HOPF ALGEBRAS IN COMBINATORICS

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Introduction

The concept of a Hopf algebra crystallized out of algebraic topology and the study of algebraic groups in the 1940s and 1950s (see [8] and [28] for its history). Being a fairly elementary algebraic notion itself, it subsequently found applications in other mathematical disciplines, and is now particularly commonplace in representation theory\(^1\).

These notes concern themselves (after a brief introduction into the algebraic foundations of Hopf algebra theory in Chapter 1) with the Hopf algebras that appear in combinatorics. These Hopf algebras tend to have bases naturally parametrized by combinatorial objects (partitions, compositions, permutations, tableaux, graphs, trees, posets, polytopes, etc.), and their Hopf-algebraic operations often encode basic operations on these objects\(^2\). Combinatorial results can then be seen as particular cases of general algebraic properties of Hopf algebras (e.g., the multiplicativity of the Möbius function can be recovered from the fact that the antipode of a Hopf algebra is an algebra anti-endomorphism), and many interesting invariants of combinatorial objects turn out to be evaluations of Hopf morphisms. In some cases (particularly that of symmetric functions), the rigidity in the structure of a Hopf algebra can lead to enlightening proofs.

One of the most elementary interesting examples of a combinatorial Hopf algebra is that of the symmetric functions. We will devote all of Chapter 2 to studying it, deviating from the usual treatments (such as in Stanley [154, Ch. 7], Sagan [138] and Macdonald [106]) by introducing the Hopf-algebraic structure early on and using it to obtain combinatorial results. Chapter 3 will underpin the importance of this algebra by proving Zelevinsky’s main theorem of PSH theory, which (roughly) claims that a Hopf algebra over \(\mathbb{Z}\) satisfying a certain set of axioms must be a tensor product of copies of the Hopf algebra of symmetric functions. These axioms are fairly restrictive, so this result is far from curtailing the diversity of combinatorial Hopf algebras; but they are natural enough that, as we will see in Chapter 4, they are satisfied for a Hopf algebra of representations of symmetric groups. As a consequence, this Hopf algebra will be revealed isomorphic to the symmetric functions – this is the famous Frobenius correspondence between symmetric functions and characters of symmetric groups, usually obtained through other ways ([53, §7.3], [138, §4.7]). We will further elaborate on the representation theories of wreath products and general linear groups over finite fields; while Zelevinsky’s PSH theory does not fully explain the latter, it illuminates it significantly.

In the next chapters, we will study further examples of combinatorial Hopf algebras: the quasisymmetric functions and the noncommutative symmetric functions in Chapter 5, various other algebras (of graphs, posets, matroids, etc.) in Chapter 7, and the Malvenuto-Reutenauer Hopf algebra of permutations in Chapter 8.

\(^{1}\)where it provides explanations for similarities between group representations and Lie algebra representations

\(^{2}\)such as concatenating two compositions, or taking the disjoint union of two graphs – but, more often, operations which return a multiset of results, such as cutting a composition into two pieces at all possible places, or partitioning a poset into two subposets in every way that satisfies a certain axiom
The main prerequisite for reading these notes is a good understanding of graduate algebra, in particular multilinear algebra (tensor products, symmetric powers and exterior powers) and basic categorical language. In Chapter 4, familiarity with representation theory of finite groups (over $\mathbb{C}$) is assumed, along with the theory of finite fields and (at some places) the rational canonical form of a matrix. Only basic knowledge of combinatorics is required (except for a few spots in Chapter 7), and familiarity with geometry and topology is needed only to understand some tangential remarks. The concepts of Hopf algebras and coalgebras and the basics of symmetric function theory will be introduced as needed. We will work over a commutative base ring most of the time, but no commutative algebra (besides, occasionally, properties of modules over a PID) will be used.

These notes began as an accompanying text for Fall 2012 Math 8680 Topics in Combinatorics, a graduate class taught by the second author at the University of Minnesota. The first author has since added many exercises (and solutions) as well as Chapter 6 on Lyndon words and the polynomiality of QSym. The notes might still grow, and any comments, corrections and complaints are welcome!

The course was an attempt to focus on examples that we find interesting, but which are hard to find fully explained currently in books or in one paper. Much of the subject of combinatorial Hopf algebras is fairly recent (1990s onwards) and still spread over research papers, although sets of lecture notes do exist, such as Foissy’s [50]. A reference which we discovered late, having a great deal of overlap with these notes is Hazewinkel, Gubareni, and Kirichenko [68]. References for the purely algebraic theory of Hopf algebras are much more frequent (see the beginning of Chapter 1 for a list).

Be warned that our notes are highly idiosyncratic in choice of topics, and they steal heavily from the sources in the bibliography.

Warnings: Unless otherwise specified ...

- $k$ here usually denotes a commutative ring.
- all maps between $k$-modules are $k$-linear.
- every ring or $k$-algebra is associative and has a 1, and every ring morphism or $k$-algebra morphism preserves the 1’s.
- all $k$-algebras $A$ have the property that $(\lambda 1_A) a = a (\lambda 1_A) = \lambda a$ for all $\lambda \in k$ and $a \in A$.
- all tensor products are over $k$.
- $I$ will denote the multiplicative identity in some ring like $k$ or in some $k$-algebra (sometimes also the identity of a group written multiplicatively).
- for any set $S$, we denote by $id_S$ (or by $id$) the identity map on $S$.
- The symbols $\subset$ (for “subset”) and $<$ (for “subgroup”) don’t imply properness (so $\mathbb{Z} \subset \mathbb{Z}$ and $\mathbb{Z} < \mathbb{Z}$).
- the $n$-th symmetric group (i.e., the group of all permutations of $\{1, 2, \ldots, n\}$) is denoted $\mathfrak{S}_n$.
- The product of permutations $a \in \mathfrak{S}_n$ and $b \in \mathfrak{S}_n$ is defined by $(ab)(i) = a(b(i))$ for all $i$.
- Words over (or in) an alphabet $I$ simply mean finite tuples of elements of a set $I$. It is customary to write such a word $(a_1, a_2, \ldots, a_k)$ as $a_1 a_2 \ldots a_k$ when this is not likely to be confused for multiplication.
- $\mathbb{N} := \{0, 1, 2, \ldots\}$.
- if $i$ and $j$ are any two objects, then $\delta_{i,j}$ denotes the Kronecker delta of $i$ and $j$; this is the integer 1 if $i = j$ and 0 otherwise.

Hopefully context will resolve some of the ambiguities.

1. What is a Hopf algebra?

The standard references for Hopf algebras are Abe [1] and Sweedler [160], and some other good ones are [27, 29, 37, 85, 117, 168, 68]. See also Foissy [50] and Manchon [113] for introductions to Hopf algebras.

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3William Schmitt’s expositions [143] are tailored to a reader interested in combinatorial Hopf algebras; his notes on modules and algebras cover a significant part of what we need from abstract algebra, whereas those on categories cover all category theory we will use and much more.

4Keith Conrad’s expository notes [32] are useful, even if not comprehensive, sources for the latter.

5The version of the notes you are reading does not contain said solutions. The version that does can be downloaded from http://web.mit.edu/~darij/www/algebra/HopfComb-sols.pdf or compiled from the source code.

6As explained below, “ring” means “associative ring with 1”. The most important cases are when $k$ is a field or when $k = \mathbb{Z}$.
tailored to combinatorial applications. Most texts only study Hopf algebras over fields (with exceptions such as [29, 27, 168]). We will work over arbitrary commutative rings\(^7\), which requires some more care at certain points (but we will not go deep enough into the algebraic theory to witness the situation over commutative rings diverge seriously from that over fields).

Let’s build up the definition of Hopf algebra structure bit-by-bit, starting with the more familiar definition of algebras.

1.1. Algebras. The following definition of \(k\)-algebras may look unfamiliar, but it is merely a restatement of their classical definition using tensors\(^8\) and \(k\)-linear maps:

**Definition 1.1.** An associative \(k\)-algebra \(A\) is a \(k\)-module with a \(k\)-linear associative operation \(A \otimes A \rightarrow A\), and a \(k\)-linear unit \(k \rightarrow A\) sending 1 in \(k\) to the two-sided multiplicative identity element 1 in \(A\). One can rephrase this by saying that these diagrams commute:

\[
\begin{align*}
&\begin{array}{c}
\begin{array}{ccc}
A \otimes A & \xrightarrow{m \otimes id} & A \\
\downarrow{m} & & \downarrow{id \otimes m}
\end{array}
\end{array} \\
&\begin{array}{ccc}
A \otimes A & \xrightarrow{id} & A \\
\downarrow{id} & & \downarrow{id}
\end{array}
\end{align*}
\]

where the maps \(A \rightarrow A \otimes k\) and \(A \rightarrow k \otimes A\) are the isomorphisms sending \(a \mapsto a \otimes 1\) and \(a \mapsto 1 \otimes a\).

We abbreviate “associative \(k\)-algebra” as “\(k\)-algebra” (associativity is assumed unless otherwise specified) or as “algebra” (when \(k\) is clear from the context).

Well-known examples of \(k\)-algebras are tensor and symmetric algebras, which we can think of as algebras of words and multisets, respectively.

**Example 1.2.** If \(V\) is a \(k\)-module and \(n \in \mathbb{N}\), then the \(n\)-fold tensor power \(V^\otimes n\) of \(V\) is the \(k\)-module \(\bigotimes_{n \geq 0} V^\otimes n\). (For \(n = 0\), this is the \(k\)-module \(k\), spanned by the “empty tensor” \(1_k\).)

The tensor algebra \(T(V) = \bigoplus_{n \geq 0} V^\otimes n\) on a \(k\)-module \(V\) is an associative \(k\)-algebra spanned (as \(k\)-module) by decomposable tensors \(v_1 v_2 \cdots v_k := v_1 \otimes v_2 \otimes \cdots \otimes v_k\) with \(k \in \mathbb{N}\) and \(v_1, v_2, \ldots, v_k \in V\). Its multiplication is defined \(k\)-linearly by

\[
m(v_1 v_2 \cdots v_k \otimes w_1 w_2 \cdots w_\ell) := v_1 v_2 \cdots v_k w_1 w_2 \cdots w_\ell
\]

for all \(k, \ell \in \mathbb{N}\) and \(v_1, v_2, \ldots, v_k, w_1, w_2, \ldots, w_\ell \in V\). The unit map \(u : k \rightarrow T(V)\) sends \(1_k\) to the empty tensor \(1_{T(V)} = 1_k \in k = V^\otimes 0\).

---

\(^7\)and we will profit from this generality in Chapters 3 and 4, where we will be applying the theory of Hopf algebras to \(k = \mathbb{Z}\) in a way that would not be possible over \(k = \mathbb{Q}\).

\(^8\)i.e., the multiplication is encoded by a \(k\)-linear map \(A \otimes A \rightarrow A\) instead of a \(k\)-bilinear map \(A \times A \rightarrow A\).

\(^9\)Some remarks about our notation (which we are using here and throughout these notes) are in order. Since we are working with tensor products of \(k\)-modules like \(T(V)\) – which themselves are made of tensors – here, we must specify what the \(\otimes\) sign means in expressions like \(a \otimes b\) where \(a\) and \(b\) are elements of \(T(V)\). Our convention is the following: When \(a\) and \(b\) are elements of a tensor algebra \(T(V)\), we always understand \(a \otimes b\) to mean the pure tensor \(a \otimes b \in T(V) \otimes T(V)\) rather than the product of \(a\) and \(b\) inside the tensor algebra \(T(V)\). The latter product will plainly be written \(a b\).

The operator precedence between \(\otimes\) and multiplication in \(T(V)\) is such that multiplication in \(T(V)\) binds more tightly than the \(\otimes\) sign; e.g., the term \(ab \otimes cd\) means \((ab) \otimes (cd)\).
If \( V \) is a free \( k \)-module, say with \( k \)-basis \( \{x_i\}_{i \in I} \), then \( T(V) \) has a \( k \)-basis of decomposable tensors \( x_{i_1} \cdots x_{i_k} := x_{i_1} \otimes \cdots \otimes x_{i_k} \) indexed by words \( (i_1, \ldots, i_k) \) in the alphabet \( I \), and the multiplication on this basis is given by concatenation of words:

\[
m(x_{i_1} \cdots x_{i_k} \otimes x_{j_1} \cdots x_{j_l}) = x_{i_1} \cdots x_{i_k} x_{j_1} \cdots x_{j_l}.
\]

Recall that in an algebra \( A \), when one has a two-sided ideal \( J \subset A \), meaning a \( k \)-submodule with \( m(J \otimes A), m(A \otimes J) \subset J \), then one can form a quotient algebra \( A/J \).

**Example 1.3.** The symmetric algebra \( \text{Sym}(V) = \bigoplus_{n \geq 0} \text{Sym}^n(V) \) is the quotient of \( T(V) \) by the two-sided ideal generated by all elements \( xy - yx \) with \( x, y \in V \). When \( V \) is a free \( k \)-module with basis \( \{x_i\}_{i \in I} \), this symmetric algebra \( S(V) \) can be identified with a (commutative) polynomial algebra \( k[x_i]_{i \in I} \), having a \( k \)-basis of (commutative) monomials \( x_{i_1} \cdots x_{i_k} \) as \( (i_1, \ldots, i_k) \) runs through all finite multisubsets \(^{10}\) of \( I \), and with multiplication defined \( k \)-linearly via multiset union.

Note that the \( k \)-module \( k \) itself canonically becomes a \( k \)-algebra. Its associative operation \( m : k \otimes k \rightarrow k \) is the canonical isomorphism \( k \otimes k \rightarrow k \), and its unit \( u : k \rightarrow k \) is the identity map.

Topology and group theory give more examples.

**Example 1.4.** The cohomology algebra \( H^*(X; k) = \bigoplus_{i \geq 0} H^i(X; k) \) with coefficients in \( k \) for a topological space \( X \) has an associative cup product. Its unit \( k = H^0(\text{pt}; k) \xrightarrow{pt} H^*(X; k) \) is induced from the unique (continuous) map \( X \rightarrow \text{pt} \), where \( \text{pt} \) is a one-point space.

**Example 1.5.** For a group \( G \), the group algebra \( kG \) has \( k \)-basis \( \{t_g\}_{g \in G} \) and multiplication defined \( k \)-linearly by \( t_g t_h = t_{gh} \), and unit defined by \( u(1) = t_e \), where \( e \) is the identity element of \( G \).

1.2. **Coalgebras.** If we are to think of the multiplication \( A \otimes A \rightarrow A \) in an algebra as putting together two basis elements of \( A \) to get a sum of basis elements of \( A \), then coalgebra structure should be thought of as taking basis elements apart.

**Definition 1.6.** A co-associative \( k \)-coalgebra \( C \) is a \( k \)-module \( C \) with a comultiplication, that is, a \( k \)-linear map \( C \xrightarrow{\Delta} C \otimes C \), and a \( k \)-linear counit \( C \xrightarrow{\epsilon} k \) making commutative the diagrams as in (1.1), (1.2) but with all arrows reversed:

(1.3)

\[
\begin{array}{ccc}
C \otimes C & \xrightarrow{\Delta \otimes \text{id}} & C \otimes C \otimes C \\
\text{id} \otimes \Delta & & \Delta \otimes \text{id} \\
\downarrow \Delta & & \downarrow \Delta \\
C \otimes C & \xleftarrow{\Delta} & C \\
\end{array}
\]

(1.4)

\[
\begin{array}{ccc}
C \otimes k & \xrightarrow{\text{id} \otimes \epsilon} & C \\
\downarrow \epsilon \otimes \text{id} & & \downarrow \epsilon \otimes \text{id} \\
C \otimes C & \xleftarrow{\epsilon} & C \otimes C \\
\end{array}
\]

Here the maps \( C \otimes k \rightarrow C \) and \( k \otimes C \rightarrow C \) are the isomorphisms sending \( c \otimes 1 \mapsto c \) and \( 1 \otimes c \mapsto c \).

We abbreviate “co-associative \( k \)-coalgebra” as “\( k \)-coalgebra” (co-associativity, i.e., the commutativity of the diagram (1.3), is assumed unless otherwise specified) or as “coalgebra” (when \( k \) is clear from the context). Sometimes, the word “coproduct” is used as a synonym for “comultiplication”\(^{11}\).

One often uses the Sweedler notation

\[ \Delta(c) = \sum_{(c)} c_1 \otimes c_2 = \sum_{(c)} c_1 \otimes c_2 \]

\(^{10}\)By a multiset of a set \( S \), we mean a multiset each of whose elements belongs to \( S \) (but can appear arbitrarily often).

\(^{11}\)although the word “coproduct” already has a different meaning in algebra
to abbreviate formulas involving $\Delta$. For example, commutativity of the left square in (1.4) asserts that
$$\sum_{(i)} c_i \epsilon(c_2) = c. \quad (\text{Example 1.7})$$

The $k$-module $k$ itself canonically becomes a $k$-coalgebra, with its comultiplication $\Delta : k \to k \otimes k$ being the canonical isomorphism $k \to k \otimes k$, and its counit $\epsilon : k \to k$ being the identity map.

**Example 1.7.** Let $k$ be a field. The homology $H_*(X;k) = \bigoplus_{i \geq 0} H_i(X;k)$ for a topological space $X$ is naturally a coalgebra: the (continuous) diagonal embedding $X \to X \times X$ sending $x \mapsto (x,x)$ induces a coassociative map

$$H_*(X;k) \to H_*(X \times X;k) \cong H_*(X;k) \otimes H_*(X;k)$$

in which the last isomorphism comes from the *K"unneth theorem* with field coefficients $k$. As before, the unique (continuous) map $X \to pt$ induces the counit $H_*(X;k) \xrightarrow{\epsilon} H_*(pt;k) \cong k$.

**Exercise 1.8.** Given a $k$-module $C$ and a $k$-linear map $\Delta : C \to C \otimes C$. Prove that there exists at most one $k$-linear map $\epsilon : C \to k$ such that the diagram (1.4) commutes.

1.3. Morphisms, tensor products, and bialgebras.

**Definition 1.9.** A morphism of algebras $A \xrightarrow{\varphi} B$ makes these diagrams commute:

$$A \xrightarrow{m_A} A \otimes A \xrightarrow{\varphi \otimes \varphi} B \otimes B \xleftarrow{m_B} B$$

Here the subscripts on $m_A, m_B, u_A, u_B$ indicate for which algebra they are part of the structure— we will occasionally use such conventions from now on.

Similarly a morphism of coalgebras is a $k$-linear map $C \xrightarrow{\varphi} D$ making the reverse diagrams commute:

$$C \xleftarrow{\Delta_C} C \otimes C \xrightarrow{\varphi \otimes \varphi} D \otimes D \xrightarrow{(\epsilon \otimes \epsilon)} k$$

**Example 1.10.** Let $k$ be a field. Continuous maps $X \xrightarrow{f} Y$ of topological spaces induce algebra morphisms $H^*(Y;k) \to H^*(X;k)$, and coalgebra morphisms $H_*(X;k) \to H_*(Y;k)$.

**Definition 1.11.** Given two $k$-algebras $A, B$, their tensor product $A \otimes B$ also becomes a $k$-algebra defining the multiplication bilinearly via

$$m((a \otimes b) \otimes (a' \otimes b')) := aa' \otimes bb'$$

or in other words $m_{A \otimes B}$ is the composite map

$$A \otimes B \otimes A \otimes B \xrightarrow{id \otimes T \otimes id} A \otimes A \otimes B \otimes B \xrightarrow{m_A \otimes m_B} A \otimes B$$

where $T$ is the *twist* map $B \otimes A \to A \otimes B$ that sends $b \otimes a \mapsto a \otimes b$.

Here we are omitting the topologist’s sign in the twist map which should be present for graded algebras and coalgebras that come from cohomology and homology: for homogeneous elements $a$ and $b$ the topologist’s twist map sends

$$b \otimes a \mapsto (-1)^{\deg(a) \deg(b)} a \otimes b.$$

This means that, if one is using the topologists’ conventions, most of our examples which we later call *graded* should actually be considered to live in only *even* degrees (which can be achieved, e.g., by artificially doubling their grading). We will, however, keep to our own definitions (so the twist map sends $b \otimes a \mapsto a \otimes b$) unless otherwise noted.

The unit element of $A \otimes B$ is $1_A \otimes 1_B$, meaning that the unit map $k \xrightarrow{u_{A \otimes B}} A \otimes B$ is the composite

$$k \xrightarrow{u_A \otimes u_B} A \otimes B.$$
Similarly, given two coalgebras $C, D$, one can make $C \otimes D$ a coalgebra in which the comultiplication and counit maps are the composites of
\[
C \otimes D \xrightarrow{\Delta_C \otimes \Delta_D} C \otimes C \otimes D \otimes D \xrightarrow{id \otimes T \otimes id} C \otimes D \otimes C \otimes D
\]
and
\[
C \otimes D \xrightarrow{\epsilon_C \otimes \epsilon_D} k \otimes k \xrightarrow{\epsilon} k.
\]

One of the first signs that these definitions interact nicely is the following straightforward proposition.

**Proposition 1.12.** When $A$ is both a $k$-algebra and a $k$-coalgebra, the following are equivalent:

- $(\Delta, \epsilon)$ are morphisms for the algebra structure $(m, u)$.
- $(m, u)$ are morphisms for the coalgebra structure $(\Delta, \epsilon)$.
- These four diagrams commute:

\[
\begin{array}{ccc}
A \otimes A & \xrightarrow{\Delta \otimes \Delta} & A \otimes A \\
\downarrow m & & \downarrow m \\
A \otimes A & \xrightarrow{\Delta} & A \\
\downarrow m \otimes m & & \downarrow m \otimes m \\
A & \xrightarrow{\epsilon \otimes \epsilon} & k \otimes k \\
\downarrow m & & \downarrow m \\
A & \xrightarrow{\epsilon} & k
\end{array}
\]

\[
\begin{array}{ccc}
A \otimes A & \xrightarrow{\epsilon \otimes m} & k \otimes k \\
\downarrow m & & \downarrow m \\
A & \xrightarrow{\epsilon} & k
\end{array}
\]

\[
\begin{array}{ccc}
k & \xrightarrow{u} & A \\
\downarrow \Delta & & \downarrow \Delta \\
k \otimes k & \xrightarrow{u \otimes u} & A \otimes A
\end{array}
\]

**Exercise 1.13.**

(a) If $A, A', B$ and $B'$ are four $k$-algebras, and $f : A \to A'$ and $g : B \to B'$ are two $k$-algebra homomorphisms, then show that $f \otimes g : A \otimes B \to A' \otimes B'$ is a $k$-algebra homomorphism.

(b) If $C, C', D$ and $D'$ are four $k$-coalgebras, and $f : C \to C'$ and $g : D \to D'$ are two $k$-coalgebra homomorphisms, then show that $f \otimes g : C \otimes D \to C' \otimes D'$ is a $k$-coalgebra homomorphism.

**Definition 1.14.** Call the $k$-module $A$ a $k$-bialgebra if it is a $k$-algebra and $k$-coalgebra satisfying the three equivalent conditions in Proposition 1.12.

**Example 1.15.** For a group $G$, one can make the group algebra $kG$ a coalgebra with counit $kG \xrightarrow{\epsilon} k$ mapping $t_g \mapsto 1$ for all $g$ in $G$, and with comultiplication $kG \xrightarrow{\Delta} kG \otimes kG$ given by $\Delta(t_g) := t_g \otimes t_g$. Checking the various diagrams in (1.8) commute is easy. For example, one can check the pentagonal diagram on each
basis element $t_g \otimes t_h$:

![Diagram](image)

**Remark 1.16.** In fact, one can think of adding a bialgebra structure to a $k$-algebra $A$ as a way of making $A$-modules $M, N$ have an $A$-module structure on their tensor product $M \otimes N$: the algebra $A \otimes A$ already acts naturally on $M \otimes N$, so one can let $a$ in $A$ act via $\Delta(a)$ in $A \otimes A$. In the theory of group representations over $k$, that is, $kG$-modules $M$, this is how one defines the diagonal action of $G$ on $M \otimes N$, namely $t_g$ acts as $t_g \otimes t_g$.

**Definition 1.17.** An element $x$ in a coalgebra for which $\Delta(x) = x \otimes x$ and $\epsilon(x) = 1$ is called group-like.

An element $x$ in a coalgebra for which $\Delta(x) = 1 \otimes x + x \otimes 1$ is called primitive.

**Example 1.18.** The tensor algebra $T(V) = \bigoplus_{n \geq 0} V^\otimes n$ is a coalgebra, with counit $\epsilon$ equal to the identity on $V^\otimes 0 = k$ and the zero map on $V^\otimes n$ for $n > 0$, and with comultiplication defined to make the elements $x$ in $V^\otimes 1 = V$ all primitive:

$$\Delta(x) := 1 \otimes x + x \otimes 1 \text{ for } x \in V^\otimes 1.$$ 

Since the elements of $V$ generate $T(V)$ as a $k$-algebra, and since $T(V) \otimes T(V)$ is also an associative $k$-algebra, the universal property of $T(V)$ as the free associative $k$-algebra on the generators $V$ allows one to define $T(V) \xrightarrow{\Delta} T(V) \otimes T(V)$ arbitrarily on $V$, and extend it as an algebra morphism.

It may not be obvious that this $\Delta$ is coassociative, but one can note that

$$((\text{id} \otimes \Delta) \circ \Delta)(x) = x \otimes 1 \otimes 1 + 1 \otimes x \otimes 1 + 1 \otimes 1 \otimes x = ((\Delta \otimes \text{id}) \circ \Delta)(x)$$

for every $x$ in $V$. Hence the two maps $(\text{id} \otimes \Delta) \circ \Delta$ and $(\Delta \otimes \text{id}) \circ \Delta$, considered as algebra morphisms $T(V) \to T(V) \otimes T(V) \otimes T(V)$, must coincide on every element of $T(V)$ since they coincide on $V$. We leave it as an exercise to check the map $\epsilon$ defined as above satisfies the counit axioms (1.4).

Here is a sample calculation in $T(V)$ when $V$ has basis $\{x, y, z\}$:

$$\Delta(xyz) = \Delta(x)\Delta(y)\Delta(z)$$

$$= (1 \otimes x + x \otimes 1)(1 \otimes y + y \otimes 1)(1 \otimes z + z \otimes 1)$$

$$= (1 \otimes xy + x \otimes y + y \otimes x + xy \otimes 1)(1 \otimes z + z \otimes 1)$$

$$= 1 \otimes xyz + x \otimes yz + y \otimes zx + z \otimes xy$$

$$+ xy \otimes z + xz \otimes y + yz \otimes x + xyz \otimes 1$$

This illustrates the idea that comultiplication “takes basis elements apart”. Here for any $v_1, v_2, \ldots, v_n$ in $V$ one has

$$\Delta(v_1 v_2 \cdots v_n) = \sum v_{j_1} \cdots v_{j_r} \otimes v_{k_1} \cdots v_{k_{n-r}}$$

where the sum is over ordered pairs $(j_1, j_2, \ldots, j_r), (k_1, k_2, \ldots, k_{n-r})$ of complementary subwords of the word $(1, 2, \ldots, n)$.  

Equivalently (and in a more familiar language),

$$\Delta(v_1 v_2 \cdots v_n) = \sum_{I \subseteq \{1, 2, \ldots, n\}} v_I \otimes v_{(1, 2, \ldots, n) \setminus I},$$

\footnote{More formally speaking, the sum is over all permutations $(j_1, j_2, \ldots, j_r, k_1, k_2, \ldots, k_{n-r})$ of $(1, 2, \ldots, n)$ satisfying $j_1 < j_2 < \cdots < j_r$ and $k_1 < k_2 < \cdots < k_{n-r}$.}
where $v_J$ (for $J$ a subset of $\{1, 2, \ldots, n\}$) denotes the product of all $v_j$ with $j \in J$ in the order of increasing $j$.

Recall one can quotient a $k$-algebra $A$ by a two-sided ideal $J$ to obtain a quotient algebra $A/J$.

**Definition 1.19.** In a coalgebra $C$, a *two-sided coideal* is a $k$-submodule $J \subset C$ for which

$$
\Delta(J) \subset J \otimes C + C \otimes J
$$

$$
\epsilon(J) = 0
$$

The quotient $k$-module $C/J$ then inherits a coalgebra structure. Similarly, in a bialgebra $A$, a subset $J \subset A$ which is both a two-sided ideal and two-sided coideal gives rise to a quotient bialgebra $A/J$.

**Exercise 1.20.** Let $A$ and $C$ be two $k$-coalgebras, and $f : A \to C$ a surjective coalgebra homomorphism.

(a) If $f$ is surjective, then show that ker $f$ is a two-sided coideal of $A$.

(b) If $k$ is a field, then show that ker $f$ is a two-sided coideal of $A$.

**Example 1.21.** The symmetric algebra $\text{Sym}(V)$ was the quotient of $T(V)$ by the two-sided ideal $J$ generated by all commutators $[x, y] = xy - yx$ for $x, y$ in $V$. Note that $x, y$ are primitive elements in $T(V)$, and the following very reusable calculation shows that *the commutator of two primitives is primitive*:

$$
\Delta[x, y] = \Delta(xy - yx)
= (1 \otimes x + x \otimes 1)(1 \otimes y + y \otimes 1) - (1 \otimes y + y \otimes 1)(1 \otimes x + x \otimes 1)
= 1 \otimes xy - 1 \otimes yx + xy \otimes 1 - yx \otimes 1
+ x \otimes y + y \otimes x - x \otimes y - y \otimes x
= 1 \otimes (xy - yx) + (xy - yx) \otimes 1
= 1 \otimes [x, y] + [x, y] \otimes 1.
$$

(1.10)

In particular, the commutators $[x, y]$ have $\Delta[x, y]$ in $J \otimes T(V) + T(V) \otimes J$. They also satisfy $\epsilon([x, y]) = 0$. Since they are generators for $J$ as a two-sided ideal, it is not hard to see this implies $\Delta(J) \subset J \otimes T(V) + T(V) \otimes J$, and $\epsilon(J) = 0$. Thus $J$ is also a two-sided coideal, and $\text{Sym}(V) = T(V)/J$ inherits a bialgebra structure.

In fact we will see in Section 3.1 that symmetric algebras are the universal example of bialgebras which are *graded, connected, commutative, cocommutative*. But first we should define some of these concepts.

**Definition 1.22.** A *graded $k$-module* $V$ is one with a $k$-module direct sum decomposition $V = \bigoplus_{n \geq 0} V_n$. Elements $x$ in $V_n$ are called *homogeneous* of degree $n$, or $\text{deg}(x) = n$.

One endows tensor products $V \otimes W$ of graded $k$-modules $V, W$ with graded module structure in which $(V \otimes W)_n := \bigoplus_{i+j=n} V_i \otimes W_j$.

A $k$-linear map $V \xrightarrow{\varphi} W$ between two graded $k$-modules is called *graded* if $\varphi(V_n) \subset W_n$ for all $n$. Say that a $k$-algebra (coalgebra, bialgebra) is *graded* if it is a graded $k$-module and all of the relevant structure maps $(u, \epsilon, m, \Delta)$ are graded.

Say that a graded module $V$ is *connected* if $V_0 \cong k$.

**Example 1.23.** Let $k$ be a field. A path-connected space $X$ has its homology and cohomology

$$
H_\ast(X; k) = \bigoplus_{i \geq 0} H_i(X; k)
$$

$$
H^\ast(X; k) = \bigoplus_{i \geq 0} H^i(X; k)
$$

carrying the structure of graded connected coalgebras and algebras, respectively. If in addition, $X$ is a topological group, or even less strongly, a *homotopy-associative* $H$-space (e.g. the loop space $\Omega Y$ on some other space $Y$), the continuous multiplication map $X \times X \to X$ induces an algebra structure on $H_\ast(X; k)$ and a coalgebra structure on $H^\ast(X; k)$, so that each become bialgebras in the topologist’s sense (i.e., with the twist as in (1.7)), and these bialgebras are dual to each other in a sense soon to be discussed. This was Hopf’s motivation: the (co-)homology of a compact Lie group carries bialgebra structure that explains why it takes a certain form; see Cartier [28, §2].

\[ \text{since } J \otimes C + C \otimes J \text{ is contained in the kernel of the canonical map } C \otimes C \to (C/J) \otimes (C/J) \]
**Example 1.24.** Tensor algebras $T(V)$ and symmetric algebras $\text{Sym}(V)$ are graded, once one picks a graded module structure for $V$; then

$$\deg(x_{i_1} \cdots x_{i_k}) = \deg(x_{i_1}) + \cdots + \deg(x_{i_k})$$

if $\{x_i\}_{i \in I}$ is a graded basis (and, more generally, $v_1 v_2 \cdots v_k$ is homogeneous of degree $i_1 + i_2 + \cdots + i_k$ if each $v_j$ is a homogeneous element of $V$ of degree $i_j$). Assuming that $V_0 = 0$, the algebras $T(V)$ and $\text{Sym}(V)$ are connected. For example, we will often say that all elements of $V$ are homogeneous of degree 1, but at other times, it will make sense to have $V$ live in different (positive) degrees.

**Exercise 1.25.** Let $A = \bigoplus_{n \geq 0} A_n$ be a graded $k$-bialgebra. We denote by $p$ the set of all primitive elements of $A$.

(a) Show that $p$ is a graded $k$-submodule of $A$ (that is, we have $p = \bigoplus_{n \geq 0} (p \cap A_n)$).

(b) Show that $p$ is a two-sided coideal of $A$.

**Exercise 1.26.** Check that for a graded connected $k$-bialgebra $A$, we have that

(a) the $k$-submodule $k = k \cdot 1_A$ of $A$ lies in $A_0$,

(b) $u$ is an isomorphism $k \xrightarrow{u} A_0$, while

(c) the two-sided ideal $\ker \epsilon$ is the $k$-module of positive degree elements $I = \bigoplus_{n > 0} A_n$.

(d) $\epsilon$ restricted to $A_0$ is the inverse isomorphism $A_0 \xrightarrow{\epsilon} k$ to $u$.

(e) for every $x \in A$, we have

$$\Delta(x) \in x \otimes 1 + A \otimes I;$$

(f) every $x$ in $I$ has comultiplication of the form

$$\Delta(x) = 1 \otimes x + x \otimes 1 + \Delta_+(x)$$

where $\Delta_+(x)$ lies in $I \otimes I$.

(Use only the gradedness of the unit $u$ and counit $\epsilon$ maps, along with commutativity of diagrams (1.4), and (1.8) and the connectedness of $A$.)

1.4. **Antipodes and Hopf algebras.** There is one more piece of structure needed to make a bialgebra a Hopf algebra, although it will come for free in the graded connected case.

**Definition 1.27.** For any coalgebra $C$ and algebra $A$, one can endow the $k$-linear maps $\text{Hom}(C, A)$ with an associative algebra structure called the convolution algebra: send $f, g$ in $\text{Hom}(C, A)$ to $f \star g$ defined by $(f \star g)(c) = \sum f(c_1)g(c_2)$, using the Sweedler notation $\Delta(c) = \sum c_1 \otimes c_2$. Equivalently, $f \star g$ is the composite

$$C \xrightarrow{\Delta} C \otimes C \xrightarrow{f \otimes g} A \otimes A \xrightarrow{m} A.$$

The associativity of this multiplication $\star$ is easy to check (see Exercise 1.28 below).

The map $u \circ \epsilon$ is a two-sided identity element for $\star$, meaning that every $f \in \text{Hom}(C, A)$ satisfies

$$\sum f(c_1)\epsilon(c_2) = f(c) = \sum \epsilon(c_1)f(c_2)$$

for all $c \in C$. One sees this by adding a top row to (1.4):

$$\begin{array}{cccc}
A \otimes k & \xrightarrow{f} & A & \xleftarrow{k \otimes A} \\
\downarrow f \otimes \text{id} & \quad & \downarrow \text{id} \otimes f & \\
C \otimes k & \xrightarrow{\text{id} \otimes \epsilon} & C & \xleftarrow{k \otimes C} \\
\downarrow \text{id} \otimes \epsilon & & \downarrow \epsilon \otimes \text{id} & \\
C \otimes C & \xleftarrow{\Delta} & C & \xrightarrow{\Delta} C \otimes C
\end{array}$$

In particular, when one has a bialgebra $A$, the convolution product $\star$ gives an associative algebra structure on $\text{End}(A) := \text{Hom}(A, A)$.

**Exercise 1.28.** Let $C$ be a $k$-coalgebra and $A$ be a $k$-algebra. Show that the binary operation $\star$ on $\text{Hom}(C, A)$ is associative.
**Exercise 1.29.** Let $C$ and $D$ be two $k$-coalgebras, and let $A$ and $B$ be two $k$-algebras. Prove that:

(a) If $f : C \to A$, $f' : C \to A$, $g : D \to B$ and $g' : D \to B$ are four $k$-linear maps, then

$$ (f \otimes g) \star (f' \otimes g') = (f \star f') \otimes (g \star g') $$

in the convolution algebra $\text{Hom}(C \otimes D, A \otimes B)$.

(b) Let $R$ be the $k$-linear map $\text{Hom}(C, A) \otimes \text{Hom}(D, B) \to \text{Hom}(C \otimes D, A \otimes B)$ which sends every tensor $f \otimes g \in \text{Hom}(C, A) \times \text{Hom}(D, B)$ to the map $f \otimes g : C \otimes D \to A \otimes B$. (Notice that the tensor $f \otimes g$ and the map $f \otimes g$ are different things which happen to be written in the same way.) Then, $R$ is a $k$-algebra homomorphism.

**Definition 1.30.** A bialgebra $A$ is called a Hopf algebra if there is an element $S$ (called an antipode for $A$) in $\text{End}(A)$ which is a 2-sided inverse under $\star$ for the identity map $\text{id}_A$. In other words, this diagram commutes:

$$
\begin{array}{ccc}
A \otimes A & \xrightarrow{S \otimes \text{id}_A} & A \otimes A \\
\downarrow{\Delta} & & \downarrow{m} \\
A & \xrightarrow{\epsilon} & k \\
\downarrow{\Delta} & & \downarrow{u} \\
A \otimes A & \xrightarrow{\text{id}_A \otimes S} & A \otimes A \\
\end{array}
$$

Or equivalently, if $\Delta(a) = \sum a_1 \otimes a_2$, then

$$
\sum_{(a)} S(a_1) a_2 = u(\epsilon(a)) = \sum_{(a)} a_1 S(a_2).
$$

**Example 1.31.** For a group algebra $kG$, one can define an antipode $k$-linearly via $S(t_g) = t_g^{-1}$. The top pentagon in the above diagram commutes because

$$
(S \star \text{id})(t_g) = m((S \otimes \text{id})(t_g \otimes t_g)) = S(t_g) t_g = t_g^{-1} t_g = \epsilon = (u \circ \epsilon)(t_g).
$$

Note that when it exists, the antipode $S$ is unique, as with all 2-sided inverses in associative algebras: if $S, S'$ are two 2-sided $\star$-inverses to $\text{id}_A$ then

$$
S' = (u \circ \epsilon) \star S' = (S \star \text{id}_A) \star S' = S \star (\text{id}_A \star S') = S \star (u \circ \epsilon) = S.
$$

On the other hand, the next property is not quite as obvious, but is useful when one wants to check that a certain map is the antipode in a particular Hopf algebra, by checking it on an algebra generating set.

**Proposition 1.32.** The antipode $S$ in a Hopf algebra $A$ is an algebra anti-endomorphism: $S(1) = 1$, and $S(ab) = S(b)S(a)$ for all $a, b$ in $A$.

**Proof.** (see [160, Chap. 4]) Since $\Delta$ is an algebra map, one has $\Delta(1) = 1 \otimes 1$, and therefore $1 = u(\epsilon(1)) = S(1) \cdot 1 = S(1)$.

To show $S(ab) = S(b)S(a)$, consider $A \otimes A$ as a coalgebra and $A$ as an algebra. Then $\text{Hom}(A \otimes A, A)$ is an associative algebra with a convolution product $\otimes$ (to be distinguished from the convolution $\star$ on $\text{End}(A)$), having two-sided identity element $u_A \epsilon_{A \otimes A}$. We will show below that these three elements of $\text{Hom}(A \otimes A, A)$

$$
\begin{align*}
f(a \otimes b) &= ab \\
g(a \otimes b) &= S(b)S(a) \\
h(a \otimes b) &= S(ab)
\end{align*}
$$

have the property that

$$
\tag{1.14}
h \otimes f = u_A \epsilon_{A \otimes A} = f \otimes g
$$

which would then show the desired equality $h = g$ via associativity:

$$
\begin{align*}
h &= h \otimes (u_A \epsilon_{A \otimes A}) = h \otimes (f \otimes g) = (h \otimes f) \otimes g = (u_A \epsilon_{A \otimes A}) \otimes g = g.
\end{align*}
$$
So we evaluate the three elements in (1.14) on $a \otimes b$, assuming $\Delta(a) = \sum_{(a)} a_1 \otimes a_2$ and $\Delta(b) = \sum_{(b)} b_1 \otimes b_2$, and hence $\Delta(ab) = \sum_{(a),(b)} a_1 b_1 \otimes a_2 b_2$. One has

$$(u_A \epsilon_{A \otimes A})(a \otimes b) = u_A(\epsilon_{A}(a) \epsilon_{A}(b)) = u_A(\epsilon_{A}(ab)).$$

$$
(h \otimes f)(a \otimes b) = \sum_{(a),(b)} h(a_1 \otimes b_1)f(a_2 \otimes b_2) \\
= \sum_{(a),(b)} S(a_1 b_1) a_2 b_2 \\
= (S \star \text{id}_A)(ab) = u_A(\epsilon_{A}(ab)).
$$

$$(f \otimes g)(a \otimes b) = \sum_{(a),(b)} f(a_1 \otimes b_1)g(a_2 \otimes b_2) \\
= \sum_{(a),(b)} a_1 b_1 S(b_2)S(a_2) \\
= \sum_{(a)} a_1 \cdot (\text{id}_A \star S)(b) \cdot S(a_2) \\
= u_A(\epsilon_{A}(b))\sum_{(a)} a_1 S(a_2) = u_A(\epsilon_{A}(b))u_A(\epsilon_{A}(a)) = u_A(\epsilon_{A}(ab))$$

\[\square\]

**Remark 1.33.** Recall from Remark 1.16 that the comultiplication on a bialgebra $A$ allows one to define an $A$-module structure on the tensor product $M \otimes N$ of two $A$-modules $M, N$. Similarly, the anti-endomorphism $S$ in a Hopf algebra allows one to turn left $A$-modules into right $A$-modules, or vice-versa.\footnote{Be warned that these two transformations are not mutually inverse! Turning a left $A$-module into a right one and then again into a left one using the antipode might lead to a non-isomorphic $A$-module, unless the antipode $S$ satisfies $S^2 = \text{id}$.} E.g., left $A$-modules $M$ naturally have a right $A$-module structure on the dual $k$-module $M^* := \text{Hom}(M,k)$, defined via $(fa)(m) := f(am)$ for $f \in M^*$ and $a \in A$. The antipode $S$ can be used to turn this back into a left $A$-module $M^*$, via $(af)(m) = f(S(a)m)$.

For groups $G$ and left $kG$-modules (group representations) $M$, this is how one defines the **contragredient action** of $G$ on $M^*$, namely $t_g$ acts as $(t_g f)(m) = f(t_g^{-1}m)$.

More generally, if $A$ is a Hopf algebra and $M$ and $N$ are two left $A$-modules, then $\text{Hom}(M,N)$ (the Hom here means $\text{Hom}_k$, not $\text{Hom}_A$) canonically becomes a left $A$-module by setting

$$(af)(m) = \sum_{(a)} a_1 f(S(a_2)m) \quad \text{for all } a \in A, f \in \text{Hom}(M,N) \text{ and } m \in M.$$  

\footnote{In more abstract terms, this $A$-module structure is given by the composition

$$A \xrightarrow{\Delta} A \otimes A \xrightarrow{\text{id}_A \otimes S} A \otimes A^{\text{op}} \xrightarrow{\text{End}(\text{Hom}(M,N))} ,$$

where the last arrow is the morphism

$$A \otimes A^{\text{op}} \rightarrow \text{End}(\text{Hom}(M,N)),$$

$$a \otimes b \mapsto (f \mapsto (M \rightarrow N, m \mapsto af(bm))).$$

Here, $A^{\text{op}}$ denotes the opposite algebra of $A$, which is the $k$-algebra differing from $A$ only in the multiplication being twisted (the product of $a$ and $b$ in $A^{\text{op}}$ is defined to be the product of $b$ and $a$ in $A$). As $k$-modules, $A^{\text{op}} = A$, but we prefer to use $A^{\text{op}}$ instead of $A$ here to ensure that all morphisms in the above composition are algebra maps.}
Along the same lines, whenever $A$ is a $k$-bialgebra, we are supposed to think of the counit $A \xrightarrow{\epsilon} k$ as giving a way to make $k$ into a trivial $A$-module. This $A$-module $k$ behaves as one would expect: the canonical isomorphisms $k \otimes M \to M$, $M \otimes k \to M$ and (if $A$ is a Hopf algebra) $\text{Hom}(M,k) \to M^*$ are $A$-module isomorphisms for any $A$-module $M$.

**Corollary 1.34.** Commutativity of $A$ implies the antipode is an involution: $S^2 = \text{id}_A$.

**Proof.** One checks that $S^2 = S \circ S$ is a right $*$-inverse to $S$, and hence $S^2 = \text{id}_A$:

$$(S \ast S^2)(a) = \sum (S(a_1)S^2(a_2)) = S\left( \sum (S(a_2)a_1) \right) \text{ by Proposition 1.32}$$

$= S(u(\epsilon(a))) = u(\epsilon(a))$ since $S(1) = 1$ by Proposition 1.32.

□

**Remark 1.35.** We won’t need it, but it is easy to adapt the above proof to show that $S^2 = \text{id}_A$ also holds for cocommutative Hopf algebras; see [117, Corollary 1.5.12] or [160, Chapter 4]. For a general Hopf algebra which is not finite-dimensional over a field $k$, the antipode $S$ may not even have finite order, even in the graded connected setting. E.g., Aguiar and Sottile [7] show that the Malvenuto-Reutenauer Hopf algebra of permutations has antipode of infinite order. In general, antipodes need not even be invertible [161].

In our frequent setting of graded connected bialgebras, antipodes come for free.

**Proposition 1.36.** A graded connected bialgebra $A$ has a unique antipode $S$, which is a graded map $A \xrightarrow{S} A$, endowing it with a Hopf structure.

**Proof.** Let us try to define a ($k$-linear) left $*$-inverse $S$ to $\text{id}_A$ on each homogeneous component $A_n$, via induction on $n$.

In the base case $n = 0$, Proposition 1.32 and its proof show that one must define $S(1) = 1$ so $S$ is the identity on $A_0 = k$.

