$ of 0 in $V$, there is $t_{o}$ such that $(g(x)-g(y)) /(x-y) \in t N$ for $x \neq y$ in $[a, c]$ and $t \geq t_{o}$. That is, $g(x)-g(y) \in(x-y) t N$. Given $N$ and $t_{o}$ as above, $g(x)-g(y) \in N$ for $|x-y|<\frac{1}{t_{o}}$. That is, as $x \rightarrow b$ the collection $g(x)$ is a bounded Cauchy net. By quasi-completeness, define $g(b) \in V$ as the limit of the values $g(x)$. For $x \rightarrow y$ the values $g(x)$ approach $g(y)$, so this extension of $g$ is continuous on $[a, c]$.

### 15.8 Proof that weak holomorphy implies continuity

With the above preparation, we prove that weak holomorphy implies (strong) continuity, completing the proof of [15.2.2], as another application of the fact that weak boundedness implies boundedness, by an argument parallel to that of [15.7] that weak $C^{1}$ implies $C^{o}$ for vector-valued functions on $[a, b]$.
[15.8.1] Claim: Weak holomorphy implies (strong) continuity.
Proof: To show that weak holomorphy of $f$ implies $f: D \rightarrow V$ is (strongly) continuous, without loss of generality prove continuity at $z=0$ and suppose $f(0)=0 \in V$. Since $\lambda \circ f$ is holomorphic for each $\lambda \in V^{*}$ and vanishes at 0 , each function $(\lambda \circ f)(z) / z$ initially defined on a punctured disk at 0 extends to a holomorphic function on a full disk at 0 . By Cauchy theory for the scalar-valued holomorphic function $z \rightarrow \frac{\lambda(f(z))}{z}$,

$$
\frac{(\lambda \circ f)(z)}{z}=\frac{1}{2 \pi i} \int_{\gamma} \frac{1}{w-z} \cdot \frac{(\lambda \circ f)(w)}{w} d w
$$

where $\gamma$ is a circle of radius $2 r$ centered at 0 , and $|z|<r$. With $M_{\lambda}$ the sup of $|\lambda \circ f|$ on $\gamma$,

$$
\left|\frac{(\lambda \circ f)(z)}{z}\right| \leq \frac{\text { length } \gamma}{2 \pi} \cdot \frac{1}{2 r-r} \cdot \frac{M_{\lambda}}{2 r}=\frac{1}{2 \pi} \cdot(2 \pi \cdot 2 r) \cdot \frac{1}{r} \cdot \frac{M_{\lambda}}{2 r}=\frac{M_{\lambda}}{r}
$$

Thus, the set of values

$$
S=\left\{\frac{f(z)}{z}:|z| \leq r\right\}
$$

is weakly bounded. Weak boundedness implies (strong) boundedness, so $S$ is bounded. That is, given a balanced convex neighborhood $N$ of 0 in $V$, there is $t_{o}>0$ such that for complex $w$ with $|w| \geq t_{o}$, the set $S$ lies inside $w N$. Then $f(z) \in z w N$ and $f(z) \in N$ for $|z|<|w|$. As $f(0)=0$, we have proven that, given $N$, for $z$ sufficiently near $0 f(z)-f(0) \in N$. This is (strong) continuity.

## 15.A Appendix: vector-valued power series

We should confirm that power series with values in a quasi-complete, locally compact vectorspace $V$ behave essentially as well as scalar-valued ones. First,
[15.A.1] Lemma: Let $c_{n}$ be a bounded sequence of vectors in the locally convex, quasi-complete topological vector space $V$. Let $z_{n}$ be a sequence of complex numbers, let $0 \leq r_{n}$ be real numbers such that $\left|z_{n}\right| \leq r_{n}$, and suppose that $\sum_{n} r_{n}<+\infty$. Then $\sum_{n} c_{n} z_{n}$ converges in $V$. Further, given a convex balanced neighborhood $U$ of 0 in $V$ let $t$ be a positive real such that for all complex $w$ with $|w| \geq t$ we have $\left\{c_{n}\right\} \subset t U$. Then

$$
\sum_{n} c_{n} z_{n} \in\left(\sum_{n}\left|z_{n}\right|\right) \cdot t U \subset\left(\sum_{n} r_{n}\right) \cdot t U
$$

Proof: For convex balanced neighborhood $N$ of 0 in the topological vector space, with complex numbers $z$ and $w$ such that $|z| \leq|w|$, then $z N \subset w N$, since $|z / w| \leq 1 \operatorname{implies}(z / w) N \subset N$, or $z N \subset w N$. Further, for an absolutely convergent series $\sum_{n} \alpha_{n}$ of complex numbers, for any $n_{o}$

$$
\sum_{n \leq n_{o}}\left(\alpha_{n} \cdot V\right)=\sum_{n \leq n_{o}}\left(\left|\alpha_{n}\right| \cdot V\right) \subset\left(\sum_{n \leq n_{o}}\left|\alpha_{n}\right|\right) \cdot N \subset\left(\sum_{n<\infty}\left|\alpha_{n}\right|\right) \cdot N
$$

For a balanced open $U$ containing 0 , let $t$ be large enough such that for any complex $w$ with $|w| \geq t$ the sequence $c_{n}$ is contained in $w U$. The previous discussion shows that

$$
\sum_{m \leq \ell \leq n} c_{\ell} z_{\ell} \in\left(\left|z_{m}\right|+\ldots+\left|z_{n}\right|\right) \cdot t U
$$

Given $\varepsilon>0$, invoking absolute convergence, take $m$ sufficiently large such that $\left|z_{m}\right|+\ldots+\left|z_{n}\right|<t \cdot \varepsilon$ for all $n \geq m$. Then

$$
\sum_{m \leq \ell \leq n} c_{\ell} z_{\ell} \in t \cdot(\varepsilon / t) \cdot U=U
$$

Thus, the original series is convergent. Since $X$ is quasi-complete the limit exists in $V$. The last containment assertion follows from this discussion, as well.
[15.A.2] Corollary: Let $c_{n}$ be a bounded sequence of vectors in a locally convex quasi-complete topological vector space $V$. Then on $|z|<1$ the series $f(z)=\sum_{n} c_{n} z^{n}$ converges and gives a holomorphic $V$-valued function. That is, the function is infinitely-many-times complex-differentiable.
Proof: The lemma shows that the series expressing $f(z)$ and its apparent $k^{t h}$ derivative $\sum_{n} c_{n}\binom{n}{k} z^{n-k}$ all converge for $|z|<1$. The usual direct proof of Abel's theorem on the differentiability of (scalar-valued)
power series can be adapted to prove the infinite differentiability of the $X$-valued function given by this power series, as follows. Let

$$
g(z)=\sum_{n \geq 0} n c_{n} z^{n-1}
$$

Then

$$
\frac{f(z)-f(w)}{z-w}-g(w)=\sum_{n \geq 1} c_{n}\left(\frac{z^{n}-w^{n}}{z-w}-n w^{n-1}\right)
$$

For $n=1$, the expression in the parentheses is 1 . For $n>1$, it is

$$
\begin{gathered}
\left(z^{n-1}+z^{n-2} w+\ldots+z w^{n-2}+w^{n-1}\right)-n w^{n-1} \\
=\left(z^{n-1}-w^{n-1}\right)+\left(z^{n-2} w-w^{n-1}\right)+\ldots+\left(z^{2} w^{n-3}-w^{n-1}\right)+\left(z w^{n-2}-w^{n-1}\right)+\left(w^{n-1}-w^{n-1}\right) \\
=(z-w)\left[\left(z^{n-2}+\ldots+w^{n-2}\right)+w\left(z^{n-3}+\ldots+w^{n-3}\right)+\ldots+w^{n-3}(z+w)+w^{n-2}+0\right] \\
=(z-w) \sum_{k=0}^{n-2}(k+1) z^{n-2-k} w^{k}
\end{gathered}
$$

For $|z| \leq r$ and $|w| \leq r$ the latter expression is dominated by

$$
|z-w| \cdot r^{n-2} \frac{n(n-1)}{2}<|z-w| \cdot n^{2} r^{n-2}
$$

Let $U$ be a balanced neighborhood of 0 in $X$, and $t$ a sufficiently large real number such that for all complex $w$ with $|w| \geq t$ all $c_{n}$ lie in $w U$. For $|z| \leq r<1$ and $|w| \leq r<1$, by the lemma,

$$
\frac{f(z)-f(w)}{z-w}-g(w)=(z-w) \sum_{n \geq 2} c_{n} \cdot\left(\sum_{k=0}^{n-2}(k+1) z^{n-2-k} w^{k}\right) \in(z-w) \cdot\left(\sum_{n} n^{2} r^{n-2}\right) \cdot t U
$$

Thus, as $z \rightarrow w$, eventually $\frac{f(z)-f(w)}{z-w}-g(w)$ lies in $U$.
[15.A.3] Corollary: Let $c_{n}$ be a sequence of vectors in a Banach space $X$ such that for some $r>0$ the series $\sum\left|c_{n}\right| \cdot r^{n}$ converges in $X$. Then for $|z|<r$ the series $f(z)=\sum c_{n} z^{n}$ converges and gives a holomorphic (infinitely-many times complex-differentiable) $X$-valued function.

## 15.B Appendix: two forms of the Baire category theorem

This standard result is indispensable and mysterious. We give the more-typical version for complete metric spaces in parallel with the argument for locally compact Hausdorff spaces.

A set $E$ in a topological space $X$ is nowhere dense if its closure $\bar{E}$ contains no non-empty open. A countable union of nowhere dense sets is said to be of first category, while every other subset is of second category. The idea, not suggested by this traditional terminology, is that first category sets are small, while second category sets are large. The theorem asserts that (non-empty) complete metric spaces and locally compact Hausdorff spaces are of second category.
[15.B.1] Theorem: For $X$ be a complete metric space or a locally compact Hausdorff topological space, the intersection of a countable collection $U_{1}, U_{2}, \ldots$ of dense open subsets $U_{i}$ of $X$ is still dense in $X$.
Proof: Let $B_{o}$ be a non-empty open set in $X$, and show that $\bigcap_{i} U_{i}$ meets $B_{o}$. Suppose that we have inductively chosen an open ball $B_{n-1}$. By the denseness of $U_{n}$, there is an open ball $B_{n}$ whose closure $\overline{B_{n}}$ satisfies

$$
\overline{B_{n}} \subset B_{n-1} \cap U_{n}
$$

Further, for complete metric spaces, take $B_{n}$ to have radius less than $1 / n$ (or any other sequence of reals going to 0 ), and in the locally compact Hausdorff case take $B_{n}$ to have compact closure.

Let

$$
K=\bigcap_{n \geq 1} \overline{B_{n}} \subset B_{o} \cap \bigcap_{n \geq 1} U_{n}
$$

For complete metric spaces, the centers of the nested balls $B_{n}$ form a Cauchy sequence (since they are nested and the radii go to 0 ). By completeness, this Cauchy sequence converges, and the limit point lies inside each closure $\overline{B_{n}}$, so lies in the intersection. In particular, $K$ is non-empty. For locally compact Hausdorff spaces, the intersection of a nested family of non-empty compact sets is non-empty, so $K$ is non-empty, and $B_{o}$ necessarily meets the intersection of the $U_{n}$.

## 15.C Appendix: Hartogs' theorem on joint analyticity

This proof roughly follows that in [Hörmander 1973] which roughly follows Hartogs' original argument [Hartogs 1906].
[15.C.1] Theorem: Let $f$ be a $\mathbb{C}$-valued function defined in a non-empty open set $U \subset \mathbb{C}^{n}$. If $f$ is analytic in each variable $z_{j}$ when the other coordinates $z_{k}$ for $k \neq j$ are fixed, then $f$ is analytic as a function of all $n$ coordinates.
[15.C.2] Remark: It is striking that no additional hypothesis on $f$ is used beyond its separate analyticity: there is no assumption of continuity, nor even of measurability. Indeed, the beginning of the proof illustrates the fact that an assumption of continuity trivializes things. The strength of the theorem is that no hypothesis whatsoever is necessary.
Proof: The assertion is local, so it suffices to prove it when the open set $U$ is a polydisk. The argument approaches the full assertion in stages.

First, suppose that $f$ is continuous on the closure $\bar{U}$ of a polydisk $U$, and separately analytic. Even without continuity, simply by separate analyticity, an $n$-fold iterated version of Cauchy's one-variable integral formula is valid, namely

$$
f(z)=\frac{1}{(2 \pi i)^{n}} \int_{C_{1}} \ldots \int_{C_{n}} \frac{f(\zeta)}{\left(\zeta_{1}-z_{1}\right) \ldots\left(\zeta_{n}-z_{n}\right)} d \zeta_{1} \ldots d \zeta_{n}
$$

where $C_{j}$ is the circle bounding the disk in which $z_{j}$ lies, traversed in the positive direction. The integral is a compactly supported integral of the function

$$
\left(\zeta_{1}, \ldots, \zeta_{n}\right) \rightarrow \frac{f\left(\zeta_{1}, \ldots, \zeta_{n}\right)}{\left(\zeta_{1}-z_{1}\right) \ldots\left(\zeta_{n}-z_{n}\right)}
$$

For $\left|z_{j}\right|<\left|\zeta_{j}\right|$, the geometric series expansion

$$
\frac{1}{\zeta_{j}-z_{j}}=\sum_{n \geq 0} \frac{z_{j}^{n}}{\zeta_{j}^{n+1}}
$$

can be substituted into the latter integral. Fubini's theorem justifies interchange of summation and integration, yielding a (convergent) power series for $f(z)$. Thus, continuity of $f(z)$ (with separate analyticity) implies joint continuity.

Note that if we could be sure that every conceivable integral of analytic functions were analytic, then this iterated one-variable Cauchy formula would prove (joint) analyticity immediately. However, it is not obvious that separate analyticity implies continuity, for example.

Next we see that boundedness of a separately analytic function on a closed polydisk implies continuity, using Schwarz' lemma and its usual corollary:
[15.C.3] Lemma: (Schwarz) Let $g(z)$ be a holomorphic function on $\{z \in \mathbb{C}:|z|<1\}$, with $g(0)=0$ and $|g(z)| \leq 1$. Then $|g(z)| \leq|z|$ and $\left|g^{\prime}(0)\right| \leq 1$. (Proof: Apply the maximum modulus principle to $f(z) / z$ on disks of radius less than 1.)
[15.C.4] Corollary: Let $g(z)$ be a holomorphic function on $\{z \in \mathbb{C}:|z|<r\}$, with $|g(z)| \leq B$ for a bound $B$. Then for $z, \zeta$ in that disk,

$$
|g(z)-g(\zeta)| \leq 2 \cdot B \cdot\left|\frac{r(z-\zeta)}{r^{2}-\bar{\zeta} z}\right|
$$

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Proof: (of corollary) The linear fractional transformation

$$
\mu: z \rightarrow r \cdot\left(\begin{array}{cc}
1 & \zeta / r \\
\bar{\zeta} / r & 1
\end{array}\right)(r z)=r \cdot \frac{z+r \zeta}{\bar{\zeta} z+r}
$$

sends the disk of radius 1 to the disk of radius $r$, and sends 0 to $\zeta$. Then the function

$$
z \rightarrow \frac{g(\mu(z))-g(\zeta)}{2 B}
$$

is normalized to match Schwarz' lemma, namely that it vanishes at 0 , and is bounded by 1 on the open unit disk. Thus, we conclude that for $|z|<1$

$$
\left|\frac{g(\mu(z))-g(\zeta)}{2 B}\right| \leq|z|
$$

Replace $z$ by

$$
\mu^{-1}(z)=\frac{r(z-\zeta)}{r^{2}-\bar{\zeta} z}
$$

to obtain

$$
\left|\frac{g(z)-g(\zeta)}{2 B}\right| \leq\left|\frac{r(z-\zeta)}{r^{2}-\bar{\zeta} z}\right|
$$

as asserted in the corollary.
Now let $f$ be separately analytic and bounded on the closure of the polydisk $\left\{\left(z_{1}, \ldots, z_{n}\right):\left|z_{j}\right|<r_{j}\right\}$. We show that $f$ is (jointly) analytic by proving it is continuous, invoking the first part of the proof (above). Let $B$ be a bound for $|f|$ on the closed polydisk. We claim that the inequality

$$
|f(z)-f(\zeta)| \leq 2 B \sum_{1 \leq j \leq n} \frac{r_{j}\left|z_{j}-\zeta_{j}\right|}{\left|r_{j}^{2}-\bar{\zeta}_{j} z_{j}\right|}
$$

holds, which would prove continuity. Because of the telescoping expression

$$
f(z)-f(\zeta)=\sum_{1 \leq j \leq n}\left(f\left(\zeta_{1}, \ldots, \zeta_{j-1}, z_{j}, \ldots, z_{n}\right)-f\left(\zeta_{1}, \ldots, \zeta_{j}, \zeta_{j}, z_{j+1}, \ldots, z_{n}\right)\right)
$$

it suffices to prove the inequality in the single-variable case, which is the immediate corollary to Schwarz' lemma as above. Thus, a bounded separately analytic $f$ is continuous, and (from above) jointly analytic.

Now we do induction on the dimension $n$ : suppose that Hartogs' theorem is proven on $\mathbb{C}^{n-1}$, and prove it for $\mathbb{C}^{n}$. Here the Baire Category Theorem intervenes, getting started on the full statement of the theorem by first showing that a separately analytic function must be bounded on some polydisk, hence (from above) continuous on that polydisk, hence (from above) analytic on that polydisk.

Let $f$ be separately analytic on a (non-empty) closed polydisk $D=\prod_{1 \leq j \leq n} D_{j}$, where $D_{j}$ is a disk in $\mathbb{C}$. We claim that there exist non-empty closed disks $E_{j} \subset D_{j}$ with $E_{n}=D_{n}$ such that $f$ is bounded on $E=\prod_{1 \leq j \leq n} E_{j}$ (and, hence, $f$ is analytic in $\left.E\right)$.

To see this, for each bound $B>0$ let

$$
\Omega_{B}=\left\{z^{\prime} \in \prod_{1 \leq j \leq n-1} E_{j}:\left|f\left(z^{\prime}, z_{n}\right)\right| \leq B \text { for all } z_{n} \in E_{n}\right\}
$$

By induction, for fixed $z_{n}$ the function $z^{\prime} \rightarrow f\left(z^{\prime}, z_{n}\right)$ is analytic, so continuous, so $\Omega_{B}$ is closed. For any fixed $z^{\prime}$, the function $z_{n} \rightarrow f\left(z^{\prime}, z_{n}\right)$ is assumed analytic, so is continuous on the closed disk $E_{n}=D_{n}$, hence bounded. Thus

$$
\bigcup_{B=1}^{\infty} \Omega_{B}=\prod_{1 \leq j \leq n-1} D_{j}
$$

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Then the Baire Category Theorem shows that some $\Omega_{B}$ must have non-empty interior, so must contain a (non-empty) closed polydisk, as claimed

Now let $f$ be separately analytic in a polydisk

$$
D=\left\{\left(z_{1}, \ldots, z_{n}\right):\left|z_{j}\right|<r\right\} \subset \mathbb{C}^{n}
$$

analytic in $z^{\prime}=\left(z_{1}, \ldots, z_{n}-1\right)$ for fixed $z_{n}$, and suppose that $f$ is analytic in a smaller (non-empty) polydisk

$$
E=\left(\prod_{1 \leq j \leq n-1}\left\{z_{j} \in \mathbb{C}:\left|z_{j}\right|<\varepsilon\right\}\right) \times\left\{z_{n} \in \mathbb{C}:\left|z_{n}\right|<r\right\}
$$

inside $D$. Then we claim that $f$ is analytic on the original polydisk $D$.
By the iterated form of Cauchy's formula, the function $z^{\prime} \rightarrow f\left(z^{\prime}, z_{n}\right)$ has a Taylor expansion in $z^{\prime}$

$$
f\left(z^{\prime}, z\right)=\sum_{\alpha} c_{\alpha}\left(z_{n}\right) z^{\prime \alpha}
$$

where the coefficients depend upon $z_{n}$, given by the usual formula

$$
c_{\alpha}\left(z_{n}\right)=\frac{\partial^{\alpha}}{\partial z^{\alpha \alpha}} f\left(0, z_{n}\right) / \alpha!
$$

using multi-index notation. Cauchy's integral formula in $z^{\prime}$ for derivatives

$$
\frac{\partial^{\alpha}}{\partial z^{\prime \alpha}} f\left(0, z_{n}\right)=\alpha!\frac{1}{(2 \pi i)^{n-1}} \int_{C_{1}} \ldots \int_{C_{n-1}} \frac{f(\zeta)}{\left(\zeta_{1}-z_{1}\right)^{\alpha_{1}+1} \ldots\left(\zeta_{n-1}-z_{n-1}\right)^{\alpha_{n-1}+1}} d \zeta_{1} \ldots d \zeta_{n-1}
$$

shows that $c_{\alpha}\left(z_{n}\right)$ is analytic in $z_{n}$, again by expanding convergent geometric series and their derivatives, and interchanging summation and integration.

Fix $0<r_{1}<r_{2}<r$ and fix $z_{n}$ with $\left|z_{n}\right|<r$. Then

$$
\left|c_{\alpha}\left(z_{n}\right)\right| \cdot r_{2}^{|\alpha|} \rightarrow 0
$$

as $|\alpha| \rightarrow \infty$, by the convergence of the power series. Let $B$ be a bound for $|f|$ on the smaller polydisk $E$. Then on that smaller polydisk the Cauchy integral formula for the derivative gives

$$
\left|c_{\alpha}\left(z_{n}\right)\right| \leq B / \varepsilon^{|\alpha|}
$$

Therefore, the subharmonic functions

$$
u_{\alpha}\left(z_{n}\right)=\frac{1}{|\alpha|} \log \left|c_{\alpha}\left(z_{n}\right)\right|
$$

are uniformly bounded from above for $\left|z_{n}\right|<r$. And the property $\left|c_{\alpha}\left(z_{n}\right)\right| \cdot r_{2}^{|\alpha|} \rightarrow 0$ shows that for fixed $z_{n}$ $\log \left(1 / r_{2}\right)$ is an upper bound for these subharmonic functions as $|\alpha| \rightarrow \infty$. Thus, Hartogs' lemma (recalled below) on subharmonic functions implies that for large $|\alpha|$, uniformly in $\left|z_{n}\right|<r_{1}$

$$
\frac{1}{|\alpha|} \log \left|c_{\alpha}\left(z_{n}\right)\right| \leq \log \left(1 / r_{1}\right)
$$

Thus, for large $|\alpha|$

$$
\left|c_{\alpha}\left(z_{n}\right)\right| \cdot r_{1}^{|\alpha|} \leq 1
$$

uniformly in $\left|z_{n}\right|<r_{1}$. Therefore, since the summands $c_{\alpha}\left(z_{n}\right) z^{\prime \alpha}$ are analytic, the series

$$
f\left(z^{\prime}, z\right)=\sum_{\alpha} c_{\alpha}\left(z_{n}\right) z^{\prime \alpha}
$$

converges to a function analytic in the polydisk $D$.
Thus, in summary, given $z \in U$, choose $r>0$ so that the polydisk of radius $2 r$ centered at $z$ is contained in $U$. The Baire category argument above shows that there is $w$ such that $z$ is inside a polydisk $D$ of radius $r$ centered at $w$, and such that $f$ is holomorphic on some smaller polydisk $E$ inside $D$ (still centered at $w$ ). Finally one uses Hartogs' lemma on subharmonic functions (below) to see that the power series for $f$ on the small polydisk $E$ at $w$ converges on the larger polydisk $D$ at $w$. Since $D$ contains the given point $z, f$ is analytic on a neighborhood of $z$. Thus, $f$ is analytic throughout $U$.
[15.C.5] Lemma: (Hartogs) Let $u_{j}$ be a sequence of real-valued subharmonic functions in an open set $U$ in $\mathbb{C}$. Suppose that the functions are uniformly bounded from above, and that

$$
\limsup _{k} u_{k}(z) \leq C
$$

for every $z \in U$. Then, given $\varepsilon>0$ and compact $K \subset U$ there exists $k_{o}$ such that for $z \in K$ and $k \geq k_{o}$

$$
u_{k}(z) \leq C+\varepsilon
$$

Proof: Without loss of generality, replacing $U$ by an open subset with compact closure contained inside $U$, we may suppose that the functions $u_{k}$ are uniformly bounded in $U$, for example $u_{k}(z) \leq 0$ for all $z \in U$. Let $r>0$ be small enough so that the distance from $K$ to every point of the complement of $U$ is more than $3 r$. Using the proposition below characterizing subharmonic functions, we have, for every $z \in K$,

$$
\pi r^{2} u_{k}(z) \leq \int_{|z-\zeta|<r} u_{k}(\zeta) d \zeta
$$

By Fatou's lemma, the $\lim$ sup of the right hand side is at most $\pi r^{2} C$ as $k \rightarrow \infty$. Thus, for every $z \in K$ there is $k_{o}$ such that for $k \geq k_{o}$

$$
\int_{|z-\zeta|<r} u_{k}(\zeta) d \zeta \leq \pi r^{2}(C+\varepsilon / 2)
$$

Since $u_{k}(z) \leq 0$, for $|z-w|<\delta<r$

$$
\pi(r+\delta)^{2} u_{k}(w) \leq \int_{|\zeta-w|<r+\delta \mid} u_{k}(\zeta) d \zeta \leq \int_{|\zeta-z| \leq r} u_{k}(\zeta) d \zeta
$$

Thus, for $\delta>0$ sufficiently small, for $k \geq k_{o}$ and $|w-z|<\delta$,

$$
u_{k}(w)<C+\varepsilon
$$

Since $K$ is compact the lemma follows.
For convenience, recall the following basic property of subharmonic functions.
[15.C.6] Proposition: For a real-valued subharmonic function $u$ bounded above on an open set $U$, for every positive measure $\mu$ on $[0, \delta]$, and for $z \in U$ of distance more than $\delta$ from the complement of $U$,

$$
u(z) \cdot 2 \pi \cdot \int d \mu \leq \int_{0}^{2 \pi} \int u\left(z+r e^{i \theta}\right) d \theta d \mu(r)
$$

Proof: The definition of a function $u$ being subharmonic on an open set $\Omega$ is that $u$ is upper semicontinuous (that is, $\{z \in \Omega: u(z)<c\}$ is open for every constant $c$ ), and for every compact $K \subset \Omega$, for every continuous function $h$ on $K$ harmonic on $K$ and $h(\beta) \geq u(\beta)$ for $\beta$ on the boundary of $K, u(z) \leq h(z)$ throughout $K$. The condition may be vacuous unless $u$ is assumed bounded from above.

Let $z \in U$ be distance more than $\delta$ away from the complement of $U$, and fix $r$ with $0<r \leq \delta$. Let $D$ be the closed disk of radius $r$ about $z$. Since $r \leq \delta, D \subset U$. For a trigonometric polynomial

$$
g(\theta)=\sum_{k} c_{k} e^{i \theta}
$$

with real coefficients $c_{k}$ with $u\left(z+r e^{i \theta}\right) \leq g(\theta)$, the polynomial

$$
G(\zeta)=c_{0}+\sum_{k>0}\left(c_{k}+c_{-k}\right) \frac{(\zeta-z)^{k}}{r^{k}}
$$

has real part $\operatorname{Re} G$ which is an upper bound for $u$ on the boundary of the disk $D$. Thus, $u \leq \operatorname{Re} G$ on $D$ by the subharmonicness of $u$, and in particular at the center of $D$, at $z$,

$$
u(z) \leq c_{o}+\frac{1}{2 \pi} \int_{0}^{2 \pi} g(\theta) d \theta
$$

Then for an arbitrary continuous real-valued function $h$ on the boundary of $D$ and with $u\left(z+r e^{i \theta}\right) \leq h(\theta)$, (by Weierstrass approximation, for example) given $\varepsilon>0$ we can find a trigonometric polynomial $g$ so that $\sup |g(\theta)-h(\theta)|<\varepsilon$. Thus, for every $\varepsilon>0$,

$$
u(z) \leq c_{o}+\frac{1}{2 \pi} \int_{0}^{2 \pi} h(\theta) d \theta+\varepsilon
$$

Thus, the latter inequality must hold with $\varepsilon=0$, for continuous $h$. Since the integral of an uppersemicontinuous function is the infimum of the integrals of continuous functions dominating it, we have the same inequality with $u$ in place of $h$. Integration with respect to the radius $r$ gives the result. ///

In fact, suppose that for every $\delta>0$ and for every $z$ at distance more than $\delta$ from the complement of $U$ there exists a positive measure $\mu$ on $[0, \delta]$ with support not just $\{0\}$ and

$$
u(z) \cdot 2 \pi \cdot \int d \mu \leq \int_{0}^{2 \pi} \int u\left(z+r e^{i \theta}\right) d \theta d \mu(r)
$$

Then $u$ is subharmonic. To see this, let $K$ be a compact subset of $U, h$ a continuous function on $K$ which is harmonic in the interior of $K$ and such that $u \leq g$ on the boundary of $K$. If the supremum of $u-h$ over $K$ is strictly positive, the upper semicontinuity of $u-h$ implies that $u-h$ attains its sup $S$ on a non-empty compact subset $M$ of the interior of $K$. Let $z_{o}$ be a point of $M$ closest to the boundary of $K$. If the distance is greater than $\delta$, then every circle $\left|z-z_{o}\right|=r$ with $0<r \leq \delta$ contains a non-empty arc of points where $u-h<S$. Then

$$
\left.\int(u-h)\left(z_{o}+r e^{i \theta}\right)\right) d \theta d \mu(r)<S \cdot 2 \pi \cdot \int d \mu(r)=(u-h)\left(z_{o}\right) \cdot 2 \pi \cdot \int d \mu(r)
$$

when $\mu$ is a measure not supported just at $\{0\}$. The mean value property for harmonic functions gives

$$
\left.\int h\left(z_{o}+r e^{i \theta}\right)\right) d \theta d \mu(r)=h\left(z_{o}\right) \cdot 2 \pi \cdot \int d \mu(r)
$$

Thus,

$$
\left.\int u\left(z_{o}+r e^{i \theta}\right)\right) d \theta d \mu(r)<u\left(z_{o}\right) \cdot 2 \pi \cdot \int d \mu(r)
$$

contradicting the hypothesis. Thus, $\sup _{K}(u-h) \leq 0$, which proves that $u$ is subharmonic.

## 16. Asymptotic expansions

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Appendix A: manipulation of asymptotic expansions
Appendix B: ordinary points
The simplest notion of asymptotic of $f(s)$ as $s$ goes to $+\infty$ on $\mathbb{R}$ is a simpler function $F(s)$ such that $\lim _{s} f(s) / F(s)=1$, written $f \sim F$. One might require an error estimate, for example,

$$
f \sim F \Longleftrightarrow f(s)=F(s) \cdot\left(1+O\left(\frac{1}{|s|}\right)\right)
$$

A more precise form is to say that

$$
f(s) \sim f_{0}(s) \cdot\left(\frac{c_{0}}{s^{\alpha}}+\frac{c_{1}}{s^{\alpha+1}}+\frac{c_{2}}{s^{\alpha+2}}+\ldots\right)
$$

(with any auxiliary function $f_{0}$ ) is an asymptotic expansion for $f$ when

$$
f=f_{0}(s) \cdot\left(\frac{c_{0}}{s^{\alpha}}+\frac{c_{1}}{s^{\alpha+1}}+\ldots+\frac{c_{n}}{s^{\alpha+n}}+O\left(\frac{1}{|s|^{\alpha+n+1}}\right)\right)
$$

The exposition is revisionist: Laplace's method is proven by reducing it to Watson's lemma.

### 16.1 Heuristic for Stirling's asymptotic

First we give a heuristic and mnemonic for the main term of the Laplace-Stirling asymptotic, namely

$$
\Gamma(s) \sim e^{-s} \cdot s^{s-\frac{1}{2}} \cdot \sqrt{2 \pi}
$$

Using Euler's integral,

$$
s \cdot \Gamma(s)=\Gamma(s+1)=\int_{0}^{\infty} e^{-u} u^{s+1} \frac{d u}{u}=\int_{0}^{\infty} e^{-u} u^{s} d u=\int_{0}^{\infty} e^{-u+s \log u} d u
$$

The trick is to replace the exponent $-u+s \log u$ by the quadratic polynomial in $u$ best approximating it near its maximum, and evaluate the resulting integral. This replacement is justified via Watson's lemma and Laplace's method, below, but the heuristic is simpler than the justification. The exponent takes its maximum where its derivative vanishes, at the unique solution $u_{o}=s$ of

$$
-1+\frac{s}{u}=0
$$

The second derivative in $u$ of the exponent is $-s / u^{2}$, which takes value $-1 / s$ at $u_{o}=s$. Thus, near $u_{o}=s$, the quadratic Taylor-Maclaurin polynomial in $u$ approximating the exponent is

$$
-s+s \log s-\frac{1}{2!s} \cdot(u-s)^{2}
$$

We imagine that

$$
s \cdot \Gamma(s) \sim \int_{0}^{\infty} e^{-s+s \log s-\frac{1}{2 s} \cdot(u-s)^{2}} d u=e^{-s} \cdot s^{s} \cdot \int_{-\infty}^{\infty} e^{-\frac{1}{2 s} \cdot(u-s)^{2}} d u
$$

The latter integral is taken over the whole real line. Evaluation of the integral over the whole line, and simple estimates on the integral over $(-\infty, 0]$, show that the integral over $(-\infty, 0]$ is of a lower order of magnitude than the whole. Thus, the leading term of the asymptotics of the integral over the whole line is the same than the integral from 0 to $+\infty$. To simplify the remaining integral, replace $u$ by su and cancel a factor of $s$ from both sides,

$$
\Gamma(s) \sim e^{-s} \cdot s^{s} \cdot \int_{-\infty}^{\infty} e^{-s(u-1)^{2} / 2} d u
$$

Replace $u$ by $u+1$, and $u$ by $u \cdot \sqrt{2 \pi / s}$, obtaining

$$
\int_{-\infty}^{\infty} e^{-s(u-1)^{2} / 2} d u=\int_{-\infty}^{\infty} e^{-s u^{2} / 2} d u=\frac{\sqrt{2 \pi}}{\sqrt{s}} \int_{-\infty}^{\infty} e^{-\pi u^{2}} d u=\frac{\sqrt{2 \pi}}{\sqrt{s}}
$$

and

$$
\Gamma(s) \sim e^{-s} \cdot s^{s-\frac{1}{2}} \cdot \sqrt{2 \pi}
$$

It is striking that this heuristic can be made rigorous, as below.

### 16.2 Watson's lemma

Watson's lemma gives an asymptotic expansion for certain Laplace transforms, valid in half-planes in $\mathbb{C}$. For example, let $h$ be a smooth function on $(0,+\infty)$ all whose derivatives are of polynomial growth, and expressible for small $x>0$ as

$$
h(x)=x^{\alpha} \cdot g(x)
$$

for some $\alpha \in \mathbb{C}$, where $g(x)$ is differentiable on $\mathbb{R}$ near 0 . Thus, $h(x)$ has an expression

$$
h(x)=x^{\alpha} \cdot \sum_{n=0}^{\infty} c_{n} x^{n} \quad(\text { for } 0<x \text { sufficiently small) }
$$

Then there is an asymptotic expansion of the Laplace transform of $h$,

$$
\int_{0}^{\infty} e^{-x s} h(x) \frac{d x}{x} \sim \frac{\Gamma(\alpha) c_{0}}{s^{\alpha}}+\frac{\Gamma(\alpha+1) c_{1}}{s^{\alpha+1}}+\frac{\Gamma(\alpha+2) c_{2}}{s^{\alpha+2}}+\ldots \quad(\text { for } \operatorname{Re}(s)>0)
$$

A simple corollary of the error estimates given below is that, letting $\operatorname{Re}(\alpha)+1-\varepsilon$ be the greatest integer less than or equal $\operatorname{Re}(\alpha)+1$,

$$
\int_{0}^{\infty} e^{-x s} h(x) \frac{d x}{x}=\int_{0}^{\infty} e^{-x s} x^{\alpha} g(x) \frac{d x}{x}=\frac{\Gamma(\alpha) g(0)}{s^{\alpha}}+O\left(\frac{1}{|s|^{\operatorname{Re}(\alpha)+1-\varepsilon}}\right)
$$

Since

$$
\operatorname{Re}(\alpha)+1-\varepsilon>\operatorname{Re}(\alpha)
$$

the error term is of strictly smaller order of magnitude in $s$.