In the inductive step, recall from Exercise 1.26 that a homogeneous element $a$ of degree $n > 0$ has $\Delta(a) = a \otimes 1 + \sum a_1 \otimes a_2$, with each $\deg(a_1) < n$. Hence in order to have $S \ast \text{id}_A = \epsilon$, one must define $S(a)$ in such a way that $S(a) \cdot 1 + \sum S(a_1)a_2 = \epsilon(a) = 0$ and hence $S(a) := -\sum S(a_1)a_2$, where $S(a_1)$ have already been uniquely defined by induction. This does indeed define such a left $*$-inverse $S$ to $\text{id}_A$, by induction. It is also a graded map by induction.

The same argument shows how to define a right $*$-inverse $S'$ to $\text{id}_A$. Then $S = S'$ is a two-sided $*$-inverse to $\text{id}_{A}$ by the associativity of $*$.

□

Here is another consequence of the fact that $S(1) = 1$.

**Proposition 1.37.** In bialgebras, primitive elements $x$ have $\epsilon(x) = 0$, and in Hopf algebras, they have $S(x) = -x$.

**Proof.** In a bialgebra, $\epsilon(1) = 1$. Hence $\Delta(x) = 1 \otimes x + x \otimes 1$ implies via (1.4) that $1 \cdot \epsilon(x) + \epsilon(1)x = x$, so $\epsilon(x) = 0$. It also implies via (1.12) that $S(x)1 + S(1)x = u(x) = u(0) = 0$, so $S(x) = -x$.

□

Thus whenever $A$ is a Hopf algebra generated as an algebra by its primitive elements, $S$ is the unique anti-endomorphism that negates all primitive elements.

**Example 1.38.** The tensor and symmetric algebras $T(V)$ and $\text{Sym}(V)$ are each generated by $V$, which contains only primitive elements in either case. Hence one has in $T(V)$ that

$$(1.15) \quad S(x_{i_1}x_{i_2} \cdots x_{i_k}) = (-x_{i_k}) \cdots (-x_{i_2})(-x_{i_1}) = (-1)^k x_{i_k} \cdots x_{i_2} x_{i_1}$$

for each word $(i_1,\ldots,i_k)$ in the alphabet $I$ if $V$ is a free $k$-module with basis $\{x_i\}_{i \in I}$. The same holds in $\text{Sym}(V)$ for each multiset $(i_1,\ldots,i_k)$, recalling that the monomials are now commutative. In other words, for a commutative polynomial $f(x)$ in $\text{Sym}(V)$, the antipode $S$ sends $f$ to $f(-x)$, negating all the variables.
The antipode for a graded connected Hopf algebra has an interesting formula due to Takeuchi [161], reminiscent of P. Hall’s formula for the Möbius function of a poset. For the sake of stating this, consider (for every $k \in \mathbb{N}$) the $k$-fold tensor power $A^\otimes k = A \otimes \cdots \otimes A$ (defined in Example 1.2) and define iterated multiplication and comultiplication maps

$$A^\otimes k \xrightarrow{m^k} A$$

by induction over $k$, setting $m^{(-1)} = u$, $\Delta^{(-1)} = \epsilon$, $m^{(0)} = \Delta^{(0)} = \text{id}_A$, and

$m^{(k)} = m \circ (\text{id}_A \otimes m^{(k-1)})$ for every $k \geq 1$;

$\Delta^{(k)} = (\text{id}_A \otimes \Delta^{(k-1)}) \circ \Delta$ for every $k \geq 1$.

Using associativity and coassociativity, one can see that for $k \geq 1$ these maps also satisfy

$m^{(k)}(a^{(1)} \otimes \cdots \otimes a^{(k+1)}) = a^{(1)} \cdots a^{(k+1)}$

$\Delta^{(k)}(b) = \sum b_1 \otimes \cdots \otimes b_{k+1}$ in Sweedler notation.

**Exercise 1.39.** Let $A$ be a $k$-algebra. Let us define, for every $k \in \mathbb{N}$, a $k$-linear map $m^{(k)} : A^\otimes (k+1) \to A$. Namely, we define these maps by induction over $k$, with the induction base $m^{(0)} = \text{id}_A$, and with the induction step $m^{(k)} = m \circ (\text{id}_A \otimes m^{(k-1)})$ for every $k \geq 1$. (This generalizes our definition of $m^{(k)}$ for Hopf algebras $A$ given above, except for $m^{(-1)}$ which we have omitted.)

(a) Show that $m^{(k)} = m \circ (m^{(i)} \otimes m^{(k-1-i)})$ for every $k \geq 0$ and $0 \leq i \leq k - 1$.

(b) Show that $m^{(k)} = m \circ (m^{(k-1)} \otimes \text{id}_A)$ for every $k \geq 1$.

(c) Show that $m^{(k)} = m^{(k-1)} \circ (\text{id}_A \otimes \Delta)$ for every $k \geq 0$ and $0 \leq i \leq k - 1$.

(d) Show that $m^{(k)} = m^{(k-1)} \circ (\text{id}_A \otimes m^{(k-1-i)})$ for every $k \geq 1$.

**Exercise 1.40.** Let $C$ be a $k$-coalgebra. Let us define, for every $k \in \mathbb{N}$, a $k$-linear map $\Delta^{(k)} : C \to C^\otimes (k+1)$. Namely, we define these maps by induction over $k$, with the induction base $\Delta^{(0)} = \text{id}_C$, and with the induction step $\Delta^{(k)} = (\text{id}_C \otimes \Delta^{(k-1)}) \circ \Delta$ for every $k \geq 1$. (This generalizes our definition of $\Delta^{(k)}$ for Hopf algebras $A$ given above, except for $\Delta^{(-1)}$ which we have omitted.)

(a) Show that $\Delta^{(k)} = (\Delta^{(i)} \otimes \Delta^{(k-1-i)}) \circ \Delta$ for every $k \geq 0$ and $0 \leq i \leq k - 1$.

(b) Show that $\Delta^{(k)} = (\Delta^{(k-1)} \otimes \text{id}_C) \circ \Delta$ for every $k \geq 1$.

(c) Show that $\Delta^{(k)} = (\text{id}_C \otimes \Delta \otimes \text{id}_C^{(k-1-i)}) \circ \Delta^{(k-1)}$ for every $k \geq 0$ and $0 \leq i \leq k - 1$.

(d) Show that $\Delta^{(k)} = (\text{id}_C \otimes \Delta^{(k-1)} \otimes \Delta) \circ \Delta^{(k-1)} = (\Delta \otimes \text{id}_C^{(k-1)} \otimes \Delta) \circ \Delta^{(k-1)}$ for every $k \geq 1$.

**Remark 1.41.** Exercise 1.39 holds more generally for nonunital associative algebras $A$ (that is, $k$-modules $A$ equipped with a $k$-linear map $m : A \otimes A \to A$ such that the diagram (1.1) is commutative, but not necessarily admitting a unit map $u$). Similarly, Exercise 1.40 holds for non-counital coassociative coalgebras $C$. The existence of a unit in $A$, respectively a counit in $C$, allows slightly extending these two exercises by additionally introducing maps $m^{(-1)} = u : k \to A$ and $\Delta^{(-1)} = \epsilon : C \to k$; however, not much is gained from this extension.

**Exercise 1.42.** For every $k \in \mathbb{N}$ and every $k$-bialgebra $H$, consider the map $\Delta_H^{(k)} : H \to H^{\otimes (k+1)}$ (this is the map $\Delta^{(k)}$ defined as in Exercise 1.40 for $C = H$), and the map $m_H^{(k)} : H^{\otimes (k+1)} \to H$ (this is the map $m^{(k)}$ defined as in Exercise 1.39 for $A = H$).
Let $H$ be a $k$-bialgebra. Let $k \in \mathbb{N}$. Show that:\footnote{The following statements are taken from [124]; specifically, part (c) is [124, Lem. 1.8].}

(a) The map $m^{(k)}_H : H^k \to H$ is a $k$-coalgebra homomorphism.
(b) The map $\Delta^{(k)}_H : H \to H^k$ is a $k$-algebra homomorphism.
(c) We have $m^{(\ell)}_H \circ (\Delta^{(k)}_H)^{(\ell+1)} = \Delta^{(k)}_H \circ m^{(\ell)}_H$ for every $\ell \in \mathbb{N}$.
(d) We have $m^{(\ell)}_H \circ \Delta^{(k)}_{H^\ell} = \Delta^{(k)}_H \circ m^{(\ell)}_H$ for every $\ell \in \mathbb{N}$.

The iterated multiplication and comultiplication maps allow explicitly computing the convolution of multiple maps; the following formula will often be used without explicit mention:

**Exercise 1.43.** Let $C$ be a $k$-coalgebra, and $A$ be a $k$-algebra. Let $k \in \mathbb{N}$. Let $f_1, f_2, \ldots, f_k$ be $k$ elements of $\text{Hom}(C,A)$. Show that

$$f_1 \star f_2 \cdots \star f_k = m^{(k-1)}_A \circ (f_1 \otimes f_2 \otimes \cdots \otimes f_k) \circ \Delta^{(k-1)}_C.$$

We are now ready to state Takeuchi’s formula for the antipode:

**Proposition 1.44.** In a graded connected Hopf algebra $A$, the antipode has formula

$$S = \sum_{k \geq 0} (-1)^k m^{(k-1)} f^\otimes k \Delta^{(k-1)}$$

(1.16)

$$= u\epsilon - f + m \circ f^\otimes 2 \circ \Delta - m^{(2)} \circ f^\otimes 3 \circ \Delta^{(2)} + \cdots$$

where $f := \text{id}_A - u\epsilon$ in $\text{End}(A)$.

**Proof.** We argue as in [161, proof of Lemma 14] or [7, §5]. For any $f$ in $\text{End}(A)$ one has this explicit formula for its $k$-fold convolution power $f^{*k} := f \ast \cdots \ast f$ in terms of its tensor powers $f^\otimes n := f \otimes \cdots \otimes f$ (according to Exercise 1.43):

$$f^{*k} = m^{(k-1)}_A \circ f^\otimes k \circ \Delta^{(k-1)}.$$ 

Therefore any $f$ annihilating $A_0$ will be locally $\ast$-nilpotent on $A$, meaning that for each $n$ one has that $A_n$ is annihilated by $f^{*n}$ for every $m > n$: homogeneity forces that for $a$ in $A_n$, every summand of $\Delta^{(m-1)}(a)$ must contain among its $m$ tensor factors at least one factor lying in $A_0$, so each summand is annihilated by $f^\otimes m$, and $f^{*m}(a) = 0$.

In particular such $f$ have the property that $u\epsilon + f$ has as two-sided $\ast$-inverse

$$(u\epsilon + f)^{\ast(-1)} = u\epsilon - f + f \ast f - f \ast f \ast f + \cdots$$

$$= \sum_{k \geq 0} (-1)^k f^{*k} = \sum_{k \geq 0} (-1)^k m^{(k-1)}_A \circ f^\otimes k \circ \Delta^{(k-1)}.$$ 

The proposition follows upon taking $f := \text{id}_A - u\epsilon$, which annihilates $A_0$. \hfill \Box

**Remark 1.45.** In fact, one can see that Takeuchi’s formula applies more generally to define an antipode $A \xrightarrow{S} A$ in any (not necessarily graded) bialgebra $A$ where the map $\text{id}_A - u\epsilon$ is locally $\ast$-nilpotent.

It is also worth noting that the proof of Proposition 1.44 gives an alternate proof of Proposition 1.36.

To finish our discussion of antipodes, we mention some properties (taken from [160, Chap. 4]) relating antipodes to convolutional inverses. It also shows that a bialgebra morphism between Hopf algebras automatically respects the antipodes.

**Proposition 1.46.** Let $H$ be a Hopf algebra with antipode $S$.

(a) For any algebra $A$ and algebra morphism $H \xrightarrow{\alpha} A$, one has $\alpha \circ S = \alpha^{*-1}$, the convolutional inverse to $\alpha$ in $\text{Hom}(H, A)$.
(b) For any coalgebra $C$ and coalgebra morphism $C \xrightarrow{\gamma} H$, one has $S \circ \gamma = \gamma^{*-1}$, the convolutional inverse to $\gamma$ in $\text{Hom}(C, H)$.
(c) If $H_1, H_2$ are Hopf algebras with antipodes $S_1, S_2$, then any bialgebra morphism $H_1 \xrightarrow{\beta} H_2$ is a Hopf morphism, that is, it commutes with the antipodes, since $\beta \circ S_1^{(a)} = \beta^{*-1} \circ S_2 \circ \beta$. 

$$\beta$$

The following statements are taken from [124]; specifically, part (c) is [124, Lem. 1.8].
Proof. We prove (a); the proof of (b) is similar, and (c) follows immediately from (a),(b) as indicated in its statement. Begin with the following: given an algebra morphism \( A \xrightarrow{\phi} A' \) and coalgebra morphism \( C \xrightarrow{\gamma} C' \), the pre- and post-composition \( f \mapsto \alpha \circ f \circ \gamma \) is a convolution algebra morphism \( \text{Hom}(C', A) \xrightarrow{\phi} \text{Hom}(C, A') \), as one can check that \( \varphi(u_A \circ \epsilon_{C'}) = \alpha \circ u_A \circ \epsilon_{C'} \circ \gamma = u_{A'} \circ \epsilon_C \), and that every \( f \in \text{Hom}(C', A) \) and \( g \in \text{Hom}(C', A) \) satisfy

\[
\varphi(f \ast g) = \alpha \circ (f \ast g) \circ \gamma = \alpha \circ m \circ (f \otimes g) \circ \Delta \circ \gamma \\
= m \circ ((\alpha \circ f \circ \gamma) \otimes (\alpha \circ g \circ \gamma)) \circ \Delta \\
= \varphi(f) \ast \varphi(g).
\]

For assertion (a), note that a special case of the above observation shows \( f \mapsto \alpha \circ f \) gives a convolution-algebra morphism \( \text{Hom}(H, H) \xrightarrow{\epsilon} \text{Hom}(H, A) \), and hence

\[
\alpha \circ S = \varphi(S) = \varphi((\text{id}_H)^{-1}) = (\varphi(\text{id}_H))^{*-1} = (\alpha \circ \text{id}_H)^{-1} = \alpha'^{-1}.
\]

Exercise 1.47. Prove that the antipode \( S \) of a Hopf algebra \( A \) is a coalgebra anti-endomorphism, i.e., that it satisfies \( \epsilon \circ S = \epsilon \) and \( \Delta \circ S = T \circ (S \otimes S) \circ \Delta \), where \( T : A \otimes A \rightarrow A \otimes A \) is the twist map sending every \( a \otimes b \) to \( b \otimes a \).

Exercise 1.48. If \( C \) is a \( k \)-coalgebra and if \( A \) is a \( k \)-algebra, then a \( k \)-linear map \( f : C \rightarrow A \) is said to be \( ** \)-invertible if it is invertible as an element of the \( k \)-algebra \( \text{Hom}(C, A), \ast \). In this case, the multiplicative inverse \( f^{-1} \) of \( f \) in \( \text{Hom}(C, A) \) is called the \( \ast \)-inverse of \( f \).

For any two \( k \)-modules \( U \) and \( V \), let \( T_{U, V} : U \otimes V \rightarrow V \otimes U \) be the twist map (i.e., the \( k \)-linear map \( U \otimes V \rightarrow V \otimes U \) sending every \( u \otimes v \) to \( v \otimes u \)).

(a) If \( C \) is a \( k \)-bialgebra, if \( A \) is a \( k \)-algebra, and if \( r : C \rightarrow A \) is a \( \ast \)-invertible \( k \)-algebra homomorphism, then prove that the \( \ast \)-inverse \( r^{-1} \) of \( r \) is a \( k \)-algebra anti-homomorphism.\(^{19}\)

(b) If \( C \) is a \( k \)-bialgebra, if \( A \) is a \( k \)-coalgebra, and if \( r : C \rightarrow A \) is a \( \ast \)-invertible \( k \)-coalgebra anti-homomorphism, then prove that the \( \ast \)-inverse \( r^{-1} \) of \( r \) is a \( k \)-coalgebra anti-homomorphism\(^{20}\).

(c) Derive Proposition 1.32 from Exercise 1.48(a), and derive Exercise 1.47 from Exercise 1.48(b).

(d) Prove Corollary 1.34 again using Proposition 1.46.

(e) Prove that the antipode \( S \) of a cocommutative Hopf algebra \( A \) satisfies \( S^2 = \text{id}_A \). (This was a statement made in Remark 1.35.)

Exercise 1.49. (a) Let \( A \) be a Hopf algebra. If \( P : A \rightarrow A \) is a \( k \)-linear map such that every \( a \in A \) satisfies

\[
\sum_{(a)} P(a_2) \cdot a_1 = u(\epsilon(a)),
\]

then prove that the antipode \( S \) of \( A \) is invertible and its inverse is \( P \).

(b) Let \( A \) be a Hopf algebra. If \( P : A \rightarrow A \) is a \( k \)-linear map such that every \( a \in A \) satisfies

\[
\sum_{(a)} a_2 \cdot P(a_1) = u(\epsilon(a)),
\]

then prove that the antipode \( S \) of \( A \) is invertible and its inverse is \( P \).

(c) Show that the antipode of a graded connected Hopf algebra is invertible.

(Compare this exercise to [117, Lemma 1.5.11].)

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\(^{19}\) A \( k \)-algebra anti-homomorphism means a \( k \)-linear map \( \varphi : P \rightarrow Q \) between two \( k \)-algebras \( P \) and \( Q \) which satisfies \( \varphi \circ m_P = m_Q \circ (\varphi \otimes \varphi) \circ T_{P, P} \) and \( \varphi \circ u_P = u_Q \).

\(^{20}\) A \( k \)-coalgebra anti-homomorphism means a \( k \)-linear map \( \varphi : P \rightarrow Q \) between two \( k \)-coalgebras \( P \) and \( Q \) which satisfies \( \Delta_Q \circ \varphi = T_{Q, Q} \circ (\varphi \otimes \varphi) \circ \Delta_P \) and \( \varepsilon_Q \circ \varphi = \varepsilon_P \).
Definition 1.50. Let $C$ be a $k$-coalgebra. A subcoalgebra of $C$ means a $k$-coalgebra $D$ such that $D \subset C$ and such that the canonical inclusion map $D \to C$ is a $k$-coalgebra homomorphism\footnote{In this definition, we follow [120, p. 55] and [168, §6.7]; other authors may use other definitions.}. When $k$ is a field, we can equivalently define a subcoalgebra of $C$ as a $k$-submodule $D$ of $C$ such that $\Delta_C(D)$ is a subset of the $k$-submodule $D \otimes D$ of $C \otimes C$; however, this might no longer be equivalent when $k$ is not a field\footnote{This is because the $k$-submodule $D \otimes D$ of $C \otimes C$ is generally not isomorphic to the $k$-module $D \otimes D$. See [120, p. 56] for specific counterexamples for the non-equivalence of the two notions of a subcoalgebra. Notice that the equivalence is salvaged if $D$ is a direct summand of $C$ as a $k$-module (see Exercise 1.51 for this).}.

Similarly, a subbialgebra of a bialgebra $C$ is a $k$-bialgebra $D$ such that $D \subset C$ and such that the canonical inclusion map $D \to C$ is a $k$-bialgebra homomorphism. Also, a Hopf subalgebra of a Hopf algebra $C$ is a $k$-Hopf algebra $D$ such that $D \subset C$ and such that the canonical inclusion map $D \to C$ is a $k$-Hopf algebra homomorphism\footnote{By Proposition 1.46(c), we can also define it as a subbialgebra of $C$ which happens to be a Hopf algebra.}.

Exercise 1.51. Let $C$ be a $k$-coalgebra. Let $D$ be a $k$-submodule of $C$ such that $D$ is a direct summand of $C$ as a $k$-module (i.e., there exists a $k$-submodule $E$ of $C$ such that $C = D \oplus E$). (This is automatically satisfied if $k$ is a field.) Assume that $\Delta(D) \subset C \otimes D$ and $\Delta(D) \subset D \otimes C$. (Here, we are abusing the notation $C \otimes D$ to denote the $k$-submodule of $C \otimes C$ spanned by tensors of the form $c \otimes d$ with $c \in C$ and $d \in D$; similarly, $D \otimes C$ should be understood.) Show that there is a canonically defined $k$-coalgebra structure on $D$ which makes $D$ a subcoalgebra of $C$.

The next exercise is implicit in [4, §5]:

Exercise 1.52. Let $k$ be a field. Let $C$ be a $k$-coalgebra, and let $U$ be any $k$-module. Let $f : C \to U$ be a $k$-linear map. Recall the map $\Delta^{(2)} : C \to C^{\otimes 3}$ from Exercise 1.40. Let $K = \ker ((\text{id}_C \otimes f \otimes \text{id}_C) \circ \Delta^{(2)})$.

(a) Show that $K$ is a $k$-subcoalgebra of $C$.
(b) Show that every $k$-subcoalgebra of $C$ which is a subset of $\ker f$ must be a subset of $K$.

Exercise 1.53. (a) Let $C = \bigoplus_{n \geq 0} C_n$ be a graded $k$-coalgebra, and $A$ be any $k$-algebra. Notice that $C_0$ itself is a $k$-subcoalgebra of $C$. Let $h : C \to A$ be a $k$-linear map such that the restriction $h\mid_{C_0}$ is a $*$-invertible map in $\text{Hom}(C_0, A)$. Prove that $h$ is a $*$-invertible map in $\text{Hom}(C, A)$. (This is a weaker version of Takeuchi’s [161, Lemma 14].)
(b) Let $A = \bigoplus_{n \geq 0} A_n$ be a graded $k$-bialgebra. Notice that $A_0$ is a sub-bialgebra of $A$. Assume that $A_0$ is a Hopf algebra. Show that $A$ is a Hopf algebra.
(c) Obtain yet another proof of Proposition 1.36.

Exercise 1.54. Let $A = \bigoplus_{n \geq 0} A_n$ be a graded connected $k$-bialgebra. Let $p$ be the $k$-submodule of $A$ consisting of the primitive elements of $A$.

(a) If $I$ is a two-sided coideal of $A$ such that $I \cap p = 0$ and such that $I = \bigoplus_{n \geq 0} (I \cap A_n)$, then prove that $I = 0$.
(b) Let $f : A \to C$ be a graded surjective coalgebra homomorphism from $A$ to a graded $k$-coalgebra $C$. If $f\mid_p$ is injective, then prove that $f$ is injective.
(c) Assume that $k$ is a field. Show that the claim of Exercise 1.54(b) is valid even without requiring $f$ to be surjective.

Remark 1.55. Exercise 1.54(b) and (c) are often used in order to prove that certain coalgebra homomorphisms are injective.

The word “bialgebra” can be replaced by “coalgebra” in Exercise 1.54, provided that the notion of a graded connected coalgebra is defined correctly (namely, as a graded coalgebra such that the restriction of $\epsilon$ to the 0-th graded component is an isomorphism), and the notion of the element 1 of a graded connected coalgebra is defined accordingly (namely, as the preimage of 1 $\in k$ under the restriction of $\epsilon$ to the 0-th graded component).

1.5. Commutativity, cocommutativity.
**Definition 1.56.** Say that the k-algebra $A$ is **commutative** if $ab = ba$, that is, this diagram commutes:

\[
\begin{array}{c}
A \otimes A \\
\downarrow m \\
A
\end{array} \quad \begin{array}{c}
A \otimes A \\
\downarrow T \\
A \otimes A
\end{array} \quad \begin{array}{c}
A \\
\downarrow m \\
A
\end{array}
\]

Say that the k-coalgebra $C$ is **cocommutative** if this diagram commutes:

\[
\begin{array}{c}
C \otimes C \\
\downarrow \Delta \\
C
\end{array} \quad \begin{array}{c}
C \otimes C \\
\downarrow T \\
C \otimes C
\end{array} \quad \begin{array}{c}
C \\
\downarrow \Delta \\
C
\end{array}
\]

**Example 1.57.** Group algebras $kG$ are always cocommutative, but commutative if and only if $G$ is abelian.

Tensor algebras $T(V)$ are always cocommutative, but not generally commutative\(^{24}\).

Symmetric algebras $\text{Sym}(V)$ are always cocommutative and commutative.

Homology and cohomology of $H$-spaces are always cocommutative and commutative **in the topologist’s sense** where one reinterprets that twist map $A \otimes A \xrightarrow{T} A \otimes A$ to have the extra sign as in (1.7).

Note how the cocommutative Hopf algebras $(T(V), \text{Sym}(V))$ have much of their structure controlled by their $k$-submodules $V$, which consist of primitive elements only (although, in general, not of all their primitive elements). This is not far from the truth in general, and closely related to Lie algebras.

**Exercise 1.58.** Recall that a **Lie algebra** over $k$ is a $k$-module $\mathfrak{g}$ with a $k$-bilinear map $[\cdot, \cdot] : \mathfrak{g} \times \mathfrak{g} \to \mathfrak{g}$ that satisfies $[x, x] = 0$ for $x$ in $\mathfrak{g}$, and the **Jacobi identity**

\[
[x, [y, z]] = [[x, y], z] + [y, [x, z]], \text{ or equivalently } [x, [y, z]] + [z, [x, y]] + [y, [z, x]] = 0.
\]

This $k$-bilinear map $[\cdot, \cdot]$ is called the **Lie bracket** of $\mathfrak{g}$.

(a) Check that any associative algebra $A$ gives rise to a Lie algebra by means of the commutator operation 

$[a, b] := ab - ba$.

(b) If $A$ is also a bialgebra, show that the $k$-submodule of primitive elements $p \subset A$ is closed under the Lie bracket, that is, $[p, p] \subset p$, and hence forms a Lie subalgebra.

Conversely, given a Lie algebra $p$, one constructs the **universal enveloping algebra** $U(p) := T(p)/J$ as the quotient of the tensor algebra $T(p)$ by the two-sided ideal $J$ generated by all elements $xy - yx - [x, y]$ for $x, y$ in $p$.

(c) Show that $J$ is also a two-sided coideal in $T(p)$ for its usual coalgebra structure, and hence the quotient $U(p)$ inherits the structure of a cocommutative bialgebra.

(d) Show that the antipode $S$ on $T(p)$ preserves $J$, meaning that $S(J) \subset J$, and hence $U(p)$ inherits the structure of a Hopf algebra.

There are theorems, discussed in [28, §3.8], [117, Chap. 5], [44, §3.2] giving various mild hypotheses in addition to cocommutativity which imply that the inclusion of the $k$-module $p$ of primitives in a Hopf algebra $A$ extends to a Hopf isomorphism $U(p) \cong A$.

**Exercise 1.59.**

(a) Let $C$ be a $k$-coalgebra. Show that $C$ is cocommutative if and only if its comultiplication $\Delta_C : C \to C \otimes C$ is a $k$-coalgebra homomorphism.

(b) Let $A$ be a $k$-algebra. Show that $A$ is commutative if and only if its multiplication $m_A : A \otimes A \to A$ is a $k$-algebra homomorphism.

**Remark 1.60.** If $C$ is a $k$-coalgebra, then $\epsilon_C : C \to k$ is always a $k$-coalgebra homomorphism. Similarly, $u_A : k \to A$ is a $k$-algebra homomorphism whenever $A$ is a $k$-algebra.

\(^{24}\)If $k$ is a field, then $T(V)$ is commutative if and only if $\dim_k V \leq 1$. 
Exercise 1.61. Let $A$ be a commutative $k$-algebra, and let $k \in \mathbb{N}$. The symmetric group $S_k$ acts on the $k$-fold tensor power $A^\otimes k$ by permuting the tensor factors: $\sigma(v_1 \otimes v_2 \otimes \cdots \otimes v_k) = v_{\sigma^{-1}(1)} \otimes v_{\sigma^{-1}(2)} \otimes \cdots \otimes v_{\sigma^{-1}(k)}$ for all $v_1, v_2, \ldots, v_k \in A$ and $\sigma \in S_k$. For every $\pi \in S_k$, denote by $\rho(\pi)$ the action of $\pi$ on $A^\otimes k$ (this is an endomorphism of $A^\otimes k$). Show that every $\pi \in S_k$ satisfies $\rho(\pi) \circ (\rho(\pi)) = \rho(k-1)$. (Recall that $m^{(k-1)} : A^\otimes k \to A$ is defined as in Exercise 1.39 for $k \geq 1$, and by $m^{(-1)} = u : k \to A$ for $k = 0$.)

Exercise 1.62. State and solve the analogue of Exercise 1.61 for cocommutative $k$-coalgebras.

Exercise 1.63. (a) If $H$ is a $k$-bialgebra and $A$ is a commutative $k$-algebra, and if $f$ and $g$ are two $k$-algebra homomorphisms $H \to A$, then prove that $f \ast g$ also is a $k$-algebra homomorphism $H \to A$.

(b) If $H$ is a $k$-bialgebra and $A$ is a commutative $k$-algebra, and if $f_1, f_2, \ldots, f_k$ are several $k$-algebra homomorphisms $H \to A$, then prove that $f_1 \ast f_2 \ast \cdots \ast f_k$ also is a $k$-algebra homomorphism $H \to A$.

(c) If $H$ is a Hopf algebra and $A$ is a commutative $k$-algebra, and if $f : H \to A$ is a $k$-algebra homomorphism, then prove that $f \circ S : H \to A$ (where $S$ is the antipode of $H$) is again a $k$-algebra homomorphism, and is a $*$-inverse to $f$.

(d) If $A$ is a commutative $k$-algebra, then show that $m^{(k)}$ is a $k$-algebra homomorphism for every $k \in \mathbb{N}$. (The map $m^{(1)} : H \otimes 1 \to A$ is defined as in Exercise 1.39.)

(e) If $C'$ and $C$ are two $k$-coalgebras, if $\gamma : C \to C'$ is a $k$-coalgebra homomorphism, if $A$ and $A'$ are two $k$-algebras, if $\alpha : A \to A'$ is a $k$-algebra homomorphism, and if $f_1, f_2, \ldots, f_k$ are several $k$-linear maps $C' \to A$, then prove that

$$\alpha \circ (f_1 \ast f_2 \ast \cdots \ast f_k) \circ \gamma = (\alpha \circ f_1 \circ \gamma) \ast (\alpha \circ f_2 \circ \gamma) \ast \cdots \ast (\alpha \circ f_k \circ \gamma).$$

(f) If $H$ is a commutative $k$-bialgebra, and $k$ and $\ell$ are two nonnegative integers, then prove that $id^{(k)}_H \circ id^{(\ell)}_H = id^{(k\ell)}_H$.

(g) If $H$ is a commutative $k$-Hopf algebra, and $k$ and $\ell$ are two integers, then prove that $id^{(k)}_H \circ id^{(\ell)}_H = id^{(k\ell)}_H$. (These powers $id^{(k)}_H$, $id^{(\ell)}_H$ and $id^{(k\ell)}_H$ are well-defined since $id_H$ is $*$-invertible.)

(h) State and prove the duals of parts (a)–(g) of this exercise.

Remark 1.64. The maps $id^{(k)}_H$ for $k \in \mathbb{N}$ are known as the Adams operators of the bialgebra $H$; they are studied, inter alia, in [5]. Particular cases of Exercise 1.63(e) appear in [124, Corollaire II.9] and [58, Theorem 1]. Exercise 1.63(f) and its dual are [100, Prop. 1.6].

Exercise 1.65. Let $A$ be a cocommutative graded Hopf algebra with antipode $S$. Define a $k$-linear map $E : A \to A$ by having $E(a) = (\deg a) \cdot a$ for every homogeneous element $a$ of $A$.

(a) Prove that for every $a \in A$, the elements $(S \ast E)(a)$ and $(E \ast S)(a)$ (where $\ast$ denotes convolution in $\text{Hom}(A, A)$) are primitive.

(b) Prove that for every primitive $p \in A$, we have $(S \ast E)(p) = (E \ast S)(p) = E(p)$.

(c) Prove that for every $a \in A$ and every primitive $p \in A$, we have $(S \ast E)(ap) = [(S \ast E)(a), p] + \epsilon(a) E(p)$, where $[u, v]$ denotes the commutator $uv - vu$ of $u$ and $v$.

(d) If $A$ is connected and $Q$ is a subring of $k$, prove that the $k$-algebra $A$ is generated by the $k$-submodule $p$ consisting of the primitive elements of $A$.

(e) Assume that $A$ is the tensor algebra $T(V)$ of a $k$-module $V$, and that the $k$-submodule $V = V^{\otimes 1}$ of $T(V)$ is the 1-degree 1 homogeneous component of $A$. Show that $(S \ast E)(x_1 x_2 \ldots x_n) = \ldots [[x_1, x_2], x_3], \ldots, x_n]$ for any $n \geq 1$ and any $x_1, x_2, \ldots, x_n \in V$.

Remark 1.66. Exercise 1.65 gives rise to a certain idempotent map $A \to A$ when $k$ is a commutative $Q$-algebra and $A$ is a cocommutative connected graded $k$-Hopf algebra. Namely, the $k$-linear map $A \to A$ sending every homogeneous $a \in A$ to $\frac{1}{\deg a} (S \ast E)(a)$ (or 0 if $\deg a = 0$) is idempotent and is a projection on the $k$-module of primitive elements of $A$. It is called the Dynkin idempotent; see [125] for more of its properties. Part (c) of the exercise is more or less Baker’s identity.

1.6. Duals. Recall that for $k$-modules $V$, taking the dual $k$-module $V^* := \text{Hom}(V, k)$ reverses $k$-linear maps. That is, every $k$-linear map $V \to W$ induces an adjoint map $W^* \to V^*$ defined uniquely by

$$(f, \varphi(v)) = (\varphi^*(f), v)$$

\footnote{We will see another such idempotent in Exercise 5.38.}
in which \((f,v)\) is the bilinear pairing \(V^* \times V \to k\) sending \((f,v) \mapsto f(v)\). If \(V\) and \(W\) are finite free \(k\)-modules\(^{26}\), more can be said: When \(\varphi\) is expressed in terms of a basis \(\{v_i\}_{i \in I}\) for \(V\) and a basis \(\{w_j\}_{j \in J}\) for \(W\) by some matrix, the map \(\varphi^*\) is expressed by the transpose matrix in terms of the dual bases of these two bases\(^{27}\).

The correspondence \(\varphi \mapsto \varphi^*\) between \(k\)-linear maps \(V \to W\) and \(k\)-linear maps \(W^* \to V^*\) is one-to-one when \(W\) is finite free. However, this is not the case in many combinatorial situations (in which \(W\) is usually free but not finite free). Fortunately, many of the good properties of finite free modules carry over to a when \(V\) is finite free. However, this is not the case in many combinatorial situations (in which \(W\) is usually free but not finite free). Fortunately, many of the good properties of finite free modules carry over to a dual situation. Let us first introduce the latter:

When \(V = \bigoplus_{n \geq 0} V_n\) is a graded \(k\)-module, note that the dual \(V^* = \prod_{n \geq 0} (V_n)^*\) can contain functionals \(f\) supported on infinitely many \(V_n\). However, we can consider the \(k\)-submodule \(V^o := \bigoplus_{n \geq 0} (V_n)^* \subset \prod_{n \geq 0} (V_n)^* = V^*,\) sometimes called the graded dual\(^{28}\), consisting of the functions \(f\) that vanish on all but finitely many \(V_n\). Notice that \(V^o\) is graded, whereas \(V^*\) (in general) is not. If \(V \to W\) is a \(k\)-linear map, then the adjoint map \(W^* \to V^*\) restricts to a graded \(k\)-linear map \(W^o \to V^o\), which we (abusively) still denote by \(\varphi^*\).

When the graded \(k\)-module \(V = \bigoplus_{n \geq 0} V_n\) is of finite type, meaning that each \(V_n\) is a finite free \(k\)-module\(^{29}\), the graded \(k\)-module \(V^o\) is again of finite type\(^{30}\) and satisfies \((V^o)^o \cong V\). Many other properties of finite free modules are salvaged in this situation; most importantly: The correspondence \(\varphi \mapsto \varphi^*\) between graded \(k\)-linear maps \(V \to W\) and graded \(k\)-linear maps \(W^o \to V^o\) is one-to-one when \(W\) is of finite type\(^{31}\).

Reversing the diagrams should then make it clear that, in the finite free or finite-type situation, duals of algebras are coalgebras, and vice-versa, and duals of bialgebras or Hopf algebras are bialgebras or Hopf algebras. For example, the product in a Hopf algebra \(A\) of finite type uniquely defines the coproduct of \(A^o\) via adjointness:

\[(\Delta_{A^o}(f), a \otimes b)_{A^o \otimes A} = (f, ab)_A.\]

Thus if \(A\) has a basis \(\{a_i\}_{i \in I}\) with product structure constants \(\{c^i_{j,k}\}\), meaning

\[a_ja_k = \sum_{i \in I} c^i_{j,k} a_i,
\]

then the dual basis \(\{f_i\}_{i \in I}\) has the same \(\{c^i_{j,k}\}\) as its coproduct structure constants:

\[\Delta_{A^o}(f_i) = \sum_{(j,k) \in I \times I} c^i_{j,k} f_j \otimes f_k.\]

The assumption that \(A\) be of finite type was indispensable here; in general, the dual of a \(k\)-algebra does not become a \(k\)-coalgebra. However, the dual of a \(k\)-coalgebra is still a \(k\)-algebra, as shown in the following exercise:

**Exercise 1.67.** For any two \(k\)-modules \(U\) and \(V\), let \(\rho_{UV} : U^* \otimes V^* \to (U \otimes V)^*\) be the \(k\)-linear map which sends every tensor \(f \otimes g \in U^* \otimes V^*\) to the composition \(U \otimes V \overset{f \otimes g}{\longrightarrow} k \otimes k \overset{m_k}{\longrightarrow} k\) of the map\(^{32}\) \(f \otimes g\) with the canonical isomorphism \(k \otimes k \overset{m_k}{\longrightarrow} k\). When \(k\) is a field and \(U\) is finite-dimensional, this map \(\rho_{UV}\) is a

\(^{26}\)A \(k\)-module is said to be finite free if it has a finite basis. If \(k\) is a field, then a finite free \(k\)-module is the same as a finite-dimensional \(k\)-vector space.

\(^{27}\)If \(\{v_i\}_{i \in I}\) is a basis of a finite free \(k\)-module \(V\), then the dual basis of this basis is defined as the basis \(\{f_i\}_{i \in I}\) of \(V^*\) that satisfies \((f_i, v_j) = \delta_{i,j}\) for all \(i\) and \(j\). (Recall that \(\delta_{i,j}\) is the Kronecker delta; \(\delta_{i,j} = 1\) if \(i = j\) and \(0\) else.)

\(^{28}\)Do not mistake this for the coalgebraic restricted dual \(A^o\) of \([160, \S 6.0]\).

\(^{29}\)This meaning of “finite type” can differ from the standard one.

\(^{30}\)More precisely, let \(V = \bigoplus_{n \geq 0} V_n\) be of finite type, and let \(\{v_i\}_{i \in I}\) be a basis \(\text{basis}\) of \(V\), that is, a basis of the \(k\)-module \(V\) such that the indexing set \(I\) is partitioned into subsets \(I_0, I_1, I_2, \ldots\) (which are allowed to be empty) with the property that, for every \(n \in \mathbb{N}\), the subfamily \(\{v_i\}_{i \in I_n}\) is a basis of the \(k\)-module \(V_n\). Then, we can define a family \(\{f_i\}_{i \in J}\) of elements of \(V^o\) by setting \((f_i, v_j) = \delta_{i,j}\) for all \(i, j \in I\). This family \(\{f_i\}_{i \in I}\) is a graded basis of the graded \(k\)-module \(V^o\). (Actually, for every \(n \in \mathbb{N}\), the subfamily \(\{f_i\}_{i \in I_n}\) is a basis of the \(k\)-submodule \((V^o)^n\) of \(V^o\) — in fact, the dual basis to the basis \(\{v_i\}_{i \in I_n}\) of \(V_n\).)

This basis \(\{f_i\}_{i \in I}\) is said to be the dual basis to the basis \(\{v_i\}_{i \in I}\) of \(V\).

\(^{31}\)Only in the case of finite free modules, \(V\) can be any graded \(k\)-module.

\(^{32}\)Keep in mind that the tensor \(f \otimes g \in U^* \otimes V^*\) is not the same as the map \(U \otimes V \overset{f \otimes g}{\longrightarrow} k \otimes k\).
k-vector space isomorphism (and usually regarded as the identity); more generally, it is injective whenever k is a field. Also, let \( s : k \to k^* \) be the canonical isomorphism. Prove that:

(a) If \( C \) is a k-coalgebra, then \( C^* \) becomes a k-algebra if we define its associative operation by \( m_{C^*} = \Delta_C \circ \rho_{C,C} : C^* \otimes C^* \to C^* \) and its unit map to be \( \epsilon_C' \circ s : k \to C^* \).

(b) The k-algebra structure defined on \( C^* \) in part (a) is precisely the one defined on \( \text{Hom} (C, k) = C^* \) in Definition 1.27 applied to \( A = k \).

(c) If \( C \) is a graded k-coalgebra, then \( C^o \) is a k-subalgebra of the k-algebra \( C^* \) defined in part (a).

(d) If \( f : C \to D \) is a homomorphism of k-coalgebras, then \( f^* : D^* \to C^* \) is a homomorphism of k-algebras.

(e) Let \( U \) be a graded k-module (not necessarily of finite type), and let \( V \) be a graded k-module of finite type. Then, there is a 1-to-1 correspondence between graded k-linear maps \( U \to V \) and graded k-linear maps \( V^o \to U^o \) given by \( f \mapsto f^* \).

(f) Let \( C \) be a graded k-coalgebra (not necessarily of finite type), and let \( D \) be a graded k-coalgebra of finite type. Part (e) of this exercise shows that there is a 1-to-1 correspondence between graded k-linear maps \( C \to D \) and graded k-linear maps \( D^o \to C^o \) given by \( f \mapsto f^* \). This correspondence has the property that a given graded k-linear map \( f : C \to D \) is a k-coalgebra map if and only if \( f^* : D^o \to C^o \) is a k-algebra map.

Another example of a Hopf algebra is provided by the so-called shuffle algebra. Before we introduce it, let us define the *shuffles* of two words:

**Definition 1.68.** Given two words \( a = (a_1, a_2, \ldots, a_n) \) and \( b = (b_1, b_2, \ldots, b_m) \), the *multiset of shuffles of \( a \) and \( b \)* is defined as the multiset

\[
\left\{ (c_{w(1)}, c_{w(2)}, \ldots, c_{w(n+m)}) : w \in \text{Sh}_{n,m} \right\}_{\text{multiset}},
\]

where \( (c_1, c_2, \ldots, c_{n+m}) \) is the concatenation \( a \cdot b = (a_1, a_2, \ldots, a_n, b_1, b_2, \ldots, b_m) \), and where \( \text{Sh}_{n,m} \) is the subset

\[
\left\{ w \in \mathfrak{S}_{n+m} : w^{-1}(1) < w^{-1}(2) < \cdots < w^{-1}(n); \ w^{-1}(n+1) < w^{-1}(n+2) < \cdots < w^{-1}(n+m) \right\}
\]

of the symmetric group \( \mathfrak{S}_{n+m} \). Informally speaking, the shuffles of the two words \( a \) and \( b \) are the words obtained by overlaying the words \( a \) and \( b \), after first moving their letters apart so that no letters get superimposed when the words are overlayed. In particular, any shuffle of \( a \) and \( b \) contains \( a \) and \( b \) as subsequences. The multiset of shuffles of \( a \) and \( b \) has \( \binom{m+n}{n} \) elements (counted with multiplicity) and is denoted by \( a \shuffle b \). For instance, the shuffles of \((1,2,1)\) and \((3,2)\) are

\[
(1,2,1,3,2), (1,2,3,2,1), (1,2,3,2,1), (1,3,2,2,1), (1,3,2,2,1),
\]

listed here as often as they appear in the multiset \((1,2,1) \shuffle (3,2)\). Here we have underlined the letters taken from \( a \) – that is, the letters at positions \( w^{-1}(1), w^{-1}(2), \ldots, w^{-1}(n) \).

**Example 1.69.** When \( A = T(V) \) is the tensor algebra for a finite free k-module \( V \), having k-basis \( \{x_i\}_{i \in I} \), its graded dual \( A^o \) is another Hopf algebra whose basis \( \{y_{(i_1, \ldots, i_t)}\} \) (the dual basis of the basis \( \{x_{i_1, \ldots, i_t}\} \) of \( A = T(V) \)) is indexed by words in the alphabet \( I \). This Hopf algebra \( A^o \) could be called the *shuffle algebra* of \( V^* \). (To be more precise, it is isomorphic to the shuffle algebra of \( V^* \) introduced in Proposition 1.73.

---

34Over arbitrary rings it does not have to be even that!

35If \( C \) is a finite free k-module, then this k-algebra structure is the same as the one defined above by adjointness. But the advantage of the new definition is that it works even if \( C \) is not a finite free k-module.

36Warning: This definition of \( \text{Sh}_{n,m} \) is highly nonstandard, and many authors define \( \text{Sh}_{n,m} \) to be the set of the inverses of the permutations belonging to what we call \( \text{Sh}_{n,m} \).

37For instance, if \( a = (1,3,2,1) \) and \( b = (2,4) \), then the shuffle \((1,2,3,2,4,1)\) of \( a \) and \( b \) can be obtained by moving the letters of \( a \) and \( b \) apart as follows:

\[
a = 1 \ 3 \ 2 \ \ 1 \\
b = 2 \ 4
\]

and then overlaying them to obtain \( 1 \ 2 \ 3 \ 2 \ 4 \ 1 \). Other ways of moving letters apart lead to further shuffles (not always distinct).
Equivalently, one has
\[ \Delta y(i_1, \ldots, i_\ell) = \sum_{j=0}^{\ell} y(i_1, \ldots, i_j) \otimes y(i_{j+1}, \ldots, i_\ell). \]

For example,
\[ \Delta y_{abc} = y \otimes y_{abc} + y_a \otimes y_{cb} + y_{ab} \otimes y_{bc} + y_{abc} \otimes y_b + y_{abc} \otimes y_w. \]

Duality also shows that the shuffle product in \( A^o \) is defined by
\[ y(i_1, \ldots, i_\ell) y(j_1, \ldots, j_m) = \sum_{k=(k_1, \ldots, k_{\ell+m}) \in i \sqcup j} y(k_1, \ldots, k_{\ell+m}) \]
where \( i \sqcup j \) (as in Definition 1.68) denotes the multiset of the \( \binom{\ell+m}{\ell} \) words obtained as shuffles of the two words \( i = (i_1, \ldots, i_\ell) \) and \( j = (j_1, \ldots, j_m) \). For example,
\[ y_{abc} y_{cb} = y_{abc} + y_{acb} + y_{cab} + y_{caba} + y_{cab} \]
\[ = y_{abc} + 2y_{acb} + 2y_{cab} + y_{cab} \]

Equivalently, one has
\[ y(i_1, i_2, \ldots, i_\ell) y(i_{\ell+1}, i_{\ell+2}, \ldots, i_{\ell+m}) = \sum_{\sigma \in S_{h_{\ell,m}}} y(i_{\sigma(1)}, i_{\sigma(2)}, \ldots, i_{\sigma(\ell+m)}) \]
\[ = \sum_{\sigma \in S_{h_{\ell,m}} \ell, m} y(i_{\sigma(1)}, i_{\sigma(2)}, \ldots, i_{\sigma(\ell+m)}) \]

(using the notations of Definition 1.68 again). Lastly, the antipode \( S \) of \( A^o \) is the adjoint of the antipode of \( A = T(V) \) described in (1.15):
\[ S y(i_1, i_2, \ldots, i_\ell) = (-1)^\ell y(i_\ell, \ldots, i_1). \]

Since the coalgebra \( T(V) \) is cocommutative, its graded dual \( T(V)^o \) is commutative.