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The idea of the proof is straightforward: the expansion is obtained from

$$
\int_{0}^{\infty} e^{-x s} h(x) \frac{d x}{x}=\int_{0}^{\infty} e^{-x s} x^{\alpha}\left(c_{0}+\ldots+c_{n} x^{n}\right) \frac{d x}{x}+\int_{0}^{\infty} e^{-x s} x^{\alpha}\left(g(x)-\left(c_{0}+\ldots+c_{n} x^{n}\right)\right) \frac{d x}{x}
$$

The first integral gives the asymptotic expansion, and for $\operatorname{Re}(s)>0$ the second integral can be integrated by parts essentially $\operatorname{Re}(\alpha)+n$ times and trivially bounded to give a $O\left(1 / s^{\alpha+n-\varepsilon}\right)$ error term for some small $\varepsilon \geq 0$. Note that for the integration by parts the denominator $x$ in the measure must be moved into the integrand proper, accounting for a slight reduction of the order of vanishing of the integrand at 0 .

To understand the error, let $\varepsilon \geq 0$ be the smallest such that

$$
N=\operatorname{Re}(\alpha)+n-\varepsilon \in \mathbb{Z}
$$

The subtraction of the initial polynomial and re-allocation of the $1 / x$ from the measure makes $x^{\alpha-1}\left(g(x)-\left(c_{0}+\ldots+c_{n} x^{n}\right)\right.$ vanish to order $N$ at 0 . This, with the exponential $e^{-s x}$ and the presumed polynomial growth of $h$ and its derivatives, allows integration by parts $N$ times without boundary terms, giving

$$
\begin{gathered}
\int_{0}^{\infty} e^{-x s} h(x) d x=\frac{\Gamma(\alpha) c_{0}}{s^{\alpha}}+\frac{\Gamma(\alpha+1) c_{1}}{s^{\alpha+1}}+\ldots+\frac{\Gamma(\alpha+n) c_{n}}{s^{\alpha+n}} \\
\quad+\frac{1}{s^{N}} \int_{0}^{\infty} e^{-s x}\left(\frac{\partial}{\partial x}\right)^{N}\left(x^{\alpha} \cdot\left(g(x)-\left(c_{0}+\ldots+c_{n} x^{n}\right)\right)\right) d x
\end{gathered}
$$

The last error-like term is $O\left(s^{-[\operatorname{Re}(\alpha)+n-\varepsilon]}\right)$. That is, computing in this fashion, the error term swallows up the last term in the asymptotic expansion. Visibly, this argument applies to more general sorts of expansions near 0 .

### 16.3 Watson's lemma illustrated on the Beta function

Here is an important example of an asymptotic result non-trivial to derive from Stirling's formula for $\Gamma(s)$, but easy to obtain from Watson's lemma. Euler's beta integral is

$$
B(s, a)=\int_{0}^{1} x^{s-1}(1-x)^{a-1} d x=\frac{\Gamma(s) \Gamma(a)}{\Gamma(s+a)}
$$

We recall how to express Beta in terms of Gamma: with $x=u /(u+1)$ in the beta integral,

$$
\begin{gathered}
B(s, a)=\int_{0}^{\infty} u^{s-1}(u+1)^{-(s-1)-(a-1)-2} d u=\int_{0}^{\infty} u^{s-1}(u+1)^{-s-a} d u \\
=\frac{1}{\Gamma(s+a)} \int_{0}^{\infty} \int_{0}^{\infty} u^{s} e^{-v(u+1)} v^{s+a} \frac{d v}{v} \frac{d u}{u}
\end{gathered}
$$

using $\int_{0}^{\infty} e^{-v y} v^{b} d v / v=\Gamma(b) / y^{b}$. Replacing $u$ by $u / v$ gives $B(s, a)=\Gamma(s) \Gamma(a) / \Gamma(s+a)$.
Fix $a$ with $\operatorname{Re}(a)>0$, and consider this integral as a function of $s$. Letting $x=e^{-u}$ gives an integrand fitting Watson's lemma,

$$
\begin{gathered}
B(s, a)=\int_{0}^{\infty} e^{-s u}\left(1-e^{-u}\right)^{a-1} d u=\int_{0}^{\infty} e^{-s u}\left(u-\frac{u^{2}}{2!}+\ldots\right)^{a-1} d u \\
=\int_{0}^{\infty} e^{-s u} u^{a} \cdot\left(1-\frac{u}{2!}+\ldots\right)^{a-1} \frac{d u}{u} \sim \frac{\Gamma(a)}{s^{a}}
\end{gathered}
$$

taking just the first term in an asymptotic expansion, using Watson's lemma. Thus,

$$
\frac{\Gamma(s) \Gamma(a)}{\Gamma(s+a)} \sim \frac{\Gamma(a)}{s^{a}}
$$

giving

$$
\left.\frac{\Gamma(s)}{\Gamma(s+a)} \sim \frac{1}{s^{a}} \quad \text { (for } a \text { fixed }\right)
$$

### 16.4 Simple form of Laplace's method

A simple version of Laplace's method obtains asymptotics in $s$ for certain integrals of the form

$$
\int_{0}^{\infty} e^{-s \cdot f(u)} d u
$$

with $f$ real-valued. The idea is that the minimum values of $f(u)$ should dominate, and the leading term of the asymptotics should be

$$
\int_{0}^{\infty} e^{-s \cdot f(u)} d u \sim e^{-s f\left(u_{o}\right)} \cdot \frac{\sqrt{2 \pi}}{\sqrt{f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}} \quad(\text { for }|s| \rightarrow \infty, \text { with } \operatorname{Re}(s) \geq \delta>0)
$$

To reduce this to Watson's lemma, break the integral at points where the derivative $f^{\prime}$ changes sign, and change variables to convert each fragment to a Watson-lemma integral. For Watson's lemma to be legitimately applied, we will find that $f$ must be smooth with all derivatives of at most polynomial growth and at most polynomial decay, as $u \rightarrow+\infty$. For simplicity assume that there is exactly one point $u_{o}$ at which $f^{\prime}\left(u_{o}\right)=0$, and that $f^{\prime \prime}\left(u_{o}\right)>0$. Further, assume that $f(u)$ goes to $+\infty$ at $0^{+}$and at $+\infty$. Since $f^{\prime}(u)>0$ for $u>u_{o}$ and $f^{\prime}(u)<0$ for $0<u<u_{o}$, on each of these two intervals there is a smooth square root $\sqrt{f(u)-f\left(u_{o}\right)}$ and there are smooth functions $F, G$ such that

$$
\left\{\begin{array}{lll}
F\left(\sqrt{f(u)-f\left(u_{o}\right)}\right) & =u & \left(\text { for } u_{o}<u<+\infty\right) \\
G\left(\sqrt{f(u)-f\left(u_{o}\right)}\right) & =u & \\
\left(\text { for } 0<u<u_{o}\right)
\end{array}\right.
$$

Then

$$
\begin{aligned}
\int_{0}^{\infty} e^{-s f(u)} d u & =e^{-s f\left(u_{o}\right)} \int_{0}^{u_{o}} e^{-s\left(f(u)-f\left(u_{o}\right)\right)} d u+e^{-s f\left(u_{o}\right)} \int_{u_{o}}^{\infty} e^{-s\left(f(u)-f\left(u_{o}\right)\right)} d u \\
& =e^{-s f\left(u_{o}\right)}\left(\int_{0}^{\infty} e^{-s x^{2}} F^{\prime}(x) d x+\int_{0}^{\infty} e^{-s x^{2}} G^{\prime}(x) d x\right)
\end{aligned}
$$

by letting $x=\sqrt{f(u)-f\left(u_{o}\right)}$ in the two intervals. In both integrals, replacing $x$ by $\sqrt{x}$ gives Watson's-lemma integrals

$$
\int_{0}^{\infty} e^{-s f(u)} d u=e^{-s f\left(u_{o}\right)}\left(\int_{0}^{\infty} e^{-s x} \frac{1}{2} x^{1 / 2} F^{\prime}(\sqrt{x}) \frac{d x}{x}+\int_{0}^{\infty} e^{-s x} \frac{1}{2} x^{1 / 2} G^{\prime}(\sqrt{x}) \frac{d x}{x}\right)
$$

At this point the needed conditions on $F$, hence, on $f$, become clear: since $F$ must be smooth with all derivatives of at most polynomial growth, direct chain-rule computations show that it suffices that no derivative of $f$ increases or decreases faster than polynomially as $u \rightarrow+\infty$. The assumptions $f^{\prime}\left(u_{o}\right)=0$ and $f^{\prime \prime}\left(u_{o}\right)>0$ assure that $F$ has a Taylor series expansion near 0 , giving a suitable expansion

$$
\frac{1}{2} x^{1 / 2} F^{\prime}(x)=\frac{1}{2} F^{\prime}(0) x^{1 / 2}+\frac{\frac{1}{2} F^{(2)}(0)}{1!} x^{3 / 2}+\frac{\frac{1}{2} F^{(3)}(0)}{2!} x^{5 / 2}+\frac{\frac{1}{2} F^{(4)}(0)}{3!} x^{7 / 2}+\ldots \quad(\text { small } x>0)
$$

From this, the main term of the Watson's lemma asymptotics for the integral involving $F$ would be

$$
\int_{0}^{\infty} e^{-s x} \frac{1}{2} x^{1 / 2} F^{\prime}(\sqrt{x}) \frac{d x}{x} \sim \frac{\Gamma\left(\frac{1}{2}\right) F^{\prime}(0)}{2} \cdot \frac{1}{\sqrt{s}}
$$

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To determine $F^{\prime}(0)$, or any higher coefficients, from $F(x)=u$, we have $F^{\prime}(x) \cdot \frac{d x}{d u}=1$. Since

$$
x=\sqrt{f(u)-f\left(u_{o}\right)}=\sqrt{\left(u-u_{o}\right)^{2} \cdot \frac{f^{\prime \prime}\left(u_{o}\right)}{2!}+\ldots}=\sqrt{\frac{f^{\prime \prime}\left(u_{o}\right)}{2}} \cdot\left(\left(u-u_{o}\right)+\ldots\right)
$$

the derivative is

$$
\frac{d x}{d u}=\sqrt{\frac{f^{\prime \prime}\left(u_{o}\right)}{2}} \cdot\left(1+O\left(u-u_{o}\right)\right)
$$

Thus,

$$
F^{\prime}(x)=\frac{1}{\frac{d x}{d u}}=\sqrt{\frac{2}{f^{\prime \prime}\left(u_{o}\right)}} \cdot\left(1+O\left(u-u_{o}\right)\right)
$$

which allows evaluation at $x=0$, namely

$$
F^{\prime}(0)=\sqrt{\frac{2}{f^{\prime \prime}\left(u_{o}\right)}}
$$

The same argument applied to $G$ gives $G^{\prime}(0)=F^{\prime}(0)$. Thus,

$$
\int_{0}^{\infty} e^{-s f(u)} d u \sim e^{-s f\left(u_{o}\right)} \cdot \frac{\Gamma\left(\frac{1}{2}\right) \cdot 2 \cdot \sqrt{\frac{2}{f^{\prime \prime}\left(u_{o}\right)}}}{2 \sqrt{s}}=e^{-s f\left(u_{o}\right)} \cdot \frac{\sqrt{2 \pi}}{\sqrt{f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}}
$$

Last, we verify that this outcome is what would be obtained by replacing $f(u)$ by its quadratic approximation

$$
f\left(u_{o}\right)+\frac{f^{\prime \prime}(0)}{2!} \cdot\left(u-u_{o}\right)^{2}
$$

in the exponent in the original integral, integrated over the whole line. The latter would be

$$
\begin{gathered}
\int_{-\infty}^{\infty} e^{s \cdot\left(f\left(u_{o}\right)+\frac{1}{2} f^{\prime \prime}\left(u_{o}\right)\left(u-u_{o}\right)^{2}\right)} d u=e^{s f\left(u_{o}\right)} \int_{-\infty}^{\infty} e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right)\left(u-u_{o}\right)^{2}} d u \\
=e^{s f\left(u_{o}\right)} \int_{-\infty}^{\infty} e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right) u^{2}} d u=e^{s f\left(u_{o}\right)} \cdot \frac{\sqrt{\pi}}{\sqrt{\frac{1}{2} f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}}=e^{s f\left(u_{o}\right)} \cdot \frac{\sqrt{2 \pi}}{\sqrt{f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}}
\end{gathered}
$$

This does indeed agree. Last, verify that the integral of the exponentiated quadratic approximation over $(-\infty, 0]$ is of a lower order of magnitude. Indeed, for $u \leq 0$ and $u_{o}>0$ we have $\left(u-u_{o}\right)^{2} \geq u^{2}+u_{o}^{2}$, and $f^{\prime \prime}\left(u_{o}\right)<0$ by assumption, so

$$
\begin{aligned}
& e^{s f\left(u_{o}\right)} \int_{-\infty}^{0} e^{s \cdot\left(\frac{1}{2} f^{\prime \prime}\left(u_{o}\right)\left(u-u_{o}\right)^{2}\right)} d u \leq e^{s f\left(u_{o}\right)} \cdot e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right) \cdot u_{o}^{2}} \int_{-\infty}^{0} e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right) u^{2}} d u \\
\leq & e^{s f\left(u_{o}\right)} \cdot e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right) \cdot u_{o}^{2}} \int_{-\infty}^{\infty} e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right) u^{2}} d u=e^{s f\left(u_{o}\right)} \cdot e^{s \cdot \frac{1}{2} f^{\prime \prime}\left(u_{o}\right) \cdot u_{o}^{2}} \cdot \frac{\sqrt{2 \pi}}{\sqrt{f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}}
\end{aligned}
$$

Thus, the integral over $(-\infty, 0]$ has an additional exponential decay by comparison to the integral over the whole line, so the leading-term of the asymptotics of the integral from 0 to $+\infty$ is the same as those of the integral from $-\infty$ to $+\infty$.

The case of $\Gamma(s)$ can be converted to this situation as follows. For real $s>0$, in the integral

$$
s \cdot \Gamma(s)=\Gamma(s+1)=\int_{0}^{\infty} e^{-u} u^{s} d u=\int_{0}^{\infty} e^{-u+s \log u} d u
$$

can replace $u$ by $s u$, to put the integral into the desired form

$$
s \cdot \Gamma(s)=\int_{0}^{\infty} e^{-s u+s \log u+s \log s} s d u=s \cdot e^{s \log s} \int_{0}^{\infty} e^{-s(u+\log u)} d u
$$

For complex $s$ with $\operatorname{Re}(s)>0$, both $s \cdot \Gamma(s)$ and the integral $s \cdot e^{s \log s} \int_{0}^{\infty} e^{-s(u+\log u)} d u$ are holomorphic in $s$, and they agree for real $s$. Thus, by the identity principle, they are equal for $\operatorname{Re}(s)>0$.

### 16.5 Laplace's method illustrated on Bessel functions

Consider the standard integral

$$
K_{\nu}(y)=\frac{1}{2} \int_{0}^{\infty} e^{\left(u+\frac{1}{u}\right) y / 2} u^{\nu} \frac{d u}{u}
$$

The function $K_{\nu}$ is variously called a Bessel function of imaginary argument or MacDonald's function or modified Bessel function of third kind. Being interested mainly in the case that the parameter $\nu$ here is purely imaginary, in the text we replace $\nu$ by $i \nu$. The leading factor of $\sqrt{y}$ arises in applications, where such an integral appears as a Whittaker function.

$$
f(y)=\sqrt{y} \int_{0}^{\infty} e^{-\left(u+\frac{1}{u}\right) y} u^{i \nu} \frac{d u}{u}
$$

The exponent $-\left(u+\frac{1}{u}\right) y$ is of the desired form, with the earlier $s$ replaced by $y$, but the $u^{i \nu}$ in the integrand does not fit into the simpler Laplace' method. Thus, consider integrals

$$
\int_{0}^{\infty} e^{-s f(u)} g(u) d u
$$

where $f$ is real-valued, but $g$ may be complex-valued. The idea still is that the minimum values of $f(u)$ should dominate, and the leading term of the asymptotics should be (assuming a unique minimum at $u_{o}$ )

$$
\int_{0}^{\infty} e^{-s \cdot f(u)} g(u) d u \sim e^{-s f\left(u_{o}\right)} \cdot \frac{\sqrt{2 \pi} \cdot g\left(u_{o}\right)}{\sqrt{f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}}
$$

As in the simpler case, reduce this to Watson's lemma by breaking the integral where $f^{\prime}$ changes sign, and change variables to convert each fragment to a Watson-lemma integral. Thus, the first part of the story is much as the simple case of Laplace's method. As in the simple case of Laplace's method, the course of the argument reveals conditions on $f$ and $g$.

For simplicity assume that there is exactly one point $u_{o}$ at which $f^{\prime}\left(u_{o}\right)=0$, that $f^{\prime \prime}\left(u_{o}\right)>0$, and that $f(u)$ goes to $+\infty$ at $0^{+}$and at $+\infty$. Thus, on these intervals there are smooth square roots $\sqrt{f(u)-f\left(u_{o}\right)}$ and smooth functions $F, G$ such that

$$
\left\{\begin{array}{lll}
F\left(\sqrt{f(u)-f\left(u_{o}\right)}\right) & =u & \\
\left(\text { for } u_{o}<u<+\infty\right) \\
G\left(\sqrt{f(u)-f\left(u_{o}\right)}\right) & =u & \\
\left(\text { for } 0<u<u_{o}\right)
\end{array}\right.
$$

Then, letting $x=\sqrt{f(u)-f\left(u_{o}\right)}$ in each of the two intervals, so that $F(x)=u$ and $G(x)=u$, respectively,

$$
\begin{gathered}
\int_{0}^{\infty} e^{-s f(u)} g(u) d u=e^{-s f\left(u_{o}\right)} \int_{0}^{u_{o}} e^{-s\left(f(u)-f\left(u_{o}\right)\right)} g(u) d u+e^{-s f\left(u_{o}\right)} \int_{u_{o}}^{\infty} e^{-s\left(f(u)-f\left(u_{o}\right)\right)} g(u) d u \\
=e^{-s f\left(u_{o}\right)}\left(\int_{0}^{\infty} e^{-s x^{2}} g(F(x)) F^{\prime}(x) d x+\int_{0}^{\infty} e^{-s x^{2}} g(G(x)) G^{\prime}(x) d x\right)
\end{gathered}
$$

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In both integrals, replacing $x$ by $\sqrt{x}$ gives Watson's-lemma integrals

$$
\int_{0}^{\infty} e^{-s f(u)} g(u) d u=e^{-s f\left(u_{o}\right)}\left(\int_{0}^{\infty} e^{-s x} \frac{\sqrt{x}}{2} g(F(\sqrt{x})) F^{\prime}(\sqrt{x}) \frac{d x}{x}+\int_{0}^{\infty} e^{-s x} \frac{\sqrt{x}}{2} g(G(\sqrt{x})) G^{\prime}(\sqrt{x}) \frac{d x}{x}\right)
$$

The assumptions $f^{\prime}\left(u_{o}\right)=0$ and $f^{\prime \prime}\left(u_{o}\right)>0$ assure that $F$ has a Taylor series expansion near 0 , which gives an expansion

$$
\frac{\sqrt{x}}{2} g(F(\sqrt{x})) F^{\prime}(\sqrt{x})=x^{1 / 2} \cdot\left(\frac{g(F(0)) F^{\prime}(0)}{2}+\ldots\right) \quad(\operatorname{small} x>0)
$$

From this, the main term of the Watson's lemma asymptotics for the integral involving $F$ would be

$$
\int_{0}^{\infty} e^{-s x} \frac{\sqrt{x}}{2} g(F(\sqrt{x})) F^{\prime}(\sqrt{x}) \frac{d x}{x} \sim \frac{\Gamma\left(\frac{1}{2}\right) g(F(0)) F^{\prime}(0)}{2} \cdot \frac{1}{\sqrt{s}}
$$

Note that $F(0)=u_{o}$. Determine $F^{\prime}(0)$ from $F(x)=u$. First, $F^{\prime}(x) \cdot \frac{d x}{d u}=1$. Since

$$
x=\sqrt{f(u)-f\left(u_{o}\right)}=\sqrt{\left(u-u_{o}\right)^{2} \cdot \frac{f^{\prime \prime}\left(u_{o}\right)}{2!}+\ldots}=\sqrt{\frac{f^{\prime \prime}\left(u_{o}\right)}{2}} \cdot\left(\left(u-u_{o}\right)+\ldots\right)
$$

the derivative is

$$
\frac{d x}{d u}=\sqrt{\frac{f^{\prime \prime}\left(u_{o}\right)}{2}} \cdot\left(1+O\left(u-u_{o}\right)\right)
$$

Thus,

$$
F^{\prime}(x)=\frac{1}{\frac{d x}{d u}}=\sqrt{\frac{2}{f^{\prime \prime}\left(u_{o}\right)}} \cdot\left(1+O\left(u-u_{o}\right)\right)
$$

which allows evaluation at $x=0$, namely

$$
F^{\prime}(0)=\sqrt{\frac{2}{f^{\prime \prime}\left(u_{o}\right)}}
$$

The same argument applied to $G$ gives $G^{\prime}(0)=F^{\prime}(0)$. Thus,

$$
\int_{0}^{\infty} e^{-s f(u)} d u \sim e^{-s f\left(u_{o}\right)} \cdot \frac{\Gamma\left(\frac{1}{2}\right) \cdot 2 \cdot g\left(u_{o}\right) \cdot \sqrt{\frac{2}{f^{\prime \prime}\left(u_{o}\right)}}}{2 \sqrt{s}}=e^{-s f\left(u_{o}\right)} \cdot \frac{\sqrt{2 \pi} \cdot g\left(u_{o}\right)}{\sqrt{f^{\prime \prime}\left(u_{o}\right)}} \cdot \frac{1}{\sqrt{s}}
$$

Returning to

$$
f(y)=\sqrt{y} \int_{0}^{\infty} e^{-\left(u+\frac{1}{u}\right) y} u^{i \nu} \frac{d u}{u}
$$

we have critical point $u_{o}=1$, and $f\left(u_{o}\right)=2$ and $f^{\prime \prime}\left(u_{o}\right)=2$. Applying the just-derived asymptotic,

$$
f(y) \sim \sqrt{y} \cdot\left(e^{-2 y} \cdot \frac{\sqrt{2 \pi} \cdot 1^{i \nu}}{\sqrt{2}} \cdot \frac{1}{\sqrt{y}}\right)=\sqrt{\pi} \cdot e^{-2 y} \quad(\text { as } y \rightarrow+\infty)
$$

Even though the exponent is plausible, it would be easy to lose track of the power of $y$, which might matter. Also, the leading constant does not depend upon the index $\nu$.

### 16.6 Regular singular points heuristic: freezing coefficients

Differential equations

$$
x^{2} u^{\prime \prime}+b x u^{\prime}+c u=0 \quad \text { (with constants } b, c \text { ) }
$$

have easy-to-understand solutions on $(0,+\infty)$ : linear combinations of $x^{\alpha}, x^{\beta}$ for $\alpha, \beta$ solutions of the indicial equation

$$
X(X-1)+b X+c=0
$$

when the roots are distinct. Therefore, it is reasonable to imagine that a differential equation

$$
x^{2} u^{\prime \prime}+x b(x) u^{\prime}+c(x) u=0
$$

with $b, c$ analytic near 0 has solutions asymptotic, as $x \rightarrow 0^{+}$, to solutions of the differential equation $x^{2} u^{\prime \prime}+b(0) x u^{\prime}+c(0) u=0$ obtained by freezing the coefficients $b(x), c(x)$ of the original at $x=0^{+}$. That is, solutions of the variable-coefficient equation should be asymptotic to $x^{\alpha}$ for solutions $\alpha$ to the indicial equation $X(X-1)+b(0) X+c(0)=0$. An equation of that form, with $b, c$ analytic near 0 , is said to have a regular singular point at 0 . Discussion below explains the behavior of solutions to such equations.

We give a useful example from the non-Euclidean geometry on the upper half-plane. Recall the $S L_{2}(\mathbb{R})$ invariant Laplacian on the upper half-plane $\mathfrak{H}$ is

$$
\Delta^{\mathfrak{H}}=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)
$$

[16.6.1] Translation-equivariant eigenfunctions We ask for $\Delta^{\mathfrak{H}}$-eigenfunctions $f(z)$ of the special form

$$
f(x+i y)=e^{2 \pi i x} u(y)
$$

That is, such an eigenfunction is equivariant under translations:

$$
f(z+t)=e^{2 \pi i(x+t)} u(y)=e^{2 \pi i t} \cdot\left(e^{2 \pi i x} u(y)\right)=e^{2 \pi i t} \cdot f(z) \quad(\text { with } t \in \mathbb{R} \text { and } z \in \mathfrak{H})
$$

The eigenfunction condition is the partial differential equation

$$
\left(\Delta^{\mathfrak{H}}-\lambda\right) e^{2 \pi i x} u(y)=0
$$

Since the dependence on $x$ is completely specified, this partial differential equation simplifies to the ordinary differential equation ${ }^{[84]}$

$$
y^{2} u^{\prime \prime}-\left(4 \pi^{2} y^{2}+\lambda\right) u=0
$$

The point $y=0$ is not an ordinary point for this equation, because in the form

$$
u^{\prime \prime}-\left(4 \pi^{2}+\frac{\lambda}{y^{2}}\right) u=0
$$

the coefficient of $u$ has a pole at 0 . But $y=0$ is a regular singular point, because that pole is of order at most 2. Thus, following the idea to freeze $y^{2} u^{\prime \prime}+y b(y) u^{\prime}+c(y)$ to $y^{2} u^{\prime \prime}+y b(0) u^{\prime}+c(0) u$, the indicial equation of the frozen equation is

$$
X(X-1)-\lambda=0
$$

Expressing $\lambda$ as $\lambda=s(s-1)$, the roots of the indicial equation are $s, 1-s$. The frozen equation has distinct solutions $y^{s}$ and $y^{1-s}$ for $s \neq \frac{1}{2}$. Thus, we could hope that solutions would have asymptotics as $y \rightarrow 0^{+}$ beginning

$$
u(y)=A y^{s}(1+O(y))+B y^{1-s}(1+O(y)) \quad\left(\text { as } y \rightarrow 0^{+}\right)
$$

[84] This equation is a type of Bessel equation, with solutions which are $K$-type and $I$-type Bessel functions.

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Indeed, this is the case, as we see below. It seems more difficult to obtain the asymptotics at $0^{+}$from integral representations of solutions of the differential equation.
[16.6.2] Remark: As we discuss below, $y^{2} u^{\prime \prime}-\left(4 \pi^{2} y^{2}+\lambda\right) u=0$ has an irregular singular point at $+\infty$, so other methods are needed to obtain asymptotics for solutions as $y \rightarrow+\infty$.
[16.6.3] Remark: Up to choices of normalizations, the function $u$ above, depending on the spectral parameter $\lambda$ or $s$, is called a Whittaker function or Bessel function, enjoying an enormous literature. Here, wish to have direct access to their properties, as instances of general phenomena.
[16.6.4] An irregular singular point For the translation-equivariant eigenfunctions on $\mathfrak{H}$, we check that $y=+\infty$ is not an ordinary point nor a regular singular point: given

$$
u^{\prime \prime}-\left(4 \pi^{2}+\frac{\lambda}{y^{2}}\right) u=0
$$

again let $u(x)=v(1 / x)$ and put $z=1 / x$, obtaining

$$
\left(z^{4} v^{\prime \prime}+2 z^{3} v^{\prime}\right)-\left(4 \pi^{2}+\lambda z^{2}\right) v=0
$$

or

$$
z^{2} v^{\prime \prime}+2 z v^{\prime}-\left(\frac{4 \pi^{2}}{z^{2}}+\lambda\right) v=0
$$

Since the coefficient of $v$ has a pole at $z=0$, this equation falls outside the present discussion. Instead, a different freezing idea succeeds: letting $y \rightarrow+\infty$ freezes the original equation at $+\infty$, giving a constantcoefficient equation

$$
u^{\prime \prime}-4 \pi^{2} u=0
$$

with easily-understood solutions $e^{ \pm 2 \pi y}$. Happily the solutions to the original equation do have asymptotics with main terms $e^{ \pm 2 \pi y}$. Details and proofs are given below, in a discussion of irregular singular points.

### 16.7 Regular singular points

A homogeneous ordinary differential equation of the form

$$
\left.x^{2} u^{\prime \prime}+x b(x) u^{\prime}+c(x) u=0 \quad \text { (with } b, c \text { analytic near } 0\right)
$$

is said to have a regular singular point at 0 . Similarly,

$$
\left(x-x_{o}\right)^{2} u^{\prime \prime}+\left(x-x_{o}\right) b(x) u^{\prime}+c(x) u=0 \quad \text { (with } b, c \text { analytic near } x_{o} \text { ) }
$$

has a regular singular point at $x_{o}$. Obviously it suffices to treat $x_{o}=0$, and is notationally convenient. The coefficients in an expansion of the form

$$
u(x)=x^{\alpha} \cdot \sum_{n=0}^{\infty} a_{n} x^{n} \quad\left(\text { with } a_{0} \neq 0, \alpha \in \mathbb{C}\right)
$$

are determined recursively, but we see below that this recursion succeeds only when $\alpha$ satisfies the indicial equation

$$
\alpha(\alpha-1)+b(0) \alpha+c(0)=0
$$

Further, when the two roots $\alpha, \alpha^{\prime}$ of the indicial equation have a relation $n+\alpha-\alpha^{\prime}=0$ for $0<n \in \mathbb{Z}$, the recursion for $\alpha$ may fail, although the recursion for $\alpha^{\prime}$ will succeed. These conditions are easily discovered, as in the following discussion. The convergence of the recursively defined series is important both because it produces a genuine function, and because it can be differentiated termwise, by Abel's theorem.
[16.7.1] The recursion The equation is

$$
x^{\alpha+2} \cdot \sum_{n=0}^{\infty}(n+\alpha)(n+\alpha-1) a_{n} x^{n-2}+b(x) x^{\alpha+1} \sum_{n=0}^{\infty}(n+\alpha) a_{n} x^{n-1}+c(x) x^{\alpha} \sum_{n=0}^{\infty} a_{n} x^{n}=0
$$

Dividing through by $x^{\alpha}$ and grouping,

$$
\sum_{n=0}^{\infty}(n+\alpha)(n+\alpha-1) a_{n} x^{n}+b(x) \sum_{n=0}^{\infty}(n+\alpha) a_{n} x^{n}+c(x) \sum_{n=0}^{\infty} a_{n} x^{n}=0
$$

The vanishing of the sum of coefficients of $x^{0}$, and $a_{0} \neq 0$, give the indicial equation. The coefficients $a_{n}$ with $n>0$ are obtained recursively, from the expected

$$
[(n+\alpha)(n+\alpha-1)+b(0)(n+\alpha)+c(0)] \cdot a_{n}=\left(\text { in terms of } a_{0}, a_{1}, \ldots, a_{n-1}\right)
$$

The coefficient of $a_{n}$ simplifies by invoking the indicial equation and the fact that the sum of the two roots $\alpha, \alpha^{\prime}$ is $1-b(0)$ :

$$
(n+\alpha)(n+\alpha-1)+b(0)(n+\alpha)+c(0)=n(n+(2 \alpha-1)+b(0))=n\left(n+\alpha-\alpha^{\prime}\right)
$$

That is,

$$
n\left(n+\alpha-\alpha^{\prime}\right) \cdot a_{n}=\left(\text { in terms of } a_{0}, a_{1}, \ldots, a_{n-1}\right) \quad(\text { for } n>0)
$$

Since $n>0$, the recursion can fail only when

$$
n+\alpha-\alpha^{\prime}=0 \quad(\text { for some } 0<n \in \mathbb{Z})
$$

[16.7.2] Convergence To complete the proof of existence, we prove convergence. Let $A, M \geq 1$ be large enough so that

$$
\begin{array}{ll}
b(x)=\sum_{n \geq 0} b_{n} x^{n} & \left(\text { with }\left|b_{n}\right| \leq A \cdot M^{n}\right) \\
c(x)=\sum_{n \geq 0} c_{n} x^{n} & \left(\text { with }\left|c_{n}\right| \leq A \cdot M^{n}\right)
\end{array}
$$

Inductively, suppose that $\left|a_{\ell}\right| \leq(C M)^{\ell}$, with a constant $C \geq 1$ to be determined in the following. Then

$$
\left|n\left(n+\alpha-\alpha^{\prime}\right) \cdot a_{n}\right| \leq A \sum_{i=1}^{n}|n-i+\alpha| M^{i} \cdot(C M)^{n-i}+A \sum_{i=1}^{n} M^{i} \cdot(C M)^{n-i} \leq A M^{n} C^{n-1}\left(\frac{n(n+1)}{2}+n|\alpha|+n\right)
$$

Dividing through by $n\left|n+\alpha-\alpha^{\prime}\right|$, this is

$$
\left|a_{n}\right| \leq A M^{n} \cdot C^{n-1} \frac{(n+1)+2|\alpha|+2}{2\left|n+\alpha-\alpha^{\prime}\right|}
$$

This motivates the choice

$$
C \geq \sup _{1 \leq n \in \mathbb{Z}} \frac{(n+1)+2|\alpha|+2}{2\left|n+\alpha-\alpha^{\prime}\right|}
$$

which gives $\left|a_{n}\right| \leq A(C M)^{n}$, and a positive radius of convergence.

### 16.8 Regular singular points at infinity

With $u(x)=v(1 / x)$,

$$
u^{\prime}(x)=\frac{-1}{x^{2}} v^{\prime}(1 / x) \quad \text { and } \quad u^{\prime \prime}(x)=\frac{1}{x^{4}} v^{\prime \prime}(1 / x)+\frac{2}{x^{3}} v^{\prime}(1 / x)
$$

Putting $z=1 / x$, this is

$$
u^{\prime}=-z^{2} v^{\prime} \quad \text { and } \quad u^{\prime \prime}=z^{4} v^{\prime \prime}+2 z^{3} v^{\prime} \quad(\text { with } u=u(x), v=v(z), z=1 / x)
$$

A differential equation $u^{\prime \prime}+p(x) u^{\prime}+q(x) u=0$ becomes

$$
\left(z^{4} v^{\prime \prime}+2 z^{3} v^{\prime}\right)+p(x)\left(-z^{2} v^{\prime}\right)+q(x) v=0
$$

or

$$
z^{2} v^{\prime \prime}+z\left(2-\frac{p(1 / z)}{z}\right) v^{\prime}+\frac{q(1 / z)}{z^{2}} v=0
$$

The point $z=0$ is a regular singular point when the coefficients

$$
2-\frac{p(1 / z)}{z} \quad \frac{q(1 / z)}{z^{2}}
$$

are analytic at 0 . That is, $z=0$ is a regular singular point when $p, q$ have expansions of the forms

$$
\left\{\begin{array} { l } 
{ p ( \frac { 1 } { z } ) = p _ { 1 } z + p _ { 2 } z ^ { 2 } + \ldots } \\
{ q ( \frac { 1 } { z } ) = q _ { 2 } z ^ { 2 } + q _ { 3 } z ^ { 3 } + \ldots }
\end{array} \quad \text { or, equivalently } \quad \left\{\begin{array}{l}
p(x)=\frac{p_{1}}{x}+\frac{p_{2}}{x^{2}}+\ldots \\
q(x)=\frac{q_{2}}{x^{2}}+\frac{q_{3}}{x^{3}}+\ldots
\end{array}\right.\right.
$$

### 16.9 Example revisited

Returning to the earlier example from the upper half-plane: we ask for $\Delta=\Delta^{\mathfrak{H}}$ eigenfunctions $f(z)$ of the special form

$$
f(x+i y)=e^{2 \pi i x} u(y)
$$

The equation $(\Delta-\lambda) f=0$ simplifies to the ordinary differential equation

$$
y^{2} u^{\prime \prime}-\left(4 \pi^{2} y^{2}+\lambda\right) u=0
$$

with regular singular point at $y=0$. The indicial equation is

$$
X(X-1)-\lambda=0
$$

With $\lambda=s(s-1)$, the roots of the indicial equation are $s, 1-s$. By now we know that, unless $s-(1-s)$ is an integer, the equation has solutions of the form

$$
u_{s}(y)=y^{s} \cdot \sum_{\ell \geq 0} a_{\ell} y^{\ell} \quad u_{1-s}(y)=y^{1-s} \cdot \sum_{\ell \geq 0} b_{\ell} y^{\ell}
$$

with coefficients $a_{\ell}$ and $b_{\ell}$ determined by the natural recursions. We emphasize that these power series have positive radius of convergence, so certainly give asymptotics as $y \rightarrow 0^{+}$. Further, convergent series can be differentiated termwise, by Abel's theorem.