**Exercise 1.70.** Let \( V \) be a 1-dimensional free \( k \)-module with basis element \( x \), so \( \text{Sym}(V) \cong k[x] \), with \( k \)-basis \( \{1 = x^0, x^1, x^2, \ldots \} \).

(a) Check that the powers \( x^i \) satisfy
\[ x^i \cdot x^j = x^{i+j} \]
\[ \Delta(x^n) = \sum_{i+j=n} \binom{n}{i} x^i \otimes x^j \]
\[ S(x^n) = (-1)^n x^n \]

(b) Check that the dual basis elements \( \{f^{(0)}, f^{(1)}, f^{(2)}, \ldots\} \) for \( \text{Sym}(V)^o \), defined by \( f^{(i)}(x^j) = \delta_{i,j} \), satisfy
\[ f^{(i)} f^{(j)} = \binom{i+j}{i} f^{(i+j)} \]
\[ \Delta(f^{(n)}) = \sum_{i+j=n} f^{(i)} \otimes f^{(j)} \]
\[ S(f^{(n)}) = (-1)^n f^{(n)} \]

(c) Show that if \( Q \) is a subring of \( k \), then the \( k \)-linear map \( \text{Sym}(V)^o \to \text{Sym}(V) \) sending \( f^{(n)} \mapsto \frac{x^n}{m} \) is a graded Hopf isomorphism.

For this reason, the Hopf structure on \( \text{Sym}(V)^o \) is called a divided power algebra.

(d) Show that when \( k \) is a field of characteristic \( p > 0 \), one has \( f^{(1)} p = 0 \), and hence why there can be no Hopf isomorphism \( \text{Sym}(V)^o \to \text{Sym}(V) \).

---

\[ ^{37} \text{Specifically, } A^o \text{ has the disadvantages of being defined only when } V^* \text{ is the dual of a finite free } k \text{-module } V, \text{ and depending on a choice of basis, whereas Proposition 1.73 will define shuffle algebras in full generality and canonically.} \]
Exercise 1.71. Let $V$ have $k$-basis $\{x_1, \ldots, x_n\}$, and let $V \oplus V$ have $k$-basis $\{x_1, \ldots, x_n, y_1, \ldots, y_n\}$, so that one has isomorphisms

$$\text{Sym}(V \oplus V) \cong k[x, y] \cong k[x] \otimes k[y] \cong \text{Sym}(V) \otimes \text{Sym}(V).$$

(a) Show that our usual coproduct on $\text{Sym}(V)$ can be re-expressed as follows:

$$\xymatrix{ \text{Sym}(V) \ar[d] \ar[r]^(0.6){\Delta} & \text{Sym}(V) \otimes \text{Sym}(V) \ar[d] \\
\mathbb{k}[x] & \mathbb{k}[x, y] \ar[l]_f }$$

$$f(x_1, \ldots, x_n) \mapsto f(x_1 + y_1, \ldots, x_n + y_n)$$

In other words, it is induced from the diagonal map

$$V \longrightarrow V \oplus V \quad x_i \longmapsto x_i + y_i$$

(1.24)

(b) One can similarly define a coproduct on the exterior algebra $\wedge V$, which is the quotient $T(V)/J$ where $J$ is the two-sided ideal generated by the elements $\{x^2 (= x \otimes x)\}_{x \in V}$ in $T^2(V)$. This becomes a graded commutative algebra

$$\wedge V = \bigoplus_{d=0}^{\infty} \wedge^d V \quad (\cong \bigoplus_{d=0}^{\infty} \wedge^d V),$$

if one views the elements of $V = \wedge^1 V$ as having odd degree, and uses the topologist’s sign conventions (as in (1.7)). One again has $\wedge(V \oplus V) = \wedge V \otimes \wedge V$ as graded algebras. Show that one can again let the diagonal map (1.24) induce a map

$$\wedge(V) \longrightarrow \wedge V \otimes \wedge V \quad f(x_1, \ldots, x_n) \longmapsto f(x_1 + y_1, \ldots, x_n + y_n)$$

(1.25)

$$\sum c_{i_1, \ldots, i_d} \cdot x_{i_1} \wedge \cdots \wedge x_{i_d} \longrightarrow \sum c_{i_1, \ldots, i_d} \cdot (x_{i_1} + y_{i_1}) \wedge \cdots \wedge (x_{i_d} + y_{i_d})$$

which makes $\wedge V$ into a graded connected Hopf algebra.

(c) Show that in the tensor algebra $T(V)$, if one views the elements of $V = V^{\otimes 1}$ as having odd degree, and uses the convention (1.7) in the twist map when defining $T(V)$, then for any $x$ in $V$ one has $\Delta(x^2) = 1 \otimes x^2 + x^2 \otimes 1$.

(d) Use part (c) to show that the two-sided ideal $J \subset T(V)$ generated by $\{x^2\}_{x \in V}$ is also a two-sided coideal, and hence the quotient $\wedge V = T(V)/J$ inherits the structure of a bialgebra. Check that the coproduct on $\wedge V$ inherited from $T(V)$ is the same as the one defined in part (b).

[Hint: The ideal $J$ in part (b) is a graded $k$-submodule of $T(V)$, but this is not completely obvious (not all elements of $V$ have to be homogeneous!).]

Exercise 1.72. Let $C$ be a $k$-coalgebra. As we know from Exercise 1.67(a), this makes $C^*$ into a $k$-coalgebra.

Let $A$ be a $k$-algebra which is finite free as $k$-module. This makes $A^*$ into a $k$-coalgebra.

Let $f : C \rightarrow A$ and $g : C \rightarrow A$ be two $k$-linear maps. Show that $f^* \star g^* = (f \star g)^*$.

The above arguments might have created the impression that duals of bialgebras have good properties only under certain restrictive conditions (e.g., the dual of a bialgebra $H$ does not generally become a bialgebra unless $H$ is of finite type), and so they cannot be used in proofs and constructions unless one is willing to sacrifice some generality (e.g., we had to require $V$ to be finite free in Example 1.69). While the first part of this impression is true, the second is not always true; often there is a way to gain back the generality lost from using duals. As an example of this, let us define the shuffle algebra of an arbitrary $k$-module (not just of a dual of a finite free $k$-module) as in Example 1.69):

Proposition 1.73. Let $V$ be a $k$-module. Define a $k$-linear map $\Delta_{\shuffle} : T(V) \rightarrow T(V) \otimes T(V)$ by setting

$$\Delta_{\shuffle}(v_1 v_2 \cdots v_n) = \sum_{k=0}^{n} (v_1 v_2 \cdots v_k) \otimes (v_{k+1} v_{k+2} \cdots v_n) \quad \text{for all } n \in \mathbb{N} \text{ and } v_1, v_2, \ldots, v_n \in V.$$
Define a \( k \)-bilinear map \( \shuffle : T(V) \times T(V) \to T(V) \), which will be written in infix notation (that is, we will write \( a \shuffle b \) instead of \( \shuffle(a,b) \)), by setting \(^{39}\)
\[
(v_1 v_2 \cdots v_\ell) \shuffle (v_{\ell+1} v_{\ell+2} \cdots v_{\ell+m}) = \sum_{\sigma \in \text{Sh}_{\ell,m}} v_{\sigma(1)} v_{\sigma(2)} \cdots v_{\sigma(\ell+m)}
\]
for all \( \ell, m \in \mathbb{N} \) and \( v_1, v_2, \ldots, v_{\ell+m} \in V \).

Consider also the comultiplication \( \epsilon \) of the Hopf algebra \( T(V) \).

Then, the \( k \)-module \( T(V) \), endowed with the multiplication \( \shuffle \), the unit \( 1_{T(V)} \in V^{\otimes 0} \subset T(V) \), the comultiplication \( \Delta_{\shuffle} \) and the counit \( \epsilon \), becomes a commutative Hopf algebra. This Hopf algebra is called the shuffle algebra of \( V \), and denoted by \( \text{Sh}(V) \). The antipode of the Hopf algebra \( \text{Sh}(V) \) is precisely the antipode \( S \) of \( T(V) \).

**Exercise 1.74.** Prove Proposition 1.73.

[**Hint:**] When \( V \) is a finite free \( k \)-module, Proposition 1.73 follows from Example 1.69. The trick is to derive the general case from this specific one. Every \( k \)-linear map \( f : W \to V \) between two \( k \)-modules \( W \) and \( V \) induces a map \( T(f) : T(W) \to T(V) \) which preserves \( \Delta_{\shuffle} \), \( \epsilon \), \( 1_{T(V)} \), \( \epsilon \) and \( S \) (in the appropriate meanings – e.g., preserving \( \Delta_{\shuffle} \) means \( \Delta_{\shuffle} \circ T(f) = (T(f) \otimes T(f)) \circ \Delta_{\shuffle} \)). Show that each of the equalities that need to be proven in order to verify Proposition 1.73 can be “transported” along such a map \( T(f) \) from a \( T(W) \) for a suitably chosen finite free \( k \)-module \( W \).]

It is also possible to prove Proposition 1.73 “by foot”, as long as one is ready to make combinatorial arguments about cutting shuffles.

**Remark 1.75.** (a) Let \( V \) be a finite free \( k \)-module. The Hopf algebra \( T(V)^\circ \) (studied in Example 1.69) is naturally isomorphic to the shuffle algebra \( \text{Sh}(V^\star) \) (defined as in Proposition 1.73 but for \( V^\star \) instead of \( V \) as Hopf algebras, by the obvious isomorphism (namely, the direct sum of the isomorphisms \((V^{\otimes n})^* \to (V^\star)^{\otimes n} \) over all \( n \in \mathbb{N} \)).

(b) The same statement applies to the case when \( V \) is a graded \( k \)-module of finite type satisfying \( V_0 = 0 \) rather than a finite free \( k \)-module, provided that \( V^\star \) and \((V^{\otimes n})^* \) are replaced by \( V^\circ \) and \((V^{\otimes n})^\circ \).

We shall return to shuffle algebras in Section 6.3, where we will show that under certain conditions (\( \mathbb{Q} \) being a subring of \( k \), and \( V \) being a free \( k \)-module) the algebra structure on a shuffle algebra \( \text{Sh}(V) \) is a polynomial algebra in an appropriately chosen set of generators.\(^{42}\)

---

\(^{38}\)This is well-defined, because the right hand side is \( n \)-multilinear in \( v_1, v_2, \ldots, v_n \), and because any \( n \)-multilinear map \( V^{\times n} \to M \) into a \( k \)-module \( M \) gives rise to a unique \( k \)-linear map \( V^{\otimes n} \to M \).

\(^{39}\)Many authors use the symbol \( \shuffle \) instead of \( \shuffle \) here, but we prefer to reserve the former notation for the shuffle product of words.

\(^{40}\)Again, this is well-defined by the \( \ell + m \)-multilinearity of the right hand side.

\(^{41}\)This can be verified by comparing (1.20) with the definition of \( \Delta_{\shuffle} \), and comparing (1.23) with the definition of \( \shuffle \).

\(^{42}\)This says nothing about the coalgebra structure on \( \text{Sh}(V) \) – which is much more complicated in these generators.
2. Review of symmetric functions \( \Lambda \) as Hopf algebra

Here we review the ring of symmetric functions, borrowing heavily from standard treatments, such as Macdonald [106, Chap. I], Sagan [138, Chap. 4], and Stanley [154, Chap. 7], but emphasizing the Hopf structure early on.

2.1. Definition of \( \Lambda \). As before, \( k \) here is a commutative ring (hence could be a field or the integers \( \mathbb{Z} \); these are the usual choices). Given an infinite variable set \( x = (x_1, x_2, \ldots) \), a monomial \( x^\alpha := x_1^{\alpha_1} x_2^{\alpha_2} \cdots \) is indexed by an element \( \alpha = (\alpha_1, \alpha_2, \ldots) \) in \( \mathbb{N}^\infty \) having finite support; such \( \alpha \) are called weak compositions.

The nonzero ones among the integers \( \alpha_1, \alpha_2, \ldots \) are called the parts of the weak composition \( \alpha \). We will consider the ring \( R(\mathbf{x}) \) of formal power series \( f(\mathbf{x}) = \sum c_\alpha x^\alpha \) with \( c_\alpha \) in \( k \) of bounded degree, that is, where there exists some bound \( d = d(f) \) for which \( \deg(x^\alpha) := \sum \alpha_i > d \) implies \( c_\alpha = 0 \). It is easy to see that the product of two such power series is well-defined, and also has bounded degree.

The symmetric group \( S_n \) permuting the first \( n \) variables \( x_1, \ldots, x_n \) acts as a group of automorphisms on \( R(\mathbf{x}) \), as does the union \( S_\infty = \bigcup_{n \geq 0} S_n \) of the infinite ascending chain \( S_0 \subset S_1 \subset S_2 \subset \cdots \) of symmetric groups. This group \( S_\infty \) can also be described as the group of all permutations of the set \( \{1, 2, 3, \ldots\} \) which leave all but finitely many elements invariant.

**Definition 2.1.** The ring of symmetric functions in \( \mathbf{x} \) with coefficients in \( k \), denoted \( \Lambda = \Lambda_k = \Lambda(\mathbf{x}) = \Lambda_k(\mathbf{x}) \), is the \( S_\infty \)-invariant subalgebra \( R(\mathbf{x})^{S_\infty} \) of \( R(\mathbf{x}) \):

\[
\Lambda := \left\{ f = \sum c_\alpha x^\alpha \in R(\mathbf{x}) : c_\alpha = c_\beta \text{ if } \alpha, \beta \text{ lie in the same } S_\infty \text{-orbit} \right\}.
\]

We refer to the elements of \( \Lambda \) as symmetric functions (over \( k \)); however, despite this terminology, they are not functions in the usual sense.\(^{43}\)

Note that \( \Lambda \) is a graded \( k \)-algebra, since \( \Lambda = \bigoplus_{n \geq 0} \Lambda_n \) where \( \Lambda_n \) are the symmetric functions \( f = \sum c_\alpha x^\alpha \) which are homogeneous of degree \( n \), meaning \( \deg(x^\alpha) = n \) for all \( c_\alpha \neq 0 \).

**Exercise 2.2.** Let \( f \in R(\mathbf{x}) \). Let \( A \) be a commutative \( k \)-algebra, and \( a_1, a_2, \ldots, a_k \) be finitely many elements of \( A \). Show that substituting \( a_1, a_2, \ldots, a_k, 0, 0, 0, \ldots \) for \( x_1, x_2, x_3, \ldots \) in \( f \) yields an infinite sum in which all but finitely many addends are zero. Hence, this sum has a value in \( A \), which is commonly denoted by \( f(a_1, a_2, \ldots, a_k) \).

**Definition 2.3.** A partition \( \lambda = (\lambda_1, \lambda_2, \ldots, \lambda_\ell, 0, 0, \ldots) \) is a weak composition whose parts weakly decrease: \( \lambda_1 \geq \cdots \geq \lambda_\ell > 0 \). The (uniquely defined) \( \ell \) is said to be the length of the partition \( \lambda \) and denoted by \( \ell(\lambda) \). One sometimes omits trailing zeroes from a partition: e.g., one can write the partition \((3, 1, 0, 0, 0, \ldots)\) as \((3, 1)\). We will often (but not always) write \( \lambda_i \) for the \( i \)-th entry of the partition \( \lambda \) (for instance, \( \lambda_2 = 3 \) if \( \lambda = (5, 3, 1, 1) \)) and call it the \( i \)-th part of \( \lambda \). The sum \( \lambda_1 + \lambda_2 + \cdots + \lambda_\ell = \lambda_1 + \lambda_2 + \cdots \) (where \( \ell = \ell(\lambda) \)) of all parts of \( \lambda \) is called the size of \( \lambda \) and denoted by \( |\lambda| \); for a given integer \( n \), the partitions of size \( n \) are referred to as the partitions of \( n \). The empty partition \( () \) is denoted by \( \emptyset \).

Every weak composition \( \alpha \) lies in the \( S_\infty \)-orbit of a unique partition \( \lambda = (\lambda_1, \lambda_2, \ldots, \lambda_\ell, 0, 0, \ldots) \) with \( \lambda_1 \geq \cdots \geq \lambda_\ell > 0 \). For any partition \( \lambda \), define the monomial symmetric function

\[
m_\lambda := \sum_{\alpha \in S_\infty \lambda} x^\alpha.
\]

Letting \( \lambda \) run through the set \( \text{Par} \) of all partitions, this gives the monomial \( k \)-basis \( \{m_\lambda\} \) of \( \Lambda \). Letting \( \lambda \) run only through the set \( \text{Par}_n \) of partitions of \( n \) gives the monomial \( k \)-basis for \( \Lambda_n \).

**Example 2.4.** For \( n = 3 \), one has

\[
m_{(3)} = x_1^3 + x_2^3 + x_3^3 + \cdots
\]
\[
m_{(2,1)} = x_1^2 x_2 + x_1 x_2^2 + x_1^2 x_3 + x_1 x_3^2 + \cdots
\]
\[
m_{(1,1,1)} = x_1 x_2 x_3 + x_1 x_2 x_4 + x_1 x_3 x_4 + x_2 x_3 x_4 + \cdots
\]

\(^{43}\)Being power series, they can be evaluated at appropriate families of variables. But this does not make them functions (no more than polynomials are functions). The terminology “symmetric function” is thus not well-chosen; but it is standard.
Remark 2.5. We have defined the symmetric functions as the elements of $R(x)$ invariant under the group $S(\infty)$. However, they also are the elements of $R(x)$ invariant under the group $S_{\infty}$ of all permutations of the set $\{1, 2, 3, \ldots\}$. 

Remark 2.6. It is sometimes convenient to work with finite variable sets $x_1, \ldots, x_n$, which one justifies as follows. Note that the algebra homomorphism

$$R(x) \rightarrow R(x_1, \ldots, x_n) = k[x_1, \ldots, x_n]$$

which sends $x_{n+1}, x_{n+2}, \ldots$ to 0 restricts to an algebra homomorphism

$$\Lambda_k(x) \rightarrow \Lambda_k(x_1, \ldots, x_n) = k[x_1, \ldots, x_n]^{S_n}.$$

Furthermore, this last homomorphism is a $k$-linear isomorphism when restricted to $\Lambda_i$ for $0 \leq i \leq n$, since it sends the monomial basis elements $m_\lambda(x)$ to the monomial basis elements $m_\lambda(x_1, \ldots, x_n)$. Thus when one proves identities in $\Lambda_k(x_1, \ldots, x_n)$ for all $n$, they are valid in $\Lambda$, that is, $\Lambda$ is the inverse limit of the $\Lambda(x_1, \ldots, x_n)$ in the category of graded $k$-algebras.

One can also define a comultiplication on $\Lambda$ as follows. Note that when one decomposes the variables into two sets $(x, y) = (x_1, x_2, \ldots, y_1, y_2, \ldots)$, one has a ring homomorphism

$$R(x) \otimes R(x) \rightarrow R(x, y)$$

$$f(x) \otimes g(x) \rightarrow f(x)g(y).$$

This restricts to an isomorphism

$$(2.2) \quad \Lambda \otimes \Lambda = R(x)^{S(\infty)} \otimes R(x)^{S(\infty)} \rightarrow R(x, y)^{S(\infty) \times S(\infty)}$$

where $S(\infty) \times S(\infty)$ denotes permutations of (finite subsets of) the $x$ and separate permutations of (finite subsets of) the $y$, because $R(x, y)^{S(\infty) \times S(\infty)}$ has $k$-basis $\{m_\lambda(x)m_\mu(y)\}_{\lambda, \mu \in \text{Par}}$. As $S(\infty) \times S(\infty)$ is a subgroup of the group $S(\infty, \infty)$ (the group of all permutations of $\{x_1, x_2, \ldots, y_1, y_2, \ldots\}$ leaving all but finitely many variables invariant) acting on all of $(x, y)$, one gets an inclusion of rings

$$\Lambda(x, y) = R(x, y)^{S(\infty, \infty)} \rightarrow R(x, y)^{S(\infty) \times S(\infty)} \cong \Lambda \otimes \Lambda$$

where the last isomorphism is the inverse of the one in (2.2). This gives a comultiplication

$$\Lambda = \Lambda(x) \quad \Delta \mapsto \Lambda(x, y) \quad \Lambda \otimes \Lambda$$

$$f(x) = f(x_1, x_2, \ldots) \quad \mapsto \quad f(x, y) = f(x_1, x_2, \ldots, y_1, y_2, \ldots).$$

Here, $f(x_1, x_2, \ldots, y_1, y_2, \ldots)$ means the result of choosing some bijection $\phi : \{x_1, x_2, x_3, \ldots\} \rightarrow \{x_1, x_2, \ldots, y_1, y_2, \ldots\}$ and substituting $\phi(x_i)$ for every $x_i$ in $f$. (The choice of $\phi$ is irrelevant since $f$ is symmetric.)

\textit{Proof.} We need to show that $\Lambda = R(x)^{S_{\infty}}$. Since

$$\Lambda = \left\{ f = \sum_\alpha c_\alpha x^\alpha \in R(x) : c_\alpha = c_\beta \text{ if } \alpha, \beta \text{ lie in the same $S(\infty)$-orbit} \right\}$$

and

$$R(x)^{S_{\infty}} = \left\{ f = \sum_\alpha c_\alpha x^\alpha \in R(x) : c_\alpha = c_\beta \text{ if } \alpha, \beta \text{ lie in the same $S_{\infty}$-orbit} \right\},$$

this will follow immediately if we can show that two weak compositions $\alpha$ and $\beta$ lie in the same $S(\infty)$-orbit if and only if they lie in the same $S_{\infty}$-orbit. But this is straightforward to check (in fact, two weak compositions $\alpha$ and $\beta$ lie in the same orbit under either group if and only if they have the same multiset of nonzero entries).}

\textit{To be more precise, the choice of $\phi$ is irrelevant because $f$ is $S_{\infty}$-invariant, with the notations of Remark 2.5.}
Example 2.7. One has
\[
\Delta m_{(2,1)} = m_{(2,1)}(x_1, x_2, \ldots, y_1, y_2, \ldots)
\]
\[
= x_1^2 x_2 + x_1 x_2^2 + \cdots
\]
\[
+ x_1^2 y_1 + x_1 y_2^2 + \cdots
\]
\[
+ x_1 y_1^2 + x_1^2 y_2 + \cdots
\]
\[
+ y_1^2 y_2 + y_1 y_2^2 + \cdots
\]
\[
= m_{(2,1)}(x) + m_{(2)}(x)m_{(1)}(y) + m_{(1)}(x)m_{(2)}(y) + m_{(2,1)}(y)
\]
\[
= m_{(2,1)} \otimes 1 + m_{(2)} \otimes m_{(1)} + m_{(1)} \otimes m_{(2)} + 1 \otimes m_{(2,1)}.
\]
This example generalizes easily to the following formula
\[
(2.3) \quad \Delta m_\lambda = \sum_{(\mu, \nu) : \mu \sqcup \nu = \lambda} m_\mu \otimes m_\nu,
\]
in which \( \mu \sqcup \nu \) is the partition obtained by taking the multiset union of the parts of \( \mu \) and \( \nu \), and then reordering them to make them weakly decreasing.

Checking that \( \Delta \) is coassociative amounts to checking that
\[
(\Delta \otimes \text{id}) \circ \Delta f = f(x, y, z) = (\text{id} \otimes \Delta) \circ \Delta f
\]
inside \( \Lambda(x, y, z) \) as a subring of \( \Lambda \otimes \Lambda \). The counit \( \Lambda \to k \) is defined in the usual fashion for graded connected coalgebras, namely \( \epsilon \) annihilates \( I = \bigoplus_{n>0} \Lambda_n \), and \( \epsilon \) is the identity on \( \Lambda_0 = k \); alternatively \( \epsilon \) sends a symmetric function \( f(x) \) to its constant term \( f(0, 0, \ldots) \).

Note that \( \Delta \) is an algebra morphism \( \Lambda \to \Lambda \otimes \Lambda \) because it is a composition of maps which are all algebra morphisms. As the unit and counit axioms are easily checked, \( \Lambda \) becomes a graded connected \( k \)-bialgebra of finite type, and hence also a Hopf algebra by Proposition 1.36. We will identify its antipode more explicitly in Section 2.4 below.

2.2. Other Bases. We introduce the usual other bases of \( \Lambda \), and explain their significance later.

Definition 2.8. Define the families of power sum symmetric functions \( p_n \), elementary symmetric functions \( e_n \), and complete homogeneous symmetric functions \( h_n \), for \( n = 1, 2, 3, \ldots \) by
\[
(2.4) \quad p_n := x_1^n + x_2^n + \cdots = m_{(n)}
\]
\[
(2.5) \quad e_n := \sum_{i_1 < \cdots < i_n} x_{i_1} \cdots x_{i_n} = m_{(1^n)}
\]
\[
(2.6) \quad h_n := \sum_{i_1 \leq \cdots \leq i_n} x_{i_1} \cdots x_{i_n} = \sum_{\lambda \in \text{Par}_n} m_\lambda
\]
Here, we are using the multiplicative notation for partitions: whenever \((m_1, m_2, m_3, \cdots)\) is a weak composition, \((1^{m_1}2^{m_2}3^{m_3} \cdots)\) denotes the partition \( \lambda \) such that for every \( i \), the multiplicity of the part \( i \) in \( \lambda \) is \( m_i \). The \( i^{m_i} \) satisfying \( m_i = 0 \) are often omitted from this notation, and so the \((1^n)\) in \((2.5)\) means \((\underbrace{1,1, \ldots, 1}_n)\). (For another example, \((1^22^34^3) = (1^22^03^14^35^06^07^0 \cdots)\) means the partition \((4, 4, 4, 3, 1, 1, 1)\).)

By convention, also define \( h_0 = e_0 = 1 \), and \( h_n = e_n = 0 \) if \( n < 0 \). Extend these multiplicatively to partitions \( \lambda = (\lambda_1, \lambda_2, \ldots, \lambda_\ell) \) with \( \lambda_1 \geq \cdots \geq \lambda_\ell > 0 \):
\[
p_\lambda := p_{\lambda_1} p_{\lambda_2} \cdots p_{\lambda_\ell}
\]
\[
e_\lambda := e_{\lambda_1} e_{\lambda_2} \cdots e_{\lambda_\ell}
\]
\[
h_\lambda := h_{\lambda_1} h_{\lambda_2} \cdots h_{\lambda_\ell}
\]
Also define the Schur function
\[
(2.7) \quad s_\lambda := \sum_T x^{\text{cont}(T)}
\]
where $T$ runs through all column-strict tableaux of shape $\lambda$, that is, $T$ is an assignment of entries in $\{1, 2, 3, \ldots\}$ to the cells of the Ferrers diagram for $\lambda$, weakly increasing left-to-right in rows, and strictly increasing top-to-bottom in columns. Here $\cont(T)$ denotes the weak composition $([T^{-1}(1)], [T^{-1}(2)], [T^{-1}(3)], \ldots)$, so that $x^{\cont(T)} = \prod_{i} x^{|T^{-1}(i)|}$. For example,

\[
T = \begin{pmatrix}
1 & 1 & 1 & 4 & 7 \\
2 & 3 & 3 \\
4 & 4 & 6 \\
6 & 7
\end{pmatrix}
\]

is a column-strict tableau of shape $\lambda = (5, 3, 3, 2)$ with $x^{\cont(T)} = x_1^3 x_2^1 x_3^2 x_4^0 x_5^2 x_6^2$.

**Example 2.9.** One has

\[
m(1) = p(1) = e(1) = h(1) = s(1) = x_1 + x_2 + x_3 + \cdots \\
s(n) = h_n \\
s(1^n) = e_n
\]

**Example 2.10.** One has for $\lambda = (2, 1)$ that

\[
p_{(2,1)} = p_{2,1} = (x_1^2 + x_2^2 + \cdots)(x_1 + x_2 + \cdots) = m_{(2,1)} + m_{(3)} \\
e_{(2,1)} = e_2 e_1 = (x_1 x_2 + x_1 x_3 + \cdots)(x_1 + x_2 + \cdots) = m_{(2,1)} + 3m_{(1,1,1)} \\
h_{(2,1)} = h_2 h_1 = (x_1^2 + x_2^2 + \cdots + x_1 x_2 + x_1 x_3 + \cdots)(x_1 + x_2 + \cdots) = m_{(2,1)} + 2m_{(2,1)} + 3m_{(1,1,1)}
\]

and

\[
s_{(2,1)} = x_1^2 x_2 + x_1^2 x_3 + x_1 x_2^2 + x_1^2 x_3 + x_1 x_2 x_3 + x_1 x_2 x_3 + x_1 x_2 x_4 + \cdots \\
\begin{pmatrix}
1 & 1 & 12 & 13 & 12 & 12 & 12 \\
2 & 3 & 2 & 3 & 3 & 2 & 4
\end{pmatrix}
\]

\[
= m_{(2,1)} + 2m_{(1,1,1)}
\]

The Ferrers diagram of a partition $\lambda$ is defined as the set of all pairs $(i, j) \in \{1, 2, 3, \ldots\}^2$ satisfying $j \leq \lambda_i$. This is a set of cardinality $|\lambda|$. Usually, one visually represents a Ferrers diagram by drawing its elements $(i, j)$ as points on the plane, although (unlike the standard convention for drawing points on the plane) one lets the x-axis go top-to-bottom (i.e., the point $(i + 1, j)$ is one step below the point $(i, j)$), and the y-axis go left-to-right (i.e., the point $(i, j + 1)$ is one step to the right of the point $(i, j)$). This is the so-called English notation, also known as the matrix notation because it is precisely the way one labels the entries of a matrix. Other notations appear in literature, such as the French notation used, e.g., in Malvenuto’s [109], and the Russian notation used, e.g., in parts of Kerov’s [79]. These points are drawn either as dots or as square boxes; in the latter case, the boxes are centered at the points they represent, and they have sidelength 1 so that the boxes centered around $(i, j)$ and $(i, j + 1)$ touch each other along a sideline. For example, the Ferrers diagram of the partition $(3, 2, 2)$ is represented as

\[
\begin{array}{cccc}
\bullet & \bullet & \bullet \\
\bullet & \bullet & \bullet & \bullet
\end{array}
\]

or as

\[
\begin{array}{cccc}
\square & \square & \square \\
\square & \square & \square & \square
\end{array}
\]

The Ferrers diagrams of a partition $\lambda$ uniquely determines $\lambda$. One refers to the elements of the Ferrers diagram of $\lambda$ as the cells (or boxes) of this diagram (which is particularly natural when one represents them by boxes) or, briefly, as the cells of $\lambda$. Notation like “west”, “north”, “left”, “right”, “row” and “column” concerning cells of Ferrers diagrams normally refers to their visual representation.

Ferrers diagrams are also known as Young diagrams.

One can characterize the Ferrers diagrams of partitions as follows: A finite subset $S$ of $\{1, 2, 3, \ldots\}^2$ is the Ferrers diagram of some partition if and only if for every $(i, j) \in S$ and every $(i', j') \in \{1, 2, 3, \ldots\}^2$ satisfying $i' \leq i$ and $j' \leq j$, we have $(i', j') \in S$. In other words, a finite subset $S$ of $\{1, 2, 3, \ldots\}^2$ is the Ferrers diagram of some partition if and only if it is a lower set of the poset $\{1, 2, 3, \ldots\}^2$ with respect to the componentwise order.

\[
\text{To visually represent a column-strict tableau } T \text{ of shape } \lambda, \text{ we draw the same picture as when representing the Ferrers diagram of } \lambda, \text{ but with a little difference: a cell } (i, j) \text{ is no longer represented by a dot or box, but instead is represented by the entry of } T \text{ assigned to this cell. Accordingly, the entry of } T \text{ assigned to a given cell } c \text{ is often referred to as the entry of } T \text{ in } c.
\]

\[
\text{Example 2.10. For } \lambda = (2, 1), \text{ we have }
\]

\[
p_{(2,1)} = p_{2,1} = (x_1^2 + x_2^2 + \cdots)(x_1 + x_2 + \cdots) = m_{(2,1)} + m_{(3)} \\
e_{(2,1)} = e_2 e_1 = (x_1 x_2 + x_1 x_3 + \cdots)(x_1 + x_2 + \cdots) = m_{(2,1)} + 3m_{(1,1,1)} \\
h_{(2,1)} = h_2 h_1 = (x_1^2 + x_2^2 + \cdots + x_1 x_2 + x_1 x_3 + \cdots)(x_1 + x_2 + \cdots) = m_{(2,1)} + 2m_{(2,1)} + 3m_{(1,1,1)}
\]

and

\[
s_{(2,1)} = x_1^2 x_2 + x_1^2 x_3 + x_1 x_2^2 + x_1^2 x_3 + x_1 x_2 x_3 + x_1 x_2 x_3 + x_1 x_2 x_4 + \cdots \\
\begin{pmatrix}
1 & 1 & 12 & 13 & 12 & 12 & 12 \\
2 & 3 & 2 & 3 & 3 & 2 & 4
\end{pmatrix}
\]

\[
= m_{(2,1)} + 2m_{(1,1,1)}
\]
In fact, one has these transition matrices for \( n = 3 \) expressing elements in terms of the monomial basis \( m_\lambda \):

\[
\begin{pmatrix}
\begin{array}{ccc}
1 & 1 & 1 \\
0 & 1 & 3 \\
0 & 0 & 6 \\
\end{array}
\end{pmatrix}
\]

\[
\begin{pmatrix}
\begin{array}{ccc}
1 & 1 & 1 \\
1 & 2 & 3 \\
1 & 3 & 6 \\
\end{array}
\end{pmatrix}
\]

\[
\begin{pmatrix}
\begin{array}{ccc}
1 & 0 & 0 \\
1 & 1 & 0 \\
1 & 2 & 1 \\
\end{array}
\end{pmatrix}
\]

Our next goal is to show that \( e_\lambda, s_\lambda, h_\lambda \) (and, under some conditions, the \( p_\lambda \) as well) all give bases for \( \Lambda \). However at the moment it is not yet even clear that \( s_\lambda \) are symmetric!

**Proposition 2.11.** Schur functions \( s_\lambda \) are symmetric, that is, they lie in \( \Lambda \).

**Proof.** It suffices to show \( s_\lambda \) is symmetric under swapping the variables \( x_i, x_{i+1} \), by providing an involution \( \iota \) on the set of all column-strict tableaux \( T \) of shape \( \lambda \) which switches the \( \text{cont}(T) \) for \( (i, i+1) \text{cont}(T) \).

Restrict attention to the entries \( i, i+1 \) in \( T \), which must look something like this:

\[
\begin{array}{cccccccc}
\hline
i & i & i & i & i & i & i & i + 1 \\
\hline
i + 1 & i + 1 & i + 1 & i + 1 & i + 1 & i + 1 & i + 1 \\
\hline
\end{array}
\]

One finds several vertically aligned pairs \( \frac{i}{i + 1} \). If one were to remove all such pairs, the remaining entries would be a sequence of rows, each looking like this:

\[
\begin{array}{cccccccc}
\hline
i & . & . & . & i & . & . & . \\
\hline
\end{array}
\]

An involution due to Bender and Knuth tells us to leave fixed all the vertically aligned pairs \( \frac{i}{i + 1} \) but change each sequence as in (2.8) to this:

\[
\begin{array}{cccccccc}
\hline
i, i, . . . , i & i + 1, i + 1, . . . , i + 1 \\
\hline
\end{array}
\]

For example, the above configuration in \( T \) would change to

\[
\begin{array}{cccccccc}
\hline
i & i & i & i & i & i & i & i + 1 \\
\hline
i & i + 1 & i + 1 & i + 1 & i + 1 & i + 1 & i + 1 \\
\hline
\end{array}
\]

It is easily checked that this map is an involution, and that it has the effect of swapping \( (i, i + 1) \) in \( \text{cont}(T) \). \qed

**Remark 2.12.** The symmetry of Schur functions allows one to reformulate them via column-strict tableaux defined with respect to any total ordering \( \mathcal{L} \) on the positive integers, rather than the usual \( 1 < 2 < 3 < \ldots \).

For example, one can use the reverse order \( \leq R \) or even more exotic orders, such as

\[
1 < 3 < 5 < 7 < \ldots < 2 < 4 < 6 < 8 < \ldots
\]

Say that an assignment \( T \) of entries in \( \{1, 2, 3, \ldots\} \) to the cells of the Ferrers diagram of \( \lambda \) is an \( \mathcal{L} \)-column-strict tableau if it is weakly \( \mathcal{L} \)-increasing left-to-right in rows, and strictly \( \mathcal{L} \)-increasing top-to-bottom in columns.

---

\(^{48}\) This reverse order is what one uses when one defines a Schur function as a generating function for reverse semistandard tableaux or column-strict plane partitions; see Stanley [154, Proposition 7.10.4].
Proposition 2.13. For any total order $\mathcal{L}$ on the positive integers,
\begin{equation}
    s_\lambda = \sum_T x^{\text{cont}(T)}
\end{equation}
as $T$ runs through all $\mathcal{L}$-column-strict tableaux of shape $\lambda$.

Proof. Given a weak composition $\alpha = (\alpha_1, \alpha_2, \ldots)$ with $\alpha_{n+1} = \alpha_{n+2} = \cdots = 0$, assume that the integers $1, 2, \ldots, n$ are totally ordered by $\mathcal{L}$ as $w(1) <_\mathcal{L} \cdots <_\mathcal{L} w(n)$ for some $w \in \mathcal{S}_n$. Then the coefficient of $x^\alpha = x_1^{\alpha_1} \cdots x_n^{\alpha_n}$ on the right side of (2.9) is the same as the coefficient of $x^{w^{-1}(\alpha)}$ on the right side of (2.7) defining $s_\lambda$, which by symmetry of $s_\lambda$ is the same as the coefficient of $x^\alpha$ on the right side of (2.7).

It is now not hard to show that $p_\lambda, e_\lambda, s_\lambda$ give bases by a triangularity argument. For this purpose, let us introduce a useful partial order on partitions.

Definition 2.14. The dominance or majorization partial order on $\text{Par}_n$, written $\lambda \succ \mu$, is defined by
\begin{equation}
    \lambda_1 + \lambda_2 + \cdots + \lambda_k \geq \mu_1 + \mu_2 + \cdots + \mu_k \quad \text{for } k = 1, 2, \ldots, n.
\end{equation}

(The definition of dominance would not change if we would replace “for $k = 1, 2, \ldots, n$” by “for every positive integer $k$” or by “for every $k \in \mathbb{N}$”.)

Definition 2.15. For a partition $\lambda$, its conjugate or transpose partition $\lambda'$ is the one whose Ferrers diagram is obtained from that of $\lambda$ by exchanging rows for columns. Alternatively, one has this formula for its $i^{th}$ part:
\begin{equation}
    (\lambda')_i := |\{j : \lambda_j \geq i\}|.
\end{equation}

Exercise 2.16. Let $\lambda, \mu \in \text{Par}_n$. Show that $\lambda \succ \mu$ if and only if $\mu' \succ \lambda'$.

Proposition 2.17. The sets $\{e_\lambda\}, \{s_\lambda\}$ as $\lambda$ runs through all partitions give $k$-bases for $\Lambda_k$ for any commutative ring $k$. The same holds for $\{p_\lambda\}$ when $\mathbb{Q}$ is a subring of $k$.

Proof. One can restrict attention to each homogeneous component $\Lambda_n$ and partitions $\lambda$ of $n$. We check that in each case, the proposed basis expands triangularly in the $\{m_\lambda\}$ with some choice of orderings on $\text{Par}_n$ indexing the rows and columns, as illustrated in Example 2.10. Moreover, we check that the diagonal entries of this triangular expansion are invertible, which then yields that the proposed basis is indeed a basis.

One has $s_\lambda = \sum_{\mu} K_{\lambda,\mu} m_\mu$ where the coefficient $K_{\lambda,\mu}$ is the Kostka number counting the column-strict tableaux $T$ of shape $\lambda$ having $\text{cont}(T) = \mu$; this follows because both sides are symmetric functions, and $K_{\lambda,\mu}$ is the coefficient of $x^\mu$ on both sides. Since for each positive integer $k$, the entries $1, 2, \ldots, k$ in $T$ must all lie within the first $k$ rows of $\lambda$, one has that $\lambda_1 + \lambda_2 + \cdots + \lambda_k \geq \mu_1 + \mu_2 + \cdots + \mu_k$, that is, $\lambda \succ \mu$ (if $K_{\lambda,\mu} \neq 0$). Hence, the expansion $s_\lambda = \sum_{\mu} K_{\lambda,\mu} m_\mu$ becomes $s_\lambda = \sum_{\mu, \lambda \succ \mu} K_{\lambda,\mu} m_\mu$. One can also check that $K_{\lambda,\lambda} = 1$, so this expansion is unimodular with appropriate ordering of rows and columns.

Before we show that $\{e_\lambda\}$ is a basis, we define a few notations regarding integer matrices. A $\{0,1\}$-matrix means a matrix whose entries belong to the set $\{0,1\}$. If $A \in \mathbb{N}^{m \times n}$ is a matrix, then the row sums of $A$ means the $\ell$-tuple $(r_1, r_2, \ldots, r_\ell)$, where each $r_i$ is the sum of all entries in the $i^{th}$ row of $A$; similarly, the column sums of $A$ means the $m$-tuple $(c_1, c_2, \ldots, c_m)$, where each $c_j$ is the sum of all entries in the $j^{th}$ column of $A$.

One has $e_\lambda = \sum_{\mu} a_{\lambda,\mu} m_\mu$ where $a_{\lambda,\mu}$ counts the $\{0,1\}$-matrices having row sums $\lambda$ and column sums $\mu$: when one expands $e_{\lambda_1} e_{\lambda_2} \cdots$, choosing the monomial $x_{j_1} \cdots x_{j_\lambda_1}$ in the $e_{\lambda_1}$ factor corresponds to putting $1$’s in the $i^{th}$ row and columns $j_1, \ldots, j_{\lambda_1}$ of the $\{0,1\}$-matrix. It is not hard to check that $a_{\lambda,\mu}$ vanishes unless $\lambda' \succ \mu$. One can also check that $a_{\lambda_{\lambda'} \lambda} = 1$, so this expansion is again unimodular with appropriate ordering of rows and columns.

Assume now that $\mathbb{Q}$ is a subring of $k$. One has $p_\lambda = \sum_{\mu} b_{\lambda,\mu} m_\mu$ where $b_{\lambda,\mu}$ counts the ways to partition the nonzero parts $\lambda_1, \ldots, \lambda_\ell$ into blocks such that the sums of the blocks give $\mu$; more formally, $b_{\lambda,\mu}$ is the number of maps $\varphi : \{1, 2, \ldots, \ell\} \to \{1, 2, 3, \ldots\}$ having $\mu_j = \sum_{i : \varphi(i) = j} \lambda_i$ for $j = 1, 2, \ldots$. Again it is not hard to check that $b_{\lambda,\mu}$ vanishes unless $\lambda \preceq \mu$, and hence this expansion is triangular, for appropriate ordering of rows and columns (but not unimodular, as $b_{\lambda,\lambda} \neq 1$ in general). The diagonal entries $b_{\lambda,\lambda}$ are positive integers and thus invertible in $k$, so $\{p_\lambda\}$ is a basis.

\[49\text{This is the easy implication in the Gale-Ryser Theorem.}\]
Remark 2.18. When \( \mathbb{Q} \) is not a subring of \( \mathbb{k} \), the family \( \{ p_\lambda \} \) is not (in general) a basis of \( \Lambda_k \); for instance, \( e_2 = \frac{1}{2} (p_{(1,1)} - p_2) \in \mathbb{Q} \) is not in the \( \mathbb{Z} \)-span of this family. However, if we define \( b_{\lambda,\mu} \) as in the above proof, then the \( \mathbb{Z} \)-linear span of all \( p_\lambda \) equals the \( \mathbb{Z} \)-linear span of all \( b_{\lambda,\lambda} m_\lambda \). Indeed, if \( \mu = (\mu_1, \mu_2, \ldots, \mu_k) \) with \( k = \ell(\mu) \), then \( b_{\lambda,\mu} \) is the size of the subgroup of \( \mathcal{S}_k \) consisting of all permutations \( \sigma \in \mathcal{S}_k \) having each \( i \) satisfy \( \mu_{\sigma(i)} = \mu_i \). As a consequence, \( b_{\mu,\mu} \) divides \( b_{\lambda,\mu} \) for every partition \( \mu \) of the same size as \( \lambda \) (because this group acts\(^50\) freely on the set which is enumerated by \( b_{\lambda,\mu} \)), so that by rescaling \( m_\lambda \) with the factor \( b_{\lambda,\lambda} \) we obtain a unitriangular integer change-of-basis matrix from the family \( \{ p_\lambda \}_{\lambda \in \text{Par}_n} \) to the family \( \{ b_{\lambda,\lambda} m_\lambda \}_{\lambda \in \text{Par}_n} \).

2.3. Comultiplications. Thinking about comultiplication \( \Lambda \xrightarrow{\Delta} \Lambda \otimes \Lambda \) on Schur functions forces us to immediately confront the following.

Definition 2.19. For partitions \( \mu, \lambda \) say that \( \mu \subseteq \lambda \) if \( \mu_i \leq \lambda_i \) for \( i = 1, 2, \ldots \), so the Ferrers diagram for \( \mu \) is a subset of the cells for the Ferrers diagram of \( \lambda \). In this case, define the \( \text{skew (Ferrers) diagram} \) \( \lambda/\mu \) to be their set difference.\(^51\)

Then define the \( \text{skew Schur function} \) \( s_{\lambda/\mu}(x) \) to be the sum \( s_{\lambda/\mu} := \sum_T x^{\text{cont}(T)} \), where the sum ranges over all \( \text{column-strict tableaux} \) \( T \) of shape \( \lambda/\mu \); that is, assignments of a value in \( \{1, 2, 3, \ldots\} \) to each cell of \( \lambda/\mu \), weakly increasing left-to-right in rows, and strictly increasing top-to-bottom in columns.