We execute a few steps of the recursion for the coefficients for $y^{s}$. The equation

$$
\sum_{\ell \geq 0}(\ell+s)(\ell+s-1) a_{\ell} y^{\ell}-\left(4 \pi^{2} y^{2}+\lambda\right) \sum_{\ell \geq 0} a_{\ell} y^{\ell}=0
$$

simplifies to

$$
\ell(\ell+2 s-1) a_{\ell}=4 \pi^{2} a_{\ell-2} \quad(\text { for } \ell \geq 1)
$$

with $a_{-1}=0$ by convention, and $a_{0}=1$. Thus, the odd-degree terms are all 0 , and

$$
u_{s}(y)=y^{s} \cdot\left(1+\frac{4 \pi^{2} y^{2}}{2(1+2 s)}+\frac{\left(4 \pi^{2}\right)^{2} y^{4}}{2(1+2 s) \cdot 4(3+2 s)}+\ldots\right)
$$

Similarly, replacing $s$ by $1-s$,

$$
u_{1-s}(y)=y^{1-s} \cdot\left(1+\frac{4 \pi^{2} y^{2}}{2(3-2 s)}+\frac{\left(4 \pi^{2}\right)^{2} y^{4}}{2(3-2 s) \cdot 4(5-2 s)}+\ldots\right)
$$

For $\operatorname{Re}(s) \neq \frac{1}{2}$, one of these solutions is obviously asymptotically larger than the other. For $\operatorname{Re}(s)=\frac{1}{2}$, they are the same size, so some cancellation can occur. Write $s=\frac{1}{2}+i \nu$, so $1-s=\frac{1}{2}-i \nu$, and rewrite the expansions in those coordinates:

$$
\left\{\begin{array}{l}
u_{\frac{1}{2}+i \nu}(y)=y^{\frac{1}{2}+i \nu} \cdot\left(1+\frac{\pi^{2} y^{2}}{(1+i \nu)}+\frac{\pi^{4} y^{4}}{(1+i \nu) \cdot 2(2+i \nu)}+\ldots\right) \\
u_{\frac{1}{2}-i \nu}(y)=y^{\frac{1}{2}-i \nu} \cdot\left(1+\frac{\pi^{2} y^{2}}{(1-i \nu)}+\frac{\pi^{4} y^{4}}{(1-i \nu) \cdot 2(2-i \nu)}+\ldots\right)
\end{array}\right.
$$

For example,

$$
\left\{\begin{array}{l}
u_{\frac{1}{2}+i \nu}+u_{\frac{1}{2}-i \nu}=2 y^{\frac{1}{2}} \cos (\log y)+O\left(y^{\frac{3}{2}}\right) \\
u_{\frac{1}{2}+i \nu}-u_{\frac{1}{2}-i \nu}=2 y^{\frac{1}{2}} \sin (\log y)+O\left(y^{\frac{3}{2}}\right)
\end{array}\right.
$$

Further, behavior of the higher terms as functions of $\nu$ is clear.

### 16.10 Irregular singular points

According to Erdélyi, Thomé found that differential equations with finite rank irregular singular points have asymptotic expansions given by the expected recursions. Thus, although the irregularity typically precludes convergence of the series expression for solutions, the series is still a legitimate asymptotic expansion.

We approximately follow Erdélyi in treating a rank-one irregular singular point of a second-order differential equation: after normalization to get rid of the first-derivative term, these are of the form

$$
u^{\prime \prime}-q(x) u=0 \quad \text { with } \quad q(x) \sim q_{o}+\frac{q_{1}}{x}+\frac{q_{2}}{x^{2}}+\ldots \quad\left(\text { as } x \rightarrow+\infty, \text { with } q_{o} \neq 0\right)
$$

with $q$ continuous in some range $x \geq a$. The series expression for $q(x)$ need not be convergent: it suffices that it be an asymptotic expansion of $q(x)$ at $+\infty$. Freezing the coefficient $q$ to its value at $+\infty$, gives the constant-coefficient equation

$$
u^{\prime \prime}-q_{o} u=0
$$

and suggests that the solutions $e^{ \pm \sqrt{q_{o}} x}$ of the constant-coefficient equation should give the leading term in the asymptotics of solutions of the original equation. This is approximately true: there is an adjustment by a power of $x$. Solutions have asymptotics of the form

$$
u(x) \sim e^{ \pm \sqrt{q_{o}} x} \cdot x^{\rho} \cdot\left(1+\sum_{n \geq 1} \frac{a_{n}}{x^{n}}\right) \quad\left(\text { with } \rho=\frac{q_{1}}{ \pm 2 \sqrt{q_{o}}}, \text { as } x \rightarrow+\infty\right)
$$

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with coefficients $a_{n}$ obtained by a natural recursion. However, the series rarely converges.
The loss of convergence is not a trifling matter. The term-wise differentiability of convergent power series is extremely useful. In contrast, term-wise differentiation of asymptotic series

$$
f(x) \sim \sum_{n \geq 0} \frac{a_{n}}{x^{n}} \quad(\text { as } x \rightarrow+\infty)
$$

for differentiable $f$ produces an asymptotic series for $f^{\prime}$ only under additional hypotheses, for example, that $f^{\prime}$ admits such an asymptotic series. (See the appendix.) While a function admitting an asymptotic expansion of this form determines that expansion uniquely, the expansion does not uniquely determine the function. For example, as $x \rightarrow+\infty, e^{-x}=o\left(x^{-N}\right)$ for all $N$, so $e^{-x}$ has the 0 asymptotic expansion, but is not the 0 function.
[16.10.1] Example: rotationally symmetric eigenfunctions on $\mathbb{R}^{n}$ A natural example arises from the eigenvalue equation for the radial component of the Euclidean Laplacian on $\mathbb{R}^{n}$ :

$$
v^{\prime \prime}+\frac{n-1}{r} v^{\prime}-\lambda v=0
$$

For large $r$, this equation resembles the constant-coefficient equation $v^{\prime \prime}-\lambda v=0$, with solutions $e^{ \pm r \sqrt{\lambda}}$. Heuristically, we should have solutions with behavior $v \sim e^{ \pm r \sqrt{\lambda}}$ as $v \rightarrow+\infty$. This is not quite right: the true asymptotic expansions have main terms

$$
v \sim \frac{e^{ \pm r \sqrt{\lambda}}}{r^{\frac{n-1}{2}}}
$$

That is, the differences between the actual equation and the constant-coefficient approximation do not alter the constant in the exponential, but do have a significant impact, as we see below.

A natural recursion, carried out just below, produces an apparent solution to differential equations in this class, of the form

$$
e^{\omega x} x^{-\rho} \sum_{n \geq 0} \frac{c_{n}}{x^{n}}
$$

However, unlike the regular singular point situation, the series is not convergent! The relation of this nonconvergent series to any genuine solution is a priori unclear. It is natural to suppose that this non-convergent series is an asymptotic expansion, but this is not obvious. A genuine solution must be identified by other means, must be proven to have an asymptotic expansion, and the latter must be compared with the series obtained by the recursion. All this will occupy us in following sections.

## [16.10.2] Recursion

In more detail, the heuristic recursion is as follows, as applied to the eigenvalue equation for the radial component of the Laplacian on $\mathbb{R}^{n}$. First, simplify by employing the standard device to eliminate the $v^{\prime}$ term: ${ }^{[85]}$ take $v=u / r^{(n-1) / 2}$, and then

$$
u^{\prime \prime}-\left(\lambda+\frac{(n-1)(n-3)}{4 r^{2}}\right) u=0
$$

The singular point at infinity is irregular, unless $n=1$ or 3 . Nevertheless, intuitively, for $x \rightarrow+\infty$, a differential equation of the shape

$$
u^{\prime \prime}-\left(\lambda+\frac{C}{x^{2}}\right) u=0
$$

[85] Let $v=u \cdot w$ and set the $u^{\prime}$ term equal to 0 in the left-hand side. This gives $2 u^{\prime} w^{\prime}+\frac{n-1}{r} u^{\prime} w=0$, which gives the differential equation $2 w^{\prime}+\frac{n-1}{r} w=0$ for $w$.
is approximately a constant-coefficient differential equation, suggesting a solution of the form ${ }^{[86]}$

$$
u(x)=e^{ \pm x \sqrt{\lambda}} \cdot \sum_{\ell=0}^{+\infty} \frac{c_{\ell}}{x^{\ell}} \quad \quad \text { (with } c_{0}=1, \text { without loss of generality) }
$$

Substituting the latter into the differential equation and dividing through by $e^{ \pm x \sqrt{\lambda}}$, letting $s= \pm \sqrt{\lambda}$, simplifies to

$$
\sum_{\ell=0}^{+\infty}\left((\ell-2)(\ell-1) c_{\ell-2}-2 s(\ell-1) c_{\ell-1}\right) \frac{1}{x^{\ell}}-\sum_{\ell=0}^{+\infty} c_{\ell-2} \frac{C}{x^{\ell}}=0
$$

where by convention $c_{-1}=c_{-2}=0$. The case $\ell=0$ is vacuous, as is $\ell=1$. The case $\ell \geq 2$ determines $c_{\ell-1}$, assuming $s \neq 0$ :

$$
(\ell-2)(\ell-1) c_{\ell-2}-2 s(\ell-1) c_{\ell-1}-C \cdot c_{\ell-2}=0
$$

or

$$
c_{\ell+1}=\frac{\ell(\ell+1)-C}{2 s(\ell+1)} \cdot c_{\ell}
$$

[16.10.3] Remark: That recursion causes the coefficients to grow approximately as factorials, and the resulting series does not converge for any finite non-zero value of $1 / x$, unless the constant $C$ happens to be of the form $(\ell-1)(\ell-2)$ for some positive integer $\ell$, in which case the series terminates, and is convergent.

Our later discussion will show that the above recursion does correctly determine asymptotic expansions for solutions. In particular, the leading part of the asymptotic is

$$
v=\frac{e^{ \pm r \sqrt{\lambda}}}{r^{\frac{n-1}{2}}} \cdot\left(1+O\left(\frac{1}{r}\right)\right) \quad\left(\text { as } r \rightarrow+\infty, \text { in } \mathbb{R}^{n}\right)
$$

The denominator $r^{(n-1) / 2}$ might be hard to anticipate. The symmetry $r \rightarrow-r$ imposes a further requirement, and for $\sqrt{\lambda}$ not purely imaginary one of the two solutions swamps the other. Indeed, for $\sqrt{\lambda}$ not purely imaginary, the asymptotic components of the large solution are all larger than the main part of the smaller solution. Further, this is an asymptotic and not merely a bound.

In fact, for $n$ odd, the asymptotic is finite: the recursion for coefficients terminates, so gives a convergent series: we obtain not merely asymptotics, but equalities. Thus, in odd-dimensional $\mathbb{R}^{n}$ the solutions to the differential equation for rotationally-invariant $\lambda$-eigenfunctions have elementary expressions. For example,

$$
\begin{cases}v=e^{ \pm r \sqrt{\lambda}} & \left(\text { on } \mathbb{R}^{1}\right) \\ v=\frac{e^{ \pm r \sqrt{\lambda}}}{r} & \left(\text { on } \mathbb{R}^{3}\right) \\ v=e^{ \pm r \sqrt{\lambda}}\left(\frac{1}{r^{2}}-\frac{1}{ \pm r^{3} \sqrt{\lambda}}\right) & \left(\text { on } \mathbb{R}^{5}\right) \\ v=e^{ \pm r \sqrt{\lambda}}\left(\frac{1}{r^{3}}-\frac{3}{ \pm r^{5} \sqrt{\lambda}}+\frac{3}{r^{7} \lambda}\right) & \left(\text { on } \mathbb{R}^{7}\right)\end{cases}
$$

[16.10.4] Remark: The same technique applies to differential equations

$$
u^{\prime \prime}-q(x) u=0
$$

with $q(x)$ continuous in some range $x>a$ and admitting an asymptotic expansion at infinity of the form

$$
q(x) \sim \sum_{\ell \geq 0} \frac{q_{\ell}}{x^{\ell}} \quad\left(\text { with } q_{o} \neq 0\right)
$$

[86] Anticipating the adjustment by $x^{\rho}$ in general, with $\rho$ determined by the asymptotics $q_{o}+\frac{q_{1}}{x}+\ldots$ of the coefficient of $u$ by $\rho=q_{1} / 2 \sqrt{q_{o}}$, in the present example we are fortunate that $q_{1}=0$, so the idea of freezing is exactly right.

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The condition $q_{o} \neq 0$ is essential ${ }^{[87]}$ for the recursion to succeed. Adjustment by $x^{\rho}$ with $\rho=q_{1} / 2 \sqrt{q_{o}}$ would be found necessary when $q_{1} \neq 0$. In any case, the recursion rarely produces a convergent power series!

## [16.10.5] Comparison to regular singular points

The behavior of the above recursion is is much different from that resulting from a regular singular point. A power series in $z=1 / x$ behaves differently under $d / d x$ than under $d / d z$. Indeed, as in the example above, the power series in $1 / x$ often diverges, while at a regular singular point the analogous power series has a positive radius of convergence. For $u^{\prime \prime}-q(x) u=0$ to have a regular singular point at infinity, changing variables to $u(x)=v(1 / x)$ and $z=1 / x$,

$$
u^{\prime}(x)=\frac{-1}{x^{2}} v^{\prime}(1 / x) \quad \text { and } \quad u^{\prime \prime}(x)=\frac{1}{x^{4}} v^{\prime \prime}(1 / x)+\frac{2}{x^{3}} v^{\prime}(1 / x)
$$

Putting $z=1 / x$, this is

$$
u^{\prime}=-z^{2} v^{\prime} \quad \text { and } \quad u^{\prime \prime}=z^{4} v^{\prime \prime}+2 z^{3} v^{\prime} \quad(\text { with } u=u(x), v=v(z), z=1 / x)
$$

Thus, in the coordinate $z$ at infinity, the differential equation becomes

$$
\left(z^{4} v^{\prime \prime}+2 z^{3} v^{\prime}\right)-q\left(\frac{1}{z}\right) v=0
$$

or

$$
v^{\prime \prime}+\frac{2}{z} v^{\prime}-\frac{q(1 / z)}{z^{4}} v=0
$$

The point $z=0$ is never an ordinary point, because of the pole in the coefficient of $v^{\prime}$. The point $z=0$ is a regular singular point only when $q(1 / z) / z^{2}$ is analytic at $z=0$, that is, when $x^{2} q(x)$ is analytic at $\infty$. This requires that $q(x)$ have the form

$$
q(x)=\frac{q_{2}}{x^{2}}+\frac{q_{3}}{x^{3}}+\ldots
$$

### 16.11 Example: translation-equivariant eigenfunctions on $\mathfrak{H}$

Another example of irregular singular point arises from the $S L_{2}(\mathbb{R})$-invariant Laplacian on the upper half-plane $\mathfrak{H}$ :

$$
\Delta^{\mathfrak{H}}=y^{2}\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right)
$$

We ask for $\Delta^{\mathfrak{H}}$-eigenfunctions $f(z)$ of the special form

$$
f(x+i y)=e^{2 \pi i x} u(y)
$$

that is, equivariant under translations:

$$
f(z+t)=e^{2 \pi i(x+t)} u(y)=e^{2 \pi i t} \cdot\left(e^{2 \pi i x} u(y)\right)=e^{2 \pi i t} \cdot f(z) \quad(\text { with } t \in \mathbb{R} \text { and } z \in \mathfrak{H})
$$

The eigenfunction condition

$$
\left(\Delta^{\mathfrak{H}}-\lambda\right) e^{2 \pi i x} u(y)=0
$$

simplifies to the ordinary differential equation

$$
y^{2} u^{\prime \prime}-\left(4 \pi^{2} y^{2}+\lambda\right) u=0
$$

[87] The condition $q_{o} \neq 0$ and the assumption that $q$ has the indicated asymptotic at $+\infty$ together imply that there is $x_{o}$ such that $q(x) \neq 0$ for $x \geq x_{o}$. That is, in the regime $x \geq x_{o}$ there are no transition points.

This equation has an irregular singular point at $+\infty$, seen by changing coordinates, as follows. Let $u(x)=v(1 / x)$ and put $z=1 / x$, obtaining

$$
\left(z^{4} v^{\prime \prime}+2 z^{3} v^{\prime}\right)-\left(4 \pi^{2}+\lambda z^{2}\right) v=0
$$

or

$$
z^{2} v^{\prime \prime}+2 z v^{\prime}-\left(\frac{4 \pi^{2}}{z^{2}}+\lambda\right) v=0
$$

Since the coefficient of $v$ has a pole at $z=0$, the singular point of this equation in the new coordinate $z$ at 0 is irregular.
[16.11.1] Recursion Happily, following our present prescription, in the form

$$
u^{\prime \prime}-\left(4 \pi^{2}+\frac{\lambda}{y^{2}}\right) u=0
$$

the coefficient

$$
q(y)=q_{o}+\frac{q_{1}}{y}+\frac{q_{2}}{y^{2}}+\ldots=4 \pi^{2}+\frac{\lambda}{y^{2}}
$$

is analytic at $y=\infty$, and $q(\infty)=q_{o}=4 \pi^{2} \neq 0$ while $q_{1}=0$, so our later discussion justifies freezing $y$ at $+\infty$, obtaining the constant-coefficient equation

$$
u^{\prime \prime}-4 \pi^{2} u=0
$$

with solutions $e^{ \pm 2 \pi y}$, and assuring existence of solutions of the original equation with asymptotics of the form

$$
u(y)=e^{ \pm 2 \pi y} \cdot \sum_{\ell \geq 0} \frac{a_{\ell}}{y^{\ell}}
$$

Substituting this into the differential equation and dividing through by $e^{ \pm 2 \pi y}$ gives

$$
\sum_{\ell=0}^{+\infty}\left((\ell-2)(\ell-1) a_{\ell-2} \mp 2 \pi(\ell-1) a_{\ell-1}\right) \frac{1}{y^{\ell}}-\sum_{\ell=0}^{+\infty} a_{\ell-2} \frac{\lambda}{y^{\ell}}=0
$$

or

$$
\pm 2 \pi(\ell-1) a_{\ell-1}=((\ell-2)(\ell-1)-\lambda) a_{\ell-2}
$$

or

$$
a_{\ell}=\frac{(\ell-1) \ell-\lambda}{ \pm 2 \pi \ell} a_{\ell-1}=\left(\ell-1-\frac{\lambda}{\ell}\right) \frac{a_{\ell}}{ \pm 2 \pi}
$$

As usual, $a_{-2}=a_{-1}=0$ by convention, and $a_{0}=1$. The cases $\ell \leq 0$ are vacuous. With $a_{0}=1$, the recursion begins

$$
\begin{aligned}
& a_{1}=\frac{-\lambda}{ \pm 2 \pi} \\
& a_{2}=\left(1-\frac{\lambda}{2}\right) \frac{a_{1}}{ \pm 2 \pi}=\left(1-\frac{\lambda}{2}\right)(-\lambda) \frac{1}{( \pm 2 \pi)^{2}} \\
& a_{3}=\left(2-\frac{\lambda}{3}\right) \frac{a_{2}}{ \pm 2 \pi}=\left(2-\frac{\lambda}{3}\right)\left(1-\frac{\lambda}{2}\right)(-\lambda) \frac{1}{( \pm 2 \pi)^{3}}
\end{aligned}
$$

[16.11.2] Remark: If $\lambda$ is of the form $\lambda=\ell(\ell-1)$ for $0<\ell \in \mathbb{Z}$, the recursion terminates. Then the asymptotic expansion is convergent, and produces an elementary solution to the eigenfunction equation. [88]
[88] These elementary solutions arise from the finite-dimensional representations of $S L_{2}(\mathbb{R})$.

### 16.12 Beginning of construction of solutions

[16.12.1] Heuristic for asymptotic expansion Consider the equation

$$
u^{\prime \prime}-q(x) u=0
$$

as $x \rightarrow+\infty$, where $q$ is continuous in some range $x>a$ and itself admits an asymptotic expansion

$$
q(x) \sim \sum_{n \geq 0} \frac{q_{n}}{x^{n}} \quad\left(\text { as } x \rightarrow+\infty, \text { with } q_{o} \neq 0\right)
$$

The $q_{o} \neq 0$ condition is essential. We look for a solution of the form

$$
u(x) \sim e^{\omega x} \cdot x^{-\rho} \cdot \sum_{n \geq 0} \frac{c_{o}}{x^{n}} \quad \quad \text { (with } c_{o} \text { non-zero) }
$$

Substituting this expansion in the differential equation and dividing through by $e^{\omega x} x^{-\rho}$, setting the coefficient of $1 / x^{n}$ to 0 ,

$$
\left((\rho+n-2)(\rho+n-1) c_{n-2}-2 \omega(\rho+n-1) c_{n-1}+\omega^{2} c_{n}\right)-\left(q_{o} c_{n}+q_{1} c_{n-1}+\ldots+q_{n-1} c_{1}+q_{n} c_{o}\right)=0
$$

By convention, $c_{-2}=c_{-1}=0$ and $q_{-2}=q_{-1}=0$. For $n=0$, the relation is

$$
\omega^{2} c_{o}-q_{o} c_{o}=0
$$

so $\omega= \pm \sqrt{q_{o}} \neq 0$, since $c_{o} \neq 0$. For $n=1$,

$$
\left(-2 \omega \rho c_{o}+\omega^{2} c_{1}\right)-\left(q_{o} c_{1}+q_{1} c_{o}\right)=0
$$

so, using $\omega^{2}=q_{o}$ and $\omega \neq 0$, this is

$$
-2 \omega \rho-q_{1}=0
$$

so $\rho=-q_{1} /(2 \omega)$. Thus, the choice of $\pm \omega$ is reflected in the choice of $\pm \rho$. For $n \geq 2$, using $\omega^{2}=q_{o}$,

$$
\left(-2 \omega(\rho+n-1)-q_{1}\right) c_{n-1}=-(\rho+n-2)(\rho+n-1) c_{n-2}+\left(q_{2} c_{n-2}+\ldots+q_{n-1} c_{1}+q_{n} c_{o}\right)
$$

and using $-2 \omega \rho-q_{1}=0$,

$$
-2 \omega(n-1) c_{n-1}=-(\rho+n-2)(\rho+n-1) c_{n-2}+\left(q_{2} c_{n-2}+\ldots+q_{n-1} c_{1}+q_{n} c_{o}\right)
$$

Since $\omega \neq 0$, this gives a successful recursion. The following discussion will show that the two asymptotics, with $\pm \omega$ and corresponding $\pm \rho$, are asymptotic expansions of two solutions of the differential equation $u^{\prime \prime}-q(x) u=0$.
[16.12.2] Remark: However, since the above expansions usually do not converge, genuine solutions must be constructed by other means, and must be shown to have asymptotic expansions at $+\infty$.
[16.12.3] Small renormalization For a solution $u$ to $u^{\prime \prime}-q(x) u=0$, let

$$
u(x)=e^{\omega x} \cdot x^{-\rho} \cdot v(x)
$$

with $\omega$ and $\rho$ determined as above. Then

$$
\left\{\begin{array}{l}
u^{\prime}=e^{\omega x} x^{-\rho}\left(\left(\omega-\frac{\rho}{x}\right) v+v^{\prime}\right) \\
u^{\prime \prime}=e^{\omega x} x^{-\rho}\left(\omega-\frac{\rho}{x}\right)^{2} v+e^{\omega x} x^{-\rho} \frac{\rho}{x^{2}} v+2 e^{\omega x} x^{-\rho}\left(\omega-\frac{\rho}{x}\right) v^{\prime}+e^{\omega x} x^{-\rho} v^{\prime \prime}
\end{array}\right.
$$

Dividing through by $e^{\omega x} x^{-\rho}$ gives the differential equation for $v$, namely,

$$
v^{\prime \prime}+2\left(\omega-\frac{\rho}{x}\right) v^{\prime}+\left(\omega^{2}-\frac{2 \omega \rho}{x}+\frac{\rho^{2}+\rho}{x^{2}}-q(x)\right) v=0
$$

Unsurprisingly, the $\omega^{2}$ and $-2 \omega \rho / x$ cancel the first two terms of $q(x)$. Thus, the function

$$
F(x)=x^{2} \cdot\left(\omega^{2}-\frac{2 \omega \rho}{x}+\frac{\rho^{2}+\rho}{x^{2}}-q(x)\right)
$$

is bounded. The differential equation is

$$
v^{\prime \prime}+2\left(\omega-\frac{\rho}{x}\right) v^{\prime}+\frac{F(x)}{x^{2}} v=0
$$

Rewrite the equation as

$$
\frac{d}{d x}\left(e^{2 \omega x} x^{-2 \rho} \frac{d v}{d x}\right)+e^{2 \omega x} x^{-2 \rho-2} F(x) v(x)=0
$$

Integrate this from $b \geq a$ to $x \geq b$, and multiply through by $e^{-2 \omega x} x^{2 \rho}$, to obtain

$$
\frac{d v}{d x}+e^{-2 \omega x} x^{2 \rho} \int_{b}^{x} e^{2 \omega t} t^{-2 \rho-2} F(t) v(t) d t=\mathrm{const} \cdot e^{-2 \omega x} x^{2 \rho}
$$

Take the constant of integration to be 0 and integrate from $a$ to $x$, to obtain

$$
v(x)+\int_{a}^{x} e^{-2 \omega s} s^{2 \rho}\left(\int_{b}^{s} e^{2 \omega t} t^{-2 \rho-2} F(t) v(t) d t\right) d s=\mathrm{const}
$$

Rearrange the double integral:

$$
\int_{a}^{x} e^{-2 \omega s} s^{2 \rho}\left(\int_{b}^{s} e^{2 \omega t} t^{-2 \rho-2} F(t) v(t) d t\right) d s=\int_{b}^{x}\left(\int_{t}^{x} e^{2 \omega(t-s)}\left(\frac{s}{t}\right)^{2 \rho} d s\right) F(t) v(t) \frac{d t}{t^{2}}
$$

Let $K(x, t)$ denote the inner integral

$$
K(x, t)=\int_{t}^{x} e^{2 \omega(t-s)}\left(\frac{s}{t}\right)^{2 \rho} d s
$$

Then the equation is

$$
v(x)-\int_{b}^{x} K(x, t) F(t) v(t) \frac{d t}{t^{2}}=\mathrm{const}
$$

Take the constant to be 1 . With $b=+\infty$, this gives an integral equation

$$
v(x)=1+\int_{x}^{\infty} K(x, t) F(t) v(t) \frac{d t}{t^{2}}
$$

We claim that this equation can be solved by successive approximations. With the obvious operator

$$
T f(x)=\int_{x}^{\infty} K(x, t) F(t) f(t) \frac{d t}{t^{2}}
$$

take $w_{o}(x)=1, w_{n+1}=T w_{n}$, and then show that the limit

$$
v(x)=w_{o}(x)+w_{1}(x)+w_{2}(x)+\ldots=\left(1+T+T^{2}+\ldots\right) w_{o}
$$

exists pointwise, is twice differentiable, and satisfies the differential equation.

### 16.13 Boundedness of $K(x, t)$

We claim that, with correct choice of $\pm \omega$, the kernel

$$
-K(x, t)=\int_{x}^{t} e^{2 \omega(t-s)}\left(\frac{s}{t}\right)^{2 \rho} d s
$$

is bounded for $t \geq x \geq a$. Choose $\pm \omega$ so that either $\operatorname{Re}(\omega)<0$, or $\operatorname{Re}(\omega)=0$ and $\operatorname{Re}(\rho) \geq 0$.
[16.13.1] Very easy case $\rho=0$ To illustrate the reasonableness of the boundedness assertion, consider the special case $\rho=0$, where the integral can be computed explicitly:

$$
-K(x, t)=\int_{x}^{t} e^{2 \omega(t-s)} d s=\frac{1}{-2 \omega}\left(1-e^{2 \omega(t-x)}\right)
$$

Since $\operatorname{Re} \omega \leq 0$ and $\omega \neq 0$, this is bounded, for $a \leq x \leq t$.
[16.13.2] Easy case $\operatorname{Re} \omega<0$ When $\operatorname{Re} \omega<0$, absolute value estimates suffice to prove boundedness of $K(x, t)$.

$$
|K(x, t)| \leq \int_{x}^{t} e^{2 \operatorname{Re} \omega(t-s)}\left(\frac{s}{t}\right)^{2 \operatorname{Re} \rho} d s \leq \int_{a}^{t} e^{2 \operatorname{Re} \omega(t-s)}\left(\frac{s}{t}\right)^{2 \operatorname{Re} \rho} d s
$$

Lighten the notation by taking $\omega, \rho$ real. For $\rho \geq 0$,

$$
\int_{a}^{t} e^{2 \omega(t-s)}\left(\frac{s}{t}\right)^{2 \rho} d s \leq \int_{0}^{t} e^{2 \omega(t-s)} d s=e^{2 \omega t} \cdot \frac{e^{-2 \omega t}-1}{2|\omega|} \leq \frac{1}{2|\omega|} \quad(\text { for } \rho \geq 0)
$$

For $\rho<0$, still with $\omega<0$,

$$
\int_{a}^{t} e^{2 \omega(t-s)}\left(\frac{s}{t}\right)^{2 \rho} d s \leq \int_{t / 2}^{t} e^{2 \omega(t-s)}\left(\frac{t / 2}{t}\right)^{2 \rho} d s+\int_{0}^{t / 2} e^{2 \omega(t-s)}\left(\frac{1}{t}\right)^{2 \rho} d s \quad(\text { for } \rho<0)
$$

The two integrals are bounded in $t \geq a$, for elementary reasons. Thus, for $\operatorname{Re}(\omega)<0$, the kernel $K(x, t)$ is bounded.

## [16.13.3] $\operatorname{Re}(\omega)=0$ and cancellation

When $\operatorname{Re}(\omega)=0$, absolute value estimates no longer suffice to prove boundedness. Cancellation must be exploited by an integration by parts. Choose $\pm \omega$ so that $\operatorname{Re}(\rho) \geq 0$. One integration by parts gives

$$
\begin{gathered}
\int_{x}^{t} e^{2 \omega(t-s)}\left(\frac{s}{t}\right)^{2 \rho} d s=\left[\frac{e^{2 \omega(t-s)}}{-2 \omega}\left(\frac{s}{t}\right)^{2 \rho}\right]_{x}^{t}+\int_{x}^{t} \frac{e^{2 \omega(t-s)}}{2 \omega} \frac{2 \rho}{s}\left(\frac{s}{t}\right)^{2 \rho} d s \\
=\frac{1}{-2 \omega}-\frac{e^{2 \omega(t-x)}}{-2 \omega}\left(\frac{x}{t}\right)^{2 \rho}+\int_{x}^{t} \frac{e^{2 \omega(t-s)}}{2 \omega} \frac{2 \rho}{s}\left(\frac{s}{t}\right)^{2 \rho} d s
\end{gathered}
$$

The leading terms are bounded for $t \geq x \geq a$. The latter integral can be estimated by absolute values, for $\operatorname{Re} \rho \neq 0$ :

$$
\left|\int_{x}^{t} e^{2 \omega(t-s)} \frac{1}{s}\left(\frac{s}{t}\right)^{2 \rho} d s\right| \leq \int_{x}^{t} \frac{1}{s}\left(\frac{s}{t}\right)^{2 \operatorname{Re} \rho} d s=\frac{1}{2 \operatorname{Re} \rho}\left(1-\left(\frac{x}{t}\right)^{2 \operatorname{Re} \rho}\right)
$$

When $\operatorname{Re} \rho=0$, a second integration by parts gives

$$
\int_{x}^{t} e^{2 \omega(t-s)} \frac{1}{s}\left(\frac{s}{t}\right)^{2 \rho} d s=\left[\frac{e^{2 \omega(t-s)}}{-2 \omega} \frac{1}{s}\left(\frac{s}{t}\right)^{2 \rho}\right]_{x}^{t}+\frac{2 \rho-1}{2 \omega} \int_{x}^{t} e^{2 \omega(t-s)} \frac{1}{s^{2}}\left(\frac{s}{t}\right)^{2 \rho} d s
$$

The latter integral is estimated by

$$
\left|\int_{x}^{t} e^{2 \omega(t-s)} \frac{1}{s^{2}}\left(\frac{s}{t}\right)^{2 \rho} d s\right| \leq \int_{x}^{t} \frac{d s}{s^{2}}=\frac{1}{x}-\frac{1}{t} \leq \frac{1}{a}
$$

Thus, in all cases, $K(x, t)$ is bounded on $t \geq x \geq a>0$.

### 16.14 End of construction of solutions

[16.14.1] Bound for $T$ As observed above, there is a bound $A$ so that $|F(x)| \leq A$ for $x \geq a$. Let $|K(x, t)| \leq B$. For $f$ satisfying a bound $|f(x)| \leq x^{-\lambda}$ for $x \geq a$, with $\lambda>-1$,

$$
|(T f)(x)| \leq \frac{A B}{\lambda+1} x^{-(\lambda+1)} \quad(\text { for } x \geq a)
$$

Indeed,

$$
|T f(x)|=\left\lvert\, \int_{x}^{\infty} K(x, t) F(t) f(t) \frac{d t}{t^{2}} \leq A B \int_{x}^{\infty} t^{-(\lambda+2)} d t\right.
$$

[16.14.2] Bound on $f_{n}$ With $f_{0}=1$ and $f_{n+1}=T f_{n}$, we claim that

$$
\left|f_{n}(x)\right| \leq \frac{(A B)^{n}}{n!} x^{-n} \quad(\text { for } n=0,1,2, \ldots \text { and } x \geq a)
$$

This holds for $n=0$, and induction using the bound on $T$ gives the result.
[16.14.3] Convergence of the series Now we show that the series

$$
f(x)=\sum_{n \geq 0} f_{n}(x)=\sum_{n \geq 0} T^{n} f_{0}(x)
$$

converges uniformly absolutely, and satisfies the integral equation

$$
f(x)=1+\int_{x}^{\infty} K(x, t) F(t) f(t) \frac{d t}{t^{2}}
$$

Uniform absolute convergence in $C^{o}[a,+\infty)$ follows from the previous estimate. This justifies interchange of summation and integration:

$$
\begin{aligned}
T f(x) & =\int_{x}^{\infty} K(x, t) F(t) f(t) \frac{d t}{t^{2}}=\sum_{n \geq 0} \int_{x}^{\infty} K(x, t) F(t) T^{n} f_{0}(t) \frac{d t}{t^{2}} \\
& =\sum_{n \geq 0} T^{n+1} f_{0}(x)=-1+\sum_{n \geq 0} T^{n} f_{0}(x)=-1+f(x)
\end{aligned}
$$

Thus, $f$ satisfies the integral equation. Since $K(x, t)$ is differentiable in $x$, and since the integral for $T$ converges well, the expression

$$
f(x)=1+\int_{x}^{\infty} K(x, t) F(t) f(t) \frac{d t}{t^{2}}
$$

demonstrates the differentiability of $f$. Further, since $K(x, x)=0$, the derivative is

$$
f^{\prime}(x)=\int_{x}^{\infty} \frac{\partial K(x, t)}{\partial x} F(t) f(t) \frac{d t}{t^{2}}=\int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} F(t) f(t) \frac{d t}{t^{2}}
$$

The integral is again continuously differentiable in $x$, so $f$ is in $C^{2}$.
[16.14.4] Back to the differential equation From the integral expression,

$$
f^{\prime \prime}(x)=-\frac{F(x)}{x^{2}} f(x)+\int_{x}^{\infty}\left(-2 \omega+\frac{2 \rho}{x}\right) e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} F(t) f(t) \frac{d t}{t^{2}}
$$

Substituting into the differential equation,

$$
\begin{gathered}
f^{\prime \prime}+2\left(\omega-\frac{\rho}{x}\right) f^{\prime}+\frac{F}{x^{2}} f= \\
-\frac{F}{x^{2}} f+\int_{x}^{\infty}\left(-2 \omega+\frac{2 \rho}{x}\right) e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} F(t) f(t) \frac{d t}{t^{2}} \\
+2\left(\omega-\frac{\rho}{x}\right) \int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} F(t) f(t) \frac{d t}{t^{2}}+\frac{F}{x^{2}} f=0
\end{gathered}
$$

Then

$$
u(x)=e^{\omega x} x^{-\rho} f(x)
$$

satisfies the original equation

$$
u^{\prime \prime}-q(x) u=0
$$

[16.14.5] Two independent solutions In the special case that $q_{o}<0$ and $q_{1} \in \mathbb{R}, \omega=\sqrt{\omega}$ has $\operatorname{Re} \omega=0$ and $\operatorname{Re} \rho=0$. In that case, the successive approximation solution to the integral equation can proceed with either values $\pm \omega, \pm \rho$, and two linearly independent solutions are obtained.