Example 2.20. \[ \begin{array}{ccc} & & 2 & 5 \\ & 1 & 1 \\ 2 & 2 & 4 \\ 4 & 5 \end{array} \]

This is a column-strict tableau of shape \( \lambda/\mu = (5, 3, 2)/(3, 1, 0, 0) \) and it has \( x^{\text{cont}(T)} = x_1^2 x_2^3 x_3^0 x_4^2 x_5^2 \).

Remark 2.21. If \( \mu \) and \( \lambda \) are partitions such that \( \mu \subseteq \lambda \), then \( s_{\lambda/\mu} \in \Lambda \). (This is proven similarly as Proposition 2.11.) Actually, if \( \mu \subseteq \lambda \), then \( s_{\lambda/\mu} \in \Lambda_{(\lambda/\mu)} \), where \( |\lambda/\mu| \) denotes the number of cells of the skew shape \( \lambda/\mu \) (so \( |\lambda/\mu| = |\lambda| - |\mu| \)).

It is customary to define \( s_{\lambda/\mu} \) to be 0 if we don’t have \( \mu \subseteq \lambda \). This can also be seen by a literal reading of the definition \( s_{\lambda/\mu} := \sum_T x^{\text{cont}(T)} \), as long as we understand that there are no column-strict tableaux of shape \( \lambda/\mu \) when \( \lambda/\mu \) is not defined.

Clearly, every partition \( \lambda \) satisfies \( s_\lambda = s_\lambda \circ \phi \).

Exercise 2.22. \( (a) \) State and prove an analogue of Proposition 2.13 for skew Schur functions.

\( (b) \) Let \( \lambda, \mu, \lambda’ \) and \( \mu’ \) be partitions such that \( \mu \subseteq \lambda \) and \( \mu’ \subseteq \lambda’ \). Assume that the skew Ferrers diagram \( \lambda’/\mu’ \) can be obtained from the skew Ferrers diagram \( \lambda/\mu \) by a \( 180^\circ \) rotation.\(^52\) Prove that \( s_{\lambda/\mu} = s_{\lambda’/\mu’} \).

Exercise 2.23. Let \( \lambda \) and \( \mu \) be two partitions, and let \( k \in \mathbb{N} \) be such that\(^53\) \( \mu_k \geq \lambda_{k+1} \). Let \( F \) be the skew Ferrers diagram \( \lambda/\mu \). Let \( F_{\text{rows} \leq k} \) denote the subset of \( F \) consisting of all \( (i, j) \) satisfying \( i \leq k \). Let \( F_{\text{rows} > k} \) denote the subset of \( F \) consisting of all \( (i, j) \in F \) satisfying \( i > k \). Let \( \alpha \) and \( \beta \) be two partitions such that \( \beta \subseteq \alpha \) and such that the skew Ferrers diagram \( \alpha/\beta \) can be obtained from \( F_{\text{rows} \leq k} \) by parallel

\(^50\)Specifically, an element \( \sigma \) of the group takes \( \varphi : \{1, 2, \ldots, \ell\} \to \{1, 2, 3, \ldots\} \) to \( \sigma \circ \varphi \).

\(^51\)In other words, the skew Ferrers diagram \( \lambda/\mu \) is the set of all \( (i, j) \in \{1, 2, 3, \ldots\}^2 \) satisfying \( \mu_i < j \leq \lambda_i \).

While the Ferrers diagram for a single partition \( \lambda \) uniquely determines \( \lambda \), the skew Ferrers diagram \( \lambda/\mu \) does not uniquely determine \( \mu \) and \( \lambda \). (For instance, it is empty whenever \( \lambda = \mu \).) When one wants to keep \( \mu \) and \( \lambda \) in memory, one speaks of the skew shape \( \lambda/\mu \); this simply means the pair \( (\mu, \lambda) \). Every notion defined for skew Ferrers diagrams also makes sense for skew shapes, because to any skew shape \( \lambda/\mu \) we can assign the skew Ferrers diagram \( \lambda/\mu \) (even if not injectively). For instance, the cells of the skew shape \( \lambda/\mu \) are the cells of the skew Ferrers diagram \( \lambda/\mu \).

One can characterize the skew Ferrers diagrams as follows: A finite subset \( S \) of \( \{1, 2, 3, \ldots\}^2 \) is a skew Ferrers diagram (i.e., there exist two partitions \( \lambda \) and \( \mu \) such that \( \mu \subseteq \lambda \) and such that \( S \) is the skew Ferrers diagram \( \lambda/\mu \) if and only if for every \( (i, j) \in S \), every \( (i’, j’) \in \{1, 2, 3, \ldots\}^2 \) and every \( (i'', j'') \in S \) satisfying \( i'' < i' \leq i \) and \( j' \leq j'' \leq j \), we have \( (i', j') \in S \).

\(^52\)For example, this happens when \( \lambda = (3, 2), \mu = (1), \lambda’ = (5, 4) \) and \( \mu’ = (3, 1) \).

\(^53\)As usual, we write \( s_k \) for the \( k \)-th entry of a partition \( \nu \).
translation. Let \( \gamma \) and \( \delta \) be two partitions such that \( \delta \subseteq \gamma \) and such that the skew Ferrers diagram \( \gamma/\delta \) can be obtained from \( F_{\text{row} \circ \delta} \) by parallel translation.\(^{54}\) Prove that \( s_{\lambda/\mu} = s_{\alpha/\beta} s_{\gamma/\delta} \).

**Proposition 2.24.** The comultiplication \( \Lambda \xrightarrow{\Delta} \Lambda \otimes \Lambda \) has the following effect on the symmetric functions discussed so far:\(^{55}\)

(i) \( \Delta p_n = 1 \otimes p_n + p_n \otimes 1 \) for every \( n \geq 1 \), that is, the power sums \( p_n \) are primitive.

(ii) \( \Delta e_n = \sum_{i+j=n} e_i \otimes e_j \) for every \( n \in \mathbb{N} \).

(iii) \( \Delta h_n = \sum_{i+j=n} h_i \otimes h_j \) for every \( n \in \mathbb{N} \).

(iv) \( \Delta s_{\lambda} = \sum_{\mu \subseteq \lambda} s_{\mu} \otimes s_{\lambda/\mu} \) for any partition \( \lambda \).

(v) \( \Delta s_{\lambda/\mu} = \sum_{\nu \subseteq \mu \subseteq \lambda} s_{\mu/\nu} \otimes s_{\lambda/\mu} \) for any partitions \( \lambda \) and \( \nu \).

**Proof.** Recall that \( \Delta \) sends \( f(x) \mapsto f(x, y) \), and one can easily check that

\[
\begin{align*}
    (i) & \quad p_n(x, y) = \sum_i x_i^n + \sum_i y_i^n = p_n(x) \cdot 1 + 1 \cdot p_n(y) \\
    (ii) & \quad e_n(x, y) = \sum_{i+j=n} e_i(x) e_j(y) \\
    (iii) & \quad h_n(x, y) = \sum_{i+j=n} h_i(x) h_j(y)
\end{align*}
\]

For assertion (iv), note that by (2.9), one has

\[
\begin{equation}
    s_{\lambda}(x, y) = \sum_T (x, y)^{\text{cont}(T)},
\end{equation}
\]

where the sum is over column-strict tableaux \( T \) of shape \( \lambda \) having entries in the linearly ordered alphabet

\[
    x_1 < x_2 < \cdots < y_1 < y_2 < \cdots.
\]

\(^{56}\) For example,

\[
    T = \begin{array}{cccc}
        x_1 & x_1 & x_2 & y_5 \\
        y_2 & y_1 & y_1 & y_1 \\
        y_4 & y_2 & y_4 & y_5
    \end{array}
\]

is such a tableau of shape \( \lambda = (5, 3, 3, 2) \). Note that the restriction of \( T \) to the alphabet \( x \) gives a column-strict tableau \( T_x \) of some shape \( \mu \subseteq \lambda \), and the restriction of \( T \) to the alphabet \( y \) gives a column-strict tableau \( T_y \) of shape \( \lambda/\mu \) (e.g. for \( T \) in the example above, the tableau \( T_y \) appeared in Example 2.20). Consequently, one has

\[
\begin{equation}
    s_{\lambda}(x, y) = \sum_T x^{\text{cont}(T_x)} \cdot y^{\text{cont}(T_y)}
\end{equation}
\]

\[
\begin{equation}
    = \sum_{\mu \subseteq \lambda} \left( \sum_{T_x} x^{\text{cont}(T_x)} \right) \left( \sum_{T_y} y^{\text{cont}(T_y)} \right) = \sum_{\mu \subseteq \lambda} s_{\mu}(x) s_{\lambda/\mu}(y).
\end{equation}
\]

Assertion (v) is obvious in the case when we don’t have \( \nu \subseteq \lambda \) (in fact, in this case, both \( s_{\lambda/\nu} \) and \( \sum_{\nu \subseteq \mu \subseteq \lambda} s_{\mu/\nu} \otimes s_{\lambda/\mu} \) are clearly zero). In the remaining case, the proof of assertion (v) is similar to that of

\(^{56}\)Here, \( (x, y)^{\text{cont}(T)} \) means the monomial \( \prod_{t \in T} a^{[T^{-1}(a)]} \), where \( \mathbb{A} \) denotes the totally ordered alphabet \( x_1 < x_2 < \cdots < y_1 < y_2 < \cdots \). In other words, \( (x, y)^{\text{cont}(T)} \) is the product of all entries of the tableau \( T \) (which is a monomial, since the entries of \( T \) are not numbers but variables).

The following rather formal argument should allay any doubts as to why (2.11) holds: Let \( L \) denote the totally ordered set which is given by the set \( \{1, 2, 3, \ldots\} \) of positive integers, equipped with the total order \( 1 < \epsilon < 3 < \epsilon < 5 < \epsilon < 7 < \epsilon < \cdots \). Then, (2.9) yields \( s_{\lambda} = \sum_T x^{\text{cont}(T)} \) as \( T \) runs through all \( L \)-column-strict tableaux of shape \( \lambda \). Substituting the variables \( x_1, y_1, x_2, y_2, x_3, y_3, \ldots \) for \( x_1, x_2, x_3, x_4, x_5, y_6, \ldots \) (that is, substituting \( x_i \) for \( x_{2i-1} \) and \( y_i \) for \( x_{2i} \)) in this equality, we obtain (2.11).
(iv). (Of course, the tableaux $T$ and $T_x$ now have skew shapes $\lambda/\nu$ and $\mu/\nu$, and instead of (2.9), we need to use the answer to Exercise 2.22(a).)

Notice that parts (ii) and (iii) of Proposition 2.24 are particular cases of part (iv), since $h_n = s(n)$ and $\epsilon_n = s(1^n)$.

**Exercise 2.25.** (a) Show that the Hopf algebra $\Lambda$ is cocommutative.

(b) Show that $\Delta s_{\lambda/\nu} = \sum_{\mu \in \text{Par}} s_{\lambda/\mu} \otimes s_{\mu/\nu}$ for any partitions $\lambda$ and $\nu$.

**Exercise 2.26.** Let $\epsilon_n$. Let $\Lambda$ be written uniquely as a polynomial in the elementary symmetric polynomials $x_i$ of Proposition 2.17. We leave the proof of (f) to the reader.

Indeed, the claims (a) and (b) are obvious, while the claims (c), (d) and (e) are proven similarly to our proof of Proposition 2.17. We leave the proof of (f) to the reader.

Claim (c) can be rewritten as follows: The elementary symmetric polynomials $e_i(x_1, x_2, \ldots, x_N)$, for $i \in \{1, 2, \ldots, N\}$, form an algebraically independent generating set of $\Lambda(x_1, x_2, \ldots, x_N)$. This is precisely the well-known theorem (due to Gauss)\(^{57}\) that every symmetric polynomial in $N$ variables $x_1, x_2, \ldots, x_N$ can be written uniquely as a polynomial in the $N$ elementary symmetric polynomials.

2.4. The antipode, the involution $\omega$, and algebra generators. Since $\Lambda$ is a graded connected $k$-bialgebra, it will have an antipode $S : \Lambda \to \Lambda$ making it a Hopf algebra by Proposition 1.36. However, several issues will be resolved by identifying $S$ more explicitly now.

**Proposition 2.28.** Each of $\{e_n\}_{n=1, 2, \ldots}$, $\{h_n\}_{n=1, 2, \ldots}$ are algebraically independent, and generate $\Lambda$ as a polynomial algebra for any commutative ring $k$. The same holds for $\{p_n\}_{n=1, 2, \ldots}$ when $Q$ is a subring of $k$.

Furthermore, the antipode $S$ acts as follows:

(i) $S(p_n) = -p_n$
(ii) $S(e_n) = (-1)^n h_n$
(iii) $S(h_n) = (-1)^n e_n$

Proof. The assertions that $\{e_n\}, \{p_n\}$ are algebraically independent and generate $\Lambda$ are equivalent to Proposition 2.17 asserting $\{e_\lambda\}, \{p_\lambda\}$ give bases for $\Lambda$. The assertion $S(p_n) = -p_n$ follows from Proposition 1.37 since $p_n$ is primitive by Proposition 2.24(i).

\(^{57}\)See, e.g., [32, Symmetric Polynomials, Theorem 5 and Remark 17].
For the remaining assertions, start with the easy generating function identities

\[ H(t) := \prod_{i=1}^{\infty} (1 - x_it)^{-1} = 1 + h_1(x)t + h_2(x)t^2 + \cdots = \sum_{n \geq 0} h_n(x)t^n \]
\[ E(t) := \prod_{i=1}^{\infty} (1 + x_it) = 1 + e_1(x)t + e_2(x)t^2 + \cdots = \sum_{n \geq 0} e_n(x)t^n \]

which shows that

\[ 1 = E(-t)H(t) = \left( \sum_{n \geq 0} e_n(x)(-t)^n \right) \left( \sum_{n \geq 0} h_n(x)t^n \right) \]

and hence, equating coefficients of powers of \( t \), that for \( n = 0, 1, 2, \ldots \) one has

\[ \sum_{i+j=n} (-1)^i e_i h_j = \delta_{0,n}. \]

This lets one recursively express the \( e_n \) in terms of \( h_n \) and vice-versa:

\[ e_0 := 1 := h_0 \]
\[ e_n = e_{n-1}h_1 - e_{n-2}h_2 + e_{n-3}h_3 - \cdots \]
\[ h_n = h_{n-1}e_1 - h_{n-2}e_2 + h_{n-3}e_3 - \cdots \]

for \( n = 1, 2, \ldots \). Thus if one uses the algebraic independence of the generators \( \{e_n\} \) for \( \Lambda \) to define an algebra endomorphism as follows

\[ \Lambda \xrightarrow{\omega} \Lambda, \quad e_n \mapsto h_n, \]

then the identical form of the two recursions in (2.18) shows that \( \omega \) also sends \( h_n \mapsto e_n \). Therefore \( \omega \) is an involutive automorphism of \( \Lambda \), and the \( \{h_n\} \) are another algebraically independent generating set for \( \Lambda \).

For the assertion about the antipode \( S \) applied to \( e_n \) or \( h_n \), note that the coproduct formulas for \( e_n, h_n \) in Proposition 2.24(ii),(iii) show that the defining relations for their antipodes (1.13) will in this case be

\[ \sum_{i+j=n} S(e_i) e_j = \delta_{0,n} = \sum_{i+j=n} e_i S(e_j) \]
\[ \sum_{i+j=n} S(h_i) h_j = \delta_{0,n} = \sum_{i+j=n} h_i S(h_j) \]

because \( u\epsilon(e_n) = u\epsilon(h_n) = \delta_{0,n} \). Comparing these to (2.17), one concludes via induction on \( n \) that \( S(e_n) = (-1)^n h_n \) and \( S(h_n) = (-1)^n e_n \). \( \square \)

The \( k \)-algebra endomorphism \( \omega \) of \( \Lambda \) defined in the proof of Proposition 2.28 is known as the fundamental involution on \( \Lambda \). We record for future use two results that were shown in the above proof: that \( \omega \) is an involution, and that it sends \( e_n \) to \( h_n \) and vice versa for every positive integer \( n \).

Proposition 2.28 shows that the antipode \( S \) on \( \Lambda \) is, up to sign, the same as the fundamental involution \( \omega \): one has

\[ S(f) = (-1)^n \omega(f) \quad \text{for} \quad f \in \Lambda_n \]

since this formula holds for all elements of the generating set \( \{e_n\} \) (or \( \{h_n\} \)).

Remark 2.29. Up to now we have not yet derived how the involution \( \omega \) and the antipode \( S \) act on (skew) Schur functions, which is quite beautiful: If \( \lambda \) and \( \mu \) are partitions satisfying \( \mu \subseteq \lambda \), then

\[ \omega(s_{\lambda/\mu}) = s_{\lambda'/\mu'} \]
\[ S(s_{\lambda/\mu}) = (-1)^{\lambda/\mu} s_{\lambda'/\mu'} \]

where recall that \( \lambda' \) is the transpose or conjugate partition to \( \lambda \), and \( |\lambda/\mu| \) is the number of squares in the skew diagram \( \lambda/\mu \), that is, \( |\lambda/\mu| = n - k \) if \( \lambda, \mu \) lie in \( \text{Par}_n, \text{Par}_k \) respectively.

We will deduce this later in three ways (once as an exercise using the Pieri rules in Exercise 2.69, once again using skewing operators in Exercise 2.76, and for the third time from the action of the antipode in
Theorem 2.31. Skew Schur functions are the following polynomials in $\{h_n\}, \{e_n\}$:

$$s_{\lambda\mu} = \det(h_{\lambda_i - \mu_j + i - j})_{i,j=1,2,\ldots,l}$$

(2.23) $$s_{\lambda\mu}^t = \det(e_{\lambda_i - \mu_j + i - j})_{i,j=1,2,\ldots,l}$$

for any two partitions $\lambda$ and $\mu$ and any $\ell \in \mathbb{N}$ satisfying $\ell(\lambda) \leq \ell$ and $\ell(\mu) \leq \ell$.

Since we appear not to need these formulas in the sequel, we will not prove them right away. However, a proof is sketched in the solution to Exercise 2.71, and various proofs are well-explained in [92, (39) and (41)], [106, §I.5], [136, Thm. 7.1], [138, §4.5], [154, §7.16], [166, Thms. 3.5 and 3.5*]; also, a simultaneous generalization of both formulas is shown in [62, Theorem 11], and two others in [133, 1.9] and [63, Thm. 3.1]. An elegant treatment of Schur polynomials taking the Jacobi-Trudi formula (2.22) as the definition of $s_\lambda$ is given by Tamvakis [162].

2.5. Cauchy product, Hall inner product, self-duality. The Schur functions, although a bit unmotivated right now, have special properties with regard to the Hopf structure. One property is intimately connected with the following Cauchy identity.

Theorem 2.31. In the power series ring $k[[x, y]] := k[[x_1, x_2, \ldots, y_1, y_2, \ldots]]$, one has the following expansion:

$$\prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} = \sum_{\lambda \in \text{Par}} s_\lambda(x)s_\lambda(y).$$

(2.24)

Remark 2.32. The left hand side of (2.24) is known as the Cauchy product, or Cauchy kernel.

An equivalent version of the equality (2.24) is obtained by replacing each $x_i$ by $x_i t$, and writing the resulting identity in the power series ring $R(x, y)[[t]]$:

$$\prod_{i,j=1}^{\infty} (1 - t x_i y_j)^{-1} = \sum_{\lambda \in \text{Par}} t^{\ell(\lambda)} s_\lambda(x)s_\lambda(y).$$

(2.25) (Recall that $|\lambda| = \lambda_1 + \lambda_2 + \cdots + \lambda_\ell$ for any partition $\lambda = (\lambda_1, \lambda_2, \ldots, \lambda_\ell)$.)

Proof of Theorem 2.31. We follow the standard combinatorial proof (see [138, §4.8],[154, §7.11,7.12]), which rewrites the left and right sides of (2.25), and then compares them with the Robinson-Schensted-Knuth (RSK) bijection. On the left side, expanding out each geometric series

$$(1 - t x_i y_j)^{-1} = 1 + t x_i y_j + (t x_i y_j)^2 + (t x_i y_j)^3 + \cdots$$

and thinking of $(x_i y_j)^m$ as $m$ occurrences of a biletter [60] $i \choose j$, we see that the left hand side can be rewritten as the sum of $t^\ell(x_{i_1}y_{j_1}) (x_{i_2}y_{j_2}) \cdots (x_{i_\ell}y_{j_\ell})$ over all multisets $\{i_1 \choose j_1}, \ldots, i_\ell \choose j_\ell \}$ of biletters. Order the biletters in such a multiset in the lexicographic order $\preceq_{\text{lex}}$, which is the total order on the set of all biletters defined by

$$i_1 \choose j_1 \preceq_{\text{lex}} i_2 \choose j_2 \iff \text{(we have } i_1 \leq i_2, \text{ and if } i_1 = i_2, \text{ then } j_1 \leq j_2).$$

Defining a biword to be an array $i \choose j = i_1 \ldots i_\ell \choose j_1 \ldots j_\ell$ in which the biletters are ordered $i_1 \choose j_1 \preceq_{\text{lex}} \cdots \preceq_{\text{lex}} i_\ell \choose j_\ell$, then the left side of (2.25) is the sum $\sum t^\ell x^{\text{cont}(i)} y^{\text{cont}(j)}$ over all biwords $i \choose j$, where $\ell$ stands for the number

58 The second of the following identities is also known as the von N"agelsbach-Kostka identity.

59 The RSK bijection has been introduced by Knuth [81], where what we call “biletters” is referred to as “two-line arrays.”

The most important ingredient of this algorithm – the RS-insertion operation – however goes back to Schensted. The special case of the RSK algorithm where the biword has to be a permutation (written in two-line notation) and the two tableaux have to be standard (i.e., each of them has content (1n), where n is the size of their shape) is the famous Robinson-Schensted correspondence [96]. More about these algorithms can be found in [138, Chapter 3], [154, §7.11-7.12], [53, Chapters 1 and A], [22, §3, §6] and various other places.

60 A biletter here simply means a pair of letters, written as a column vector. A letter means a positive integer.
of biletters in the biword. On the right side, expanding out the Schur functions as sums of tableaux gives \( \sum_{(P,Q)} t^i x^{\mathrm{cont}(Q)} y^{\mathrm{cont}(P)} \) in which the sum is over all ordered pairs \((P,Q)\) of column-strict tableaux having the same shape\(^{61}\), with \(\ell\) cells. (We shall refer to such pairs as tableau pairs from now on.)

The Robinson-Schensted-Knuth algorithm gives us a bijection between the biwords \((\underline{1})\) and the tableau pairs \((P,Q)\), which has the property that

\[
\begin{align*}
\mathrm{cont}(i) &= \mathrm{cont}(Q), \\
\mathrm{cont}(j) &= \mathrm{cont}(P)
\end{align*}
\]

(and that the length \(\ell\) of the biword \((\underline{1})\) equals the size \(|\lambda|\) of the common shape of \(P\) and \(Q\); but this follows automatically from \(\mathrm{cont}(i) = \mathrm{cont}(Q)\)). Clearly, once such a bijection is constructed, the equality (2.25) will follow.

Before we define this algorithm, we introduce a simpler operation known as RS-insertion (short for Robinson-Schensted insertion). RS-insertion takes as input a column-strict tableau \(P\) and a letter \(j\), and returns a new column-strict tableau \(P'\) along with a corner cell\(^{62}\) \(c\) of \(P'\), which is constructed as follows: Start out by setting \(P' = P\). The letter \(j\) tries to insert itself into the first row of \(P'\) by either bumping out the leftmost letter in the first row strictly larger than \(j\), or else placing itself at the right end of the row if no such larger letter exists. If a letter was bumped from the first row, this letter follows the same rules to insert itself into the second row, and so on\(^{63}\). This series of bumps must eventually come to an end\(^{64}\). At the end of the bumping, the tableau \(P'\) created has an extra corner cell not present in \(P\). If we call this corner cell \(c\), then \(P'\) (in its final form) and \(c\) are what the RS-insertion operation returns. One says that \(P'\) is the result of inserting\(^{65}\) \(c\) into the tableau \(P\). It is straightforward to see that this resulting filling \(P'\) is a column-strict tableau\(^{66}\).

**Example 2.33.** To give an example of this operation, let us insert the letter \(j = 3\) into the column-strict tableau

\[
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 6 \\
3 & 4 & 7 & \\
5 & \\
\end{array}
\]

(where we are showing all intermediate states of \(P'\); the underlined letter is always the one that is going to be bumped out at the next step):  

\[
\begin{array}{cccc}
1 & 1 & 3 & 4 \\
2 & 2 & 4 & 6 \\
3 & 4 & 7 & \\
5 & \\
\end{array}
\overset{\text{insert 3;}}{\longrightarrow}
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 6 \\
3 & 4 & 7 & \\
5 & \\
\end{array}
\overset{\text{bump out 4;}}{\longrightarrow}
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 4 \\
3 & 4 & 7 & \\
5 & \\
\end{array}
\overset{\text{insert 4;}}{\longrightarrow}
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 4 \\
3 & 4 & 7 & \\
5 & \\
\end{array}
\overset{\text{bump out 6;}}{\longrightarrow}
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 4 \\
3 & 4 & 6 & \\
5 & \\
\end{array}
\overset{\text{insert 7;}}{\longrightarrow}
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 4 \\
3 & 4 & 6 & \\
5 & \\
\end{array}
\overset{\text{done}}{\longrightarrow}
\begin{array}{cccc}
1 & 1 & 3 & 3 \\
2 & 2 & 4 & 4 \\
3 & 4 & 6 & \\
5 & \\
\end{array}
\]

The last tableau in this sequence is the column-strict tableau that is returned. The corner cell that is returned is the second cell of the fourth row (the one containing 7).

\(^{61}\)And this shape should be the Ferrers diagram of a partition (not just a skew partition).

\(^{62}\)A **corner cell** of a tableau or a Ferrers diagram is defined to be a cell \(c\) which belongs to the tableau (resp. diagram) but whose immediate neighbors to the east and to the south don’t.

\(^{63}\)Here, rows are allowed to be empty – so it is possible that a letter is bumped from the last nonempty row of \(P'\) and settles in the next, initially empty, row.

\(^{64}\)Since we can only bump out entries from nonempty rows

\(^{65}\)This terminology is reminiscent of insertion into binary search trees, a basic operation in theoretical computer science. This is more than superficial similarity; there are, in fact, various analogies between Ferrers diagrams (and their fillings) and unlabelled plane binary trees (resp. their labellings), and one of them is the analogy between RS-insertion and binary search tree insertion. See [71, §4.1].

\(^{66}\)Indeed, the reader can check that \(P'\) remains a column-strict tableau throughout the algorithm that defines RS-insertion. (The only part of this that isn’t obvious is showing that when a letter \(t\) bumped out of some row \(k\) is inserted into row \(k + 1\), the property that the letters increase strictly down columns is preserved. Argue that the bumping-out of \(t\) from row \(k\) was caused by the insertion of another letter \(u < t\), and that the cell of row \(k + 1\) into which \(t\) is then being inserted is in the same column as this \(u\), or in a column further left than it.)
RS-insertion will be used as a step in the RSK algorithm; the construction will rely on a simple fact known as the row bumping lemma. Let us first define the notion of a bumping path (or bumping route): If \( P \) is a column-strict tableau, and \( j \) is a letter, then some letters are inserted into some cells when RS-insertion is applied to \( P \) and \( j \). The sequence of these cells (in the order in which they see letters inserted into them) is called the bumping path for \( P \) and \( j \). This bumping path always ends with the corner cell \( c \) which is returned by RS-insertion. As an example, when \( j = 1 \) is inserted into the tableau \( P \) shown below, the result \( P' \) is shown with all entries on the bumping path underlined:

\[
P = \begin{array}{cccc}
1 & 1 & 2 & 3 \\
2 & 2 & 4 & 4 \\
3 & 4 & 5 & \text{insert} \\
4 & 6 & 6 & j=1
\end{array}
\quad \quad \quad \quad \quad \quad
P' = \begin{array}{cccc}
1 & 1 & 1 & 2 & 3 \\
2 & 2 & 2 & 4 \\
4 & 5 & 6 & \text{insert}
\end{array}
\]

A first simple observation about bumping paths is that bumping paths trend weakly left — that is, if the bumping path of \( P \) and \( j \) is \((c_1, c_2, \ldots, c_k)\), then, for each \( 1 \leq i < k \), the cell \( c_{i+1} \) lies in the same column as \( c_i \) or in a column further left.\(^67\) A subtler property of bumping paths is the following row bumping lemma (\([53, \text{p. 9}]\)):

**Row bumping lemma:** Let \( P \) be a column-strict tableau, and let \( j \) and \( j' \) be two letters. Applying RS-insertion to the tableau \( P \) and the letter \( j \) yields a new column-strict tableau \( P' \) and a corner cell \( c \). Applying RS-insertion to the tableau \( P' \) and the letter \( j' \) yields a new column-strict tableau \( P'' \) and a corner cell \( c' \).

(a) Assume that \( j \leq j' \). Then, the bumping path for \( P' \) and \( j' \) stays strictly to the right, within each row, of the bumping path for \( P \) and \( j \). The cell \( c' \) (in which the bumping path for \( P' \) and \( j' \) ends) is in the same row as the cell \( c \) (in which the bumping path for \( P \) and \( j \) ends) or in a row further up; it is also in a column further right than \( c \).

(b) Assume instead that \( j > j' \). Then, the bumping path for \( P' \) and \( j' \) stays weakly to the left, within each row, of the bumping path for \( P \) and \( j \). The cell \( c' \) (in which the bumping path for \( P' \) and \( j' \) ends) is in a row further down than the cell \( c \) (in which the bumping path for \( P \) and \( j \) ends); it is also in the same column as \( c \) or in a column further left.

This lemma can be easily proven by induction over the row.\(^68\)

We can now define the actual RSK algorithm. Let \( (i, j) \) be a biword. Starting with the pair \( (P_0, Q_0) = (\emptyset, \emptyset) \) and \( m = 0 \), the algorithm applies the following steps (see Example 2.34 below):

- If \( i_{m+1} \) does not exist (that is, \( m \) is the length of \( i \)), stop.
- Apply RS-insertion to the column-strict tableau \( P_m \) and the letter \( j_{m+1} \) (the bottom letter of \( \binom{j}{j_{m+1}} \)). Let \( P_{m+1} \) be the resulting column-strict tableau, and let \( c_{m+1} \) be the resulting corner cell.
- Create \( Q_{m+1} \) from \( Q_m \) by adding the top letter \( i_{m+1} \) of \( \binom{i}{j_{m+1}} \) to \( Q_m \) in the cell \( c_{m+1} \) (which, as we recall, is the extra corner cell of \( P_{m+1} \) not present in \( P_m \)).
- Set \( m \) to \( m + 1 \).

After all of the biletters have been thus processed, the result of the RSK algorithm is \( (P_T, Q_T) =: (P, Q) \).

---

\(^67\) This follows easily from the preservation of column-strictness during RS-insertion.

\(^68\) We leave the details to the reader, only giving the main idea for (a) (the proof of (b) is similar). To prove the first claim of (a), it is enough to show that for every \( i \), if any letter is inserted into row \( i \) during RS-insertion for \( P' \) and \( j' \), then some letter is also inserted into row \( i \) during RS-insertion for \( P \) and \( j \), and the former insertion happens in a cell strictly to the right of the cell where the latter insertion happens. This follows by induction over \( i \). In the induction step, we need to show that if, for a positive integer \( i \), we try to consecutively insert two letters \( k \) and \( k' \), in this order, into the \( i \)-th row of a column-strict tableau, possibly bumping out existing letters in the process, and if we have \( k < k' \), then the cell into which \( k \) is inserted is strictly to the left of the cell into which \( k' \) is inserted, and the letter bumped out by the insertion of \( k \) is \( \leq \) to the letter bumped out by the insertion of \( k' \) (or else the insertion of \( k' \) bumps out no letter at all — but it cannot happen that \( k' \) bumps out a letter but \( k \) does not). This statement is completely straightforward to check (by only studying the \( i \)-th row). This way, the first claim of (a) is proven, and this entails that the cell \( c' \) (being the last cell of the bumping path for \( P' \) and \( j' \)) is in the same row as the cell \( c \) or in a row further up. It only remains to show that \( c' \) is in a column further right than \( c \). This follows by noticing that, if \( k \) is the row in which the cell \( c' \) lies, then \( c' \) is in a column further right than the entry of the bumping path for \( P \) and \( j \) in row \( k \) (by the first claim of (a)), and this latter entry is further right than or in the same column as the ultimate entry \( c \) of this bumping path (since bumping paths trend weakly left).
Example 2.34. The term in the expansion of the left side of (2.24) corresponding to 

\[(x_1 y_2)^i (x_1 y_3)^i (x_2 y_1)^i (x_3 y_1)^i (x_4 y_3)^2 (x_5 y_2)^1\]

is the biword \(\binom{i}{j} = (1124445)_{(2411332)}\), whose RSK algorithm goes as follows:

\[
\begin{aligned}
P_0 &= \emptyset & Q_0 &= \emptyset \\
P_1 &= 2 & Q_1 &= 1 \\
P_2 &= 24 & Q_2 &= 11 \\
P_3 &= 14 & Q_3 &= 11 \\
P_4 &= 114 & Q_4 &= 114 \\
P_5 &= 113 & Q_5 &= 114 \\
P_6 &= 1133 & Q_6 &= 1144 \\
P := P_7 &= 1123 & Q := Q_7 &= 1144 \\
\end{aligned}
\]

The bumping rule obviously maintains the property that \(P_m\) is a column-strict tableau of some Ferrers shape throughout. It should be clear that \((P_m, Q_m)\) have the same shape at each stage. Also, the construction of \(Q_m\) shows that it is at least weakly increasing in rows and weakly increasing in columns throughout. What is perhaps least clear is that \(Q_m\) remains strictly increasing down columns. That is, when one has a string of equal letters on top \(i_m = i_{m+1} = \cdots = i_{m+r}\), so that on bottom one bumps in \(j_m \leq j_{m+1} \leq \cdots \leq j_{m+r}\), one needs to know that the new cells form a horizontal strip, that is, no two of them lie in the same column. This follows from (the last claim of) part (a) of the row bumping lemma. Hence, the result \((P, Q)\) of the RSK algorithm is a tableau pair.

To see that the RSK map is a bijection, we show how to recover \(\binom{i}{j}\) from \((P, Q)\). This is done by reverse bumping from \((P_{m+1}, Q_{m+1})\) to recover both the biletter \(\binom{i_{m+1}}{j_{m+1}}\) and the tableaux \((P_m, Q_m)\), as follows. Firstly, \(i_{m+1}\) is the maximum entry of \(Q_{m+1}\), and \(Q_m\) is obtained by removing the rightmost occurrence of this letter \(i_{m+1}\) from \(Q_{m+1}\). To produce \(P_m\) and \(j_{m+1}\), find the position of the rightmost occurrence of \(i_{m+1}\) in \(Q_{m+1}\), and start reverse bumping in \(P_{m+1}\) from the entry in this same position, where reverse bumping an entry means inserting it into one row higher by having it bump out the rightmost entry which is strictly smaller. The entry bumped out from the first row is \(j_{m+1}\), and the resulting tableau is \(P_m\).

69 Actually, each of these new cells (except for the first one) is in a column further right than the previous one. We will use this stronger fact further below.

70 It necessarily has to be the rightmost occurrence, since (according to the previous footnote) the cell into which \(i_{m+1}\) was filled at the step from \(Q_m\) to \(Q_{m+1}\) lies further right than any existing cell of \(Q_m\) containing the letter \(i_{m+1}\).

71 Let us give a few more details on this “reverse bumping” procedure. Reverse bumping (also known as RS-deletion or reverse RS-insertion) is an operation which takes a column-strict tableau \(P^e\) and a corner cell \(c\) of \(P^e\), and constructs a column-strict tableau \(P\) and a letter \(j\) such that RS-insertion for \(P\) and \(j\) yields \(P^e\) and \(c\). It starts by setting \(P = P^e\), and removing the entry in the cell \(c\) from \(P\). This removed entry is then denoted by \(k\), and is inserted into the row of \(P\) above \(c\), bumping out the rightmost entry which is smaller than \(k\). The letter which is bumped out – say, \(\ell\), in turn, is inserted into the row above it, bumping out the rightmost entry which is smaller than \(\ell\). This procedure continues in the same way until an entry is bumped out of the first row (which will eventually happen). The reverse bumping operation returns the resulting tableau \(P\) and the entry which is bumped out of the first row.

It is straightforward to check that the reverse bumping operation is well-defined (i.e., \(P\) does stay a column-strict tableau throughout the procedure) and is the inverse of the RS-insertion operation. (In fact, these two operations undo each other step by step.)
Finally, to see that the RSK map is surjective, one needs to show that the reverse bumping procedure can be applied to any pair \((P, Q)\) of column-strict tableaux of the same shape, and will result in a (lexicographically ordered) biword \(\binom{1}{i}\). We leave this verification to the reader.\(^\circ\)

This is by far not the only known proof of Theorem 2.31. Two further proofs will be sketched in Exercise 2.68 and Exercise 2.66.

**Corollary 2.35.** In the Schur function basis \(\{s_\lambda\}\) for \(\Lambda\), the structure constants for multiplication and comultiplication are the same, that is, if one defines \(c_\lambda^{\mu,\nu}, \hat{c}_\lambda^{\mu,\nu}\) via the unique expansions

\[
\begin{align*}
s_\mu s_\nu &= \sum_\lambda c_\lambda^{\mu,\nu} s_\lambda \\
\Delta(s_\lambda) &= \sum_{\mu,\nu} \hat{c}_\lambda^{\mu,\nu} s_\mu \otimes s_\nu
\end{align*}
\]

(2.26)

then \(c_\lambda^{\mu,\nu} = \hat{c}_\lambda^{\mu,\nu}\).

**Proof.** Work in the ring \(k[[x, y, z]]\), where \(y = (y_1, y_2, y_3, \ldots)\) and \(z = (z_1, z_2, z_3, \ldots)\) are two new sets of variables. The identity (2.24) lets one interpret both \(c_\lambda^{\mu,\nu}, \hat{c}_\lambda^{\mu,\nu}\) as the coefficient\(^72\) of \(s_\mu(x)s_\nu(y)s_\lambda(z)\) in the product

\[
\prod_{i,j=1}^\infty (1-x_i z_j)^{-1} \prod_{i,j=1}^\infty (1-y_i z_j)^{-1} \left(\sum_\mu s_\mu(x)s_\mu(z)\right) \left(\sum_\nu s_\nu(y)s_\nu(z)\right) = \sum_{\mu,\nu} s_\mu(x)s_\nu(y) s_\lambda(z)
\]

since, regarding \(x_1, x_2, \ldots, y_1, y_2, \ldots\) as lying in a single variable set \((x, y)\), separate from the variables \(z\), the Cauchy identity (2.24) expands the same product as

\[
\prod_{i,j=1}^\infty (1-x_i z_j)^{-1} \prod_{i,j=1}^\infty (1-y_i z_j)^{-1} = \sum_\lambda s_\lambda(x, y)s_\lambda(z)
\]

\[
= \sum_\lambda \left(\sum_{\mu,\nu} c_\lambda^{\mu,\nu} s_\mu(x)s_\nu(y)\right) s_\lambda(z).
\]

\(\square\)

**Definition 2.36.** The coefficients \(c_\lambda^{\mu,\nu} = \hat{c}_\lambda^{\mu,\nu}\) appearing in the expansions (2.26) are called Littlewood-Richardson coefficients.

**Remark 2.37.** We will interpret \(c_\lambda^{\mu,\nu}\) combinatorially in Section 2.6. By now, however, we can already prove some properties of these coefficients:

We have

\[
c_\lambda^{\mu,\nu} = c_\nu^{\lambda,\mu}
\]

(2.27)

(by comparing coefficients in \(\sum_\lambda c_\lambda^{\mu,\nu} s_\lambda = s_\mu s_\nu = s_\nu s_\mu = \sum_\lambda \hat{c}_\lambda^{\mu,\nu} s_\lambda\)). Also, noting on one hand the expansion

\[
s_\lambda(x, y) = \Delta(s_\lambda) = \sum_{\mu,\nu} c_\lambda^{\mu,\nu} s_\mu(x)s_\nu(y)
\]

\(\square\)

\(\text{It is easy to see that repeatedly applying reverse bumping to } (P, Q) \text{ will result in a sequence } \binom{t_1}{j_{i_1}}, \binom{t_{i_1}}{j_{i_1-1}}, \ldots, \binom{t_r}{j_1} \text{ of billetters such that applying the RSK algorithm to } \binom{t_1}{j_{i_1}}, \ldots, \binom{t_r}{j_1} \text{ gives back } (P, Q). \text{ The question is why we have } \binom{t_1}{j_1} \leq \text{lex} \cdots \leq \text{lex} \binom{t_r}{j_r}. \text{ Since the chain of inequalities } i_1 \leq i_2 \leq \cdots \leq i_r \text{ is clear from the choice of entry to reverse-bump, it only remains to show that for every string } i_m = i_{m+1} = \cdots = i_{m+r} \text{ of equal top letters, the corresponding bottom letters weakly increase (that is, } j_m \leq j_{m+1} \leq \cdots \leq j_{m+r}). \text{ One way to see it is to apply part (b) of the row bumping lemma.}\)

\(\text{Speaking of coefficients makes sense here since } \Lambda (x) \otimes \Lambda (y) \otimes \Lambda (z) \text{ embeds into } k[[x, y, z]]\).
and on the other hand
\[ s_{\lambda}(x, y) = \sum_{\mu \subseteq \lambda} s_{\mu}(x)s_{\lambda/\mu}(y) = \sum_{\mu \in \text{Par}} s_{\mu}(x)s_{\lambda/\mu}(y) \]

one concludes another standard interpretation for \( c_{\mu,\nu}^\lambda \):
\[ s_{\lambda/\mu} = \sum_{\nu} c_{\mu,\nu}^\lambda s_{\nu} \]

whenever \( \lambda \) and \( \mu \) are two partitions (not necessarily satisfying \( \mu \subseteq \lambda \)). In particular, \( c_{\mu,\nu}^\lambda \) vanishes unless \( \mu \subseteq \lambda \). Consequently, \( c_{\mu,\nu}^\lambda \) vanishes unless \( \nu \subseteq \lambda \) as well (since \( c_{\mu,\nu}^\lambda = c_{\nu,\mu}^\lambda \)) and furthermore vanishes unless the equality \( |\mu| + |\nu| = |\lambda| \) holds\(^\text{75}\). Altogether, we conclude that \( c_{\mu,\nu}^\lambda \) vanishes unless \( \mu, \nu \subseteq \lambda \) and \( |\mu| + |\nu| = |\lambda| \).

**Exercise 2.38.** Show that any four partitions \( \kappa, \lambda, \varphi \) and \( \psi \) satisfy
\[ \sum_{\rho \in \text{Par}} c_{\kappa,\lambda,\rho}^\varphi = \sum_{(\alpha,\beta,\gamma,\delta) \in \text{Par}^4} c_{\alpha,\beta}^\gamma c_{\beta,\gamma}^\delta c_{\gamma,\delta}^\psi, \]

**Exercise 2.39.** (a) For any partition \( \mu \), prove that
\[ \sum_{\lambda \in \text{Par}} s_{\lambda}(x) s_{\lambda/\mu}(y) = s_{\mu}(x) \cdot \prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} \]

in the power series ring \( k[[x, y]] = k[[x_1, x_2, x_3, \ldots, y_1, y_2, y_3, \ldots]] \).

(b) Let \( \alpha \) and \( \beta \) be two partitions. Show that
\[ \sum_{\lambda \in \text{Par}} s_{\lambda/\alpha}(x) s_{\lambda/\beta}(y) = \left( \sum_{\rho \in \text{Par}} s_{\rho/\beta}(y) s_{\alpha/\rho}(x) \right) \cdot \prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} \]

in the power series ring \( k[[x, y]] = k[[x_1, x_2, x_3, \ldots, y_1, y_2, y_3, \ldots]] \).

[Hint: For (b), expand the product
\[ \prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} \prod_{i,j=1}^{\infty} (1 - x_i w_j)^{-1} \prod_{i,j=1}^{\infty} (1 - z_i y_j)^{-1} \prod_{i,j=1}^{\infty} (1 - z_i w_j)^{-1} \]

in the power series ring \( k[[x_1, x_2, x_3, \ldots, y_1, y_2, y_3, \ldots, z_1, z_2, z_3, \ldots, w_1, w_2, w_3, \ldots]] \) in two ways: once by applying Theorem 2.31 to the two variable sets \( (z, x) \) and \( (w, y) \) and then using (2.13); once again by applying (2.24) to the two variable sets \( z \) and \( w \) and then applying Exercise 2.39(a) twice.]

The statement of Exercise 2.39(b) is known as the *skew Cauchy identity*, and appears in Sagan-Stanley [139, Cor. 6.12], Stanley [154, exercise 7.27(c)] and Macdonald [106, §15.5, example 26]; it seems to be due to Zelevinsky. It generalizes the statement of Exercise 2.39(a), which in turn is a generalization of Theorem 2.31.

**Definition 2.40.** Define the *Hall inner product* on \( \Lambda \) to be the \( k \)-bilinear form \( \langle \cdot, \cdot \rangle \) which makes \( \{ s_{\lambda} \} \) an orthonormal basis, that is, \( \langle s_{\lambda}, s_{\nu} \rangle = \delta_{\lambda,\nu} \).

**Exercise 2.41.** (a) If \( n \) and \( m \) are two distinct nonnegative integers, and if \( f \in \Lambda_n \) and \( g \in \Lambda_m \), then show that \( \langle f, g \rangle = 0 \).

(b) If \( n \in \mathbb{N} \) and \( f \in \Lambda_n \), then prove that \( \langle h_n, f \rangle = f(1) \) (where \( f(1) \) is defined as in Exercise 2.2).

The Hall inner product induces a \( k \)-module homomorphism \( \Lambda \rightarrow \Lambda^o \) (sending every \( f \in \Lambda \) to the \( k \)-linear map \( \Lambda \rightarrow k \), \( g \mapsto \langle f, g \rangle \)). This homomorphism is invertible (since the Hall inner product has an orthonormal basis), so that \( \Lambda^o \cong \Lambda \) as \( k \)-modules. But in fact, more can be said:

**Corollary 2.42.** The isomorphism \( \Lambda^o \cong \Lambda \) induced by the Hall inner product is an isomorphism of Hopf algebras.