In all other cases, the successive approximation argument succeeds for only one choice of sign, producing a solution $u$ as above. Nevertheless, a second solution can be constructed as follows, by a standard device. Since $f(x)=1+O(1 / x)$, there is $b \geq a$ large enough so that $u(x) \neq 0$ for $x \geq b$. Then let $v=u \cdot w$, require that $v$ satisfy $v^{\prime \prime}-q v=0$, and see what condition this imposes on $w$. From

$$
v^{\prime \prime}-q v=u^{\prime \prime} w+2 u^{\prime} w^{\prime}+u w^{\prime \prime}-q u w=0
$$

using $u^{\prime \prime}-q u=0$, we obtain

$$
\frac{w^{\prime \prime}}{w^{\prime}}=\frac{-2 u^{\prime}}{u}
$$

Then

$$
\log w^{\prime}=-2 \log u+C
$$

and

$$
w(x)=\int_{b}^{x} u(t)^{-2} d t
$$

Thus, a second solution is

$$
u(x) \cdot \int_{b}^{x} u(t)^{-2} d t
$$

That integral is not constant, so the two solutions are linearly independent.

### 16.15 Asymptotics of solutions

We show that the solutions on $x \geq a$ have the same asymptotics as the heuristic indicated earlier.
[16.15.1] Some elementary asymptotics Use the standard device $(\rho)_{\ell}=\rho(\rho+1) \ldots(\rho+\ell-1)$ and $(\rho)_{0}=1$. Let $0 \neq \omega \in \mathbb{C}$ with $\operatorname{Re} \omega \leq 0$. If $\operatorname{Re} \omega=0$, require that $\operatorname{Re} \rho>1$. Repeated integration by parts and easy estimates yield asymptotic expansions,

$$
\left\{\begin{array}{l}
\int_{x}^{\infty} e^{\omega t} t^{-\rho} d t \sim e^{\omega x} \cdot \sum_{\ell \geq 0} \frac{(\rho)_{\ell}}{(-\omega)^{\ell+1}} \frac{1}{x^{\rho+\ell}} \\
\int_{b}^{x} e^{-\omega t} t^{-\rho} d t \sim e^{-\omega x} \cdot \sum_{\ell \geq 0} \frac{(\rho)_{\ell}}{\omega^{\ell+1}} \frac{1}{x^{\rho+\ell}}
\end{array}\right.
$$

Since the sup of $\left|e^{\omega t} t^{-\rho}\right|$ occurs farther to the right for larger $\operatorname{Re}(\rho)<0$, these asymptotics are not uniform in $\rho$. Note that the boundedness of the kernel $K(x, t)$ proven earlier has a weaker hypothesis than the second asymptotic assertion, requires a slightly more complicated argument, and has a weaker conclusion.
[16.15.2] Asymptotics of $T^{n} f_{0}$ With $f_{0}=1$, we claim that $f_{n}=T^{n} f_{0}$ has an asymptotic expansion at $+\infty$, of the form

$$
f_{n} \sim \sum_{\ell \geq n} c_{n \ell} x^{-\ell}
$$

This holds for $f_{0}=1$. To do the induction step, assume $f_{n}$ has such an asymptotic expansion. Then $F(x) \cdot f_{n}(x)$ has a similar expansion

$$
F f_{n} \sim \sum_{\ell \geq n} b_{\ell} x^{-\ell}
$$

because (as the product of two asymptotic expansions in $1 / x^{n}$ is readily shown to be an asymptotic expansion for the product function)

$$
F(x)=x^{2} \cdot\left(\omega^{2}-\frac{2 \omega \rho}{x}+\frac{\rho^{2}+\rho}{x^{2}}-q(x)\right)
$$

and $q$ is assumed to have an asymptotic expansion in the functions $1 / x^{n}$ at $+\infty$. We want to insert the asymptotic expansion for $F f_{n}$ into the integral in the differentiated form of $f_{n+1}=T f_{n}$, namely, into the equation

$$
f_{n+1}^{\prime}(x)=\int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} F(t) f_{n}(t) \frac{d t}{t^{2}}
$$

Indeed, from

$$
F(x) f_{n}(x)-\sum_{n \leq \ell \leq N} b_{\ell} x^{-\ell}=O\left(x^{-(N+1)}\right)
$$

and from the boundedness of $K(x, t)$ we have

$$
\begin{gathered}
\left|\int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho}\left(F(t) f_{n}(t)-\sum_{n \leq \ell \leq N} b_{\ell} t^{-\ell}\right) \frac{d t}{t^{2}}\right|=\left|\int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} O\left(t^{-(N+1)}\right) \frac{d t}{t^{2}}\right| \\
\ll \omega, \rho, N \quad x^{-(N+1)} \int_{x}^{\infty} \frac{d t}{t^{2}}=O\left(x^{-(N+2)}\right)=o\left(x^{-(N+1)}\right)
\end{gathered}
$$

Thus, the desired asymptotics for $f_{n+1}^{\prime}$ would follow from asymptotics for the collection

$$
\int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho}\left(\sum_{n \leq \ell \leq N} b_{\ell} t^{-\ell}\right) \frac{d t}{t^{2}} \quad(\text { for } N \geq n)
$$

As noted above,

$$
\int_{x}^{\infty} e^{\omega t} t^{-\rho} d t \sim e^{\omega x} \cdot \sum_{\ell \geq 0} \frac{(\rho)_{\ell}}{(-\omega)^{\ell+1}} \frac{1}{x^{\rho+\ell}}
$$

Note that for each $N$ only finitely-many asymptotic expansions are used, so uniformity is not an issue. After some preliminary rearrangements, this gives

$$
\begin{aligned}
& \int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho}\left(\sum_{n \leq \ell \leq N} b_{\ell} t^{-\ell}\right) \frac{d t}{t^{2}}=\sum_{n \leq \ell \leq N} b_{\ell} \int_{x}^{\infty} e^{2 \omega(t-x)}\left(\frac{x}{t}\right)^{2 \rho} t^{-\ell} \frac{d t}{t^{2}} \\
&=\sum_{n \leq \ell \leq N} b_{\ell} e^{-2 \omega x} x^{2 \rho} \int_{x}^{\infty} e^{2 \omega t} t^{-(2 \rho+\ell+2)} d t \\
&= \sum_{n \leq \ell \leq N} b_{\ell} e^{-2 \omega x} x^{2 \rho} \cdot e^{2 \omega x}\left(\sum_{0 \leq m \leq N-(2+\ell)} \frac{(\rho+\ell+2)_{m}}{(-2 \omega)^{m+1}} \frac{1}{x^{2 \rho+\ell+2+m}}+O\left(\frac{1}{x^{2 \rho+N+1}}\right)\right)
\end{aligned}
$$

$$
=\sum_{n \leq \ell \leq N} b_{\ell} \sum_{0 \leq m \leq N-(2+\ell)} \frac{(\rho+\ell+2)_{m}}{(-2 \omega)^{m+1}} \frac{1}{x^{\ell+2+m}}+O\left(\frac{1}{x^{N+1}}\right)
$$

This holds for all $N$, so we have an asymptotic expansion for $f_{n+1}^{\prime}$ :

$$
f_{n+1}^{\prime}(x) \sim \sum_{k \geq n+1}\left(\sum_{\ell: n \leq \ell \leq k} b_{\ell} \frac{(\rho+\ell+2)_{k-\ell}}{(-2 \omega)^{m+1}}\right) \frac{1}{x^{k+2}}
$$

Integrating this in $x$ gives the asymptotic expansion of $f_{n+1}$. (See the appendix.)
[16.15.3] Asymptotics of the solution $f$ Obviously we expect the asymptotic expansion of $f=\sum_{n} f_{n}$ to be the sum of those of $f_{n}$, all the more so since the $1 / x^{m}$ terms in the expansion of $f_{n}$ vanish for $m<n$. The uniform pointwise bound

$$
\left|f_{n}(x)\right| \leq \frac{(A B)^{n}}{n!} x^{-n} \quad(\text { for } n=0,1,2, \ldots \text { and } x \geq a)
$$

proven earlier legitimizes this. Thus, the solution $f$ has an asymptotic expansion of the desired type.
To prove that this asymptotic expansion is the same as the expansion obtained by a recursion earlier, we show that the coefficients satisfy the same recursion.

The integral expression for $f^{\prime}$ in terms of $f$ (above) proves that $f^{\prime}$ has an asymptotic expansion, and similarly for $f^{\prime \prime}$. As proven in the appendix, this justifies two termwise differentiations of the asymptotic for $f$.

The asymptotics for $f, f^{\prime}$, and $f^{\prime \prime}$ can be inserted in the differential equation

$$
f^{\prime \prime}+2\left(\omega-\frac{\rho}{x}\right) f^{\prime}+\left(\omega^{2}-\frac{2 \omega \rho}{x}+\frac{\rho^{2}+\rho}{x^{2}}-q(x)\right) f=0
$$

for $f$. We have assumed that the coefficient of $f$ has an asymptotic expansion, and this equation gives the expected recursive relation on the coefficients of the asymptotic for $f$. Therefore, the solution

$$
u(x)=e^{\omega x} x^{-\rho} f(x)
$$

to the original differential equation has the asymptotics inherited from $f$, which match the heuristic asymptotics from the earlier formal/heuristic solution.
[16.15.4] The second solution Now we show that the second solution

$$
v(x)=u(x) \cdot \int_{b}^{x} u(t)^{-2} d t
$$

to the original differential equation has the asymptotics given by the heuristic recursion, but with the opposite choice of $\pm \omega$ and $\pm \rho$. In terms of $f$,

$$
\begin{gathered}
v(x)=u(x) \cdot \int_{b}^{x} u(t)^{-2} d t=e^{\omega x} x^{-\rho} f(x) \int_{b}^{x} e^{-2 \omega t} x^{-2 \rho} f(t) d t \\
=e^{-\omega x} x^{\rho} f(x) \int_{b}^{x} e^{2 \omega(x-s)}\left(\frac{s}{x}\right)^{2 \rho} f(s)^{-2} d s
\end{gathered}
$$

This motivates taking

$$
g(x)=f(x) \int_{b}^{x} e^{2 \omega(x-s)}\left(\frac{s}{x}\right)^{2 \rho} f(s)^{-2} d s
$$

The lower bound $b$ has been chosen large enough so that $f(x)$ is bounded away from 0 for $x \geq b$. Since $f$ has an asymptotic expansion with leading coefficient 1 , it is elementary that there are coefficients $a_{n}$ so that $1 / f^{2}$ has asymptotics

$$
\frac{1}{f(x)^{2}}=1+\sum_{1 \leq n \leq N} \frac{a_{n}}{x^{n}}+O\left(\frac{1}{x^{N+1}}\right) \quad\left(\text { with } a_{0}=1\right)
$$

Then

$$
\frac{g(x)}{f(x)}=\sum_{0 \leq n \leq N} a_{n} \int_{b}^{x} e^{2 \omega(x-s)}\left(\frac{s}{x}\right)^{2 \rho} \frac{1}{s^{n}} d s+\int_{b}^{x} e^{2 \omega(x-s)}\left(\frac{s}{x}\right)^{-2 \rho} O\left(\frac{1}{s^{N+1}}\right) d s
$$

The last integral is $O\left(1 / x^{N+1}\right)$, from the elementary asymptotics. For each fixed $N$, the finitely-many integrals inside the summation have elementary asymptotics. Since for fixed $N$ there are only finitely-many such asymptotics, they are trivially uniform, so the asymptotics can be added. The asymptotic expansion for

$$
\int_{b}^{x} e^{2 \omega(x-s)}\left(\frac{s}{x}\right)^{2 \rho} \frac{1}{s^{n}} d s
$$

begins with $1 / x^{n}$, so the coefficient of each $1 / x^{n}$ is a finite sum, and there is no issue of convergence. Multiplying this asymptotic by that of $f(g)$ give the asymptotic expansion of $g(x)$.

As with $f$, the derivatives $g^{\prime}$ and $g^{\prime \prime}$ of $g$ have integral representations which yield asymptotic expansions. Thus, as in the appendix, the asymptotic expansion for $g$ can be twice differentiated term-wise to give those of $g^{\prime}$ and $g^{\prime \prime}$. Thus, their asymptotic expansions can be inserted in the differential equation. Their coefficients must satisfy the same recursion with some choice of $\pm \omega$ and corresponding $\pm \rho$. Arguing that the asymptotic for $g$ cannot be identical to that of $f$, we infer that the recursion for the coefficients of $g$ uses the opposite choice $-\omega,-\rho$ from the choice $\omega, \rho$ used to construct $f$.
[16.15.5] Remark: When $\omega$ and $\rho$ are both purely imaginary, $u$ and $v$ are bounded, neither approaches 0 , and they are uniquely determined up to constant factors. In all other cases, one solution approaches 0 , and is uniquely determined up to a constant, while the other is unbounded and ambiguous by multiples of the first, insofar as it depends on the choice of lower bound $b$ in the integral above.
[16.15.6] Remark: (Stokes' phenomenon) When the coefficient $q(x)$ of the differential equation $u^{\prime \prime}-q(x) u=$ 0 is analytic in a sector in $\mathbb{C}$, and when $q$ admits the same sort of asymptotic expansion

$$
\left.q\left(e^{i \theta} x\right) \sim \sum_{n \geq 0} \frac{q_{n} e^{-i n \theta}}{x^{n}} \quad \quad \quad \text { (uniformly in } \theta\right)
$$

in that sector, uniformly in the argument $\theta$, with $q_{o} \neq 0$, the above discussion still applies. In the realvariable discussion, with $\omega= \pm \sqrt{q_{o}}$, the case $\operatorname{Re} \omega=0$ was at the interface between the regimes $\operatorname{Re} \omega \leq 0$ and $\operatorname{Re} \omega \geq 0$ in which behaviors of solutions differed. Similarly, in the complex-variable situation the line $\operatorname{Re}\left(z \cdot \sqrt{q_{o}}\right)=0$ is the boundary between regimes of different behavior. On that line, the behavior is as in the $\operatorname{Re} \omega=0$ case. On either side of that line, one solution is exponentially larger than the other, etc. This is Stokes' phenomenon.

## 16.A Appendix: manipulation of asymptotic expansions

To say that $\varphi_{\ell}$ is an asymptotic sequence at $x_{o}$ means that $\varphi_{\ell+1}(x)=o\left(\varphi_{\ell}(x)\right)$ as $x \rightarrow x_{o}$, for all $\ell$. A function $f$ has an asymptotic expansion in terms of the $\varphi_{n}$, expressed with coefficients $c_{n}$ as

$$
f(x) \sim \sum_{n \geq 0} c_{n} \varphi_{n}
$$

when, for all $N \geq 0$,

$$
f(x)-\sum_{0 \leq n \leq N} c_{n} \varphi_{n}=o\left(\varphi_{N}\right)
$$

It is not surprising that a sum or integral of asymptotic expansions uniform in a parameter has the expected asymptotics. Circumstances under which an asymptotic expansion can be differentiated are more special.
[16.A.1] Summing asymptotic expansions Let functions $f_{n}$ have asymptotic expansions $f_{n} \sim$ $\sum_{\ell \geq 0} c_{n \ell} \varphi_{\ell}$, uniform in $n$, meaning that

$$
f_{n}(x)-\sum_{\ell \leq N} c_{n \ell} \varphi_{\ell}=o\left(\varphi_{N}\right) \quad \text { (implied constant and neighborhood of } x_{o} \text { uniform in } n \text { ) }
$$

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Let $a_{n}$ be coefficients such that $\sum_{n} a_{n} \cdot c_{n \ell}$ is convergent and $\sum_{n} a_{n}$ is absolutely convergent. We claim that $\sum_{n} a_{n} f_{n}$ converges in a neighborhood of $x_{o}$ and has the expected asymptotic expansion

$$
\sum_{n} a_{n} f_{n} \sim \sum_{\ell}\left(\sum_{n} a_{n} c_{n \ell}\right) \varphi_{n}
$$

The uniformity of the asymptotic expansions, and $\sum_{n}\left|a_{n}\right|<\infty$, give

$$
\left.\sum_{n \geq 1} a_{n}\left(f_{n}(x)-c_{n 1} \varphi_{1}(x)\right)=o\left(\varphi_{1}(x)\right) \quad \text { (uniformly in } x\right)
$$

In particular, the sum on the left-hand side converges for fixed $x$. Since $\sum_{n} a_{n} c_{n 1}$ converges, $\sum_{n>1} a_{n} f_{n}(x)$ converges. Similarly,

$$
\sum a_{n} f_{n}(x)-\sum_{\ell \leq N}\left(\sum_{n} a_{n} c_{n \ell}\right) \varphi_{\ell}=o\left(\varphi_{N}\right)
$$

[16.A.2] Integrals The general case is readily extrapolated from the example of an infinite sum. Namely, let $f(x, y) \sim \sum_{\ell} c_{\ell}(y) \varphi_{\ell}$ be asymptotic expansions uniform in a parameter $y \in Y$, where $Y$ is a measure space. Suppose that $y \rightarrow f(x, y)$ is measurable for each $x$, and that every $c_{\ell}(y)$ is measurable. Let $a(y)$ be absolutely integrable on $Y$, and assume that the integrals

$$
\int_{Y} a(y) c_{\ell}(y) d y
$$

converge for all $n$. Then

$$
\int_{Y} a(y) f(x, y) d y
$$

exists for $x$ close to $x_{o}$, and has asymptotic expansion

$$
\int_{Y} a(y) f(x, y) d y \sim \sum_{\ell}\left(\int_{Y} a(y) c_{\ell}(y) d y\right) \varphi_{\ell}
$$

[16.A.3] Differentiation of asymptotics in $1 / x^{n}$ Asymptotic power series are asymptotic expansions

$$
f(x) \sim c_{o}+\frac{c_{1}}{x}+\frac{c_{2}}{x^{2}}+\ldots \quad(\text { as } x \rightarrow+\infty)
$$

Unlike general situations, two such asymptotic expansions can be multiplied. A special property of asymptotic power series is the absolute integrability of $f(x)-c_{o}-c_{1} / x=O\left(x^{-2}\right)$ on intervals $[a,+\infty)$. Let

$$
F(x)=\int_{x}^{\infty}\left(f(t)-c_{o}-\frac{c_{1}}{t}\right) d t
$$

We claim that $F$ has an asymptotic expansion obtained from that of $f(x)-c_{o}-c_{1} / x$ by integrating termwise, namely,

$$
F(x) \sim \frac{c_{2}}{t}+\frac{c_{3}}{2 t^{2}}+\frac{c_{4}}{3 t^{3}}+\ldots
$$

To prove this, use

$$
f(x)-\left(c_{o}+\frac{c_{1}}{x}+\ldots+\frac{c_{N}}{x^{N}}\right)=O\left(x^{-(N+1)}\right)
$$

Then

$$
\begin{gathered}
F(x)-\left(\frac{c_{2}}{t}+\frac{c_{3}}{2 t^{2}}+\ldots+\frac{c_{N}}{N x^{N-1}}\right)=\int_{x}^{\infty}\left(f(t)-c_{o}-\frac{c_{1}}{t}\right) d t-\int_{x}^{\infty}\left(\frac{c_{2}}{t^{2}}+\frac{c_{3}}{t^{3}}+\ldots+\frac{c_{N}}{x^{N}}\right) d t \\
=\int_{x}^{\infty} O\left(t^{-(N+1)}\right) d t=O\left(x^{-N}\right)=o\left(x^{-(N-1)}\right)
\end{gathered}
$$

This has a surprising corollary about differentiation: for $f$ with an asymptotic power series at $+\infty$ as above, if $f$ is differentiable, and if $f^{\prime}$ has an asymptotic power series at $+\infty$, then the asymptotics of $f^{\prime}$ are obtained by differentiating that of $f$ termwise:

$$
f^{\prime}(x) \sim-\frac{c_{1}}{x^{2}}-\frac{2 c_{2}}{x^{3}}-\frac{3 c_{3}}{x^{4}}-\ldots
$$

When $f$ is holomorphic in a region in which the asymptotic holds uniformly in the argument of $x$, Cauchy's integral formula for $f^{\prime}$ produces an asymptotic for $f^{\prime}$ from that for $f$, thus avoiding the need to make a hypothesis that $f^{\prime}$ admits an asymptotic expansion.