\(^{74}\)In the last equality, we removed the condition \( \mu \subseteq \lambda \) on the addends of the sum; this does not change the value of the sum (because we have \( s_{\lambda/\mu} = 0 \) whenever we don’t have \( \mu \subseteq \lambda \)).

\(^{75}\)In fact, this is clear when we don’t have \( \mu \subseteq \lambda \). When we do have \( \mu \subseteq \lambda \), this follows from observing that \( s_{\lambda/\mu} \in \Lambda_{|\lambda/\mu|} \) has zero coefficient before \( s_{\mu} \) whenever \( |\mu| + |\nu| \neq |\lambda| \).
**Proposition 2.43.** One can also expand

\[(2.28) \quad \prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} = \sum_{\lambda \in \mathcal{P}} h_{\lambda}(x) m_{\lambda}(y) = \sum_{\lambda \in \mathcal{P}} z_{\lambda}^{-1} p_{\lambda}(x)p_{\lambda}(y)\]

where \(z_{\lambda} := m_1! \cdot 1^m_1 \cdot m_2! \cdot 2^m_2 \cdots \) if \(\lambda = (1^{m_1}, 2^{m_2}, \ldots)\) with multiplicity \(m_i\) for the part \(i\). (Here, we assume that \(\mathbb{Q}\) is a subring of \(\mathbb{k}\) for the last equality.)

**Remark 2.44.** It is relevant later (and explains the notation) that \(z_{\lambda}\) is the size of the \(\mathfrak{S}_n\)-centralizer subgroup for a permutation having cycle type \(\lambda\) with \(|\lambda| = n\). This is a classical (and fairly easy) result (see, e.g., [138, Prop. 1.1.1] or [154, Prop. 7.7.3] for a proof).

**Proof of Proposition 2.43.** For the first expansion, note that \((2.14)\) shows

\[\prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} = \prod_{j=1}^{\infty} \sum_{n \geq 0} h_n(x) y_j^n = \sum_{\text{weak compositions} (n_1, n_2, \ldots)} \left( h_{n_1}(x) h_{n_2}(x) \cdots \right) (y_1^{n_1} y_2^{n_2} \cdots) \]

\[= \sum_{\lambda \in \mathcal{P}} \sum_{\text{weak compositions} (n_1, n_2, \ldots) \text{ satisfying} (n_1, n_2, \ldots) \in \mathcal{P}^\infty_{\lambda}} \left( h_{n_1}(x) h_{n_2}(x) \cdots \right) (y_1^{n_1} y_2^{n_2} \cdots) = \sum_{\lambda \in \mathcal{P}} h_{\lambda}(x) \sum_{\text{weak compositions} (n_1, n_2, \ldots) \text{ satisfying} (n_1, n_2, \ldots) \in \mathcal{P}^\infty_{\lambda}} y^{(n_1, n_2, \ldots)} = m_{\lambda}(y).\]

For the second expansion (and for later use in the proof of Theorem 4.59) note that

\[(2.29) \quad \log H(t) = \log \prod_{i=1}^{\infty} (1 - x_i t)^{-1} = \sum_{i=1}^{\infty} - \log(1 - x_i t) = \sum_{i=1}^{\infty} \sum_{m=1}^{\infty} \frac{(x_i t)^m}{m} = \sum_{m=1}^{\infty} \frac{1}{m} p_m(x) t^m\]

so that taking \(\frac{d}{dt}\) then shows that

\[(2.30) \quad P(t) := \sum_{m \geq 0} p_{m+1} t^m = \frac{H'(t)}{H(t)} = H'(t) E(-t).\]

---

\(^76\)If \(\sigma\) is a permutation of a finite set \(X\), then the cycle type of \(\sigma\) is defined as the list of the lengths of all cycles of \(\sigma\) (that is, of all orbits of \(\sigma\) acting on \(X\)) written in decreasing order. This is clearly a partition of \(|X|\). (Some other authors write it in increasing order instead, or treat it as a multiset.)

For instance, the permutation of the set \(\{0, 3, 6, 9, 12\}\) which sends 0 to 3, 3 to 9, 6 to 6, 9 to 0, and 12 to 12 is \((3, 1, 1)\), since the cycles of this permutation have lengths 3, 1 and 1.

It is known that two permutations in \(\mathfrak{S}_n\) have the same cycle type if and only if they are conjugate. Thus, for a given partition \(\lambda\) with \(|\lambda| = n\), any two permutations in \(\mathfrak{S}_n\) having cycle type \(\lambda\) are conjugate and therefore their \(\mathfrak{S}_n\)-centralizer subgroups have the same size.
A similar calculation shows that

\[
\log \prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} = \sum_{m=1}^{\infty} \frac{1}{m} p_m(x) p_m(y)
\]
and hence

\[
\prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} = \exp \left( \sum_{m=1}^{\infty} \frac{1}{m} p_m(x) p_m(y) \right) = \prod_{m=1}^{\infty} \exp \left( \frac{1}{m} p_m(x) p_m(y) \right)
\]

\[
= \prod_{m=1}^{\infty} \sum_{k=0}^{\infty} \frac{1}{k!} \left( \frac{1}{m} p_m(x) p_m(y) \right)^k = \sum_{\text{weak compositions}} \prod_{m=1}^{\infty} \left( \frac{1}{k_m!} \left( \frac{1}{m} p_m(x) p_m(y) \right)^{k_m} \right)
\]

(by the product rule)

\[
= \sum_{\text{weak compositions}} \prod_{m=1}^{\infty} \frac{1}{k_m!} \left( \frac{1}{m} p_m(x) p_m(y) \right)^{k_m} = \sum_{\text{weak compositions}} \prod_{m=1}^{\infty} \frac{1}{k_m!} \left( \frac{1}{m} p_m(x) p_m(y) \right)^{k_m}
\]

\[
= \sum_{\text{weak compositions}} \frac{p_{\lambda}(x) p_{\lambda}(y)}{\lambda!}
\]
due to the fact that every partition can be uniquely written in the form \((1^{k_1} 2^{k_2} 3^{k_3} \cdots)\) with \((k_1, k_2, k_3, \ldots)\) a weak composition.

**Corollary 2.45.** (a) With respect to the Hall inner product on \(\Lambda\), one also has dual bases \(\{h_{\lambda}\}\) and \(\{m_{\lambda}\}\).

(b) If \(\mathbb{Q}\) is a subring of \(k\), then \(\{p_{\lambda}\}\) and \(\{z_{\lambda}^{-1} p_{\lambda}\}\) are also dual bases with respect to the Hall inner product on \(\Lambda\).

(c) If \(\mathbb{R}\) is a subring of \(k\), then \(\left\{ \frac{p_{\lambda}}{\sqrt{z_{\lambda}}} \right\}\) is an orthonormal basis of \(\Lambda\) with respect to the Hall inner product.

**Proof.** Since (2.24) and (2.28) showed

\[
\prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} = \sum_{\lambda \in \text{Par}} s_\lambda(x) s_\lambda(y) = \sum_{\lambda \in \text{Par}} h_\lambda(x) m_\lambda(y) = \sum_{\lambda \in \text{Par}} p_\lambda(x) z^{-1}_\lambda p_\lambda(y) = \sum_{\lambda \in \text{Par}} \frac{p_{\lambda}(x) p_{\lambda}(y)}{\sqrt{z_{\lambda}}} \sqrt{z_{\lambda}}
\]
it suffices to show that any pair of graded bases\(^\text{77}\) \(\{u_{\lambda}\}\), \(\{v_{\lambda}\}\) having

\[
\sum_{\lambda \in \text{Par}} s_\lambda(x) s_\lambda(y) = \sum_{\lambda \in \text{Par}} u_\lambda(x) v_\lambda(y)
\]
will be dual with respect to \((\cdot, \cdot)\). Write transition matrices \(A = (a_{\nu, \lambda})\), \(B = (b_{\nu, \lambda})\) uniquely expressing

\[
u_\lambda = \sum_{\nu} a_{\nu, \lambda} s_\nu,
\]

\[
u_\lambda = \sum_{\nu} b_{\nu, \lambda} s_\nu.
\]

Then orthonormality of the \(s_\lambda\) gives \((u_\alpha, v_\beta) = \sum_\nu a_{\nu, \alpha} b_{\nu, \beta}\), and hence we want \(A' B = I\), that is, \(B^{-1} = A'\). On the other hand, one has

\[
\sum_{\lambda} s_\lambda(x) s_\lambda(y) = \sum_{\lambda} u_\lambda(x) v_\lambda(y) = \sum_{\lambda} \sum_{\nu} a_{\nu, \lambda} s_\nu(x) \sum_{\rho} b_{\rho, \lambda} s_\rho(y).
\]

\(^{77}\)A basis \(\{w_{\lambda}\}\) of \(\Lambda\) indexed by the partitions \(\lambda\) is said to be graded if every \(n \in \mathbb{N}\) and \(\lambda \in \text{Par}_n\) satisfy \(w_{\lambda} \in \Lambda_n\).
Comparing coefficients\(^ {78}\) of \(s_\nu(x)s_\rho(y)\) forces \(\sum_\lambda a_\nu,\lambda b_\rho,\lambda = \delta_{\nu,\rho}\), or in other words, \(AB^t = I\). Since \(A\) is invertible (being a transition matrix between two bases\(^ {79}\)), this yields \(B^tA = I\), and hence \(A^tB = I\), as desired.

Corollary 2.45 is a known and fundamental fact\(^ {80}\). However, our definition of the Hall inner product is unusual; most authors (e.g., Macdonald in [106, §1.4], Hazewinkel/Gubareni/Kirichenko in [68, Def. 4.1.21], and Stanley in [154, (7.30)]) define the Hall inner product as the bilinear form satisfying \((h_\lambda, m_\mu) = \delta_{\lambda,\mu}\) (or, alternatively, \((m_\lambda, h_\mu) = \delta_{\lambda,\mu}\)), and only later prove that the basis \(\{s_\lambda\}\) is orthonormal with respect to this scalar product. (Of course, the fact that this definition is equivalent to our Definition 2.40 follows either from this orthonormality, or from our Corollary 2.45(a).)

The tactic applied in the proof of Corollary 2.45 can not only be used to show that certain basis of \(\Lambda\) are dual, but also, with a little help from linear algebra over rings (Exercise 2.46), it can be strengthened to show that certain families of symmetric functions are bases to begin with, as we will see in Exercise 2.47.

**Exercise 2.46.** (a) Prove that if an endomorphism of a finitely generated \(k\)-module is surjective, then this endomorphism is a \(k\)-module isomorphism.

(b) Let \(A\) be a finite free \(k\)-module with finite basis \((\gamma_i)_{i\in I}\). Let \((\beta_i)_{i\in I}\) be a family of elements of \(A\) which spans the \(k\)-module \(A\). Prove that \((\beta_i)_{i\in I}\) is a \(k\)-basis of \(A\).

**Exercise 2.47.** (a) Assume that for every partition \(\lambda\), two homogeneous elements \(u_\lambda\) and \(v_\lambda\) of \(\Lambda\), both having degree \(|\lambda|\), are given. Assume further that
\[
\sum_{\lambda \in \text{Par}} s_\lambda(x)s_\lambda(y) = \sum_{\lambda \in \text{Par}} u_\lambda(x)v_\lambda(y)
\]
in \(k[[x,y]] = k[[x_1, x_2, x_3, \ldots, y_1, y_2, y_3, \ldots]]\). Show that \((u_\lambda)_{\lambda \in \text{Par}}\) and \((v_\lambda)_{\lambda \in \text{Par}}\) are \(k\)-bases of \(\Lambda\), and actually are dual bases with respect to the Hall inner product on \(\Lambda\).

(b) Use this to give a new proof of the fact that \((h_\lambda)_{\lambda \in \text{Par}}\) is a \(k\)-basis of \(\Lambda\).

**Exercise 2.48.** Prove that \(\sum_{m\geq 0} p_{m+1}t^m = \frac{H'(t)}{H(t)}\). (This was proven in (2.30) in the case when \(Q\) is a subring of \(k\), but here we make no requirements on \(k\).)

### 2.6. Bialternants, Littlewood-Richardson: Stembridge’s concise proof.

There is a more natural way in which Schur functions arise as a \(k\)-basis for \(\Lambda\), coming from consideration of polynomials in a finite variable set, and the relation between those which are symmetric and those which are alternating.

For the remainder of this section, fix a nonnegative integer \(n\), and let \(x = (x_1, \ldots, x_n)\) be a finite variable set. This means that \(s_{\lambda/\mu}(x) = s_{\lambda/\mu}(x_i) = \sum_T x^{\text{cont}(T)}\) is a generating function for column-strict tableaux \(T\) as in Definition 2.19, but with the extra condition that \(T\) has entries in \(\{1, 2, \ldots, n\}\). As a consequence, \(s_{\lambda/\mu}\) is a polynomial in \(k[x_1, x_2, \ldots, x_n]\) (not just a power series), since there are only finitely many column-strict tableaux \(T\) of shape \(\lambda/\mu\) having all their entries in \(\{1, 2, \ldots, n\}\). We will assume without further mention that all partitions appearing in the section have at most \(n\) parts.

**Definition 2.49.** Let \(k\) be the ring \(\mathbb{Z}\) or a field of characteristic not equal to 2. (We require this to avoid certain annoyances in the discussion of alternating polynomials in characteristic 2.)

Say that a polynomial \(f(x) = f(x_1, \ldots, x_n)\) is alternating if for every permutation \(w\) in \(S_n\) one has that
\[(w f)(x) = f(x_{w(1)}, \ldots, x_{w(n)}) = \text{sgn}(w)f(x).\]

Let \(\Lambda_{\text{sgn}} \subset k[x_1, \ldots, x_n]\) denote the subset of alternating polynomials\(^ {82}\).
As with $\Lambda$ and its monomial basis $\{m_\lambda\}$, there is an obvious $k$-basis for $\Lambda^{\text{sgn}}$, coming from the fact that a polynomial $f = \sum a_\alpha x^\alpha$ is alternating if and only if $c_w(\alpha) = \text{sgn}(w)c_\alpha$ for every $w$ in $S_n$ and every $\alpha \in \mathbb{N}^n$. This means that every alternating $f$ is a $k$-linear combination of the following elements.

**Definition 2.50.** For $\alpha = (\alpha_1, \ldots, \alpha_n)$ in $\mathbb{N}^n$, define the *alternant*

$$a_\alpha := \sum_{w \in S_n} \text{sgn}(w)w(x^\alpha) = \det \begin{bmatrix} x_{1}^{\alpha_1} & \cdots & x_{n}^{\alpha_n} \\ x_{1}^{\alpha_1} & \cdots & \vdots \\ \vdots & \cdots & \vdots \\ x_{n}^{\alpha_1} & \cdots & x_{n}^{\alpha_n} \end{bmatrix}.$$  

**Example 2.51.** One has

$$a_{(1,5,0)} = x_1^1x_2^5x_3^0 - x_1^5x_2^0x_3^5 - x_1^1x_2^0x_3^5 + x_1^0x_2^5x_3^1 + x_1^0x_2^1x_3^5 = -a_{(5,1,0)}.$$  

Similarly, $a_w(\alpha) = \text{sgn}(w)a_\alpha$ for every $w \in S_n$ and every $\alpha \in \mathbb{N}^n$. Meanwhile, $a_{(5,2,2)} = 0$ since the transposition $t = (1 2 3)$ fixes $(5, 2, 2)$ and hence

$$a_{(5,2,2)} = t(a_{(5,2,2)}) = \text{sgn}(t)a_{(5,2,2)} = -a_{(5,2,2)}.$$  

83 Alternatively, $a_{(5,2,2)} = 0$ as it is a determinant of a matrix with two equal columns. Similarly, $a_\alpha = 0$ for every $n$-tuple $\alpha \in \mathbb{N}^n$ having two equal entries.

This example illustrates that, for a $k$-basis for $\Lambda^{\text{sgn}}$, one can restrict attention to alternants $a_\alpha$ in which $\alpha$ is a *strict partition*, i.e., in which $\alpha$ satisfies $\alpha_1 > \alpha_2 > \cdots > \alpha_n$. One can therefore uniquely express $\alpha = \lambda + \rho$, where $\lambda$ is a (weak) partition $\lambda_1 \geq \cdots \geq \lambda_n \geq 0$ and where $\rho := (n - 1, n - 2, \ldots, 2, 1, 0)$ is sometimes called the *staircase partition* 84. For example $\alpha = (5, 1, 0) = (3, 0, 0) + (2, 1, 0) = \lambda + \rho$.

**Proposition 2.52.** Let $k$ be the ring $\mathbb{Z}$ or a field of characteristic not equal to 2.

The alternants $\{a_\alpha\}$ as $\lambda$ runs through the partitions with at most $n$ parts form a $k$-basis for $\Lambda^{\text{sgn}}$. In addition, the *bilateral*nt $\left\{ \frac{a_{\lambda + \rho}}{a_\rho} \right\}$ as $\lambda$ runs through the same set form a $k$-basis for $\Lambda(x_1, \ldots, x_n) = k[x_1, \ldots, x_n]^S_n$.

**Proof.** The first assertion should be clear from our previous discussion: the alternants $\{a_{\lambda + \rho}\}$ span $\Lambda^{\text{sgn}}$ by definition, and they are $k$-linearly independent because they are supported on disjoint sets of monomials $x^\alpha$.

The second assertion follows from the first, after proving the following *Claim*: $f(x)$ lies in $\Lambda^{\text{sgn}}$ if and only if $f(x) = a_\rho \cdot g(x)$ where $g(x)$ lies in $k[x]\natural^{S_n}$ and where

$$a_\rho = \det(x_i^{n-j})_{i,j=1,2,\ldots,n} = \prod_{1 \leq i < j \leq n} (x_i - x_j)$$

is the *Vandermonde determinant/product*. In other words

$$\Lambda^{\text{sgn}} = a_\rho \cdot k[x]\natural^{S_n}$$

is a free $k[x]\natural^{S_n}$-module of rank one, with $a_\rho$ as its $k[x]\natural^{S_n}$-basis element.

To see the Claim, first note the inclusion

$$\Lambda^{\text{sgn}} \supset a_\rho \cdot k[x]\natural^{S_n}$$

since the product of a symmetric polynomial and an alternating polynomial is an alternating polynomial. For the reverse inclusion, note that since an alternating polynomial $f(x)$ changes sign whenever one exchanges two distinct variables $x_i, x_j$, it must vanish upon setting $x_i = x_j$, and therefore be divisible by $x_i - x_j$, so

83 One subtlety should be addressed: We want to prove that $a_{(5,2,2)} = 0$ in $k[x_1, \ldots, x_n]$ for every commutative ring $k$. It is clearly enough to prove that $a_{(5,2,2)} = 0$ in $\mathbb{Z}[x_1, \ldots, x_n]$. Since 2 is not a zero-divisor in $\mathbb{Z}[x_1, \ldots, x_n]$, we can achieve this by showing that $a_{(5,2,2)} = -a_{(5,2,2)}$. We would not be able to make this argument directly over an arbitrary commutative ring $k$.

84 The name is owed to its Ferrers shape. For instance, if $n = 5$, then the Ferrers diagram of $\rho$ (represented using dots) has the form

```
• • •
• •
•
```
divisible by the entire product $\prod_{1 \leq i < j \leq n} (x_i - x_j) = a_\rho$. But then the quotient $g(x) = \frac{f(x)}{a_\rho}$ is symmetric, as it is a quotient of two alternating polynomials.

Let us now return to the general setting, where $k$ is an arbitrary commutative ring. We are not requiring that the assumptions of Proposition 2.52 be valid; we can still study the $a_\alpha$ of Definition 2.50, but we cannot use Proposition 2.52 anymore. We will show that the fraction $\frac{a_\lambda + \rho}{a_\rho}$ is nevertheless a well-defined polynomial in $\Lambda(x_1, \ldots, x_n)$ whenever $\lambda$ is a partition\textsuperscript{[85]}, and in fact equals the Schur function $s_\lambda(x)$. As a consequence, the mysterious bialternant basis $\{ a_\lambda + \rho \}$ of $\Lambda(x_1, \ldots, x_n)$ defined in Proposition 2.52 still exists in the general setting, and is plainly the Schur functions $\{ s_\lambda(x) \}$. Stembridge [157] noted that one could give a remarkably concise proof of an even stronger assertion, which simultaneously gives one of the standard combinatorial interpretations for the Littlewood-Richardson coefficients $c_{\mu,\nu}^\lambda$. For the purposes of stating it, we introduce for a tableau $T$ the notation $T|_{\text{cols} \geq j}$ (resp. $T|_{\text{cols} \leq j}$) to indicate the subtableau which is the restriction of $T$ to the union of its columns $j, j + 1, j + 2, \ldots$ (resp. columns $1, 2, \ldots, j$).

**Theorem 2.53.** For partitions $\lambda, \mu, \nu$ with $\mu \subseteq \lambda$, one has$^{86}$

$$a_\nu + \rho s_{\lambda/\mu} = \sum_T a_{\nu + \text{cont}(T) + \rho}$$

where $T$ runs through all column-strict tableaux with entries in $\{1, 2, \ldots, n\}$ of shape $\lambda/\mu$ with the property that for all $j = 1, 2, \ldots$ one has $\nu + \text{cont}(T|_{\text{cols} \geq j})$ a partition.

Before proving Theorem 2.53, let us see some of its consequences.

**Corollary 2.54.** For any partition $\lambda$, we have$^{87}$

$$s_{\lambda}(x) = \frac{a_{\lambda + \rho}}{a_\rho}.$$  

**Proof.** Take $\nu = \emptyset$ in Theorem 2.53. Note that for any $\lambda$, there is only one column-strict tableau $T$ of shape $\lambda$ having each $\text{cont}(T|_{\text{cols} \geq j})$ a partition, namely the one having every entry in row $i$ equal to $i$:

$$
\begin{array}{cccccc}
1 & 1 & 1 & 1 & 1 \\
2 & 2 & 2 \\
3 & 3 & 3 \\
4 & & & & & \\
\end{array}
$$

Furthermore, this $T$ has $\text{cont}(T) = \lambda$, so the theorem says $a_\rho s_{\lambda} = a_{\lambda + \rho}$.

**Example 2.55.** For $n = 2$, so that $\rho = (1,0)$, if we take $\lambda = (4,2)$, then one has

$$
\begin{align*}
\frac{a_{\lambda + \rho}}{a_\rho} &= \frac{a_{(4,2) + (1,0)}}{a_{(1,0)}} = \frac{a_{(5,2)}}{a_{(1,0)}} \\
&= \frac{x_1^5 x_2^3 - x_1^2 x_2^5}{x_1 - x_2} \\
&= x_1^4 x_2 + x_1^3 x_2^3 + x_1^2 x_2^4 \\
&= x \begin{pmatrix} 1 & 1 & 1 & 1 \\ 2 & 2 & 2 & 2 \\ \end{pmatrix} + x \begin{pmatrix} 1 & 1 & 1 \\ 2 & 2 & 2 \\ \end{pmatrix} + x \begin{pmatrix} 1 & 1 \\ 2 & 2 \\ \end{pmatrix} \\
&= s_{(4,2)} = s_{\lambda}.
\end{align*}
$$

\textsuperscript{85}This can also be deduced by base change from the $k = \mathbb{Z}$ case of Proposition 2.52.

\textsuperscript{86}Again, we can drop the requirement that $\mu \subseteq \lambda$, provided that we understand that there are no column-strict tableaux of shape $\lambda/\mu$ unless $\mu \subseteq \lambda$.

\textsuperscript{87}Notice that division by $a_\rho$ is unambiguous in the ring $k[x_1, \ldots, x_n]$, since $a_\rho$ is not a zero-divisor (in fact, $a_\rho = \prod_{1 \leq i < j \leq n} (x_i - x_j)$ is the product of the binomials $x_i - x_j$, none of which is a zero-divisor).
Some authors use the equality in Corollary 2.54 to define the Schur polynomial $s_\lambda(x_1, x_2, \ldots, x_n)$ in $n$ variables; this definition, however, has the drawback of not generalizing easily to infinitely many variables or to skew Schur functions.\footnote{With some effort, it is possible to use Corollary 2.54 in order to define the Schur function $s_\lambda$ in infinitely many variables. Indeed, one can define this Schur function as the unique element of $\Lambda$ whose evaluation at $(x_1, x_2, \ldots, x_n)$ equals $s_\lambda(x_{1+\rho})$ for every $n \in \mathbb{N}$. If one wants to use such a definition, however, one needs to check that such an element exists. This is the approach to defining $s_\lambda$ taken in [92, Definition 1.4.2] and in [106, §I.3].}

Next divide through by $a_\mu$ on both sides of Theorem 2.53 (and use Corollary 2.54) to give the following.

**Corollary 2.56.** For partitions $\lambda, \mu, \nu$ having at most $n$ parts, one has

\[
(2.32) \quad s_\nu s_{\lambda/\mu} = \sum_T s_{\nu+\cont(T)}
\]

where $T$ runs through the same set as in Theorem 2.53. In particular, taking $\nu = \emptyset$, we obtain

\[
(2.33) \quad s_{\lambda/\mu} = \sum_T s_{\cont(T)}
\]

where in the sum $T$ runs through all column-strict tableaux of shape $\lambda/\mu$ for which each $\cont(T|_{\text{cols} \geq j})$ is a partition.

**Proof of Theorem 2.53.** Start by rewriting the left side of the theorem, and using the fact that $w(s_{\lambda/\mu}) = s_{\lambda/\mu}$ for any $w$ in $S_n$:

\[
\begin{align*}
a_{\nu+\rho}s_{\lambda/\mu} &= \sum_{w \in S_n} \sgn(w)x^{w(\nu+\rho)}w(s_{\lambda/\mu}) \\
&= \sum_{w \in S_n} \sgn(w)x^{w(\nu+\rho)} \sum_{\text{column-strict } T \text{ of shape } \lambda/\mu} x^{w(\cont(T))} \\
&= \sum_{\text{column-strict } T \text{ of shape } \lambda/\mu} \sum_{w \in S_n} \sgn(w)x^{w(\nu+\cont(T)+\rho)} \\
&= \sum_{\text{column-strict } T \text{ of shape } \lambda/\mu} a_{\nu+\cont(T)+\rho}
\end{align*}
\]

We wish to cancel out all the summands indexed by column-strict tableaux $T$ which fail any of the conditions that $\nu + \cont(T|_{\text{cols} \geq j})$ be a partition. Given such a $T$, find the maximal $j$ for which it fails this condition\footnote{Such a $j$ exists because $\nu + \cont(T|_{\text{cols} \geq j})$ is a partition for all sufficiently high $j$ (in fact, $\nu$ itself is a partition).}, and then find the minimal $k$ for which

\[
\nu_k + \cont_k(T|_{\text{cols} \geq j}) < \nu_{k+1} + \cont_{k+1}(T|_{\text{cols} \geq j}).
\]

Maximality of $j$ forces

\[
\nu_k + \cont_k(T|_{\text{cols} \geq j+1}) \geq \nu_{k+1} + \cont_{k+1}(T|_{\text{cols} \geq j+1}).
\]

Since column-strictness implies that column $j$ of $T$ can contain at most one occurrence of $k$ or of $k + 1$ (or neither or both), the previous two inequalities imply that column $j$ must contain an occurrence of $k + 1$ and no occurrence of $k$, so that

\[
\nu_k + \cont_k(T|_{\text{cols} \geq j}) + 1 = \nu_{k+1} + \cont_{k+1}(T|_{\text{cols} \geq j}).
\]

This implies that the adjacent transposition $t_{k,k+1}$ swapping $k$ and $k+1$ fixes the vector $\nu + \cont(T|_{\text{cols} \geq j}) + \rho$.

Now create a new tableau $T^*$ from $T$ by applying the Bender–Knuth involution (from the proof of Proposition 2.11) on letters $k, k + 1$, but only to columns $1, 2, \ldots, j - 1$ of $T$, leaving columns $j, j + 1, j + 2, \ldots$ unchanged. One should check that $T^*$ is still column-strict, but this holds because column $j$ of $T$ has no occurrences of letter $k$. Note that

\[
t_{k,k+1} \cont(T|_{\text{cols} \leq j-1}) = \cont(T^*|_{\text{cols} \leq j-1})
\]

and hence

\[
t_{k,k+1}(\nu + \cont(T) + \rho) = \nu + \cont(T^*) + \rho
\]

so that $a_{\nu+\cont(T)+\rho} = -a_{\nu+\cont(T^*)+\rho}$.
Because $T, T^*$ have exactly the same columns $j, j + 1, j + 2, \ldots$, the tableau $T^*$ is also a violator of at least one of the conditions that $\nu + \text{cont}(T^*|_{\text{cols} \geq j})$ be a partition, and has the same choice of maximal $j$ and minimal $k$ as did $T$. Hence the map $T \mapsto T^*$ is an involution on the violators that lets one cancel their summands $a_{\nu + \text{cont}(T) + \rho}$ and $a_{\nu + \text{cont}(T^*) + \rho}$ in pairs.

So far (in this section) we have worked with a finite set of variables $x_1, x_2, \ldots, x_n$ (where $n$ is a fixed nonnegative integer) and with partitions having at most $n$ parts. We now drop these conventions and restrictions; thus, partitions again mean arbitrary partitions, and $x$ again means the infinite family $(x_1, x_2, x_3, \ldots)$ of variables. In this setting, we have the following analogue of Corollary 2.56:

**Corollary 2.57.** For partitions $\lambda, \mu, \nu$ (of any lengths), one has

$$s_\mu s_{\lambda/\mu} = \sum_T s_{\nu + \text{cont}(T)}$$

where $T$ runs through all column-strict tableaux of shape $\lambda/\mu$ with the property that for all $j = 1, 2, \ldots$ one has $\nu + \text{cont}(T|_{\text{cols} \geq j})$ a partition. In particular, taking $\nu = \emptyset$, we obtain

$$s_{\lambda/\mu} = \sum_T s_{\text{cont}(T)}$$

where in the sum $T$ runs through all column-strict tableaux of shape $\lambda/\mu$ for which each cont$(T|_{\text{cols} \geq j})$ is a partition.

**Proof of Corollary 2.57.** Essentially, Corollary 2.57 is obtained from Corollary 2.56 by “letting $n$ (that is, the number of variables) tend to $\infty$”. This can be formalized in different ways: One way is to endow the ring of power series $k[[x]] = k[[x_1, x_2, x_3, \ldots]]$ with the coefficientwise topology\textsuperscript{90}, and to show that the left hand side of (2.32) tends to the left hand side of (2.34) when $n \to \infty$, and the same holds for the right hand side. A different approach proceeds by regarding $\Lambda$ as the inverse limit of the $\Lambda \langle x_1, x_2, \ldots, x_n \rangle$. □

Comparing coefficients of a given Schur function $s_\nu$ in (2.35), we obtain the following version of the Littlewood-Richardson rule.

**Corollary 2.58.** For partitions $\lambda, \mu, \nu$ (of any lengths), the Littlewood-Richardson coefficient $c_{\mu,\nu}^\lambda$ counts column-strict tableaux $T$ of shape $\lambda/\mu$ with cont$(T) = \nu$ having the property that each cont$(T|_{\text{cols} \geq j})$ is a partition.

### 2.7. The Pieri and Assaf-McNamara skew Pieri rule

The classical Pieri rule refers to two special cases of the Littlewood-Richardson rule. To state them, recall that a skew shape is called a horizontal (resp. vertical) strip if no two of its cells lie in the same column (resp. row). A horizontal (resp. vertical) $n$-strip (for $n \in \mathbb{N}$) shall mean a horizontal (resp. vertical) strip of size $n$ (that is, having exactly $n$ cells).

**Theorem 2.59.** For every partition $\lambda$ and any $n \in \mathbb{N}$, we have

$$s_{\lambda^+} h_n = \sum_{\lambda^+ = \lambda^+ / \lambda \text{ is a horizontal } n\text{-strip}} s_{\lambda^+}$$

$$s_{\lambda^+} e_n = \sum_{\lambda^+ = \lambda^+ / \lambda \text{ is a vertical } n\text{-strip}} s_{\lambda^+}$$

\textsuperscript{90}This topology is defined as follows:

We endow the ring $k$ with the discrete topology. Then, we can regard the $k$-module $k[[x]]$ as a direct product of infinitely many copies of $k$ (by identifying every power series in $k[[x]]$ with the family of its coefficients). Hence, the product topology is a well-defined topology on $k[[x]]$: this topology is denoted as the coefficientwise topology. Its name is due to the fact that a sequence $(a_n)_{n \in \mathbb{N}}$ of power series converges to a power series $a$ with respect to this topology if and only if for every monomial $m$, all sufficiently high $n \in \mathbb{N}$ satisfy

$$(\text{the coefficient of } m \text{ in } a_n) = (\text{the coefficient of } m \text{ in } a).$$
Example 2.60.
\[ s \]
\[ s \]
\[ h \]
\[ s \]
\[ s \]
\[ s \]
\[ s \]
\[ s \]
\[ s \]
\[ s \]

Proof of Theorem 2.59. For the first Pieri formula involving \( h_n \), as \( h_n = s(s) \) one has
\[ s_s^{\lambda} h_n = \sum \lambda\mu \mu^{\lambda/\mu} + \sum \lambda\mu \mu^{\lambda/\mu} = s_s^{\lambda/\mu} + s_s^{\lambda/\mu} \]

Assaf and McNamara [9] recently proved an elegant generalization.

Theorem 2.61. For any partitions \( \lambda \) and \( \mu \) and any \( n \in \mathbb{N} \), we have\(^91\)
\[ s_{\lambda/\mu} h_n = \sum_{\lambda^{\mu}, \mu \subseteq \lambda} (-1)^{|\mu/\mu^-|} s_{\lambda^{\mu}/\mu^-} \]
\[ s_{\lambda/\mu} e_n = \sum_{\lambda^{\mu}, \mu \subseteq \lambda} (-1)^{|\mu/\mu^-|} s_{\lambda^{\mu}/\mu^-} \]

\(^91\)Note that \( \mu \subseteq \lambda \) is not required. (The left hand sides are 0 otherwise, but this does not trivialize the equalities.)
Example 2.62.

\[ s \quad \begin{array}{cccc} \ota & \ota & h & \ota \\ \ota \\ \ota \end{array} \]

\[ = \quad \begin{array}{cccc} \ota & \ota \\ \ota \\ \ota \end{array} + \begin{array}{cccc} \ota & s & \ota & \ota \\ \ota \\ \ota \end{array} + \begin{array}{cccc} \ota & \ota & \ota \end{array} + \begin{array}{cccc} \ota & \ota & \ota & \ota \\ \ota \end{array} \]

\[ = \quad \begin{array}{cccc} \ota & \ota \\ \ota \\ \ota \end{array} + \begin{array}{cccc} \ota \end{array} \]

\[ + \begin{array}{cccc} \ota & \ota \end{array} \]

\[ - \begin{array}{cccc} \ota \end{array} \]

\[ + \begin{array}{cccc} \ota \end{array} \]

Theorem 2.61 is proven in the next section, using an important Hopf algebra tool.

Exercise 2.63. Let \( \lambda = (\lambda_1, \lambda_2, \lambda_3, \ldots) \) and \( \mu = (\mu_1, \mu_2, \mu_3, \ldots) \) be two partitions such that \( \mu \subseteq \lambda \).

(a) Show that \( \lambda/\mu \) is a horizontal strip if and only if every \( i \in \{1, 2, 3, \ldots\} \) satisfies \( \mu_i \geq \lambda_i + 1 \).

(b) Show that \( \lambda/\mu \) is a vertical strip if and only if every \( i \in \{1, 2, 3, \ldots\} \) satisfies \( \lambda_i \leq \mu_i + 1 \).

Exercise 2.64. Let \( \lambda \) and \( \mu \) be two partitions such that \( \mu \subseteq \lambda \). Let \( n \in \mathbb{N} \).

(a) Show that \( (h_n, s_{\lambda/\mu}) \) equals 1 if \( \lambda/\mu \) is a horizontal \( n \)-strip, and equals 0 otherwise.

(b) Use part (a) to give a new proof of (2.36).

Exercise 2.65. Prove Theorem 2.59 again using the ideas of the proof of Theorem 2.31.

Exercise 2.66. Let \( A \) be a commutative ring, and \( n \in \mathbb{N} \).

(a) Let \( a_1, a_2, \ldots, a_n \) be \( n \) elements of \( A \). Let \( b_1, b_2, \ldots, b_n \) be \( n \) further elements of \( A \). If \( a_i - b_j \) is an invertible element of \( A \) for every \( i \in \{1, 2, \ldots, n\} \) and \( j \in \{1, 2, \ldots, n\} \), then prove that

\[
\det \left( \frac{1}{a_i - b_j} \right)_{i,j=1,2,\ldots,n} = \frac{\prod_{1 \leq j < i \leq n} ((a_i - a_j)(b_j - b_i))}{\prod_{(i,j) \in \{1,2,\ldots,n\}^2} (a_i - b_j)}. \]

(b) Let \( a_1, a_2, \ldots, a_n \) be \( n \) elements of \( A \). Let \( b_1, b_2, \ldots, b_n \) be \( n \) further elements of \( A \). If \( 1 - a_ib_j \) is an invertible element of \( A \) for every \( i \in \{1, 2, \ldots, n\} \) and \( j \in \{1, 2, \ldots, n\} \), then prove that

\[
\det \left( \frac{1}{1 - a_ib_j} \right)_{i,j=1,2,\ldots,n} = \frac{\prod_{1 \leq j < i \leq n} ((a_i - a_j)(b_i - b_j))}{\prod_{(i,j) \in \{1,2,\ldots,n\}^2} (1 - a_ib_j)}. \]

\(^{92}\)In other words, \( \lambda/\mu \) is a horizontal strip if and only if \( (\lambda_2, \lambda_3, \lambda_4, \ldots) \subseteq \mu \). This simple observation has been used by Pak and Postnikov [122, §10] for a new approach to RSK-type algorithms.
(c) Use the result of part (b) to give a new proof for Theorem 2.31.\footnote{This approach to Theorem 2.31 is taken in [36, \S 4] (except that [36] only works with finitely many variables).}

The determinant on the left hand side of Exercise 2.66(a) is known as the \textit{Cauchy determinant}.

**Exercise 2.67.** Prove that \( s_{(a,b)} = h_a h_b - h_{a+1} h_{b-1} \) for any two integers \( a \geq b \geq 0 \) (where we set \( h_{-1} = 0 \) as usual).

(Note that this is precisely the Jacobi-Trudi formula (2.22) in the case when \( \lambda = (a, b) \) is a partition with at most two entries and \( \mu = \emptyset \).)

**Exercise 2.68.** If \( \lambda \) is a partition and \( \mu \) is a weak composition, let \( K_{\lambda,\mu} \) denote the number of column-strict tableaux \( T \) of shape \( \lambda \) having \( \text{cont} (T) = \mu \). (This \( K_{\lambda,\mu} \) is called the \((\lambda, \mu)\)-\textit{Kostka number}.)

(a) Use Theorem 2.59 to show that every partition \( \mu \) satisfies \( h_\mu = \sum_\lambda K_{\lambda,\mu} s_\lambda \), where the sum ranges over all partitions \( \lambda \).

(b) Use this to give a new proof for Theorem 2.31.\footnote{Of course, this gives a new proof of Theorem 2.31 only when coupled with a proof of Theorem 2.59 which does not rely on Theorem 2.31. The proof of Theorem 2.59 we gave in the text above did not rely on Theorem 2.31, whereas the proof of (2.30) given in Exercise 2.64(b) did.}

(c) Give a new proof of the fact that \((h_\lambda)_{\lambda \in \text{Par}}\) is a \( k \)-basis of \( \Lambda \).

**Exercise 2.69.** (a) Define a \( k \)-linear map \( \mathfrak{3} : \Lambda \rightarrow \Lambda \) by having it send \( s_\lambda \) to \( s_{\lambda'} \) for every partition \( \lambda \). (This is clearly well-defined, since \((s_\lambda)_{\lambda \in \text{Par}}\) is a \( k \)-basis of \( \Lambda \).) Show that
\[
\mathfrak{3} (fh_\nu) = \mathfrak{3} (f) \cdot \mathfrak{3} (h_\nu)
\]
for every \( f \in \Lambda \) and every \( n \in \mathbb{N} \).

(b) Show that \( \mathfrak{3} = \omega \).

(c) Show that \( c_{\mu,\nu}^\lambda = c_{\mu',\nu'}^\lambda \) for any three partitions \( \lambda, \mu \) and \( \nu \).

(d) Use this to prove (2.21).\footnote{The first author learned this approach to (2.21) from Alexander Postnikov.}

**Exercise 2.70.** (a) Show that
\[
\prod_{i,j=1}^{\infty} (1 + x_i y_j) = \sum_{\lambda \in \text{Par}} s_\lambda (x) s_{\lambda'} (y) = \sum_{\lambda \in \text{Par}} e_\lambda (x) m_\lambda (y)
\]
in the power series ring \( k[[x,y]] = k[[x_1,x_2,x_3,\ldots, y_1,y_2,y_3,\ldots]] \).

(b) Assume that \( Q \) is a subring of \( k \). Show that
\[
\prod_{i,j=1}^{\infty} (1 + x_i y_j) = \sum_{\lambda \in \text{Par}} (-1)^{|\lambda| - \ell(\lambda)} z_\lambda^{-1} p_\lambda (x) p_\lambda (y)
\]
in the power series ring \( k[[x,y]] = k[[x_1,x_2,x_3,\ldots, y_1,y_2,y_3,\ldots]] \), where \( z_\lambda \) is defined as in Proposition 2.43.

The first equality of Exercise 2.70(a) appears in [154, Thm. 7.14.3], [138, Thm. 4.8.6] and several other references under the name of the \textit{dual Cauchy identity}, and is commonly proven using a “dual” analogue of the Robinson-Schensted-Knuth algorithm.

**Exercise 2.71.** Prove Theorem 2.30.\footnote{This is the proof given in Stanley [154, \S 7.16, Second Proof of Thm. 7.16.1] and Macdonald [106, proof of (5.4)].}

[\textbf{Hint:}] Switch \( x \) and \( y \) in the formula of Exercise 2.39(a), and specialize the resulting equality by replacing \( y \) by a finite set of variables \((y_1, y_2, \ldots, y_\ell)\); then, set \( n = \ell \) and \( \rho = (n-1, n-2, \ldots, 0) \), and multiply with the alternant \( a_{\rho} (y_1, y_2, \ldots, y_\ell) \), using Corollary 2.54 to simplify the result; finally, extract the coefficient of \( y^{\lambda + \rho} \).]
2.8. **Skewing and Lam’s proof of the skew Pieri rule.** We codify here the operation \( s^\perp_\mu \) of skewing by \( s_\mu \), acting on Schur functions via

\[
s^\perp_\mu(s_\lambda) = s_{\lambda/\mu}
\]

(where, as before, one defines \( s_{\lambda/\mu} = 0 \) if \( \mu \nsubseteq \lambda \)). These operations play a crucial role

- in Lam’s proof of the skew Pieri rule,
- in Lam, Lauve, and Sottile’s proof \([88]\) of a more general skew Littlewood-Richardson rule that had been conjectured by Assaf and McNamara, and
- in Zelevinsky’s structure theory of PSH’s to be developed in the next chapter.

**Definition 2.72.** Given a Hopf algebra \( A \) of finite type, and its (graded) dual \( A^\circ \), let \((\cdot, \cdot)_A\) be the pairing \((f, a) := f(a)\) for \( f \) in \( A^\circ \) and \( a \) in \( A \). Then define for each \( f \) in \( A^\circ \) an operator \( A \xrightarrow{f^\perp} A \) as follows\(^97\): for \( a \) in \( A \) with \( \Delta(a) = \sum a_1 \otimes a_2 \), let

\[
f^\perp(a) = \sum (f, a_1) a_2.
\]

Note that when one takes \( A = \Lambda = A^\circ \) \(^98\), the element \( a = s_\lambda \) has \( \Delta s_\lambda = \sum s_\mu \otimes s_{\lambda/\mu} \). Hence if \( f = s_\mu \), then one has \( f^\perp(a) = s_{\lambda/\mu} = s^\perp_\mu(s_\lambda) \) as desired.

**Proposition 2.73.** The \( f^\perp \) operators \( A \to A \) have the following properties.

(i) For every \( f \in A^\circ \), the map \( f^\perp \) is adjoint to left multiplication \( A^\circ \xrightarrow{f^\perp} A^\circ \) in the sense that

\[
(g, f^\perp(a)) = (fg, a).
\]

(ii) For every \( f, g \in A^\circ \), we have \( f^\perp g = g^\perp f^\perp \), that is, \( A \) becomes a right \( A^\circ \)-module via the \( f^\perp \) action.

(iii) If an \( f \in A^\circ \) satisfies \( \Delta(f) = \sum f_1 \otimes f_2 \), then

\[
f^\perp(ab) = \sum f_1^\perp(a) f_2^\perp(b).
\]

In particular, if \( f \) is primitive in \( A^\circ \), so that \( \Delta(f) = f \otimes 1 + 1 \otimes f \), then \( f^\perp \) is a derivation:

\[
f^\perp(ab) = f^\perp(a) \cdot b + a \cdot f^\perp(b).
\]

**Proof.** For (i), note that

\[
(g, f^\perp(a)) = \sum (f, a_1)(g, a_2) = (f \otimes g, \Delta_A(a)) = (m_{A^\circ}(f \otimes g), a) = (fg, a).
\]

For (ii), using (i) and considering any \( h \) in \( A^\circ \), one has that

\[
(h, (fg)^\perp(a)) = (fgh, a) = (gh, f^\perp(a)) = (h, g^\perp(f^\perp(a)))).
\]

For (iii), noting that

\[
\Delta(ab) = \Delta(a)\Delta(b) = \left( \sum_{(a)} a_1 \otimes a_2 \right) \left( \sum_{(b)} b_1 \otimes b_2 \right) = \sum_{(a), (b)} a_1 b_1 \otimes a_2 b_2,
\]

one has that

\[
f^\perp(ab) = \sum_{(a), (b)} (f, a_1 b_1)_A a_2 b_2 = \sum_{(a), (b)} (\Delta(f), a_1 \otimes b_1)_{A \otimes A} a_2 b_2
\]

\[
= \sum_{(f), (a), (b)} (f_1, a_1)_{A} (f_2, b_1)_A a_2 b_2
\]

\[
= \sum_{(f)} \left( \sum_{(a)} (f_1, a_1)_A a_2 \right) \left( \sum_{(b)} (f_2, b_1)_A b_2 \right) = \sum_{(f)} f_1^\perp(a) f_2^\perp(b).
\]

\(^97\)This \( f^\perp(a) \) is called \( a \xrightarrow{f} \) in Montgomery \([117, \text{Example 1.6.5}]\).

\(^98\)Here, we are using the isomorphism \( \Lambda^\circ \cong \Lambda \) induced by the Hall inner product (see Corollary 2.42 for why it is an isomorphism) to identify \( \Lambda \) with \( \Lambda^\circ \).
The following interaction between multiplication and $h^\perp$ is the key to deducing the skew Pieri formula from the usual Pieri formulas.

**Lemma 2.74.** For any $f, g$ in $\Lambda$ and any $n \in \mathbb{N}$, one has

$$f \cdot h_n^\perp(g) = \sum_{k=0}^{n} (-1)^k h_{n-k}^\perp(e_k^\perp(f) \cdot g).$$

**Proof.** Starting with the right side, first apply Proposition 2.73(iii):

$$\sum_{k=0}^{n} (-1)^k h_{n-k}^\perp(e_k^\perp(f) \cdot g)$$

$$= \sum_{k=0}^{n} (-1)^k \sum_{j=0}^{n-k} h_j^\perp(e_j^\perp(f)) \cdot h_{n-k-j}(g)$$

$$= \sum_{i=0}^{n} (-1)^{n-i} \left( \sum_{j=0}^{n-i} (-1)^j h_j^\perp(e_{n-i-j}^\perp(f)) \right) \cdot h_i^\perp(g) \quad \text{(reindexing $i := n - k - j$)}$$

$$= \sum_{i=0}^{n} (-1)^{n-i} \left( \sum_{j=0}^{n-i} (-1)^j e_{n-i-j}h_j \right) \downarrow^\perp (f) \cdot h_i^\perp(g) \quad \text{(by Proposition 2.73(ii))}$$

$$= 1^\perp(f) \cdot h_i^\perp(g) = f \cdot h_n^\perp(g)$$

where the second-to-last equality used (2.17).

**Proof of Theorem 2.61.** We prove the first skew Pieri rule; the second is analogous, swapping $h_i \leftrightarrow e_i$ and swapping the words “vertical” ↔ “horizontal”. Symmetry of $(\cdot, \cdot)_{\Lambda}$, commutativity of $\Lambda$, and Proposition 2.73(i) imply for $f$ in $\Lambda$,

$$s_{\lambda/\mu, f} = (s^\perp_{\mu}(s_{\lambda}), f) = (s_{\lambda}, s_{\mu}f) \quad \text{(2.38)}$$

Hence for any $g$ in $\Lambda$, one can compute that

$$(h_n s_{\lambda/\mu} \cdot g) \overset{\text{Prop. 2.73(i)}}{=} (s^\perp_{\mu}(s_{\lambda}) \cdot h_n^\perp(g)) \overset{\text{(2.38)}}{=} (s_{\lambda}, s_{\mu} \cdot h_n^\perp(g))$$

$$\overset{\text{Lemma 2.74}}{=} \sum_{k=0}^{n} (-1)^k (s_{\lambda}, h_{n-k}^\perp(e_k^\perp(s_{\mu}) \cdot g))$$

$$\overset{\text{Prop. 2.73(i)}}{=} \sum_{k=0}^{n} (-1)^k (h_{n-k} s_{\lambda}, e_k^\perp(s_{\mu}) \cdot g). \quad \text{(2.39)}$$

The first Pieri rule in Theorem 2.59 lets one rewrite $h_{n-k}s_{\lambda} = \sum_{\lambda^+} s_{\lambda^+}$, with the sum running through $\lambda^+$ for which $\lambda^+/\lambda$ is a horizontal $(n-k)$-strip. The second Pieri rule in Theorem 2.59 lets one rewrite $e_k^\perp s_{\mu} = \sum_{\mu^-} s_{\mu^-}$, with the sum running through $\mu^-$ for which $\mu/\mu^-$ is a vertical $k$-strip, since $(s_{\mu^-}, e_k^\perp s_{\mu}) = (e_k s_{\mu^-}, s_{\mu}).$ Thus the last line of (2.39) becomes

$$\sum_{k=0}^{n} (-1)^k \left( \sum_{\lambda^+} s_{\lambda^+} \cdot \sum_{\mu^-} s_{\mu^-} \cdot g \right) \overset{\text{(2.38)}}{=} \left( \sum_{k=0}^{n} (-1)^k \sum_{(\lambda^+, \mu^-)} s_{\lambda^+/\mu^-} \cdot g \right)$$

where the sum is over the pairs $(\lambda^+, \mu^-)$ for which $\lambda^+/\lambda$ is a horizontal $(n-k)$-strip and $\mu/\mu^-$ is a vertical $k$-strip. \hfill \Box

**Exercise 2.75.** Let $n \in \mathbb{N}$.

(a) For every $k \in \mathbb{N}$, let $p(n,k)$ denote the number of partitions of $n$ of length $k$. Let $c(n)$ denote the number of *self-conjugate* partitions of $n$ (that is, partitions $\lambda$ of $n$ satisfying $\lambda' = \lambda$). Show that

$$(-1)^n c(n) = \sum_{k=0}^{n} (-1)^k p(n, k).$$
The goal of this exercise is to prove (2.21) using the skewing operators that we have developed. Recall the involution \( \omega : \Lambda \to \Lambda \) defined in (2.19).

(a) Show that \( \omega (p_\lambda) = (-1)^{|\lambda| - \ell(\lambda)} p_\lambda \) for any \( \lambda \in \text{Par} \), where \( \ell(\lambda) \) denotes the length of the partition \( \lambda \).

(b) Show that \( \omega \) is an isometry.

(c) Show that this same map \( \omega : \Lambda \to \Lambda \) is a Hopf automorphism.

(d) Prove that \( \omega (a^T b) = (\omega (a))^\perp (\omega (b)) \) for every \( a \in \Lambda \) and \( b \in \Lambda \).

(e) For any partition \( \lambda = (\lambda_1, \ldots, \lambda_\ell) \) with length \( \ell(\lambda) = \ell \), prove that
\[
e_\ell^\perp s_\lambda = s(\lambda_1 - 1, \lambda_2 - 1, \ldots, \lambda_\ell - 1).
\]

(f) For any partition \( \lambda = (\lambda_1, \lambda_2, \ldots) \), prove that
\[
h_\lambda^\perp s_\lambda = s(\lambda_2, \lambda_3, \lambda_4, \ldots).
\]

(g) Prove (2.21).

2.9. **Assorted exercises on symmetric functions.** Over a hundred exercises on symmetric functions are collected in Stanley’s [154, chapter 7], and even more (but without any hints or references) on his website. Further sources for results related to symmetric functions are Macdonald’s work, including his monograph [106] and his expository [107]. In this section, we gather a few exercises that are not too difficult to handle with the material given above.

**Exercise 2.77.**

(a) Let \( m \in \mathbb{Z} \). Prove that, for every \( f \in \Lambda \), the infinite sum \( \sum_{i \in \mathbb{N}} (-1)^i h_{m+i} e_i^\perp f \) is convergent in the discrete topology (i.e., all but finitely many addends of this sum are zero). Hence, we can define a map \( B_m : \Lambda \to \Lambda \) by setting
\[
B_m (f) = \sum_{i \in \mathbb{N}} (-1)^i h_{m+i} e_i^\perp f \quad \text{for all } f \in \Lambda.
\]

Show that this map \( B_m \) is \( \mathbb{k} \)-linear.

(b) Let \( \lambda = (\lambda_1, \lambda_2, \lambda_3, \ldots) \) be a partition, and let \( m \in \mathbb{Z} \) be such that \( m \geq \lambda_1 \). Show that
\[
\sum_{i \in \mathbb{N}} (-1)^i h_{m+i} e_i^\perp s_\lambda = s(m, \lambda_1, \lambda_2, \lambda_3, \ldots).
\]

(c) Let \( n \in \mathbb{N} \). For every \( n \)-tuple \((\alpha_1, \alpha_2, \ldots, \alpha_n) \in \mathbb{Z}^n\), we define an element \( \bar{s}_{(\alpha_1, \alpha_2, \ldots, \alpha_n)} \in \Lambda \) by
\[
\bar{s}_{(\alpha_1, \alpha_2, \ldots, \alpha_n)} = \det \left( (h_{\alpha_i - i + j})_{i,j=1,2,\ldots,n} \right).
\]

Show that
\[
s_\lambda = \bar{s}_{(\lambda_1, \lambda_2, \ldots, \lambda_n)}
\]
for every partition \( \lambda = (\lambda_1, \lambda_2, \lambda_3, \ldots) \) having at most \( n \) parts (where “part” means “nonzero part”). Furthermore, show that for every \( n \)-tuple \((\alpha_1, \alpha_2, \ldots, \alpha_n) \in \mathbb{Z}^n\), the symmetric function \( \bar{s}_{(\alpha_1, \alpha_2, \ldots, \alpha_n)} \) either is 0 or equals \( \pm s_\nu \) for some partition \( \nu \) having at most \( n \) parts.

Finally, show that for any \( n \)-tuples \((\alpha_1, \alpha_2, \ldots, \alpha_n) \in \mathbb{Z}^n\) and \((\beta_1, \beta_2, \ldots, \beta_n) \in \mathbb{N}^n\), we have
\[
\bar{s}_{(\beta_1, \beta_2, \ldots, \beta_n)} \bar{s}_{(\alpha_1, \alpha_2, \ldots, \alpha_n)} = \det \left( (h_{\alpha_i - \beta_j - i + j})_{i,j=1,2,\ldots,n} \right).
\]

\[\text{http://math.mit.edu/~rstan/ec/ch7supp.pdf}\]
(d) For every \( n \in \mathbb{N} \), every \( m \in \mathbb{Z} \) and every \( n \)-tuple \((\alpha_1, \alpha_2, \ldots, \alpha_n) \in \mathbb{Z}^n\), prove that
\[
\sum_{i \in \mathbb{N}} (-1)^i h_{m+i} e_i(s(\alpha_1, \alpha_2, \ldots, \alpha_n)) = s(m, \alpha_1, \alpha_2, \ldots, \alpha_n),
\]
where we are using the notations of Exercise 2.77(c).

(e) For every \( n \in \mathbb{N} \) and every \( n \)-tuple \((\alpha_1, \alpha_2, \ldots, \alpha_n) \in \mathbb{Z}^n\), prove that
\[
s(\alpha_1, \alpha_2, \ldots, \alpha_n) = (B_{\alpha_1} \circ B_{\alpha_2} \circ \cdots \circ B_{\alpha_n})(1),
\]
where we are using the notations of Exercise 2.77(c) and Exercise 2.77(a).

Remark 2.78. The map \( B_n \) defined in Exercise 2.77(a) is the so-called \( m \)-th Bernstein creation operator; it appears in Zelevinsky [170, §4.20(a)] and has been introduced by J.N. Bernstein, who found the result of Exercise 2.77(b). It is called a “Schur row adder” in [54]. Exercise 2.77(e) appears in Berg/Bergeron/Saliola/Zabrocki [16, Theorem 2.3], where it is used as a prototype for defining noncommutative analogues of Schur functions, the so-called immaculate functions. The particular case of Exercise 2.77(e) for \((\alpha_1, \alpha_2, \ldots, \alpha_n)\) a partition of length \( n \) (a restatement of Exercise 2.77(b)) is proven in [106, §3.5, example 29].

Exercise 2.79.  
(a) Prove that there exists a unique family \((x_n)_{n \geq 1}\) of elements of \( \Lambda \) such that
\[
H(t) = \prod_{n=1}^{\infty} \left(1 - x_n t^n\right)^{-1}.
\]
Denote this family \((x_n)_{n \geq 1}\) by \((w_n)_{n \geq 1}\). For instance,
\[
w_1 = s(1), \quad w_2 = -s(1,1), \quad w_3 = -s(2,1),
\]
\[
w_4 = -s(1,1,1,1) - s(2,1,1) - s(2,2) - s(3,1), \quad w_5 = -s(2,1,1,1) - s(2,2,1) - s(3,1,1) - s(3,2) - s(4,1).
\]

(b) Show that \( w_n \) is homogeneous of degree \( n \) for every positive integer \( n \).

(c) For every partition \( \lambda \), define \( w_{\lambda} \in \Lambda \) by \( w_{\lambda} = w_{\lambda_1} w_{\lambda_2} \cdots w_{\lambda_{\ell(\lambda)}} \) (where \( \lambda = (\lambda_1, \lambda_2, \ldots, \lambda_{\ell}) \) with \( \ell = \ell(\lambda) \)). Notice that \( w_{\lambda} \) is homogeneous of degree \( |\lambda| \). Prove that \( \sum_{\lambda \in \mathbb{P}_{\lambda}} w_{\lambda} = h_n \) for every \( n \in \mathbb{N} \).

(d) Show that \( \{w_{\lambda}\}_{\lambda \in \mathbb{P}_{\lambda}} \) is a \( k \)-basis of \( \Lambda \). (This basis is called the Witt basis\textsuperscript{106}; it is studied in [65, §9-10,\textsuperscript{101}]
\textsuperscript{107}.)

(e) Prove that \( p_n = \sum_{d|n} d w_n/d \) for every positive integer \( n \). (Here, the summation sign \( \sum_{d|n} \) means a sum over all positive divisors \( d \) of \( n \).)

(f) We are going to show that \(-w_n\) is a sum of Schur functions (possibly with repetitions, but without signs!) for every \( n \geq 2 \). (For \( n = 1 \), the opposite is true: \( w_1 \) is a single Schur function.) This proof goes back to Doran [41]\textsuperscript{102}.

For any positive integers \( n \) and \( k \), define \( f_{n,k} \in \Lambda \) by \( f_{n,k} = \sum_{\lambda \in \mathbb{P}_{\lambda}} w_{\lambda} \), where \( \min \lambda \) denotes the smallest nonzero part of \( \lambda \). Show that
\[
-f_{n,k} = s(n-1,1) + \sum_{i=2}^{k-1} f_{i,i} f_{n-i,i} \quad \text{for every } n \geq k \geq 2.
\]
Conclude that \(-f_{n,k}\) is a sum of Schur functions for every \( n \in \mathbb{N} \) and \( k \geq 2 \). Conclude that \(-w_n\) is a sum of Schur functions for every \( n \geq 2 \).

(g) For every partition \( \lambda \), define \( r_{\lambda} \in \Lambda \) by \( r_{\lambda} = \prod_{i \geq 1} h_{v_i} (x_1^i, x_2^i, x_3^i, \ldots) \), where \( v_i \) is the number of occurrences of \( i \) in \( \lambda \). Show that \( \sum_{\lambda \in \mathbb{P}_{\lambda}} w_{\lambda}(x) r_{\lambda}(y) = \prod_{i,j=1}^{\infty} (1 - x_i y_j)^{-1} \).

(h) Show that \( \{r_{\lambda}\}_{\lambda \in \mathbb{P}_{\lambda}} \) and \( \{w_{\lambda}\}_{\lambda \in \mathbb{P}_{\lambda}} \) are dual bases of \( \Lambda \).

\textsuperscript{106}This is due to its relation with Witt vectors in the appropriate sense. Most of the work on this basis has been done by Reutenauer and Hazewinkel.

\textsuperscript{107}It also implicitly appears in [11, §5]. Indeed, the \( q_n \) of [11] are our \( w_n \) for (\( k = R \)).

\textsuperscript{108}See also Stanley [154, Exercise 7.46].
Exercise 2.80. For this exercise, set $k = \mathbb{Z}$, and consider $\Lambda = \Lambda_{\mathbb{Z}}$ as a subring of $\Lambda_{\mathbb{Q}}$. Also, consider $\Lambda \otimes_{\mathbb{Z}} \Lambda$ as a subring of $\Lambda_{\mathbb{Q}} \otimes_{\mathbb{Q}} \Lambda_{\mathbb{Q}}$.  

Recall that the family $(p_n)_{n \geq 1}$ generates the $\mathbb{Q}$-algebra $\Lambda_{\mathbb{Q}}$, but does not generate the $\mathbb{Z}$-algebra $\Lambda$.

(a) Define a $\mathbb{Q}$-linear map $Z : \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}}$ by setting

$$Z(p_\lambda) = z_\lambda p_\lambda$$

for every partition $\lambda$, where $z_\lambda$ is defined as in Proposition 2.43. 

(b) Define a $\mathbb{Q}$-algebra homomorphism $\Delta_\times : \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}} \otimes_{\mathbb{Q}} \Lambda_{\mathbb{Q}}$ by setting

$$\Delta_\times (p_n) = p_n \otimes p_n$$

for every positive integer $n$.

(c) Let $r \in \mathbb{Z}$. Define a $\mathbb{Q}$-algebra homomorphism $\epsilon_r : \Lambda_{\mathbb{Q}} \to \mathbb{Q}$ by setting

$$\epsilon_r (p_n) = r$$

for every positive integer $n$.

(d) Let $r \in \mathbb{Z}$. Define a $\mathbb{Q}$-algebra homomorphism $i_r : \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}}$ by setting

$$i_r (p_n) = rp_n$$

for every positive integer $n$.

(e) Define a $\mathbb{Q}$-linear map $Sq : \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}}$ by setting

$$Sq (p_\lambda) = p_\lambda^2$$

for every partition $\lambda$.

(f) Let $r \in \mathbb{Z}$. Define a $\mathbb{Q}$-algebra homomorphism $\Delta_r : \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}} \otimes_{\mathbb{Q}} \Lambda_{\mathbb{Q}}$ by setting

$$\Delta_r (p_n) = \sum_{i=1}^{n-1} \binom{n}{i} p_i \otimes p_{n-i} + r \otimes p_n + p_n \otimes r$$

for every positive integer $n$.

(h) Define a $\mathbb{Q}$-bilinear map $* : \Lambda_{\mathbb{Q}} \times \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}}$, which will be written in infix notation (that is, we will write $a * b$ instead of $*(a,b)$), by setting

$$p_\lambda * p_\mu = \delta_{\lambda,\mu} z_\lambda p_\lambda$$

for any partitions $\lambda$ and $\mu$.

(i) Show that $\epsilon_1 (f) = f(1)$ for every $f \in \Lambda_{\mathbb{Q}}$ (where we are using the notation $\epsilon_r$ defined in Exercise 2.80(c)).

[Hint:]

Here is how this works: We have $\Lambda_{\mathbb{Q}} \cong \mathbb{Q} \otimes_{\mathbb{Z}} \Lambda$. But fundamental properties of tensor products yield

\begin{align}(2.43)\end{align}

\begin{align}
\mathbb{Q} \otimes_{\mathbb{Z}} (\Lambda \otimes_{\mathbb{Z}} \Lambda) &\cong (\mathbb{Q} \otimes_{\mathbb{Z}} \Lambda) \otimes_{\mathbb{Q}} (\Lambda \otimes_{\mathbb{Z}} \Lambda) \\
&\cong \Lambda_{\mathbb{Q}} \otimes_{\mathbb{Q}} \Lambda_{\mathbb{Q}}
\end{align}

as $\mathbb{Q}$-algebras. But $\Lambda \otimes_{\mathbb{Z}} \Lambda$ is a free $\mathbb{Z}$-module (since $\Lambda$ is a free $\mathbb{Z}$-module), and so the canonical ring homomorphism $\Lambda \otimes_{\mathbb{Z}} \Lambda \to \mathbb{Q} \otimes_{\mathbb{Z}} (\Lambda \otimes_{\mathbb{Z}} \Lambda)$ sending every $u$ to $1_{\mathbb{Q}} \otimes_{\mathbb{Z}} u$ is injective. Composing this ring homomorphism with the $\mathbb{Q}$-algebra isomorphism of (2.43) gives an injective ring homomorphism $\Lambda \otimes_{\mathbb{Z}} \Lambda \to \Lambda_{\mathbb{Q}} \otimes_{\mathbb{Q}} \Lambda_{\mathbb{Q}}$. We use this latter homomorphism to identify $\Lambda \otimes_{\mathbb{Z}} \Lambda$ with a subring of $\Lambda_{\mathbb{Q}} \otimes_{\mathbb{Q}} \Lambda_{\mathbb{Q}}$. 

4This is well-defined, since $(p_n)_{\lambda \in \mathbb{Q} \otimes_{\mathbb{Z}} \Lambda}$ is a $\mathbb{Q}$-module basis of $\Lambda_{\mathbb{Q}}$.

5This is well-defined, since the family $(p_n)_{n \geq 1}$ generates the $\mathbb{Q}$-algebra $\Lambda_{\mathbb{Q}}$ and is algebraically independent.

6This is well-defined, since the family $(p_n)_{n \geq 1}$ generates the $\mathbb{Q}$-algebra $\Lambda_{\mathbb{Q}}$ and is algebraically independent.

7This is well-defined, since the family $(p_n)_{n \geq 1}$ generates the $\mathbb{Q}$-algebra $\Lambda_{\mathbb{Q}}$ and is algebraically independent.

8This is well-defined, since $(p_n)_{\lambda \in \mathbb{Q} \otimes_{\mathbb{Z}} \Lambda}$ is a $\mathbb{Q}$-module basis of $\Lambda_{\mathbb{Q}}$.

9This is well-defined, since the family $(p_n)_{n \geq 1}$ generates the $\mathbb{Q}$-algebra $\Lambda_{\mathbb{Q}}$ and is algebraically independent.

10But unlike $\Lambda_{\mathbb{Q}}$ with the usual coalgebra structure, it is neither graded nor a Hopf algebra.

11This is well-defined, since $(p_n)_{\lambda \in \mathbb{Q} \otimes_{\mathbb{Z}} \Lambda}$ is a $\mathbb{Q}$-module basis of $\Lambda_{\mathbb{Q}}$. 
For (b), show that, for every \( f \in A_Q \), the tensor \( \Delta_X \) is the preimage of \( f \left( (x_i y_j)_{(i,j) \in \{1,2,3,...,n\}} \right) = f(x_1y_1, x_1y_2, x_1y_3, \ldots, x_2y_1, x_2y_2, x_2y_3, \ldots) \in \mathbb{Q}[x,y] \) under the canonical injection \( A_Q \otimes \mathbb{Q} A_Q \to \mathbb{Q}[x,y] \) which maps every \( f \otimes g \) to \( f(x)g(y) \). (This requires making sure that the evaluation \( f \left( (x_i y_j)_{(i,j) \in \{1,2,3,...,n\}} \right) \) is well-defined to begin with, i.e., converges as a formal power series.)

For an alternative solution to (b), compute \( \Delta_X(h_n) \) or \( \Delta_X(e_n) \).

- For (c), compute \( \epsilon_r(e_n) \) or \( \epsilon_r(h_n) \).
- Reduce (d) to (b) and (c) using Exercise 1.13.
- Reduce (e) to (b).
- (f) is the hardest part. It is tempting to try and interpret the definition of \( \Delta_r \) as a convoluted way of saying that \( \Delta_r(f) \) is the preimage of \( f \left( (x_i + y_j)_{(i,j) \in \{1,2,3,...,n\}} \right) \) under the canonical injection \( A_Q \otimes \mathbb{Q} A_Q \to \mathbb{Q}[x,y] \) which maps every \( f \otimes g \) to \( f(x)g(y) \). However, this does not make sense since the evaluation \( f \left( (x_i + y_j)_{(i,j) \in \{1,2,3,...,n\}} \right) \) is (in general) not well-defined \(^{112}\) (and even if it was, it would fail to explain the \( r \)). So we need to get down to finitely many variables. For every \( N \in \mathbb{N} \), define a \( \mathbb{Q} \)-algebra homomorphism \( \mathcal{E}_N : A_Q \otimes \mathbb{Q} A_Q \to \mathbb{Q}[x_1, x_2, \ldots, x_N, y_1, y_2, \ldots, y_N] \) by sending each \( f \otimes g \) to \( f(x_1, x_2, \ldots, x_N)g(y_1, y_2, \ldots, y_N) \). Show that \( \Delta_N(\Lambda) \subset \mathcal{E}_N^{-1}(\bigotimes \mathbb{Z}[x_1, x_2, \ldots, x_N, y_1, y_2, \ldots, y_N]) \). This shows that, at least, the coefficients of \( \Delta_r(f) \) (in front of the \( m_{\lambda} \otimes m_{\mu} \) with \( \ell(\lambda) \leq r \) and \( \ell(\mu) \leq r \) in the \( \mathbb{Q} \)-basis \( (m_{\lambda} \otimes m_{\mu})_{\lambda,\mu \in \text{Par}, \lambda \leq r, \mu \leq r} \)) are integral for \( f \in A \). Of course, we want all coefficients. Show that \( \Delta_r = \Delta_r \ast (\Delta_Q \circ \text{id}_-\) in \( \text{Hom}(A_Q, A_Q \otimes A_Q) \) for any integers \( a \) and \( b \). This allows “moving” the \( r \). This approach to (f) was partly suggested to the first author by Richard Stanley.

- For (h), notice that Definition 3.2(b) (below) allows us to construct a bilinear form \( \langle \cdot , \cdot \rangle_{A_Q \otimes A_Q} : (A_Q \otimes A_Q) \times (A_Q \otimes A_Q) \to \mathbb{Q} \) from the Hall inner product \( \langle \cdot , \cdot \rangle : A_Q \times A_Q \to \mathbb{Q} \). Show that

\[
(a \ast b, c) \equiv (a \otimes b, \Delta_X(c))_{A_Q \otimes A_Q} \quad \text{for all } a, b, c \in A_Q,
\]

and then use (b).

\[\text{Remark 2.81.} \] The map \( \Delta_X \) defined in Exercise 2.80(b) is known as the \textit{internal comultiplication} (or \textit{Kronecker comultiplication}) on \( A_Q \). Unlike the standard comultiplication \( \Delta_{A_Q} \), it is not a graded map, but rather sends every homogeneous component \( (A_Q)_n \) into \( (A_Q)_n \otimes (A_Q)_m \). The bilinear map \( \ast \) from Exercise 2.80(h) is the so-called \textit{internal multiplication} (or \textit{Kronecker multiplication}), and is similarly not graded but rather takes \( (A_Q)_m \times (A_Q)_n \) to \( (A_Q)_m \) if \( n = m \) and to 0 otherwise.

The analogy between the two internal structures is not perfect: While we saw in Exercise 2.80(g) how the internal comultiplication yields another bialgebra structure on \( A_Q \), it is not true that the internal multiplication (combined with the usual coalgebra structure of \( A_Q \)) forms a bialgebra structure as well. What is missing is a multiplicative unity; if we would take the closure of \( A_Q \) with respect to the grading, then \( 1 + h_1 + h_2 + h_3 + \cdots \) would be such a unity.

The structure constants of the internal comultiplication on the Schur basis \( (s_\lambda)_{\lambda \in \text{Par}} \) are equal to the structure constants of the internal multiplication on the Schur basis\(^{113}\), and are commonly referred to as the \textit{Kronecker coefficients}. They are known to be nonnegative integers (this follows from Exercise 4.38(c)\(^{114}\)), but no combinatorial proof is known for their nonnegativity. Combinatorial interpretations for these coefficients akin to the Littlewood-Richardson rule have been found only in special cases (cf., e.g., [135]).

The map \( \Delta_r \) of Exercise 2.80(f) also has some classical theory behind it, relating to Chern classes of tensor products ([114], [106, §1.4, example 5]).

Parts (b), (c), (d), (e) and (f) of Exercise 2.80 are instances of a general phenomenon: Many \( \mathbb{Z} \)-algebra homomorphisms \( \Lambda \to A \) (with \( A \) a commutative ring, usually torsionfree) are easiest to define by first defining a \( \mathbb{Q} \)-algebra homomorphism \( \Lambda_Q \to A \otimes \mathbb{Q} \) and then showing that this homomorphism restricts to a \( \mathbb{Z} \)-algebra homomorphism \( \Lambda \to A \). One might ask for general criteria when this is possible; specifically, for what choices

\[^{112}\text{e.g., it involves summing infinitely many } x_1 \text{'s if } f = e_1\]

\[^{113}\text{This can be obtained, e.g., from (2.44).}\]

\[^{114}\text{Their integrality can also be easily deduced from Exercise 2.80(b).}\]
of \((b_n)_{n \geq 1} \in A^{(1,2,3,\ldots)}\) does there exist a \(\mathbb{Z}\)-algebra homomorphism \(A \to A\) sending the \(p_n\) to \(b_n\)? Such choices are called \textit{ghost-Witt vectors} in Hazewinkel [65], and we can give various equivalent conditions for a family \((b_n)_{n \geq 1}\) to be a ghost-Witt vector:

**Exercise 2.82.** Let \(A\) be a commutative ring.

For every \(n \in \{1,2,3,\ldots\}\), let \(\varphi_n : A \to A\) be a ring endomorphism of \(A\). Assume that the following properties hold:

- We have \(\varphi_n \circ \varphi_m = \varphi_{nm}\) for any two positive integers \(n\) and \(m\).
- We have \(\varphi_1 = \text{id}\).
- We have \(\varphi_p(a) \equiv a^p \mod pA\) for every \(a \in A\) and every prime number \(p\).

(For example, when \(A = \mathbb{Z}\), one can set \(\varphi_n = \text{id}\) for all \(n\); this simplifies the exercise somewhat. More generally, setting \(\varphi_n = \text{id}\) works whenever \(A\) is a binomial ring.\(^{115}\) However, the results of this exercise are at their most useful when \(A\) is a multivariate polynomial ring \(\mathbb{Z}[x_1,x_2,x_3,\ldots]\) over \(\mathbb{Z}\) and the homomorphism \(\varphi_n\) sends every \(P \in A\) to \(P(x_1^n,x_2^n,x_3^n,\ldots)\).

Let \(\mu\) denote the \textit{number-theoretic Möbius function}; this is the function \(\{1,2,3,\ldots\} \to \mathbb{Z}\) defined by

\[
\mu(m) = \begin{cases} 
0, & \text{if } m \text{ is not squarefree;} \\
(\text{-1})^{\text{number of prime factors of } m}, & \text{if } m \text{ is squarefree}
\end{cases}
\]

for every positive integer \(m\).

Let \(\phi\) denote the \textit{Euler totient function}; this is the function \(\{1,2,3,\ldots\} \to \mathbb{N}\) which sends every positive integer \(m\) to the number of elements of \(\{1,2,\ldots,m\}\) coprime to \(m\).

Let \((b_n)_{n \geq 1} \in A^{(1,2,3,\ldots)}\) be a family of elements of \(A\). Prove that the following seven assertions are equivalent:

- **Assertion C:** For every positive integer \(n\) and every prime factor \(p\) of \(n\), we have
  \[
  \varphi_p(b_{n/p}) \equiv b_n \mod p^{v_p(n)}A.
  \]
  Here, \(v_p(n)\) denotes the exponent of \(p\) in the prime factorization of \(n\).

- **Assertion D:** There exists a family \((\alpha_n)_{n \geq 1} \in A^{(1,2,3,\ldots)}\) of elements of \(A\) such that every positive integer \(n\) satisfies
  \[
  b_n = \sum_{d|n} d\alpha_d^{n/d}.
  \]

- **Assertion E:** There exists a family \((\beta_n)_{n \geq 1} \in A^{(1,2,3,\ldots)}\) of elements of \(A\) such that every positive integer \(n\) satisfies
  \[
  b_n = \sum_{d|n} d\varphi_n/d(\beta_d).
  \]

- **Assertion F:** Every positive integer \(n\) satisfies
  \[
  \sum_{d|n} \mu(d) \varphi_d(b_{n/d}) \in nA.
  \]

- **Assertion G:** Every positive integer \(n\) satisfies
  \[
  \sum_{d|n} \phi(d) \varphi_d(b_{n/d}) \in nA.
  \]

\(^{115}\) A \textit{binomial ring} is defined to be a torsionfree (as an additive group) commutative ring \(A\) which has one of the following equivalent properties:

- For every \(n \in \mathbb{N}\) and \(a \in A\), we have \(a(a - 1) \cdots (a - n + 1) \in n! \cdot A\). (That is, binomial coefficients \(\binom{a}{n}\) with \(a \in A\) and \(n \in \mathbb{N}\) are defined in \(A\).
- We have \(a^p \equiv a \mod pA\) for every \(a \in A\) and every prime number \(p\).

See [169] and the references therein for studies of these rings. It is not hard to check that \(\mathbb{Z}\) and every localization of \(\mathbb{Z}\) are binomial rings, and so is any commutative \(\mathbb{Q}\)-algebra as well as the ring

\[
\{ P \in \mathbb{Q}[X] \mid P(n) \in \mathbb{Z} \text{ for every } n \in \mathbb{Z}\}
\]

(but not the ring \(\mathbb{Z}[X]\) itself).

\(^{116}\) Here and in the following, summations of the form \(\sum_{d|n}\) range over all \textbf{positive} divisors of \(n\).
• Assertion \( H \): Every positive integer \( n \) satisfies
\[
\sum_{i=1}^{n} \varphi_{n/gcd(i,n)} \left( b_{gcd(i,n)} \right) \in nA.
\]

• Assertion \( J \): There exists a ring homomorphism \( \Lambda_\mathbb{Z} \to A \) which, for every positive integer \( n \), sends \( p_n \) to \( b_n \).

[Hint: The following identities hold for every positive integer \( n \):
\begin{align*}
(2.45) & \quad \sum_{d|n} \phi(d) = n; \\
(2.46) & \quad \sum_{d|n} \mu(d) = \delta_{n,1}; \\
(2.47) & \quad \sum_{d|n} \mu(d) \frac{n}{d} = \phi(n); \\
(2.48) & \quad \sum_{d|n} d\mu(d) \phi \left( \frac{n}{d} \right) = \mu(n).
\end{align*}

Furthermore, the following simple lemma is useful: If \( k \) is a positive integer, and if \( p \in \mathbb{Z} \), \( a \in A \) and \( b \in A \) are such that \( a \equiv b \mod p^kA \), then \( a^{p^n} \equiv b^{p^n} \mod p^{k+\ell}A \) for every \( \ell \in \mathbb{N} \).]

Remark 2.83. Much of Exercise 2.82 is folklore, but it is hard to pinpoint concrete appearances in literature. The equivalence \( C \iff D \) appears in Hesselholt [70, Lemma 1] (in slightly greater generality), where it is referred to as Dwork’s lemma and used in the construction of the Witt vector functor. This equivalence is also [65, Lemma 9.93]. The equivalence \( D \iff F \iff G \iff H \) in the case \( A = \mathbb{Z} \) is [42, Corollary on p. 10], where it is put into the context of Burnside rings and necklace counting. The equivalence \( C \iff F \) for finite families \( (b_n)_{n \in \{1,2,...,m\}} \) in lieu of \( (b_n)_{n \geq 1} \) is [154, Exercise 5.2 a]. One of the likely oldest relevant sources is Schur’s [145], which proves the equivalence \( C \iff D \iff F \) for finite families \( (b_n)_{n \in \{1,2,...,m\}} \), as well as a “finite version” of \( C \iff J \) (Schur did not have \( \Lambda \), but was working with actual power sums of roots of polynomials).

Exercise 2.84. Let \( A \) denote the ring \( \mathbb{Z} \). For every \( n \in \{1,2,3,\ldots\} \), let \( \varphi_n \) denote the identity endomorphism \( \text{id} \) of \( A \). Prove that the seven equivalent assertions \( C, D, E, F, G, H \) and \( J \) of Exercise 2.82 are satisfied for each of the following families \( (b_n)_{n \geq 1} \in \mathbb{Z}^{\{1,2,3,\ldots\}} \):

- the family \( (b_n)_{n \geq 1} = (q^n)_{n \geq 1} \), where \( q \) is a given integer.
- the family \( (b_n)_{n \geq 1} = (q^n)_{n \geq 1} \), where \( q \) is a given integer.
- the family \( (b_n)_{n \geq 1} = \binom{qn}{rn}_{n \geq 1} \), where \( r \in \mathbb{Q} \) and \( q \in \mathbb{Z} \) are given. (Here, a binomial coefficient
\[
\binom{a}{b}
\]
has to be interpreted as 0 when \( b \notin \mathbb{N} \).)
- the family \( (b_n)_{n \geq 1} = \binom{qn-1}{rn-1}_{n \geq 1} \), where \( r \in \mathbb{Z} \) and \( q \in \mathbb{Z} \) are given.

Exercise 2.85. For every \( n \in \{1,2,3,\ldots\} \), define a map \( f_n : \Lambda \to \Lambda \) by setting
\[
f_n(a) = a(x_1^n, x_2^n, x_3^n, \ldots)
\]
for every \( a \in \Lambda \).

(So what \( f_n \) does to a symmetric function is replacing all variables \( x_1, x_2, x_3, \ldots \) by their \( n \)-th powers.)

(a) Show that \( f_n : \Lambda \to \Lambda \) is a \( \mathbb{k} \)-algebra homomorphism for every \( n \in \{1,2,3,\ldots\} \).
(b) Show that \( f_n \circ f_m = f_{nm} \) for any two positive integers \( n \) and \( m \).
(c) Show that \( f_1 = \text{id} \).
(d) Prove that \( f_n : \Lambda \to \Lambda \) is a Hopf algebra homomorphism for every \( n \in \{1,2,3,\ldots\} \).
(e) Prove that \( f_2(h_m) = \sum_{i=0}^{2m} \binom{m}{i} h_i h_{2m-i} \) for every \( m \in \mathbb{N} \).
(f) Assume that \( \mathbb{k} = \mathbb{Z} \). Prove that \( f_p(a) \equiv a^p \mod p\Lambda \) for every \( a \in \Lambda \) and every prime number \( p \).
(g) Use Exercise 2.82 to obtain new solutions to parts (b), (c), (d), (e) and (f) of Exercise 2.80.
Exercise 2.86. For every $n \in \{1, 2, 3, \ldots\}$, define a $k$-algebra homomorphism $v_n : \Lambda \to \Lambda$ by

$$v_n(h_m) = \begin{cases} h_{m/n}, & \text{if } n \mid m; \\ 0, & \text{if } n \nmid m \end{cases}$$

for every positive integer $m$.

(a) Show that any positive integers $n$ and $m$ satisfy

$$v_n(p_m) = \begin{cases} np_{m/n}, & \text{if } n \mid m; \\ 0, & \text{if } n \nmid m \end{cases}.$$ 

(b) Show that any positive integers $n$ and $m$ satisfy

$$v_n(e_m) = \begin{cases} (-1)^{m-m/n} ne_{m/n}, & \text{if } n \mid m; \\ 0, & \text{if } n \nmid m \end{cases}.$$ 

(c) Prove that $v_n \circ v_m = v_{nm}$ for any two positive integers $n$ and $m$.

(d) Prove that $v_1 = \text{id}$.

(e) Prove that $v_n : \Lambda \to \Lambda$ is a Hopf algebra homomorphism for every $n \in \{1, 2, 3, \ldots\}$.

Now, consider also the maps $f_n : \Lambda \to \Lambda$ defined in Exercise 2.85. Fix a positive integer $n$.

(f) Prove that the maps $f_n : \Lambda \to \Lambda$ and $v_n : \Lambda \to \Lambda$ are adjoint with respect to the Hall inner product on $\Lambda$.

(g) Show that $v_n \circ f_n = \text{id}^n$.

(h) Prove that $f_n \circ v_m = v_m \circ f_n$ whenever $m$ is a positive integer coprime to $n$.

Finally, recall the $w_m \in \Lambda$ defined in Exercise 2.79.

(i) Show that any positive integer $m$ satisfies

$$v_n(w_m) = \begin{cases} w_{m/n}, & \text{if } n \mid m; \\ 0, & \text{if } n \nmid m \end{cases}.$$ 

The homomorphisms $v_n : \Lambda \to \Lambda$ defined in Exercise 2.86 are called the Verschiebung endomorphisms of $\Lambda$; this name comes from German, where “Verschiebung” means “shift”. This terminology, as well as that of Frobenius endomorphisms, originates in the theory of Witt vectors, and the connection between the Frobenius and Verschiebung endomorphisms of $\Lambda$ and the identically named operators on Witt vectors is elucidated in [65, Chapter 13].

Exercise 2.87. Fix $n \in \mathbb{N}$. For any $n$-tuple $w = (w_1, w_2, \ldots, w_n)$ of integers, define the descent set $\text{Des}(w)$ of $w$ to be the set $\{i \in \{1, 2, \ldots, n-1\} : w_i > w_{i+1}\}$.

(a) We say that an $n$-tuple $(w_1, w_2, \ldots, w_n)$ is Smirnov if every $i \in \{1, 2, \ldots, n-1\}$ satisfies $w_i \neq w_{i+1}$.

Fix $k \in \mathbb{N}$, and let $X_{n,k} \in k[[x]]$ denote the sum of the monomials $x_{w_1}x_{w_2} \cdots x_{w_n}$ over all Smirnov $n$-tuples $w = (w_1, w_2, \ldots, w_n) \in \{1, 2, 3, \ldots\}^n$ satisfying $|\text{Des}(w)| = k$. Prove that $X_{n,k} \in \Lambda$.

(b) For any $n$-tuple $w = (w_1, w_2, \ldots, w_n)$, define the stagnation set $\text{Stag}(w)$ of $w$ to be the set $\{i \in \{1, 2, \ldots, n-1\} : w_i = w_{i+1}\}$. (Thus, an $n$-tuple is Smirnov if and only if its stagnation set is empty.)

For any $d \in \mathbb{N}$ and $s \in \mathbb{N}$, define a power series $X_{n,d,s} \in k[[x]]$ as the sum of the monomials $x_{w_1}x_{w_2} \cdots x_{w_n}$ over all $n$-tuples $w = (w_1, w_2, \ldots, w_n) \in \{1, 2, 3, \ldots\}^n$ satisfying $|\text{Des}(w)| = d$ and $|\text{Stag}(w)| = s$. Prove that $X_{n,d,s} \in \Lambda$ for any nonnegative integers $d$ and $s$. 

\footnote{In the notations of [154, (A2.160)], the value $f_n(a)$ for an $a \in \Lambda$ can be written as a $[p_a]$ or (when $k = \mathbb{Z}$) as $p_a[a]$.}

\footnote{This is well-defined, since the family $(h_m)_{m \geq 1}$ generates the $k$-algebra $\Lambda$ and is algebraically independent.}

\footnote{which is also where most of the statements of Exercises 2.85 and 2.86 come from.
(c) Assume that $n$ is positive. For any $d \in \mathbb{N}$ and $s \in \mathbb{N}$, define three further power series $U_{n,d,s}, V_{n,d,s}$ and $W_{n,d,s}$ in $k[[x]]$ by the following formulas:

\begin{align}
U_{n,d,s} &= \sum_{w=(w_1, w_2, \ldots, w_n) \in \{1, 2, 3, \ldots\}^n; \quad |\text{Des}(w)|=d; \quad |\text{Stag}(w)|=s; \quad w_i < w_n} x_{w_1} x_{w_2} \cdots x_{w_n}; \\
V_{n,d,s} &= \sum_{w=(w_1, w_2, \ldots, w_n) \in \{1, 2, 3, \ldots\}^n; \quad |\text{Des}(w)|=d; \quad |\text{Stag}(w)|=s; \quad w_1 = w_n} x_{w_1} x_{w_2} \cdots x_{w_n}; \\
W_{n,d,s} &= \sum_{w=(w_1, w_2, \ldots, w_n) \in \{1, 2, 3, \ldots\}^n; \quad |\text{Des}(w)|=d; \quad |\text{Stag}(w)|=s; \quad w_1 > w_n} x_{w_1} x_{w_2} \cdots x_{w_n}.
\end{align}

Prove that these three power series $U_{n,d,s}, V_{n,d,s}$ and $W_{n,d,s}$ belong to $\Lambda$.

**Remark 2.88.** The function $X_{n,k}$ in Exercise 2.87(a) is a simple example ([148, Example 2.5, Theorem C.3]) of a chromatic quasisymmetric function that happens to be symmetric. See Shareshian/Wachs [148] for more general criteria for such functions to be symmetric, as well as deeper results. For example, [148, Theorem 6.3] gives an expansion for a wide class of chromatic quasisymmetric functions in the Schur basis of $\Lambda$, which, in particular, shows that our $X_{n,k}$ satisfies

$$X_{n,k} = \sum_{\lambda \in \text{Par}_n} a_{\lambda,k} s_{\lambda},$$

where $a_{\lambda,k}$ is the number of all assignments $T$ of entries in $\{1, 2, \ldots, n\}$ to the cells of the Ferrers diagram of $\lambda$ such that the following four conditions are satisfied:

- Every element of $\{1, 2, \ldots, n\}$ is used precisely once in the assignment (i.e., we have $\text{cont}(T) = (1^n)$).
- Whenever a cell $y$ of the Ferrers diagram lies immediately to the right of a cell $x$, we have $T(y) - T(x) \geq 2$.
- Whenever a cell $y$ of the Ferrers diagram lies immediately below a cell $x$, we have $T(y) - T(x) \geq -1$.
- There exist precisely $k$ elements $i \in \{1, 2, \ldots, n-1\}$ such that the cell $T^{-1}(i)$ lies in a row below $T^{-1}(i+1)$.

Are there any such rules for the $X_{n,d,s}$ of part (b)?

Smirnov $n$-tuples are more usually called Smirnov words, or (occasionally) Carlitz words.

**Exercise 2.89.**

(a) Let $n \in \mathbb{N}$. Define a matrix $A_n = (a_{i,j})_{i,j=1,2,\ldots,n}$ by

$$a_{i,j} = \begin{cases} 
p_{i-j+1}, & \text{if } i \geq j; \\
p_i, & \text{if } i = j - 1; \\
0, & \text{if } i < j - 1
\end{cases} \quad \text{for all } (i,j) \in \{1,2,\ldots,n\}^2.$$

This matrix $A_n$ looks as follows:

$$A_n = \begin{pmatrix}
p_1 & 1 & 0 & \cdots & 0 & 0 \\
p_2 & p_1 & 2 & \cdots & 0 & 0 \\
p_3 & p_2 & p_1 & \cdots & 0 & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
p_n-1 & p_{n-2} & p_{n-3} & \cdots & p_1 & n-1 \\
p_n & p_{n-1} & p_{n-2} & \cdots & p_2 & p_1
\end{pmatrix}.$$

Show that $\det(A_n) = n!e_n$.

(b) Let $n$ be a positive integer. Define a matrix $B_n = (b_{i,j})_{i,j=1,2,\ldots,n}$ by

$$b_{i,j} = \begin{cases} 
i e_i, & \text{if } j = 1; \\
e_{i-j+1}, & \text{if } j > 1
\end{cases} \quad \text{for all } (i,j) \in \{1,2,\ldots,n\}^2.$$
The matrix $B_n$ looks as follows:

$$
B_n = \begin{pmatrix}
\begin{array}{ccccccc}
\epsilon_1 & \epsilon_0 & \epsilon_1 & \cdots & \epsilon_{n+2} & \epsilon_{n+3} & \epsilon_{n+4} \\
2\epsilon_2 & \epsilon_0 & \epsilon_1 & \cdots & \epsilon_{n+4} & \epsilon_{n+5} & \epsilon_{n+6} \\
3\epsilon_3 & \epsilon_2 & \epsilon_1 & \cdots & \epsilon_{n+5} & \epsilon_{n+6} & \epsilon_{n+7} \\
& \vdots & \vdots & \ddots & \vdots & \vdots & \vdots \\
(n-1)\epsilon_{n-1} & \epsilon_{n-2} & \epsilon_{n-3} & \cdots & \epsilon_1 & \epsilon_0 \\
ne_n & \epsilon_{n-1} & \epsilon_{n-2} & \cdots & \epsilon_2 & \epsilon_1
\end{array}
\end{pmatrix}
$$

Show that $\det(B_n) = p_n$.