## 16.B Appendix: ordinary points

The following discussion is well-known, although the convergence discussion is often omitted. This is the simpler case extended by the discussion of the regular singular points. A homogeneous ordinary differential equation of the form

$$
u^{\prime \prime}+b(x) u^{\prime}+c(x) u=0 \quad(\text { with } b, c \text { analytic near } 0)
$$

is said to have an ordinary point at 0 . The coefficients in a proposed expansion of the form

$$
u(x)=\sum_{n=0}^{\infty} a_{n} x^{n} \quad\left(\text { with } a_{0} \neq 0\right)
$$

are determined recursively from $a_{0}$ and $a_{1}$, as follows. The equation is

$$
\sum_{n=0}^{\infty} n(n-1) a_{n} x^{n-2}+b(x) \sum_{n=0}^{\infty} n a_{n} x^{n-1}+c(x) \sum_{n=0}^{\infty} a_{n} x^{n}=0
$$

or

$$
\sum_{n=0}^{\infty} n(n-1) a_{n} x^{n-2}+b(x) \sum_{n=0}^{\infty}(n-1) a_{n-1} x^{n-2}+c(x) \sum_{n=0}^{\infty} a_{n-2} x^{n-2}=0
$$

The coefficients $a_{n}$ with $n \geq 2$ are obtained recursively, from the expected

$$
n(n-1) \cdot a_{n}=\left(\text { in terms of } a_{0}, a_{1}, \ldots, a_{n-1}\right)
$$

To complete the proof of existence, we prove convergence. Take $A, M \geq 1$ large enough so that

$$
\begin{cases}b(x)=\sum_{n \geq 0} b_{n} x^{n} & \left(\text { with }\left|b_{n}\right| \leq A \cdot M^{n}\right) \\ c(x)=\sum_{n \geq 0} c_{n} x^{n} & \left(\text { with }\left|c_{n}\right| \leq A \cdot M^{n}\right)\end{cases}
$$

Inductively, suppose that $\left|a_{\ell}\right| \leq(C M)^{\ell}$, with a constant $C \geq 1$ to be determined in the following. Then
$n(n-1) \cdot\left|a_{n}\right| \leq A \sum_{i=1}^{n}(n-i) M^{i-1} \cdot(C M)^{n-i}+A \sum_{i=2}^{n} M^{i-2} \cdot(C M)^{n-i} \leq A M^{n-1} \cdot C^{n-1}\left(\frac{n(n+1)}{2}+n-1\right)$

Dividing through by $n(n-1)$, this is

$$
\left|a_{n}\right| \leq A M^{n-1} C^{n-1} \frac{n^{2}+3 n-2}{n(n-1)}
$$

This motivates taking

$$
C \geq A \sup _{2 \leq n \in \mathbb{Z}} \frac{n^{2}+3 n-2}{n(n-1)}
$$

which gives $\left|a_{n}\right| \leq(C M)^{n}$. In particular, for arbitrary $a_{0}$ and $a_{1}$ the resulting power series has a positive radius of convergence. In particular, these series can be differentiated termwise, by Abel's theorem.
[16.B.1] Ordinary points at infinity Let $u(x)=v(1 / x)$ and $z=1 / x$. Then

$$
u^{\prime}(x)=\frac{-1}{x^{2}} v^{\prime}(1 / x) \quad \text { and } \quad u^{\prime \prime}(x)=\frac{1}{x^{4}} v^{\prime \prime}(1 / x)+\frac{2}{x^{3}} v^{\prime}(1 / x)
$$

or

$$
u^{\prime}=-z^{2} v^{\prime} \quad \text { and } \quad u^{\prime \prime}=z^{4} v^{\prime \prime}+2 z^{3} v^{\prime} \quad(\text { with } u=u(x), v=v(z), z=1 / x)
$$

A differential equation $u^{\prime \prime}+b(x) u^{\prime}+c(x) u=0$ becomes

$$
\left(z^{4} v^{\prime \prime}+2 z^{3} v^{\prime}\right)+b(x)\left(-z^{2} v^{\prime}\right)+c(x) v=0
$$

or

$$
v^{\prime \prime}+\frac{2 z-b\left(\frac{1}{z}\right)}{z^{2}} v^{\prime}+\frac{c\left(\frac{1}{z}\right)}{z^{4}} v=0
$$

The point $z=0$ is an ordinary point when the coefficients of $v^{\prime}$ and $v$ are analytic at 0 . That is, $z=0$ is an ordinary point when $b, c$ have expansions at infinity of the form

$$
\left\{\begin{array}{l}
b\left(\frac{1}{z}\right)=2 z+b_{2} z^{2}+b_{3} z^{3} \ldots \\
c\left(\frac{1}{z}\right)=c_{4} z^{4}+c_{5} z^{5}+\ldots
\end{array}\right.
$$

[16.B.2] Not-quite-ordinary points Consider a differential equation with coefficients having poles of at most first order at 0 :

$$
u^{\prime \prime}+\frac{b(x)}{x} u^{\prime}+\frac{c(x)}{x} u=0
$$

with $b, c$ analytic at 0 . The coefficients in a proposed expansion of the form

$$
u(x)=\sum_{n=0}^{\infty} a_{n} x^{n} \quad\left(\text { with } a_{0} \neq 0\right)
$$

are determined recursively as follows. The equation is

$$
\sum_{n=0}^{\infty} n(n-1) a_{n} x^{n-2}+b(x) \sum_{n=0}^{\infty} n a_{n} x^{n-2}+c(x) \sum_{n=0}^{\infty} a_{n} x^{n-1}=0
$$

or

$$
\sum_{n=0}^{\infty} n(n-1) a_{n} x^{n-2}+b(x) \sum_{n=0}^{\infty} n a_{n} x^{n-2}+c(x) \sum_{n=0}^{\infty} a_{n-1} x^{n-2}=0
$$

We expect to determine the coefficients $a_{n}$ with $n \geq 2$ recursively, from

$$
(n(n-1)+b(0) n) \cdot a_{n}=\left(\text { in terms of } a_{0}, a_{1}, \ldots, a_{n-1}\right) \quad(\text { for } n \geq 1)
$$

For $b(0)$ not a non-positive integer, the recursion succeeds, and $a_{0}$ determines all the other coefficients $a_{n}$.
For $b(0)=0$, so that the coefficient of $v^{\prime}$ has no pole, the relation from the coefficient of $x^{-1}$,

$$
b(0) a_{1}+c(0) a_{0}=0
$$

implies that either $c(0)=0$ and the coefficient of $v$ has no pole, returning us to the ordinary-point case, or $a_{0}=0$, and there is no non-zero solution of this form.

For $b(0)$ a negative integer $-\ell$, the recursion for $a_{\ell}$ gives $a_{\ell}$ the coefficient 0 , and imposes a non-trivial relation on the prior coefficients $a_{n}$.

To complete the proof of existence, we prove convergence, assuming $b(0)$ is not a non-positive integer. Dividing through by a constant if necessary, we can take $M \geq 1$ large enough so that

$$
\begin{cases}b(x)=\sum_{n \geq 0} b_{n} x^{n} & \left(\text { with }\left|b_{n}\right| \leq M^{n}\right) \\ c(x)=\sum_{n \geq 0} c_{n} x^{n} & \left(\text { with }\left|c_{n}\right| \leq M^{n}\right)\end{cases}
$$

Inductively, suppose that $\left|a_{\ell}\right| \leq(C M)^{\ell}$, with a constant $C \geq 1$ to be determined in the following. Then

$$
(n(n-1)+b(0) n) \cdot\left|a_{n}\right|=\left|\sum_{i=1}^{n}(n-i) M^{i-1}(C M)^{n-i}+\sum_{i=1}^{n} M^{i-1}(C M)^{n-i}\right| \leq M^{n-1} C^{n-1}\left(\frac{n(n+1)}{2}+n\right)
$$

Dividing through by $n(n-1)+b(0) n$, this is

$$
\left|a_{n}\right| \leq M^{n-1} C^{n-1} \frac{n^{2}+3 n}{n(n-1)+b(0) n}
$$

This motivates taking

$$
C \geq \sup _{2 \leq n \in \mathbb{Z}} \frac{n^{2}+3 n}{n(n-1)+b(0) n}
$$

which gives $\left|a_{n}\right| \leq(C M)^{n}$. In particular, for arbitrary $a_{0}$ the resulting power series has a positive radius of convergence. For example, the series can be differentiated termwise, by Abel's theorem.

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[^0]:    [1] In slightly more sophisticated terms inessential to this discussion: the four examples $G$ immediately considered are real-rank one semi-simple Lie groups, and the discrete subgroups $\Gamma$ are unicuspidal in the sense that $\Gamma \backslash G / K$ is reasonably compactified by adding just a single cusp, where $K$ is a (maximal) compact subgroup of $G$. That is, the reduction theory of $\Gamma \backslash G$ is especially simple in these four cases. Examples with larger real rank, such as $G L_{n}$ with $n \geq 3$, will be considered later.

[^1]:    [14] There is no universal choice of normalization. Here, the choice is made so that the critical strip is $0 \leq \operatorname{Re}(s) \leq 1$, the rightmost pole is at $s=1$, and the functional equation relates $E_{s}$ and $E_{1-s}$. In more general contexts, other considerations dominate.

[^2]:    ${ }^{[16]}$ This $\widetilde{\Delta}_{b}$ is the Friedrichs self-adjoint extension [9.2] of the restriction of the unbounded operator $\Delta$ to the test functions $C_{c}^{\infty}(\Gamma \backslash G / K) \cap L_{b}^{2}(\Gamma \backslash G / K)$ in the space $L_{b}^{2}(\Gamma \backslash G / K)$ of pseudo-cuspforms.

[^3]:    [18] The potential conflict with $k$ being an element of a compact subgroup $K$ is avoidable only by other notational infelicities.

[^4]:    [21] This space $J_{\chi}$ is an instance of an induced representation, but we use no properties of such. Rather, the natural appearance of this function space explains attention to induced representations.

[^5]:    [25] Also, the terminology itself has a long and complicated history, and admits a much larger context, inessential to the present illustrative discussion.

[^6]:    [26] A more extensible form of the definition of Weyl group is as the normalizer of $M^{\text {min }}$ modulo the centralizer of $M^{\mathrm{min}}$, that is, monomial matrices (with one non-zero entry in each row and column) modulo diagonal matrices. However, for $G=G L_{r}$, it is very often convenient to fix a set of representatives in $G$.

[^7]:    [28] As with nomenclature for other objects, the terminology simple root is the correct name, but the origins, general definitions, and abstracted properties of simple roots would not help us here.

[^8]:    [30] Our description of what is needed to have cuspidal-data Eisenstein series be Hecke eigenfunctions would usually be the conclusion of a highly non-trivial chain of reasoning. That is, we have directly described what is needed to set $u p$ the proof that Eisenstein series formed from spherical Hecke-algebra eigenfunctions on Levi components are Hecke eigenfunctions. A more usual characterization, inherited from the chaotic historical order of developments, would be to require that the local representation generated by the cuspform be admissible and irreducible. Admissibility is equivalent to $K_{v}$-finiteness of every vector in $\pi_{v}$, and irreducibility has the usual meaning of having no closed $M_{v}^{P}$-stable subspaces, with the representation space suitably topologized. The Borel-Casselman-Matsumoto theorem [Borel 1976], [Casselman 1980], [Matsumoto 1977] asserts that $\pi_{v}$ is a subrepresentation of an unramified principal series. Given that, the key point is that induction in stages is legitimate, as in [6.9]. Unitariness of the representation, which follows from square-integrability of the cuspform, implies admissibility, from [Harish-Chandra 1970], for example. The global result, stated in [3.7] and proven in chapter 7, on discrete decomposition of cuspforms, in fact proves there is an orthogonal basis for everywhere-locally spherical square-integrable cuspforms generating irreducible representations of the global spherical Hecke algebra $C_{c}^{\infty}\left(K_{\mathbb{A}} \backslash G_{\mathbb{A}} / K_{\mathbb{A}}\right)$.

[^9]:    [34] Non-constant residues for $P^{n, n}$ are called Speh forms, since for $G L_{4}(\mathbb{R})$ the relevant unitary representations appear in [Speh 1981/2]. The general pattern for residual spectrum for $G L_{n}$ was conjectured in [Jacquet 1982/3] and proven in [Moeglin-Waldspurger 1989].
    [35] I first saw Bochner's lemma in the appendix of [Langlands 1967/76] treating minimal-parabolic Eisenstein series for $S L_{n}(\mathfrak{o})$ for rings of integers $\mathfrak{o}$. General accounts of several complex variables are [Bochner-Martin 1948] and [Hörmander 1973].

[^10]:    [36] When the group $G$ and the set $M$ are subsets of Euclidean spaces defined as zero sets or level sets of differentiable functions, differentiability of the action can be posed in the ambient Euclidean coordinates and the Implicit Function Theorem. In any particular example, even less is usually required to make sense of this requirement.
    [37] As in other instances of a group acting transitively on a set with additional structure, under modest hypotheses $M$ is a quotient $G_{o} \backslash G$ of $G$ by the isotropy group $G_{o}$ of a chosen point in $M$.

[^11]:    [38] The action of the Lie algebra by differentiating the action of the Lie group also applies abstractly to certain vectors $v$ in vectorspaces $V$ on which $G$ acts, namely, those $v$ such that $g \rightarrow g \cdot v$ is a differentiable $V$-valued function on $G$. Under mild hypotheses, smooth vectors are dense [14.6].

[^12]:    [43] Yes, the notation is somewhat in conflict with previous use of $X_{x}$ to denote the differential operator attached to $x \in \mathfrak{g}$.

[^13]:    ${ }^{[45]}$ The germ of a smooth function $f$ near a point $x_{o}$ on a smooth manifold $M$ is the equivalence class of $f$ under the equivalence relation $\sim$, where $f \sim g$ if $f, g$ are smooth functions defined on some neighborhoods of $x_{o}$, and which agree on some neighborhood of $x_{0}$. This is a construction, which does admit a more functional reformulation. That is, for each neighborhood $U$ of $x_{o}$, let $\mathcal{O}(U)$ be the ring of smooth functions on $U$, and for $U \supset V$ neighborhoods of $x_{o}$ let $\rho_{U V}: \mathcal{O}(U) \rightarrow \mathcal{O}(V)$ be the restriction map. Then the colimit $\operatorname{colim}_{U} \mathcal{O}(U)$ is exactly the ring of germs of smooth functions at $x_{o}$.
    ${ }^{[46]}$ This map $\varphi$ is presumably part of an atlas, meaning a maximal family of charts (homeomorphisms-to-image) $\varphi_{i}$ of opens $U_{i}$ in $M$ to subsets of a fixed $\mathbb{R}^{n}$, with the smooth manifold property that on overlaps things fit together smoothly, in the sense that

[^14]:    [51] The same ideas apply to $G$ a product $G=G_{\infty} \times G_{o}$ of a real Lie group $G_{\infty}$ and a totally disconnected group $G_{o}$, with $\Gamma$ a discrete subgroup.

[^15]:    [55] Happily, but not obviously, many important representations of such groups are admissible, which reduces the topological or analytical delicacy.
    [56] Some sources give the impression that a colimit of finite-dimensional spaces has no topology, or that topology is ignored. However, as noted here, there is a canonical topology, and every linear map from it to any topological vector space is continuous.

[^16]:    [58] The reasonable-ness of the compact-open topology is in its function. First, on a compact topological space $X$, the space $C^{o}(X)$ of continuous $\mathbb{C}$-valued functions with the sup-norm (of absolute value) is a Banach space. On non-compact $X$, the semi-norms given by sups of absolute values on compacts make $C^{o}(X)$ a Fréchet space. The compact-open topology accommodates spaces of continuous functions $C^{o}(X, Y)$ where the target space $Y$ is not a subset of a normed real or complex vector space, and is most interesting when $Y$ is a topological group. In the latter case, when the source $X$ is also a topological group, the subset of all continuous functions $f: X \rightarrow Y$ consisting of group homomorphisms is a (locally compact, Hausdorff) topological group, as proven below.

[^17]:    [61] This is a slightly abstracted version of a Sobolev norm, as in [9.5] and [9.7].

[^18]:    [67] When $R$ has a structure of ring or group that is reflected in its action on the vector space, the involution $r \rightarrow r^{*}$ should be an anti-automorphism, in the sense that $\left(r_{1} r_{2}\right)^{*}=r_{2}^{*} \cdot r_{1}^{*}$, since the adjoint map on continuous/bounded endomorphisms of a Hilbert space has that behavior.

[^19]:    [69] This is a Cartan involution.

[^20]:    [70] This trace form is a concrete instantiation of the Cartan-Killing form.

[^21]:    [73] In fact, such a diagram is a non-directed type of colimit, that is, with index set that is not a directed set.

[^22]:    [76] These subrepresentations have names, based on how they arose in other circumstances: are the sum of the holomorphic discrete series and anti-holomorphic discrete series representations.

[^23]:    [81] A topological group is usually understood to be locally compact and Hausdorff, and multiplication and inversion are continuous. To avoid measure-theoretic pathologies, a countable basis is often assumed. Perhaps oddly, the local compactness excludes most topological vector spaces.

[^24]:    [83] We have no immediate need of subtleties concerning functions of several complex variables, such as Hartogs' theorem that separate analyticity implies joint analyticity.