The formulas of Exercise 2.89, for finitely many variables, appear in Prasolov’s [127, §4.1]¹²⁰. In [127, §4.2], Prasolov gives four more formulas, which express $e_n$ as a polynomial in the $h_1, h_2, h_3, \ldots, n$ as a polynomial in the $e_1, e_2, e_3, \ldots, p_n$ as a polynomial in the $h_1, h_2, h_3, \ldots, n$ as a polynomial in the $p_1, p_2, p_3, \ldots$. These are not novel for us, since the first two of them are particular cases of Theorem 2.30, whereas the latter two can be derived from Exercise 2.89 by applying $\omega$. (Note that $\omega$ is only well-defined on symmetric functions in infinitely many indeterminates, so we need to apply $\omega$ before evaluating at finitely many indeterminates; this explains why Prasolov has to prove the latter two identities separately.)

**Exercise 2.90.** In the following, if $k \in \mathbb{N}$, we shall use the notation $1^k$ for $1, 1, \ldots, 1$ (in contexts such as $(n, 1^k)$). So, for example, $(3, 1^4)$ is the partition $(3, 1, 1, 1)$.

(a) Show that $e_n h_m = s_{(m+1, 1^{n-1})} + s_{(m, 1^n)}$ for any two positive integers $n$ and $m$.

(b) Show that

$$
\sum_{i=0}^{b} (-1)^i h_{a+i+1} e_{b-i} = s_{(a+1, 1^b)}
$$

for any $a \in \mathbb{N}$ and $b \in \mathbb{N}$.

(c) Show that

$$
\sum_{i=0}^{b} (-1)^i h_{a+i+1} e_{b-i} = (-1)^b \delta_{a+b, -1}
$$

for any negative integer $a$ and every $b \in \mathbb{N}$. (As usual, we set $h_j = 0$ for $j < 0$ here.)

(d) Show that

$$
\Delta s_{(a+1, 1^b)} = 1 \otimes s_{(a+1, 1^b)} + s_{(a+1, 1^b)} \otimes 1 \\
+ \sum_{(c,d,e,f) \in \mathbb{N}^4; \atop c+e=a-1;} s_{(c+1, 1^d)} \otimes s_{(e+1, 1^f)} + \sum_{(c,d,e,f) \in \mathbb{N}^4; \atop c+e=a;} s_{(c+1, 1^d)} \otimes s_{(e+1, 1^f)}
$$

for any $a \in \mathbb{N}$ and $b \in \mathbb{N}$.

Our next few exercises survey some results on Littlewood-Richardson coefficients.

**Exercise 2.91.** Let $m \in \mathbb{N}$ and $k \in \mathbb{N}$. Let $\lambda$ and $\mu$ be two partitions such that $\ell(\lambda) \leq k$ and $\ell(\mu) \leq k$. Assume that all parts of $\lambda$ and all parts of $\mu$ are $\leq m$ (that is, $\lambda_i \leq m$ and $\mu_i \leq m$ for every positive integer $i$)¹²¹. Let $\lambda^r$ and $\mu^r$ denote the $k$-tuples $(m - \lambda_k, m - \lambda_{k-1}, \ldots, m - \lambda_1)$ and $(m - \mu_k, m - \mu_{k-1}, \ldots, m - \mu_1)$, respectively.

---

¹²⁰ where our symmetric functions $e_k, h_k, p_k$, evaluated in finitely many indeterminates, are denoted $\sigma_k, p_k, s_k$, respectively.

¹²¹ As usual, we are denoting by $\nu_i$ the $i$-th part of a partition $\nu$ here.
(a) Show that $\lambda^\nu$ and $\mu^\nu$ are partitions, and that $s_{\lambda/\mu} = s_{\mu^\nu/\lambda^\nu}$.

(b) Show that $c_{\mu,\nu}^\lambda = c_{\lambda^\nu,\nu}^\mu$ for any partition $\nu$.

(c) Let $\nu$ be a partition such that $\ell(\nu) \leq k$, and such that all parts of $\nu$ are $\leq m$. Let $\nu^\lambda$ denote the $k$-tuple $(m - \nu_k, m - \nu_{k-1}, \ldots, m - \nu_1)$. Show that $\nu^\lambda$ is a partition, and satisfies

$$c_{\mu,\nu}^\lambda = c_{\lambda^\nu,\nu}^\mu = c_{\mu^\nu,\lambda^\nu}^\nu = c_{\lambda^\nu,\nu}^\mu.$$

(d) Show that

$$s_{\lambda^\nu}(x_1, x_2, \ldots, x_k) = (x_1 x_2 \cdots x_k)^m \cdot s_{\lambda}(x_1^{-1}, x_2^{-1}, \ldots, x_k^{-1})$$

in the Laurent polynomial ring $k[x_1, x_2, \ldots, x_k]$.

(e) Let $r$ be a nonnegative integer. Show that $(r + \lambda_1, r + \lambda_2, \ldots, r + \lambda_k)$ is a partition and satisfies

$$s(r+\lambda_1, r+\lambda_2, \ldots, r+\lambda_k)(x_1, x_2, \ldots, x_k) = (x_1 x_2 \cdots x_k)^r \cdot s_{\lambda}(x_1, x_2, \ldots, x_k)$$

in the polynomial ring $k[x_1, x_2, \ldots, x_k]$.

**Exercise 2.92.** Let $m \in \mathbb{N}$, $n \in \mathbb{N}$ and $k \in \mathbb{N}$. Let $\mu$ and $\nu$ be two partitions such that $\ell(\mu) \leq k$ and $\ell(\nu) \leq k$. Assume that all parts of $\mu$ are $\leq m$ (that is, $\mu_i \leq m$ for every positive integer $i$) and that all parts of $\nu$ are $\leq n$ (that is, $\nu_i \leq n$ for every positive integer $i$). Let $\mu^\nu(m)$ denote the $k$-tuple $(m - \mu_k, m - \mu_{k-1}, \ldots, m - \mu_1)$, and let $\nu^\lambda(n)$ denote the $k$-tuple $(n - \nu_k, n - \nu_{k-1}, \ldots, n - \nu_1)$.

(a) Show that $\mu^\nu(m)$ and $\nu^\lambda(n)$ are partitions.

Now, let $\lambda$ be a further partition such that $\ell(\lambda) \leq k$.

(b) If not all parts of $\lambda$ are $\leq m + n$, then show that $c_{\mu,\nu}^\lambda = 0$.

(c) If all parts of $\lambda$ are $\leq m + n$, then show that $c_{\mu,\nu}^\lambda = c_{\mu^\nu(m),\nu^\lambda(n)}^\mu$, where $\lambda^\nu(m+n)$ denotes the $k$-tuple $(m + n - \lambda_k, m + n - \lambda_{k-1}, \ldots, m + n - \lambda_1)$.

The results of Exercise 2.69(c) and Exercise 2.91(c) are two symmetries of Littlewood-Richardson coefficients\textsuperscript{122}; combining them yields further such symmetries. While these symmetries were relatively easy consequences of our algebraic definition of the Littlewood-Richardson coefficients, it is a much more challenging task to derive them bijectively from a combinatorial definition of these coefficients (such as the one given in Corollary 2.58). Some such derivations appear in \cite{165}, in \cite[Example 3.6, Proposition 5.11 and references therein]{15}, \cite[§5.1, §A.1, §A.4]{53} and \cite[(2.12)]{86} (though a different combinatorial interpretation of $c_{\mu,\nu}^\lambda$ is used in the latter three).

**Exercise 2.93.** Recall our usual notations: For every partition $\lambda$ and every positive integer $i$, the $i$-th entry of $\lambda$ is denoted by $\lambda_i$. The sign $\triangleright$ stands for dominance order. We let $\lambda^t$ denote the conjugate partition of a partition $\lambda$.

For any two partitions $\mu$ and $\nu$, we define two new partitions $\mu + \nu$ and $\mu \sqcup \nu$ of $|\mu| + |\nu|$ as follows:

- The partition $\mu + \nu$ is defined as $(\mu_1 + \nu_1, \mu_2 + \nu_2, \mu_3 + \nu_3, \ldots)$.
- The partition $\mu \sqcup \nu$ is defined as the result of sorting the list $(\mu_1, \mu_2, \ldots, \mu_{\ell(\mu)}, \nu_1, \nu_2, \ldots, \nu_{\ell(\nu)})$ in decreasing order.

(a) Show that any two partitions $\mu$ and $\nu$ satisfy $(\mu + \nu)^t = \mu^t + \nu^t$ and $(\mu \sqcup \nu)^t = \mu^t \sqcup \nu^t$.

(b) Show that any two partitions $\mu$ and $\nu$ satisfy $c_{\mu,\nu}^{\mu^t \sqcup \nu^t} = 1$ and $c_{\mu,\nu}^{\mu \sqcup \nu^t} = 1$.

(c) If $k \in \mathbb{N}$ and $n \in \mathbb{N}$ satisfy $k \leq n$, and if $\mu \in \text{Par}_k$, $\nu \in \text{Par}_{n-k}$ and $\lambda \in \text{Par}_n$ are such that $c_{\mu,\nu}^\lambda \neq 0$, then prove that $\mu + \nu \triangleright \lambda \triangleright \mu \sqcup \nu$.

(d) If $n \in \mathbb{N}$ and $m \in \mathbb{N}$ and $\alpha, \beta \in \text{Par}_n$ and $\gamma, \delta \in \text{Par}_m$ are such that $\alpha \triangleright \beta$ and $\gamma \triangleright \delta$, then show that $\alpha + \gamma \triangleright \beta + \delta$ and $\alpha \sqcup \gamma \triangleright \beta \sqcup \delta$.

(e) Let $m \in \mathbb{N}$ and $k \in \mathbb{N}$, and let $\lambda$ be the partition $(m^k) = \left(\frac{m, m, \ldots, m}{k \text{ times}}\right)$. Show that any two partitions $\mu$ and $\nu$ satisfy $c_{\mu,\nu}^\lambda \in \{0, 1\}$.

(f) Let $a \in \mathbb{N}$ and $b \in \mathbb{N}$, and let $\lambda$ be the partition $(a + 1, b^k)$ (using the notation of Exercise 2.90). Show that any two partitions $\mu$ and $\nu$ satisfy $c_{\mu,\nu}^\lambda \in \{0, 1\}$.

\textsuperscript{122}As usual, we are denoting by $\nu_i$ the $i$-th part of a partition $\nu$ here.

\textsuperscript{123}The result of Exercise 2.92(c) can also be regarded as a symmetry of Littlewood-Richardson coefficients; see \cite[§3.3]{10}.
Dyck words (written using 1’s and 2’s) are precisely the Yamanouchi words whose letters are 1’s and 2’s.

Exercise 2.93(g) is part of Stembridge’s [158, Thm. 2.1]; we refer to that article for further results of its kind.

The Littlewood-Richardson rule comes in many different forms, whose equivalence is not always immediate. Our version (Corollary 2.58) has the advantage of being the simplest to prove and one of the simplest to state. Other versions can be found in [154, appendix 1 to Ch. 7], Fulton’s [53, Ch. 5] and van Leeuwen’s [95]. We restrict ourselves to proving some very basic equivalences that allow us to restate parts of Corollary 2.58:

Exercise 2.94. We shall use the following notations:

- If $T$ is a column-strict tableau and $j$ is a positive integer, then we use the notation $T|_{\text{cols} \geq j}$ for the restriction of $T$ to the union of its columns $j, j+1, j+2, \ldots$. (This notation has already been used in Section 2.6.)
- If $T$ is a column-strict tableau and $S$ is a set of cells of $T$, then we write $T|_S$ for the restriction of $T$ to the set $S$ of cells.\(^{125}\)
- If $T$ is a column-strict tableau, then an NE-set of $T$ means a set $S$ of cells of $T$ such that whenever $s \in S$, every cell of $T$ which lies northeast\(^{126}\) of $s$ must also belong to $S$.
- The Semitic reading word\(^{127}\) of a column-strict tableau $T$ is the concatenation\(^{128}\) $r_1 r_2 r_3 \cdots$, where $r_j$ is the word obtained by reading the $i$-th row of $T$ from right to left.\(^{129}\)
- If $w = (w_1, w_2, \ldots, w_n)$ is a word, then a prefix of $w$ means a word of the form $(w_1, w_2, \ldots, w_i)$ for some $i \in \{0, 1, \ldots, n\}$. (In particular, both $w$ and the empty word are prefixes of $w$.)

A word $w$ over the set of positive integers is said to be Yamanouchi if for any prefix $v$ of $w$ and any positive integer $i$, there are at least as many $i$’s among the letters of $v$ as there are $(i+1)$’s among them.\(^{130}\)

Prove the following two statements:

(a) Let $\mu$ be a partition. Let $b_{i,j}$ be a nonnegative integer for every two positive integers $i$ and $j$. Assume that $b_{i,j} = 0$ for all but finitely many pairs $(i, j)$.

The following two assertions are equivalent:

- **Assertion A**: There exist a partition $\lambda$ and a column-strict tableau $T$ of shape $\lambda/\mu$ such that all $(i,j) \in \{1, 2, 3, \ldots\}^2$ satisfy

$$b_{i,j} = \text{(the number of all entries } i \text{ in the } j\text{-th row of } T) .$$

- **Assertion B**: The inequality

$$\mu_{j+1} + (b_{1,j+1} + b_{2,j+1} + \cdots + b_{i+1,j+1}) \leq \mu_j + (b_{1,j} + b_{2,j} + \cdots + b_{i,j})$$

holds for all $(i, j) \in \mathbb{N} \times \{1, 2, 3, \ldots\}$.

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\(^{124}\) A partition is called rectangular if it has the form $(m^k) = \left( m, \frac{m}{k}, m, \frac{m}{k}, \ldots, m, \frac{m}{k}, m \right)$ for some $m \in \mathbb{N}$ and $k \in \mathbb{N}$.

\(^{125}\) This restriction $T|_S$ is not necessarily a tableau of skew shape; it is just a map from $S$ to $\{1, 2, 3, \ldots\}$. The content $\text{cont}(T|_S)$ is nevertheless well-defined (in the usual way: $(\text{cont}(T|_S))_{(i)} = (T|_S)^{-1}(i)$).

\(^{126}\) A cell $(r, c)$ is said to lie northeast of a cell $(r', c')$ if and only if we have $r \leq r'$ and $c \geq c'$.

\(^{127}\) The notation comes from [95] and is a reference to the Arabic and Hebrew way of writing.

\(^{128}\) If $s_1, s_2, s_3, \ldots$ are several words (finitely or infinitely many), then the concatenation $s_1 s_2 s_3 \cdots$ is defined as the word which is obtained by starting with the empty word, then appending $s_1$ to its end, then appending $s_2$ to the end of the result, then appending $s_3$ to the end of the result, etc.

\(^{129}\) For example, the Semitic reading word of the tableau

$$\begin{array}{cccc}
3 & 4 & 4 & 5 \\
1 & 4 & 6 \\
3 & 5 \\
\end{array}$$

is 544364153.

The Semitic reading word of a tableau $T$ is what is called the reverse reading word of $T$ in [154, §A.1.3].

\(^{130}\) For instance, the words 121323132 and 1213 are Yamanouchi, while the words 132, 21 and 112132332111 are not. The Dyck words (written using 1’s and 2’s) are precisely the Yamanouchi words whose letters are 1’s and 2’s.

Yamanouchi words are often called lattice permutations.
(b) Let $\lambda$ and $\mu$ be two partitions, and let $T$ be a column-strict tableau of shape $\lambda/\mu$. Then, the following five assertions are equivalent:
- **Assertion C:** For every positive integer $j$, the weak composition $\text{cont}(T_{\text{cole} \geq j})$ is a partition.
- **Assertion D:** For every positive integers $j$ and $i$, the number of entries $i+1$ in the first $j$ rows of $T$ is $\leq$ to the number of entries $i$ in the first $j-1$ rows of $T$.
- **Assertion E:** For every NE-set $S$ of $T$, the weak composition $\text{cont}(T|_S)$ is a partition.
- **Assertion F:** The Semitic reading word of $T$ is Yamamoto-ic.
- **Assertion G:** There exists a column-strict tableau $S$ whose shape is a partition and which satisfies the following property: For any positive integers $i$ and $j$, the number of entries $i$ in the $j$-th row of $T$ equals the number of entries $j$ in the $i$-th row of $S$.

**Remark 2.95.** The equivalence of Assertions C and F in Exercise 2.94(b) is the “not-too-difficult exercise” mentioned in [157]. It yields the equivalence between our version of the Littlewood-Richardson rule (Corollary 2.58) and that in [154, A1.3.3].

**Exercise 2.96.** Let $K$ be a field. If $N \in K^{n \times n}$ is a nilpotent matrix, then the **Jordan type of $N$** is defined to be the list of the sizes of the Jordan blocks in the Jordan normal form of $N$, sorted in decreasing order. This Jordan type is a partition of $n$, and uniquely determines $N$ up to similarity (i.e., two nilpotent $n \times n$-matrices $N$ and $N'$ are similar if and only if the Jordan types of $N$ and $N'$ are equal). If $f$ is a nilpotent endomorphism of a finite-dimensional $K$-vector space $V$, then we define the **Jordan type of $f$** as the Jordan type of any matrix representing $f$ (the choice of the matrix does not matter, since the Jordan type of a matrix remains unchanged under conjugation).

(a) Let $n \in \mathbb{N}$. Let $N \in K^{n \times n}$ be a nilpotent matrix. Let $\lambda \in \text{Par}_n$. Show that the matrix $N$ has Jordan type $\lambda$ if and only if every $k \in \mathbb{N}$ satisfies
\[ \dim(\ker(N^k)) = (\lambda^t)_1 + (\lambda^t)_2 + \ldots + (\lambda^t)_k. \]
(Here, we are using the notation $\lambda^t$ for the transpose of a partition $\lambda$, and the notation $\nu_i$ for the $i$-th entry of a partition $\nu$.)

(b) Let $f$ be a nilpotent endomorphism of a finite-dimensional $K$-vector space $V$. Let $U$ be an $f$-stable $K$-vector subspace of $V$ (that is, a $K$-vector subspace of $V$ satisfying $f(U) \subset U$). Then, restricting $f$ to $U$ gives a nilpotent endomorphism $f|U$ of $U$, and the endomorphism $f$ also induces a nilpotent endomorphism $\overline{f}$ of the quotient space $V/U$. Let $\lambda, \mu$ and $\nu$ be the Jordan types of $f$, $f|U$ and $\overline{f}$, respectively. Show that $c^\lambda_{\mu,\nu} \neq 0$ (if $\mathbb{Z}$ is a subring of $K$).

**[Hint:** For (b), Exercise 2.69(c) shows that it is enough to prove that $c^\lambda_{\mu',\nu'} \neq 0$. Due to Corollary 2.58, this only requires constructing a column-strict tableau $T$ of shape $\lambda^t/\mu^t$ with $\text{cont}T = \nu^t$ which has the property that each $\text{cont}(T_{\text{cole} \geq j})$ is a partition. Construct this tableau by defining $a_{i,j} = \dim \left( (f^t)^{-1}(U) \cap \ker(f^t) \right)$ for all $(i,j) \in \mathbb{N}^2$, and requiring that the number of entries $i$ in the $j$-th row of $T$ be $a_{i,j} - a_{i,j-1} - a_{i-1,j} + a_{i-1,j-1}$ for all $(i,j) \in \{1, 2, 3, \ldots\}^2$. Use Exercise 2.94(a) to prove that this indeed defines a column-strict tableau, and Exercise 2.94(b) to verify that it satisfies the condition on $\text{cont}(T_{\text{cole} \geq j})$.]

**Remark 2.97.** Exercise 2.96 is a taste of the connections between the combinatorics of partitions and the Jordan normal form. Much more can, and has, been said. Marc van Leeuwen’s [93] is dedicated to some of these connections; in particular, our Exercise 2.96(a) is [93, Proposition 1.1], and a far stronger version of Exercise 2.96(b) appears in [93, Theorem 4.3 (2)], albeit only for the case of an infinite $K$. One can prove a converse to Exercise 2.96(b) as well: If $c^\lambda_{\mu,\nu} \neq 0$, then there exist $V$, $f$ and $U$ satisfying the premises of Exercise 2.96(b). When $K$ is a finite field, we can ask enumerative questions, such as how many $U$’s are there for given $V$, $f$, $\lambda$, $\mu$ and $\nu$; we will see a few answers in Section 4.9 (specifically, Proposition 4.58), and a more detailed treatment is given in [106, Ch. 2].

The relationship between partitions and Jordan normal forms can be exploited to provide linear-algebraic proofs of purely combinatorial facts. See [22, Sections 6 and 9] for some examples. Note that [22, Lemma

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131 The “first $j$ rows” mean the 1-st row, the 2-nd row, etc., the $j$-th row (even if some of these rows are empty).
132 This field has no relation to the ring $K$, over which our symmetric functions are defined.
133 The Jordan normal form of $N$ is well-defined even if $K$ is not algebraically closed, because $N$ is nilpotent (so the characteristic polynomial of $N$ is $X^n$).
9.10] is the statement that, under the conditions of Exercise 2.96(b), we have \( \nu \subseteq \lambda \). This is a direct consequence of Exercise 2.96(b) (since \( c_{\mu,\nu}^\lambda \neq 0 \) can happen only if \( \nu \subseteq \lambda \)).
3. Zelevinsky’s structure theory of positive self-dual Hopf algebras

Chapter 2 showed that, as a \( \mathbb{Z} \)-basis for the Hopf algebra \( \Lambda = \Lambda_\mathbb{Z} \), the Schur functions \( \{ s_\lambda \} \) have two special properties: they have the same structure constants \( c_{\lambda,\mu}^{\nu} \) for their multiplication as for their comultiplication (Corollary 2.35), and these structure constants are all nonnegative integers (Corollary 2.58). Zelevinsky [170, §2.3] isolated these two properties as crucial.

**Definition 3.1.** Say that a graded connected Hopf algebra \( A \) over \( k = \mathbb{Z} \) with a distinguished \( \mathbb{Z} \)-basis \( \{ \sigma_\lambda \} \) consisting of homogeneous elements\(^\text{134}\) is a *positive self-dual Hopf algebra* (or PSH) if it satisfies the two further axioms

- **(self-duality)** The same structure constants \( a_{\mu,\nu}^\lambda \) appear for the product \( \sigma_\mu \sigma_\nu = \sum\lambda a_{\mu,\nu}^\lambda \sigma_\lambda \) and the coproduct \( \Delta \sigma_\lambda = \sum\mu,\nu a_{\mu,\nu}^\lambda \sigma_\mu \otimes \sigma_\nu \).
- **(positivity)** The \( a_{\mu,\nu}^\lambda \) are all nonnegative (integers).

Call \( \{ \sigma_\lambda \} \) the *PSH-basis* of \( A \).

He then developed a beautiful structure theory for PSH’s, explaining how they can be uniquely expressed as tensor products of copies of PSH’s each isomorphic to \( \Lambda \) after rescaling their grading. The next few sections explain this, following his exposition closely.

3.1. **Self-duality implies polynomiality.** We begin with a property that forces a Hopf algebra to have algebra structure which is a *polynomial* algebra, specifically the symmetric algebra \( \text{Sym}(p) \), where \( p \) is the \( k \)-submodule of primitive elements.

Recall from Exercise 1.26 that for a graded connected Hopf algebra \( A = \bigoplus_{n=0}^\infty A_n \), every \( x \) in the two-sided ideal \( I := \ker \epsilon = \bigoplus_{n>0} A_n \) has the property that its comultiplication takes the form

\[
\Delta(x) = 1 \otimes x + x \otimes 1 + \Delta_+(x)
\]

where \( \Delta_+(x) \) lies in \( I \otimes I \). Recall also that the elements \( x \) for which \( \Delta_+(x) = 0 \) are called the *primitives*. Denote by \( p \) the \( k \)-submodule of primitive elements inside \( A \).

Given a PSH \( A \) (over \( k = \mathbb{Z} \)) with a PSH-basis \( \{ \sigma_\lambda \} \), we consider the bilinear form \( (\cdot, \cdot)_A : A \times A \to \mathbb{Z} \) on \( A \) that makes this basis orthonormal. Similarly, the elements \( \{ \sigma_\lambda \otimes \sigma_\mu \} \) give an orthonormal basis for an \( A \oplus A \) on \( A \times A \). The bilinear form \( (\cdot, \cdot)_A \) on the PSH \( A \) gives rise to a \( \mathbb{Z} \)-linear map \( A \to A^o \), which is easily seen to be injective and a \( \mathbb{Z} \)-algebra homomorphism. Thus we can identify \( A \) with a subalgebra of \( A^o \). When \( A \) is of finite type, this map is a Hopf algebra isomorphism, thus allowing us to identify \( A \) with \( A^o \).

This is an instance of the following notion of self-duality.

**Definition 3.2.**

(a) If \( (\cdot, \cdot) : V \times W \to k \) is a bilinear form on the product \( V \times W \) of two graded \( k \)-modules \( V = \bigoplus_{n \geq 0} V_n \) and \( W = \bigoplus_{m \geq 0} W_m \), then we say that this form \( (\cdot, \cdot) \) is *graded* if every two distinct nonnegative integers \( n \) and \( m \) satisfy \( (V_n, W_m) = 0 \) (that is, if every two homogeneous elements \( v \in V \) and \( w \in W \) having distinct degrees satisfy \( (v, w) = 0 \)).

(b) If \( (\cdot, \cdot)_V : V \times V \to k \) and \( (\cdot, \cdot)_W : W \times W \to k \) are two symmetric bilinear forms on some \( k \)-modules \( V \) and \( W \), then we can canonically define a symmetric bilinear form \( (\cdot, \cdot)_V \otimes W \) on the \( k \)-module \( V \otimes W \) by letting

\[
(v \otimes w, v' \otimes w')_V \otimes W = (v, v')_V \cdot (w, w')_W
\]

for all \( v, v' \in V \) and \( w, w' \in W \).

This new bilinear form is graded if the original two forms \((\cdot, \cdot)_V\) and \((\cdot, \cdot)_W\) were graded (presuming that \( V \) and \( W \) are graded).\(^\text{135}\)

(c) Say that a bialgebra \( A \) is *self-dual* with respect to a given symmetric bilinear form \( (\cdot, \cdot) : A \times A \to k \) if one has \( (a, m(b \otimes c))_A = (\Delta(a), b \otimes c)_{A \otimes A} \) and \( (1_A, a) = \epsilon(a) \) for \( a, b, c \) in \( A \). If \( A \) is a graded Hopf algebra of finite type, and this form \( (\cdot, \cdot) \) is graded, then this is equivalent to the \( k \)-module map \( A \to A^o \) induced by \( (\cdot, \cdot)_A \) giving a Hopf algebra homomorphism.

Thus, any PSH \( A \) is self-dual with respect to the bilinear form \( (\cdot, \cdot)_A \) that makes its PSH-basis orthonormal.

**Proposition 3.3.** Let \( A \) be a Hopf algebra over \( k = \mathbb{Z} \) or \( k = \mathbb{Q} \) which is graded, connected, and self-dual with respect to a positive definite graded\(^\text{134,135}\) bilinear form. Then:

---

\(^{134}\)not necessarily indexed by partitions

\(^{135}\)That is, \( (A_i, A_j) = 0 \) for \( i \neq j \).
(a) Within the ideal $I$, the $k$-submodule of primitives $p$ is the orthogonal complement to the $k$-submodule $I^2$.

(b) In particular, $p \cap I^2 = 0$.

(c) When $k = \mathbb{Q}$, one has $I = p \oplus I^2$.

Proof. (a) Note that $I^2 = m(I \otimes I)$. Hence an element $x$ in $I$ lies in the perpendicular space to $I^2$ if and only if one has for all $y$ in $I \otimes I$ that

$$0 = (x, m(y))_A = (\Delta(x), y)_{A \otimes A} = (\Delta_+(x), y)_{A \otimes A}$$

where the second equality uses self-duality, while the third equality uses the fact that $y$ lies in $I \otimes I$ and the form $(\cdot, \cdot)_{A \otimes A}$ makes distinct homogeneous components orthogonal. Since $y$ was arbitrary, this means $x$ is perpendicular to $I^2$ if and only if $\Delta_+(x) = 0$, that is, $x$ lies in $p$.

(b) This follows from (a), since the form $(\cdot, \cdot)_{A}$ is positive definite.

(c) This follows from (a) using some basic linear algebra\textsuperscript{136} when $A$ is of finite type (which is the only case we will ever encounter in practice). See Exercise 3.7 for the general proof. \hfill $\square$

Remark 3.4. One might wonder why we didn’t just say $I = p \oplus I^2$ even when $k = \mathbb{Z}$ in Proposition 3.3(c). However, this is false even for $A = A_2$: the second homogeneous component $(p \oplus I^2)_2$ is the index 2 sublattice of $A_2$ which is $\mathbb{Z}$-spanned by $\{p_{2,2},e_1^2\}$, containing $2e_2$, but not containing $e_2$ itself.

Already the fact that $p \cap I^2 = 0$ has a strong implication.

Lemma 3.5. A graded connected Hopf algebra $A$ over any ring $k$ having $p \cap I^2 = 0$ has algebra structure which is commutative.

Proof. The component $A_0 = k$ commutes with all of $A$. This forms the base case for an induction on $i + j$ in which one shows that any elements $x$ in $A_i$ and $y$ in $A_j$ with $i, j > 0$ will have $[x, y] := xy - yx = 0$. Since $[x, y]$ lies in $I^2$, it suffices to show that $[x, y]$ also lies in $p$:

$$\Delta[x, y] = [\Delta(x), \Delta(y)] = [1 \otimes x + x \otimes 1 + \Delta_+(x), 1 \otimes y + y \otimes 1 + \Delta_+(y)]$$

$$= [1 \otimes x + x \otimes 1, 1 \otimes y + y \otimes 1] + [1 \otimes x + x \otimes 1, \Delta_+(y)] + [\Delta_+(x), 1 \otimes y + y \otimes 1] + [\Delta_+(x), \Delta_+(y)]$$

$$= [1 \otimes x + x \otimes 1, 1 \otimes y + y \otimes 1] + [x, y] \otimes 1$$

showing that $[x, y]$ lies in $p$. Here the second-to-last equality used the inductive hypotheses: homogeneity implies that $\Delta_+(x)$ is a sum of homogeneous tensors of the form $z_1 \otimes z_2$ satisfying $\deg(z_1), \deg(z_2) < i$, so that by induction they will commute with $1 \otimes y, y \otimes 1$, thus proving that $[\Delta_+(x), 1 \otimes y + y \otimes 1] = 0$; a symmetric argument shows $[1 \otimes x + x \otimes 1, \Delta_+(y)] = 0$, and, a similar argument shows $[\Delta_+(x), \Delta_+(y)] = 0$. The last equality is an easy calculation, and was done already in (1.10). \hfill $\square$

Remark 3.6. Zelevinsky actually shows [170, Proof of A.1.3, p. 150] that the assumption of $p \cap I^2 = 0$ (along with hypotheses of unit, connit, graded, connected, and $\Delta$ being a morphism for multiplication) already implies the associativity of the multiplication in $A$! One shows by induction on $i + j + k$ that any $x, y, z$ in $A_i, A_j, A_k$ with $i, j, k > 0$ have vanishing associator $\text{assoc}(x, y, z) := x(yz) - (xy)z$. In the inductive step, one first notes that $\text{assoc}(x, y, z)$ lies in $I^2$, and then checks that $\text{assoc}(x, y, z)$ also lies in $p$, by a calculation very similar to the one above, repeatedly using the fact that $\text{assoc}(x, y, z)$ is multilinear in its three arguments.

Exercise 3.7. Prove Proposition 3.3(c) in the general case.

This leads to a general structure theorem.

Theorem 3.8. If a graded, connected Hopf algebra $A$ over a field $k$ of characteristic zero has $I = p \oplus I^2$, then the inclusion $p \hookrightarrow A$ extends to a Hopf algebra isomorphism from the symmetric algebra $\text{Sym}_k(p) \to A$. In particular, $A$ is both commutative and cocommutative.

\textsuperscript{136}Specifically, either the existence of an orthogonal projection on a subspace of a finite-dimensional inner-product space over $\mathbb{Q}$, or the fact that $\dim(W^+) = \dim V - \dim W$ for a subspace $W$ of a finite-dimensional inner-product space $V$ over $\mathbb{Q}$ can be used.
Note that these hypotheses are valid, using Proposition 3.3(c), whenever $A$ is obtained from a PSH (over $\mathbb{Z}$) by tensoring with $\mathbb{Q}$.

**Proof of Theorem 3.8.** Since Lemma 3.5 implies that $A$ is commutative, the universal property of $\text{Sym}_k(p)$ as a free commutative algebra on generators $p$ shows that the inclusion $p \hookrightarrow A$ at least extends to an algebra morphism $\text{Sym}_k(p) \xrightarrow{\varphi} A$. Since the Hopf structure on $\text{Sym}_k(p)$ makes the elements of $p$ primitive (see Example 1.21), this $\varphi$ is actually a coalgebra morphism (since $\Delta \circ \varphi = (\varphi \otimes \varphi) \circ \Delta$ and $\epsilon \circ \varphi = \epsilon$ need only to be checked on algebra generators), hence a bialgebra morphism, hence a Hopf algebra morphism (by Proposition 1.46(c)). It remains to show that $\varphi$ is surjective, and injective.

For the surjectivity of $\varphi$, note that the hypothesis $I = p \oplus I^2$ implies that the composite $p \hookrightarrow I \to I/I^2$ gives a $k$-vector space isomorphism. What follows is a standard argument to deduce that $p$ generates $A$ as a commutative graded $k$-algebra. One shows by induction on $n$ that any homogeneous element $a$ in $A_n$ lies in the $k$-subalgebra generated by $p$. The base case $n = 0$ is trivial as $a$ lies in $k$. In the inductive step where $a$ lies in $I_i$, write $a \equiv p \mod I^2$ for some $p$ in $p$. Thus $a = p + \sum b_i c_i$, where $b_i, c_i$ lie in $I$ but have strictly smaller degree, so that by induction they lie in the subalgebra generated by $p$, and hence so does $a$.

Note that the surjectivity argument did not use the assumption that $k$ has characteristic zero, but we will now use it in the injectivity argument for $\varphi$, to establish the following

(3.1) **Claim:** Every primitive element of $\text{Sym}(p)$ lies in $p = \text{Sym}^1(p)$.

Note that this claim fails in positive characteristic, e.g. if $k$ has characteristic 2 then $x^2$ lies in $\text{Sym}^2(p)$, however

$$\Delta(x^2) = 1 \otimes x^2 + 2x \otimes x + x^2 \otimes 1 = 1 \otimes x^2 + x^2 \otimes 1.$$ 

To see the claim, assume not, so that by gradedness, there must exist some primitive element $y \neq 0$ lying in some $\text{Sym}^n(p)$ with $n \geq 2$. This would mean that the composite map $f$ that follows the coproduct with a component projection

$$\text{Sym}^n(p) \xrightarrow{\Delta} \bigoplus_{i+j=n} \text{Sym}^i(p) \otimes \text{Sym}^j(p) \longrightarrow \text{Sym}^1(p) \otimes \text{Sym}^{n-1}(p)$$

has $f(y) = 0$. However, one can check on a basis that the multiplication backward $\text{Sym}^1(p) \otimes \text{Sym}^{n-1}(p) \xrightarrow{m} \text{Sym}^n(p)$ has the property that $m \circ f = n \cdot \text{id}_{\text{Sym}^n(p)}$;

$$(m \circ f)(x_1 \cdots x_n) = m \left( \sum_{j=1}^n x_j \otimes x_1 \cdots \hat{x_j} \cdots x_n \right) = n \cdot x_1 \cdots x_n$$

for $x_1, \ldots, x_n$ in $p$. Then $n \cdot y = m(f(y)) = m(0) = 0$ leads to the contradiction that $y = 0$, since $k$ has characteristic zero.

Now one can argue the injectivity of the (graded) map$^{137}$ $\varphi$ by assuming that one has a nonzero homogeneous element $u$ in $\ker(\varphi)$ of minimum degree. In particular, $\deg(u) \geq 1$. Also since $p \hookrightarrow A$, one has that $u$ is not in $\text{Sym}^1(p) = p$, and hence $u$ is not primitive by the previous Claim. Consequently $\Delta_+(u) \neq 0$, and one can find a nonzero component $u^{(i,j)}$ of $\Delta_+(u)$ lying in $\text{Sym}(p)_i \otimes \text{Sym}(p)_j$ for some $i, j > 0$. Since this forces $i, j < \deg(u)$, one has that $\varphi$ maps both $\text{Sym}(p)_i, \text{Sym}(p)_j$ injectively into $A_i, A_j$. Hence the tensor product map

$$\text{Sym}(p)_i \otimes \text{Sym}(p)_j \xrightarrow{\varphi^\otimes_\Delta} A_i \otimes A_j$$

$^{137}$The grading on $\text{Sym}(p)$ is induced from the grading on $p$, a homogeneous subspace of $I \cap A$ as it is the kernel of the graded map $I \xrightarrow{\Delta^+_p} A \otimes A$. 

is also injective\(^{138}\). This implies \((\varphi \otimes \varphi)(u^{(i,j)}) \neq 0\), giving the contradiction that
\[
0 = \Delta_A(0) = \Delta_A(\varphi(u)) = (\varphi \otimes \varphi)(\Delta^{\text{Sym}(\mathfrak{p})}(u))
\]
contains the nonzero \(A_i \otimes A_j\)-component \((\varphi \otimes \varphi)(u^{(i,j)})\).

(An alternative proof of the injectivity of \(\varphi\) proceeds as follows: By (3.1), the subspace of primitive elements of \(\text{Sym}(\mathfrak{p})\) is \(\mathfrak{p}\), and clearly \(\varphi |_{\mathfrak{p}}\) is injective. Hence, Exercise 1.54(b) (applied to the homomorphism \(\varphi\) shows that \(\varphi\) is injective.) \(\square\)

Before closing this section, we mention one nonobvious corollary of the Claim (3.1), when applied to the ring of symmetric functions \(\Lambda_\mathbb{Q}\) with \(\mathbb{Q}\)-coefficients, since Proposition 2.28 says that \(\Lambda_\mathbb{Q} = \mathbb{Q}[p_1, p_2, \ldots] = \text{Sym}(V)\) where \(V = \mathbb{Q}[p_1, p_2, \ldots]\).

**Corollary 3.9.** The subspace \(\mathfrak{p}\) of primitives in \(\Lambda_\mathbb{Q}\) is one-dimensional in each degree \(n = 1, 2, \ldots\), and spanned by \(\{p_1, p_2, \ldots\}\).

We note in passing that this corollary can also be obtained in a simpler fashion and a greater generality:

**Exercise 3.10.** Let \(k\) be any commutative ring. Show that the primitive elements of \(\Lambda\) are precisely the elements of the \(k\)-linear span of \(p_1, p_2, p_3, \ldots\).

3.2. **The decomposition theorem.** Our goal here is Zelevinsky’s theorem [170, Theorem 2.2] giving a canonical decomposition of any PSH as a tensor product into PSH’s that each have only one primitive element in their PSH-basis. For the sake of stating it, we introduce some notation:

**Definition 3.11.** Given \(A\) a PSH with PSH-basis \(\Sigma\), let \(C := \Sigma \cap \mathfrak{p}\) be the primitive elements in \(\Sigma\). For each \(\rho\) in \(C\), let \(A(\rho) \subset A\) be the \(\mathbb{Z}\)-span of
\[
\Sigma(\rho) := \{\sigma \in \Sigma: \text{ there exists } n \geq 0 \text{ with } (\sigma, \rho^n) \neq 0\}.
\]

**Theorem 3.12.** Any PSH \(A\) has a canonical tensor product decomposition\(^{139}\)
\[
A = \bigotimes_{\rho \in C} A(\rho)
\]
with \(A(\rho)\) a PSH, and \(\rho\) the only primitive element in its PSH-basis \(\Sigma(\rho)\).

Although in all the applications, \(C\) will be finite, when \(C\) is infinite one should interpret the tensor product in the theorem as the inductive limit of tensor products over finite subsets of \(C\), that is, linear combinations of basic tensors \(\bigotimes_p a_p\) in which there are only finitely many factors \(a_p \neq 1\).

The first step toward the theorem uses a certain unique factorization property.

**Lemma 3.13.** Given \(\mathcal{P}\) a set of pairwise orthogonal primitives in a PSH \(A\),
\[
(p_1 \cdots p_r, \pi_1 \cdots \pi_s) = 0
\]
for \(p_i, \pi_j\) in \(\mathcal{P}\) unless \(r = s\) and one can reindex so that \(p_1 = \pi_1\).

**Proof.** Induct on \(r\). For \(r > 0\), one has
\[
(p_1 \cdots p_r, \pi_1 \cdots \pi_s) = (p_2 \cdots p_r, \rho_1^1 (\pi_1 \cdots \pi_s))
\]
\[
= (p_2 \cdots p_r, \sum_{j=1}^s (\pi_1 \cdots \pi_{j-1} \cdot \rho_1^j (\pi_j) \cdot \pi_{j+1} \cdots \pi_s))
\]

\(^{138}\)One needs to know that for two injective maps \(V_i \xrightarrow{\varphi_i} W_i\) of \(k\)-vector spaces \(V_i, W_i\) with \(i = 1, 2\), the tensor product \(\varphi_1 \otimes \varphi_2\) is also injective. Factoring it as \(\varphi_1 \otimes \varphi_2 = (\text{id} \otimes \varphi_2) \circ (\varphi_1 \otimes \text{id})\), one sees that it suffices to show that for an injective map \(V \xrightarrow{\varphi} W\) of free \(k\)-modules, and any free \(k\)-module \(U\), the map \(V \otimes U \xrightarrow{\varphi \otimes \text{id}} W \otimes U\) is also injective. Since tensor products commute with direct sums, and \(U\) is (isomorphic to) a direct sum of copies of \(k\), this reduces to the easy-to-check case where \(U = k\).

Note that some kind of freeness or flatness hypothesis on \(U\) is needed here since, e.g. the injective \(\mathbb{Z}\)-module maps \(\mathbb{Z} \xrightarrow{\varphi_1 = (\times 2)} \mathbb{Z}\) and \(\mathbb{Z}/2\mathbb{Z} \xrightarrow{\varphi_2 = \text{id}} \mathbb{Z}/2\mathbb{Z}\) have \(\varphi_1 \otimes \varphi_2 = 0\) on \(\mathbb{Z} \otimes_{\mathbb{Z}} \mathbb{Z}/2\mathbb{Z} \cong \mathbb{Z}/2\mathbb{Z} \neq 0\).

\(^{139}\)The tensor product of two PSHs \(A_1\) and \(A_2\) with PSH-bases \(\Sigma_1\) and \(\Sigma_2\) is defined as the graded Hopf algebra \(A_1 \otimes A_2\) with PSH-basis \(\{\pi_1 \otimes \pi_2\}_{\pi_1, \pi_2 \in \Sigma_1 \times \Sigma_2}\). The tensor product of any family of PSHs is defined similarly (for the empty family, it is the graded connected Hopf algebra \(\mathbb{Z}\) with PSH-basis \(\{1\}\)).
from Proposition 2.73(iii) because \( \rho_1 \) is primitive.\(^{140}\) On the other hand, since each \( \pi_j \) is primitive, one has \( \rho_1^\perp(\pi_j) = (\rho_1, 1) \cdot \pi_j + (\rho_1, \pi_j) \cdot 1 = (\rho_1, \pi_j) \) which vanishes unless \( \rho_1 = \pi_j \). Hence \( (\rho_1 \cdots \rho_r, \pi_1 \cdots \pi_s) = 0 \) unless \( \rho_1 \in \{\pi_1, \ldots, \pi_s\} \), in which case after reindexing so that \( \pi_1 = \rho_1 \), it equals

\[
n \cdot (\rho_1, \rho_1) \cdot (\rho_2 \cdots \rho_r, \pi_2 \cdots \pi_s)
\]

if there are exactly \( n \) occurrences of \( \rho_1 \) among \( \pi_1, \ldots, \pi_s \). Now apply induction. \( \square \)

So far the positivity hypothesis for a PSH has played little role. Now we use it to introduce a certain partial order on the PSH \( A \), and then a semigroup grading.

**Definition 3.14.** For a subset \( S \) of an abelian group, let \( ZS \) (resp. \( NS \)) denote the subgroup of \( Z \)-linear combinations (resp. submonoid of \( N \)-linear combinations\(^{141}\)) of the elements of \( S \).

In a PSH \( A \) with PSH-basis \( \Sigma \), the subset \( NS \) forms a submonoid, and lets one define a partial order on \( A \) via \( a \leq b \) if \( b - a \) lies in \( NS \).

We note a few trivial properties of this partial order:

- The positivity hypothesis implies that \( NS \cdot NS \subset NS \).
- Hence multiplication by an element \( c \geq 0 \) (meaning \( c \) lies in \( NS \)) preserves the order: \( a \leq b \) implies \( ac \leq bc \) since \( (b - a)c \) lies in \( NS \).
- Thus \( 0 \leq c \leq d \) and \( 0 \leq a \leq b \) implies \( ac \leq bc \leq bd \).

This allows one to introduce a semigroup grading on \( A \).

**Definition 3.15.** Let \( N_{\text{lin}}^C \) denote the additive submonoid of \( N^C \) consisting of those \( \alpha = (\alpha_\rho)_{\rho \in C} \) with finite support.

Note that for any \( \alpha \) in \( N_{\text{lin}}^C \), one has that the product \( \prod_{\rho \in C} \rho^{\alpha_\rho} \geq 0 \). Define

\[
\Sigma(\alpha) := \{ \sigma \in \Sigma : \sigma \leq \prod_{\rho \in C} \rho^{\alpha_\rho} \},
\]

that is, the subset of \( \Sigma \) on which \( \prod_{\rho \in C} \rho^{\alpha_\rho} \) has support. Also define

\[
A(\alpha) := Z\Sigma(\alpha) \subset A.
\]

**Proposition 3.16.** The PSH \( A \) has an \( N_{\text{lin}}^C \)-semigroup-grading: one has an orthogonal direct sum decomposition

\[
A = \bigoplus_{\alpha \in N_{\text{lin}}^C} A(\alpha)
\]

for which

\[
A(\alpha)A(\beta) \subset A(\alpha + \beta)
\]

(3.2)

\[
\Delta A(\alpha) \subset \bigoplus_{\alpha = \beta + \gamma} A(\beta) \otimes A(\gamma).
\]

(3.3)

**Proof.** We will make free use of the fact that a PSH \( A \) is commutative, since it embeds in \( A \otimes_{\mathbb{Z}} \mathbb{Q} \), which is commutative by Theorem 3.8.

Note that the orthogonality \( (A(\alpha), A(\beta)) = 0 \) for \( \alpha \neq \beta \) is equivalent to the assertion that

\[
\left( \prod_{\rho \in C} \rho^{\alpha_\rho}, \prod_{\rho \in C} \rho^{\beta_\rho} \right) = 0,
\]

which follows from Lemma 3.13.

\(^{140}\)Strictly speaking, this argument needs further justification since \( A \) might not be of finite type (and if it is not, Proposition 2.73(iii) cannot be applied). It is more adequate to refer to the proof of Proposition 2.73(iii), which indeed goes through with \( \rho_1 \) taking the role of \( f \).

\(^{141}\)Recall that \( \mathbb{N} := \{0, 1, 2, \ldots\} \).
Next let us deal with the assertion (3.2). It suffices to check that when \( \tau, \omega \) in \( \Sigma \) lie in \( A_{(\alpha)}, A_{(\beta)} \), respectively, then \( \tau \omega \) lies in \( A_{(\alpha+\beta)} \). But note that any \( \sigma \) in \( \Sigma \) having \( \sigma \leq \tau \omega \) will then have

\[
\sigma \leq \tau \omega \leq \prod_{\rho \in \mathcal{C}} \rho^{a_{\rho}} \cdot \prod_{\rho \in \mathcal{C}} \rho^{b_{\rho}} = \prod_{\rho \in \mathcal{C}} \rho^{a_{\rho}+b_{\rho}}
\]

so that \( \sigma \) lies in \( A_{(\alpha+\beta)} \). This means that \( \tau \omega \) lies in \( A_{(\alpha+\beta)} \).

This lets us check that \( \bigoplus_{\alpha \in \mathbb{N}^k} A_{(\alpha)} \) exhaust \( A \). It suffices to check that any \( \sigma \) in \( \Sigma \) lies in some \( A_{(\alpha)} \). Proceed by induction on \( \deg(\sigma) \), with the case \( \sigma = 1 \) being trivial; the element 1 always lies in \( \Sigma \), and hence lies in \( A_{(\alpha)} \) for \( \alpha = 0 \). For \( \sigma \) lying in \( I \), one either has \( \langle \sigma, a \rangle \neq 0 \) for some \( a \) in \( I^2 \), or else \( \sigma \) lies in \( (I^2)^{\perp} = \mathfrak{p} \) (by Proposition 3.3(a)), so that \( \sigma \) is in \( \mathcal{C} \) and we are done. If \( \langle \sigma, a \rangle \neq 0 \) with \( a \) in \( I^2 \), then \( \sigma \) appears in the support of some \( \mathbb{Z} \)-linear combination of elements \( \tau \omega \) where \( \tau, \omega \) lie in \( \Sigma \) and have strictly smaller degree than \( \sigma \) has. There exists at least one such pair \( \tau, \omega \) for which \( \langle \sigma, \tau \omega \rangle \neq 0 \), and therefore \( \sigma \leq \tau \omega \). Then by induction \( \tau, \omega \) lie in some \( A_{(\alpha)}, A_{(\beta)} \), respectively, so \( \tau \omega \) lies in \( A_{(\alpha+\beta)} \), and hence \( \sigma \) lies in \( A_{(\alpha+\beta)} \) also.

Self-duality shows that (3.2) implies (3.3): if \( a, b, c \) lie in \( A_{(\alpha)}, A_{(\beta)}, A_{(\gamma)} \), respectively, then \( \langle \Delta a, b \otimes c \rangle_{A \otimes A} = \langle a, bc \rangle_A = 0 \) unless \( \alpha = \beta + \gamma \). \( \square \)

**Proposition 3.17.** For \( \alpha, \beta \in \mathbb{N}_{\text{fin}}^C \) with disjoint support, one has a bijection

\[
\Sigma(\alpha) \times \Sigma(\beta) \longrightarrow \Sigma(\alpha + \beta)
\]

\[
(\sigma, \tau) \longmapsto \sigma \tau.
\]

Thus, the multiplication map \( A_{(\alpha)} \otimes A_{(\beta)} \rightarrow A_{(\alpha+\beta)} \) is an isomorphism.

**Proof.** We first check that for \( \sigma_1, \sigma_2 \) in \( \Sigma(\alpha) \) and \( \tau_1, \tau_2 \) in \( \Sigma(\beta) \), one has

(3.4)

\[
(\sigma_1 \tau_1, \sigma_2 \tau_2) = \delta(\sigma_1, \tau_1, 2\tau_2).
\]

Note that this is equivalent to showing both

- that \( \sigma \tau \) lie in \( \Sigma(\alpha + \beta) \) so that the map is well-defined, since it shows \( (\sigma \tau, \sigma \tau) = 1 \), and
- that the map is injective.

One calculates

\[
(\sigma_1 \tau_1, \sigma_2 \tau_2)_A = (\sigma_1 \tau_1, m(\sigma_2 \otimes \tau_2))_A
\]

\[
= (\Delta(\sigma_1 \tau_1), \sigma_2 \otimes \tau_2)_{A \otimes A}
\]

\[
= (\Delta(\sigma_1) \Delta(\tau_1), \sigma_2 \otimes \tau_2)_{A \otimes A}
\]

Note that due to (3.3), \( \Delta(\sigma_1) \Delta(\tau_1) \) lies in \( \sum A_{(\alpha' + \beta')} \otimes A_{(\alpha'' + \beta'')} \) where

\[
\alpha' + \alpha'' = \alpha
\]

\[
\beta' + \beta'' = \beta.
\]

Since \( \sigma_2 \otimes \tau_2 \) lies in \( A_{(\alpha)} \otimes A_{(\beta)} \), the only nonvanishing terms in the inner product come from those with

\[
\alpha' + \beta' = \alpha
\]

\[
\alpha'' + \beta'' = \beta.
\]

As \( \alpha, \beta \) have disjoint support, this can only happen if

\[
\alpha' = \alpha, \ \alpha'' = 0, \ \beta' = 0, \ \beta'' = \beta,
\]

that is, the only nonvanishing term comes from \( (\sigma_1 \otimes 1)(1 \otimes \tau_1) = \sigma_1 \otimes \tau_1 \). Hence

\[
(\sigma_1 \tau_1, \sigma_2 \tau_2)_A = (\sigma_1 \otimes \tau_1, \sigma_2 \otimes \tau_2)_{A \otimes A} = \delta(\sigma_1, \tau_1, (\sigma_2, \tau_2)).
\]

To see that the map is surjective, express

\[
\prod_{\rho \in \mathcal{C}} \rho^{a_{\rho}} = \sum_i \sigma_i
\]

\[
\prod_{\rho \in \mathcal{C}} \rho^{b_{\rho}} = \sum_j \tau_j
\]
with \( \sigma_i \in \Sigma(\alpha) \) and \( \tau_j \in \Sigma(\beta) \). Then each product \( \sigma_i \tau_j \) is in \( \Sigma(\alpha + \beta) \) by (3.4), and
\[
\prod_{\rho \in \mathcal{C}} \rho^{\sigma_i \tau_j} = \sum_{i,j} \sigma_i \tau_j
\]
shows that \( \{\sigma_i \tau_j\} \) exhausts \( \Sigma(\alpha + \beta) \). This gives surjectivity. \( \square \)

**Proof of Theorem 3.12.** Recall from Definition 3.11 that for each \( \rho \) in \( \mathcal{C} \), one defines \( A(\rho) \subset A \) to be the \( \mathbb{Z} \)-span of
\[
\Sigma(\rho) := \{ \sigma \in \Sigma : \text{ there exists } n \geq 0 \text{ with } (\sigma, \rho^n) \neq 0 \}.
\]
In other words, \( A(\rho) := \bigoplus_{n \geq 0} A_{n,\rho} \) where \( e_\rho \) in \( \mathbb{N}_0^{\mathbb{N}_0} \) is the standard basis element indexed by \( \rho \). Proposition 3.16 then shows that \( A(\rho) \) is a Hopf subalgebra of \( A \). Since every \( \alpha \) in \( \mathbb{N}_0^{\mathbb{N}_0} \) can be expressed as the (finite) sum \( \sum_\rho \alpha_\rho e_\rho \), and the \( e_\rho \) have disjoint support, iterating Proposition 3.17 shows that \( A = \bigotimes_{\rho \in \mathcal{C}} A(\rho) \). Lastly, \( \Sigma(\rho) \) is clearly a PSH-basis for \( A(\rho) \), and if \( \sigma \) is any primitive element in \( \Sigma(\rho) \) then \( (\sigma, \rho^n) \neq 0 \) lets one conclude via Lemma 3.13 that \( \sigma = \rho \) (and \( n = 1 \)). \( \square \)

3.3. \( A \) is the unique indecomposable PSH. The goal here is to prove the rest of Zelevinsky’s structure theory for PSH’s. Namely, if \( A \) has only one primitive element \( \rho \) in its PSH-basis \( \Sigma \), then \( A \) must be isomorphic as a PSH to the ring of symmetric functions \( \Lambda \), after one rescales the grading of \( \rho \). Note that every \( \sigma \) in \( \Sigma \) has \( \sigma \leq \rho^m \) for some \( m \), and hence has degree divisible by the degree of \( \rho \). Thus one can divide all degrees by that of \( \rho \) and assume \( \rho \) has degree 1.

The idea is to find within \( A \) and \( \Sigma \) a set of elements that play the role of
\[
\{ h_n = s_{(n)} \}_{n=0,1,2,...}, \{ e_n = s_{(1^n)} \}_{n=0,1,2,...}
\]
within \( A = \Lambda \) and its PSH-basis of Schur functions \( \Sigma = \{ s_\lambda \} \). Zelevinsky’s argument does this by isolating some properties that turn out to characterize these elements:

(a) \( h_0 = e_0 = 1 \), and \( h_1 = e_1 =: \rho \) has \( \rho^2 \) a sum of two elements of \( \Sigma \), namely
\[
\rho^2 = h_2 + e_2.
\]
(b) For all \( n = 0,1,2,... \), there exist unique elements \( h_n, e_n \) in \( A_n \cap \Sigma \) that satisfy
\[
h_k^\perp e_n = 0, \quad e_k^\perp h_n = 0
\]
with \( h_2, e_2 \) being the two elements of \( \Sigma \) introduced in (a).

(c) For \( k = 0,1,2,..., n \) one has
\[
h_k^\perp h_n = h_{n-k} \quad \text{and} \quad \sigma_i^\perp h_n = 0 \quad \text{for} \quad \sigma \in \Sigma \setminus \{h_0, h_1, \ldots, h_n\}
\]
\[
e_k^\perp e_n = e_{n-k} \quad \text{and} \quad \sigma^\perp e_n = 0 \quad \text{for} \quad \sigma \in \Sigma \setminus \{e_0, e_1, \ldots, e_n\}.
\]
In particular, \( e_k^\perp e_n = 0 = h_k^\perp e_n \) for \( k \geq 2 \).

(d) Their coproducts are
\[
\Delta(h_n) = \sum_{i+j=n} h_i \otimes h_j,
\]
\[
\Delta(e_n) = \sum_{i+j=n} e_i \otimes e_j.
\]

We will prove Zelevinsky’s result \([170, \text{ Theorem } 3.1]\) as a combination of the following two theorems.

**Theorem 3.18.** Let \( A \) be a PSH with PSH-basis \( \Sigma \) containing only one primitive \( \rho \), and assume that the grading has been rescaled so that \( \rho \) has degree 1. Then, after renaming \( \rho = e_1 = h_1 \), one can find unique sequences \( \{h_n\}_{n=0,1,2,...}, \{e_n\}_{n=0,1,2,...} \) of elements of \( \Sigma \) having properties (a),(b),(c),(d) listed above.

The second theorem uses the following notion.
**Definition 3.19.** A PSH-morphism \( A \xrightarrow{\varphi} A' \) between two PSH's \( A, A' \) having PSH-bases \( \Sigma, \Sigma' \) is a graded Hopf algebra morphism for which \( \varphi(\Sigma \subseteq \mathbb{N} \Sigma) \subseteq \mathbb{N} \Sigma' \). If \( A = A' \) and \( \Sigma = \Sigma' \) it will be called a PSH-endomorphism. If \( \varphi \) is an isomorphism and restricts to a bijection \( \Sigma \rightarrow \Sigma' \), it will be called a PSH-isomorphism; if it is both a PSH-isomorphism and an endomorphism, it is a PSH-automorphism.\(^{142}\)

**Theorem 3.20.** The elements \( \{h_n\}_{n=0,1,2,\ldots}, \{e_n\}_{n=0,1,2,\ldots} \) in Theorem 3.18 also satisfy the following.

\( (e) \) The elements \( h_n, e_n \) in \( A \) satisfy the same relation (2.17)

\[ \sum_{i+j=n} (-1)^i e_i h_j = \delta_{0,n} \]

as their counterparts in \( \Lambda \), along with the property that

\[ A = \mathbb{Z}[h_1, h_2, \ldots] = \mathbb{Z}[e_1, e_2, \ldots]. \]

\( (f) \) There is exactly one nontrivial automorphism \( A \xrightarrow{\rho} A \) as a PSH, swapping \( h_n \leftrightarrow e_n \).

\( (g) \) There are exactly two PSH-isomorphisms \( A \rightarrow \Lambda \),

- one sending \( h_n \) to the complete homogeneous symmetric functions \( h_n(x) \), while sending \( e_n \) to the elementary symmetric functions \( e_n(x) \),
- the second one (obtained by composing the first with \( \omega \)) sending \( h_n \mapsto e_n(x) \) and \( e_n \mapsto h_n(x) \).

Before embarking on the proof, we mention one more bit of convenient terminology: say that an element \( \sigma \) in \( \Sigma \) is a constituent of \( a \) in \( \mathbb{N} \Sigma \) when \( \sigma \leq a \), that is, \( \sigma \) appears with nonzero coefficient \( c_\sigma \) in the unique expansion \( a = \sum_{\tau \in \Sigma} c_\tau \tau \).

**Proof of Theorem 3.18.** One fact that occurs frequently is this:

\[ (3.5) \quad \text{Every } \sigma \text{ in } \Sigma \cap A_n \text{ is a constituent of } \rho^n. \]

This follows from Theorem 3.12, since \( \rho \) is the only primitive element of \( \Sigma \): one has \( A = A(\rho) \) and \( \Sigma = \Sigma(\rho) \), so that \( \sigma \) is a constituent of some \( \rho^n \), and homogeneity considerations force \( m = n \).

Notice that \( A \) is of finite type (due to (3.5)). Thus, \( A^\circ \) is a graded Hopf algebra isomorphic to \( A \).

**Assertion (a).** Note that

\[ (\rho^2, \rho^2) = (\rho^\perp(\rho^2), \rho) = (2\rho, \rho) = 2 \]

using the fact that \( \rho^\perp \) is a derivation since \( \rho \) is primitive (Proposition 2.73(iii)). On the other hand, expressing \( \rho^2 = \sum_{\sigma \in \Sigma} c_\sigma \sigma \) with \( c_\sigma \) in \( \mathbb{N} \), one has \( (\rho^2, \rho^2) = \sum_\sigma c_\sigma^2 \). Hence exactly two of the \( c_\sigma = 1 \), so \( \rho^2 \) has exactly two distinct constituents. Denote them by \( h_2 \) and \( e_2 \). One concludes that \( \Sigma \cap A_2 = \{h_2, e_2\} \) from (3.5).

Note also that the same argument shows \( \Sigma \cap A_1 = \{\rho\} \), so that \( A_1 = \mathbb{Z} \rho \). Since \( \rho^\perp h_2 \) lies in \( A_1 = \mathbb{Z} \rho \) and \( (\rho^\perp h_2, \rho) = (h_2, \rho^2) = 1 \), we have \( \rho^\perp h_2 = \rho \). Similarly \( \rho^\perp e_2 = \rho \).

**Assertion (b).** We will show via induction on \( n \) the following three assertions for \( n \geq 1 \):

- There exists an element \( h_n \) in \( \Sigma \cap A_n \) with \( c_2^\perp h_n = 0 \).

\[ (3.6) \quad \text{This element } h_n \text{ is unique.} \]

- Furthermore \( \rho^\perp h_n = h_{n-1} \).

In the base cases \( n = 1, 2 \), it is not hard to check that our previously labelled elements, \( h_1, h_2 \) (namely \( h_1 := \rho \), and \( h_2 \) as named in part (a)) really are the unique elements satisfying these hypotheses.

\[^{142}\]The reader should be warned that not every invertible PSH-endomorphism is necessarily a PSH-automorphism. For instance, it is an easy exercise to check that \( \Lambda \otimes \Lambda \rightarrow \Lambda \otimes \Lambda \), \( f \otimes g \mapsto \sum_{(f)} f_1 \otimes f_2 g \) is a well-defined invertible PSH-endomorphism of the PSH \( \Lambda \otimes \Lambda \) with PSH-basis \( (s_\lambda \otimes s_\mu)_{(\lambda, \mu) \in \text{Par} \times \text{Par}} \), but not a PSH-automorphism.
In the inductive step, it turns out that we will find $h_n$ as a constituent of $\rho h_{n-1}$. Thus we again use the derivative property of $\rho^\perp$ to compute that $\rho h_{n-1}$ has exactly two constituents:

$$\begin{align*}
(\rho h_{n-1}, \rho h_{n-1}) &= (\rho^\perp(\rho h_{n-1}), h_{n-1}) \\
&= (h_{n-1} + \rho \cdot \rho^\perp h_{n-1}, h_{n-1}) \\
&= (h_{n-1} + \rho h_{n-2}, h_{n-1}) \\
&= 1 + (h_{n-2}, \rho^\perp h_{n-1}) \\
&= 1 + (h_{n-2}, h_{n-2}) = 1 + 1 = 2
\end{align*}$$

where the inductive hypothesis $\rho^\perp h_{n-1} = h_{n-2}$ was used twice. We next show that exactly one of the two constituents of $\rho h_{n-1}$ is annihilated by $e_2^\perp$. Note that since $e_2$ lies in $A_2$, and $A_1$ has $\mathbb{Z}$-basis element $\rho$, there is a constant $c$ in $\mathbb{Z}$ such that

$$\Delta(e_2) = e_2 \otimes 1 + c \rho \otimes \rho + 1 \otimes e_2.$$ 

On the other hand, (a) showed

$$1 = (e_2, \rho^2)_A = (\Delta(e_2), \rho \otimes \rho)_{A \otimes A}$$

so one must have $c = 1$. Therefore by Proposition 2.73(iii) again,

$$\begin{align*}
e_2^\perp(\rho h_{n-1}) &= e_2^\perp(\rho) h_{n-1} + \rho^\perp(\rho) \rho^\perp(h_{n-1}) + \rho e_2^\perp(h_{n-1}) \\
&= 0 + h_{n-2} + 0 \\
&= h_{n-2}.
\end{align*}$$

where the first term vanished due to degree considerations and the last term vanished by the inductive hypothesis. Bearing in mind that $\rho h_{n-1}$ lies in $\mathbb{N}\Sigma$, and in a PSH with PSH-basis $\Sigma$, any skewing operator $\sigma^\perp$ for $\sigma$ in $\Sigma$ will preserve $\mathbb{N}\Sigma$, one concludes from $(3.8)$ that

- one of the two distinct constituents of the element $\rho h_{n-1}$ must be sent by $e_2^\perp$ to $h_{n-2}$, and
- the other constituent of $\rho h_{n-1}$ must be annihilated by $e_2^\perp$; call this second constituent $h_n$.

Lastly, to see that this $h_n$ is unique, it suffices to show that any element $\sigma$ of $\Sigma \cap A_1$ which is killed by $e_2^\perp$ must be a constituent of $\rho h_{n-1}$. This holds for the following reason. We know $\sigma \leq \rho^n$ by $(3.5)$, and hence $0 \neq (\rho^n, \sigma) = (\rho^{n-1}, \rho^\perp \sigma)$, implying that $\rho^\perp \sigma \neq 0$. On the other hand, since $0 = \rho^\perp e_2^\perp \sigma = e_2^\perp \rho^\perp \sigma$, one has that $\rho^\perp \sigma$ is annihilated by $e_2^\perp$, and hence $\rho^\perp \sigma$ must be a (positive) multiple of $h_{n-1}$ by part of our inductive hypothesis. Therefore $(\sigma, \rho h_{n-1}) = (\rho^\perp \sigma, h_{n-1})$ is positive, that is, $\sigma$ is a constituent of $\rho h_{n-1}$.

The preceding argument, applied to $\sigma = h_n$, shows that $\rho^\perp h_n = ch_{n-1}$ for some $c$ in $\{1, 2, \ldots\}$. Since $(\rho^\perp h_n, h_{n-1}) = (h_n, \rho h_{n-1}) = 1$, this $c$ must be 1, so that $\rho^\perp h_n = h_{n-1}$. This completes the induction step in the proof of $(3.6)$.

One can then argue, swapping the roles of $e_n, h_n$ in the above argument, the existence and uniqueness of a sequence $\{e_n\}_{n=0}^\infty$ in $\Sigma$ satisfying the properties analogous to $(3.6)$, with $e_0 := 1$, $e_1 := \rho$.

**Assertion (c)**. Iterating the property from (b) that $\rho^\perp h_n = h_{n-1}$ shows that $(\rho^k)^\perp h_n = h_{n-k}$ for $0 \leq k \leq n$. However one also has an expansion

$$\rho^k = ch_k + \sum_{\sigma \in \Sigma \cap A_k \setminus \{h_k\}} c_\sigma \sigma$$

for some integers $c, c_\sigma > 0$, since every $\sigma$ in $\Sigma \cap A_k$ is a constituent of $\rho^k$. Hence

$$1 = (h_{n-k}, h_{n-k}) = ((\rho^k)^\perp h_n, (\rho^k)^\perp h_{n-k}) \geq c^2(h_k^\perp h_n, h_k^\perp h_{n-k})$$

using Proposition 2.73(ii). Hence if we knew that $h_k^\perp h_n \neq 0$ this would force

$$h_k^\perp h_n = (\rho^k)^\perp h_n = h_{n-k}$$

as well as $\sigma^\perp h_n = 0$ for all $\sigma \notin \{h_0, h_1, \ldots, h_n\}$. But

$$(\rho^{n-k})^\perp h_k^\perp h_n = h_k^\perp (\rho^{n-k})^\perp h_n = h_k^\perp h_k = 1 \neq 0$$

so $h_k^\perp h_n \neq 0$, as desired. The argument for $e_k^\perp e_n = e_{n-k}$ is symmetric.

The last assertion in (c) follows if one checks that $e_n \neq h_n$ for each $n \geq 2$, but this holds since $e_2^\perp(h_n) = 0$ but $e_2^\perp(e_n) = e_{n-2}$. 


Assertion (d). Part (c) implies that
\[(\Delta h_n, \sigma \otimes \tau)_{A \otimes A} = (h_n, \sigma \tau)_A = (\sigma^\perp h_n, \tau)_A = 0\]
unless \(\sigma = h_k\) for some \(k = 0, 1, 2, \ldots, n\) and \(\tau = h_{n-k}\). Also one can compute
\[(\Delta h_n, h_k \otimes h_{n-k}) = (h_n, h_k h_{n-k}) = (h_k^\perp h_n, h_{n-k}) \underbrace{= (h_{n-k}, h_{n-k})}_{(c)} = 1.\]
This is equivalent to the assertion for \(\Delta h_n\) in (d). The argument for \(\Delta e_n\) is symmetric. \(\square\)

Before proving Theorem 3.20, we note some consequences of Theorem 3.18. Define for each partition \(\lambda = (\lambda_1 \geq \lambda_2 \geq \cdots \geq \lambda_\ell)\) the elements of \(A\)
\[h_\lambda = h_{\lambda_1} h_{\lambda_2} \cdots,\]
\[e_\lambda = e_{\lambda_1} e_{\lambda_2} \cdots.\]
Also, define the lexicographic order on \(\text{Par}_n\) by saying \(\lambda <_{\text{lex}} \mu\) if \(\lambda \neq \mu\) and the smallest index \(i\) for which \(\lambda_i \neq \mu_i\) has \(\lambda_i < \mu_i\). Recall also that \(\lambda^t\) denotes the conjugate or transpose partition to \(\lambda\), obtained by swapping rows and columns in the Ferrers diagram.

The following unitriangularity lemma will play a role in the proof of Theorem 3.20(e).

**Lemma 3.21.** Under the hypotheses of Theorem 3.18, for \(\lambda, \mu\) in \(\text{Par}_n\), one has
\[e_\mu^\perp h_\lambda = \begin{cases} 1 & \text{if } \mu = \lambda^t \\ 0 & \text{if } \mu >_{\text{lex}} \lambda^t. \end{cases}\]
Consequently
\[\det [(e_\mu^\perp, h_\lambda)]_{\lambda, \mu \in \text{Par}_n} = \pm 1.\]

**Proof.** Notice that \(A\) is of finite type (as shown in the proof of Theorem 3.18). Thus, \(A^o\) is a graded Hopf algebra isomorphic to \(A\).

Induct on the length of \(\mu\). If \(\lambda\) has length \(\ell\), so that \(\lambda_\ell^t = \ell\), then
\begin{align*}
e_\mu^\perp h_\lambda &= e_\mu^\perp \left( e_{(\mu_2, \mu_3, \ldots)}^\perp (h_{(\lambda_1, \ldots, \lambda_\ell)}) \right) \\
&= e_{(\mu_2, \mu_3, \ldots)}^\perp \sum_{i_1 + \cdots + i_\ell = \mu_1} e_{i_1}^\perp (h_{(\lambda_1, \ldots, \lambda_\ell)}) \\
&= \begin{cases} 0 & \text{if } \mu_1 > \ell = \lambda_1^t \\ e_{(\mu_2, \mu_3, \ldots)}^\perp h_{(\lambda_1 - 1, \ldots, \lambda_\ell - 1)} & \text{if } \mu_1 = \ell = \lambda_1^t \end{cases}
\end{align*}
where the second-to-last equality used Proposition 2.73(iii) along with (the iterates of) the coproduct formula for \(\Delta(e_n)\) in Theorem 3.18(d), and the last equality used
\[e_{\ell}^\perp (h_n) = \begin{cases} h_{n-1} & \text{if } k = 1, \\ 0 & \text{if } k \geq 2 \end{cases}\]
Now apply induction, since \((\lambda_1 - 1, \ldots, \lambda_\ell - 1)^t = (\lambda_2^t, \lambda_3^t, \ldots).\)

For the last assertion, note \((e_\mu^\perp, h_\lambda) = (e_\mu^\perp(\lambda), 1) = e_{\mu^t}(\lambda)\) for \(\lambda, \mu\) in \(\text{Par}_n\). \(\square\)

The following proposition will be the crux of the proof of Theorem 3.20(f) and (g), and turns out to be closely related to Kerov’s asymptotic theory of characters of the symmetric groups [79].

**Proposition 3.22.** Given a PSH \(A\) with PSH-basis \(\Sigma\) containing only one primitive \(\rho\), the two maps \(A \to \mathbb{Z}\) defined on \(A = \bigoplus_{n \geq 0} A_n\) via
\[\delta_h = \bigoplus_n h_n^\perp, \]
\[\delta_e = \bigoplus_n e_n^\perp\]
are characterized as the only two \(\mathbb{Z}\)-linear maps \(A \to \mathbb{Z}\) with the three properties of being
- **positive:** \(\delta(\mathbb{N}\Sigma) \subset \mathbb{N},\)
- **multiplicative:** \(\delta(a_1 a_2) = \delta(a_1) \delta(a_2),\) and
• normalized: \( \delta(\rho) = 1 \).

**Proof.** Notice that \( A \) is of finite type (as shown in the proof of Theorem 3.18). Thus, \( A^\circ \) is a graded Hopf algebra isomorphic to \( A \).

It should be clear from their definitions that \( \delta_h, \delta_e \) are \( \mathbb{Z} \)-linear, positive and normalized. To see that \( \delta_h \) is multiplicative, by \( \mathbb{Z} \)-linearity, it suffices to check that for \( a_1, a_2 \) in \( A_{n_1}, A_{n_2} \) with \( n_1 + n_2 = n \), one has

\[
\delta_h(a_1a_2) = h_n^n(a_1a_2) = \sum_{i_1 + i_2 = n} h_{i_1}^<(a_1)h_{i_2}^>(a_2) = h_{n_1}^<(a_1)h_{n_2}^>(a_2) = \delta_h(a_1)\delta_h(a_2)
\]

in which the second equality used Proposition 2.73(iii) and Theorem 3.18(d). The argument for \( \delta_e \) is symmetric.

Conversely, given \( A \xrightarrow{\delta} \mathbb{Z} \) which is \( \mathbb{Z} \)-linear, positive, multiplicative, and normalized, note that

\[
\delta(h_2) + \delta(e_2) = \delta(h_2 + e_2) = \delta(\rho^2) = \delta(\rho)^2 = 1^2 = 1
\]

and hence positivity implies that either \( \delta(h_2) = 0 \) or \( \delta(e_2) = 0 \). Assume the latter holds, and we will show that \( \delta = \delta_h \).

Given any \( \sigma \in \Sigma \cap A_n \setminus \{h_n\} \), note that \( e_2^2\sigma \neq 0 \) by Theorem 3.18(b), and hence \( 0 \neq (e_2^2\sigma, \rho^{n-2}) = (\sigma, e_2\rho^{n-2}) \). Thus \( \sigma \) is a constituent of \( e_2\rho^{n-2} \), so positivity implies

\[
0 \leq \delta(\sigma) \leq \delta(e_2\rho^{n-2}) = \delta(e_2)\delta(\rho^{n-2}) = 0.
\]

Thus \( \delta(\sigma) = 0 \) for \( \sigma \in \Sigma \cap A_n \setminus \{h_n\} \). Since \( \delta(\rho^n) = \delta(\rho)^n = 1^n = 1 \), this forces \( \delta(h_n) = 1 \), for each \( n \geq 0 \) (including \( n = 0 \), as \( 1 = \delta(\rho) = \delta(\rho \cdot 1) = \delta(\rho)\delta(1) = 1 \cdot \delta(1) = \delta(1) \)). Thus \( \delta = \delta_h \). The argument when \( \delta(h_2) = 0 \) showing \( \delta = \delta_e \) is symmetric. \( \square \)

**Proof of Theorem 3.20.** Many of the assertions of parts (e) and (f) will come from constructing the unique nontrivial PSH-automorphism \( \omega \) of \( A \) from the antipode \( S \): for homogeneous \( a \) in \( A_n \), define \( \omega(a) := (-1)^n S(a) \). We now study some of the properties of \( S \) and \( \omega \).

Notice that \( A \) is of finite type (as shown in the proof of Theorem 3.18). Thus, \( A^\circ \) is a graded Hopf algebra isomorphic to \( A \).

Since \( A \) is a PSH, it is commutative by Theorem 3.8 (applied to \( A \otimes_{\mathbb{Z}} \mathbb{Q} \)). This implies both that \( S, \omega \) are actually algebra endomorphisms by Proposition 1.32, and that \( S^2 = \text{id}_A = \omega^2 \) by Corollary 1.34.

Since \( A \) is self-dual and the defining diagram (1.12) satisfied by the antipode \( S \) is sent to itself when one replaces \( A \) by \( A^\circ \) and all maps by their adjoints, one concludes that \( S = S^* \) (where \( S^* \) means the restricted adjoint \( S^* : A^\circ \to A^\circ \)) i.e., \( S \) is self-adjoint. Since \( S \) is an algebra endomorphism, and \( S = S^* \), in fact \( S \) is also a coalgebra endomorphism, a bialgebra endomorphism, and a Hopf endomorphism (by Proposition 1.46(c)). The same properties are shared by \( \omega \).

Since \( \text{id}_A = S^2 = SS^* \), one concludes that \( S \) is an isometry, and hence so is \( \omega \).

Since \( \rho \) is primitive, one has \( S(\rho) = -\rho \) and \( \omega(\rho) = \rho \). Therefore \( \omega(\rho^n) = \rho^n \) for \( n = 1, 2, \ldots \). Use this as follows to check that \( \omega \) is a PSH-automorphism, which amounts to checking that every \( \sigma \) in \( \Sigma \) has \( \omega(\sigma) \) in \( \Sigma \):

\[
(\omega(\sigma), \omega(\sigma)) = (\sigma, \sigma) = 1
\]

so that \( \pm \omega(\sigma) \) lies in \( \Sigma \), but also if \( \sigma \) lies in \( A_n \), then

\[
(\omega(\sigma), \rho^n) = (\sigma, \omega(\rho^n)) = (\sigma, \rho^n) > 0.
\]

In summary, \( \omega \) is a PSH-automorphism of \( A \), an isometry, and an involution.

Let us try to determine the action of \( \omega \) on the \( \{h_n\} \). By similar reasoning as in (3.7), one has

\[
\Delta(h_2) = h_2 \otimes 1 + \rho \otimes \rho + 1 \otimes h_2.
\]

Thus \( 0 = S(h_2) + S(\rho)\rho + h_2 \), and combining this with \( S(\rho) = -\rho \), one has \( S(h_2) = e_2 \). Thus also \( \omega(h_2) = (-1)^2 S(h_2) = e_2 \).
We claim that this forces \( \omega(h_n) = e_n \), because \( h_2^+ \omega(h_n) = 0 \) via the following calculation: for any \( a \) in \( A \) one has

\[
(h_2^+ \omega(h_n), a) = (\omega(h_n), h_2a) = (h_n, \omega(h_2a)) = (h_n, e_2\omega(a)) = (e_2^+ h_n, \omega(a)) = (0, \omega(a)) = 0.
\]

Consequently the involution \( \omega \) swaps \( h_n \) and \( e_n \), while the antipode \( S \) has \( S(h_n) = (-1)^n e_n \) and \( S(e_n) = (-1)^n h_n \). Thus the coproduct formulas in (d) and definition of the antipode \( S \) imply the relation (2.17) between \( \{h_n\} \) and \( \{e_n\} \).

This relation (2.17) also lets one recursively express the \( h_n \) as polynomials with integer coefficients in the \( \{e_n\} \), and vice-versa, so that \( \{h_n\} \) and \( \{e_n\} \) each generate the same \( \mathbb{Z} \)-subalgebra \( A' \) of \( A \). We wish to show that \( A' \) exhausts \( A \).

We argue that Lemma 3.21 implies the Gram matrix \([h_{\mu}, \lambda], \mu, \lambda \in \text{Par}_n\) has determinant \( \pm 1 \) as follows. Since \( \{h_n\} \) and \( \{e_n\} \) both generate \( A' \), there exists a \( \mathbb{Z} \)-matrix \((a_{\mu, \lambda})\) expressing \( e_\mu = \sum_\lambda a_{\mu, \lambda} h_\lambda \), and one has

\[
[(e_\mu, h_\lambda)] = [a_{\mu, \lambda}] \cdot [(h_\mu, h_\lambda)].
\]

Taking determinants of these three \( \mathbb{Z} \)-matrices, and using the fact that the determinant on the left is \( \pm 1 \), both determinants on the right must also be \( \pm 1 \).

Now we will show that every \( \sigma \in \Sigma \cap A_n \) lies in \( A' \). Uniquely express \( \sigma = \sigma' + \sigma'' \) in which \( \sigma' \) lies in the \( \mathbb{R} \)-span \( \mathbb{R} A'_n \) and \( \sigma'' \) lies in the real perpendicular space \((\mathbb{R} A'_n)\perp \) inside \( \mathbb{R} \otimes_{\mathbb{Z}} A_n \). One can compute \( \mathbb{R} \)-coefficients \((c_{\mu})_{\mu \in \text{Par}_n}\) that express \( \sigma' = \sum_\mu c_{\mu} h_\mu \) by solving the system

\[
\left( \sum_\mu c_{\mu} h_\mu, h_\lambda \right) = (\sigma, h_\lambda) \quad \text{for } \lambda \in \text{Par}_n.
\]

This linear system is governed by the Gram matrix \([h_{\mu}, h_\lambda], \mu, \lambda \in \text{Par}_n\) with determinant \( \pm 1 \), and its right side has \( \mathbb{Z} \)-entries since \( \sigma, h_\lambda \) lie in \( A \). Hence the solution \((c_{\mu})_{\mu \in \text{Par}_n}\) will have \( \mathbb{Z} \)-entries, so \( \sigma' \) lies in \( A' \). Furthermore, \( \sigma'' = \sigma - \sigma' \) will lie in \( A \), and hence by the orthogonality of \( \sigma', \sigma'' \),

\[
1 = (\sigma, \sigma) = (\sigma', \sigma') + (\sigma'', \sigma'').
\]

One concludes that either \( \sigma'' = 0 \), or \( \sigma' = 0 \). The latter cannot occur since it would mean that \( \sigma = \sigma'' \) is perpendicular to all of \( A' \). But \( \rho'' = h_2^\perp \) lies in \( A' \), and \((\sigma, \rho'') \neq 0 \). Thus \( \sigma'' = 0 \), meaning \( \sigma = \sigma' \) lies in \( A' \).

This completes the proof of assertion (e). Note that in the process, having shown det\((h_{\mu}, h_\lambda), \mu, \lambda \in \text{Par}_n = \pm 1 \), one also knows that \( \{h_\lambda\}, \lambda \in \text{Par}_n \) are \( \mathbb{Z} \)-linearly independent, so that \( \{h_1, h_2, \ldots\} \) are algebraically independent, and \( A = \mathbb{Z}[h_1, h_2, \ldots] \) is the polynomial algebra generated by \( \{h_1, h_2, \ldots\} \).

For assertion (f), we have seen that \( \omega \) gives such a PSH-automorphism \( A \to A \), swapping \( h_n \leftrightarrow e_n \). Conversely, given a PSH-automorphism \( A \cong A \), consider the positive, multiplicative, normalized \( \mathbb{Z} \)-linear map \( \delta := \delta_0 \circ \varphi : A \to \mathbb{Z} \). Proposition 3.22 shows that either

- \( \delta = \delta_0 \), which then forces \( \varphi(h_n) = h_n \) for all \( n \), so \( \varphi = \text{id}_A \), or
- \( \delta = \delta_e \), which then forces \( \varphi(e_n) = h_n \) for all \( n \), so \( \varphi = \omega \).

For assertion (g), given a PSH \( A \) with PSH-basis \( \Sigma \) having exactly one primitive \( \rho \), since we have seen \( A = \mathbb{Z}[h_1, h_2, \ldots] \), where \( h_n \) in \( A \) is as defined in Theorem 3.18, one can uniquely define an algebra morphism \( A \cong A \) that sends the element \( h_n \) to the complete homogeneous symmetric function \( h_n(x) \). Assertions (d) and (e) show that \( \varphi \) is a bialgebra isomorphism, and hence it is a Hopf isomorphism. To show that it is a PSH-isomorphism, we first note that it is an isometry because one can iterate Proposition 2.73(iii) together with assertions (c) and (d) to compute all inner products

\[
(h_{\mu}, h_\lambda)_A = (1, h_\mu^+ h_\lambda^+) = (1, h_\mu^+ h_\mu^+ \cdots (h_\lambda, h_\lambda \cdots))_A
\]

for \( \mu, \lambda \) in \( \text{Par}_n \). Hence

\[
(h_{\mu}, h_\lambda)_A = (h_{\mu}(x), h_\lambda(x))_A = (\varphi(h_\mu), \varphi(h_\lambda))_A
\]

Once one knows \( \varphi \) is an isometry, then elements \( \omega \) in \( \Sigma \cap A_n \) are characterized in terms of the form \((\cdot, \cdot)\) by \((\omega, \omega) = 1 \) and \((\omega, \rho'') > 0 \). Hence \( \varphi \) sends each \( \sigma \) in \( \Sigma \) to a Schur function \( s_\lambda \), and is a PSH-isomorphism. □
4. Complex representations for $\mathfrak{S}_n$, wreath products, $GL_n(F_q)$

After reviewing the basics that we will need from representation and character theory of finite groups, we give Zelevinsky’s three main examples of PSH’s arising as spaces of virtual characters for three towers of finite groups:

- symmetric groups,
- their wreath products with any finite group, and
- the finite general linear groups.

Much in this chapter traces its roots to Zelevinsky’s book [170]. The results concerning the symmetric groups, however, are significantly older and spread across the literature: see, e.g., [154, §I.7], [138, §7.3], [106, §I.7], [138, §4.7], [87], for proofs using different tools.

4.1. Review of complex character theory. We shall now briefly discuss some basics of representation (and character) theory that will be used below. A good source for this material, including the crucial Mackey formula, is Serre [146, Chaps. 1-7].

4.1.1. Basic definitions, Maschke, Schur. For a group $G$, a representation is a homomorphism $G \xrightarrow{\varphi} GL(V)$ for some vector space $V$ over a field. We will take the field to be $\mathbb{C}$ from now on, and we will also assume that $V$ is finite-dimensional over $\mathbb{C}$. Thus a representation of $G$ is the same as a finite-dimensional (left-) $\mathbb{C}G$-module $V$. (We use the notations $\mathbb{C}G$ and $\mathbb{C}[G]$ synonymously for the group algebra of $G$ over $\mathbb{C}$. More generally, if $S$ is a set, then $\mathbb{C}S = \mathbb{C}[S]$ denotes the free $\mathbb{C}$-module with basis $S$.)

We also assume that $G$ is finite, so that Maschke’s Theorem\textsuperscript{143} says that $\mathbb{C}G$ is semisimple, meaning that every $\mathbb{C}G$-module $U \subset V$ has a $\mathbb{C}G$-module complement $U'$ with $V = U \oplus U'$. Equivalently, indecomposable $\mathbb{C}G$-modules are the same thing as simple (=irreducible) $\mathbb{C}G$-modules.

Schur’s Lemma implies that for two simple $\mathbb{C}G$-modules $V_1, V_2$, one has

$$\text{Hom}_{\mathbb{C}G}(V_1, V_2) \cong \begin{cases} \mathbb{C} & \text{if } V_1 \cong V_2, \\ 0 & \text{if } V_1 \ncong V_2. \end{cases}$$

4.1.2. Characters and Hom spaces. A $\mathbb{C}G$-module $V$ is completely determined up to isomorphism by its character

$$G \xrightarrow{\chi_V} \mathbb{C}, \quad g \mapsto \chi_V(g) := \text{trace}(g : V \to V).$$

This character $\chi_V$ is a class function, meaning it is constant on $G$-conjugacy classes. The space $R_{\mathbb{C}}(G)$ of class functions $G \to \mathbb{C}$ has a Hermitian, positive definite form

$$(f_1, f_2)_G := \frac{1}{|G|} \sum_{g \in G} f_1(g) \overline{f_2(g)}.$$

For any two $\mathbb{C}G$-modules $V_1, V_2$,

$$(\chi_{V_1}, \chi_{V_2})_G = \dim_{\mathbb{C}} \text{Hom}_{\mathbb{C}G}(V_1, V_2).$$

The set of all irreducible characters

$$\text{Irr}(G) = \{\chi_V : V \text{ is a simple } \mathbb{C}G\text{-module}\}$$

forms an orthonormal basis of $R_{\mathbb{C}}(G)$ with respect to this form, and spans a $\mathbb{Z}$-sublattice

$$R(G) := \mathbb{Z}\text{Irr}(G) \subset R_{\mathbb{C}}(G)$$

sometimes called the virtual characters of $G$. For every $\mathbb{C}G$-module $V$, the character $\chi_V$ belongs to $R(G)$.

Instead of working with the Hermitian form $(\cdot, \cdot)_G$ on $G$, we could also (and some authors do) define a $\mathbb{C}$-bilinear form $(\cdot, \cdot)_{\mathbb{C}}$ on $R_{\mathbb{C}}(G)$ by

$$(f_1, f_2)_{\mathbb{C}} := \frac{1}{|G|} \sum_{g \in G} f_1(g) f_2(g^{-1}).$$

---

\textsuperscript{143}... which has a beautiful generalization to finite-dimensional Hopf algebras due to Larson and Sweedler; see Montgomery [117, §2.2].
This form is not identical with $\langle \cdot , \cdot \rangle_G$ (indeed, $\langle \cdot , \cdot \rangle_G$ is bilinear while $\langle \cdot , \cdot \rangle_G$ is Hermitian), but it still satisfies (4.1), and thus is identical with $\langle \cdot , \cdot \rangle_G$ on $R(G) \times R(G)$. Hence, for all we are going to do until Section 4.9, we could just as well use the form $\langle \cdot , \cdot \rangle_G$ instead of $\langle \cdot , \cdot \rangle_G$.

4.1.3. Tensor products. Given two groups $G_1, G_2$ and $CG_i$-modules $V_i$ for $i = 1, 2$, their tensor product $V_1 \otimes_{CG_1} V_2$ becomes a $CG_1 \times G_2$-module via $(g_1, g_2)(v_1 \otimes v_2) = g_1(v_1) \otimes g_2(v_2)$. This module is called the (outer) tensor product of $V_1$ and $V_2$. When $V_1, V_2$ are both simple, then so is $V_1 \otimes V_2$, and every simple $CG_1 \times G_2$-module arises this way (with $V_1$ and $V_2$ determined uniquely up to isomorphism).\(^{144}\) Thus one has identifications and isomorphisms

$$\text{Irr}(G_1 \times G_2) = \text{Irr}(G_1) \times \text{Irr}(G_2),$$

$$R(G_1 \times G_2) \cong R(G_1) \otimes_{\mathbb{Z}} R(G_2);$$

here, $\chi_{V_1} \otimes \chi_{V_2} \in R(G_1) \otimes_{\mathbb{Z}} R(G_2)$ is being identified with $\chi_{V_1} \otimes \chi_{V_2} \in R(G_1 \times G_2)$ for all $CG_i$-modules $V_1$ and $CG_2$-modules $V_2$.

The latter isomorphism is actually a restriction of the isomorphism $R_C(G_1 \times G_2) \cong R_C(G_1) \otimes_{\mathbb{C}} R_C(G_2)$ under which every pure tensor $\phi_1 \otimes \phi_2 \in R_C(G_1) \otimes_{\mathbb{C}} R_C(G_2)$ corresponds to the class function $G_1 \times G_2 \to \mathbb{C}$, $(g_1, g_2) \mapsto \phi_1(g_1) \otimes \phi_2(g_2)$.

Given two $CG_i$-modules $V_1$ and $W_1$ and two $CG_2$-modules $V_2$ and $W_2$, we have

$$\langle \chi_{V_1} \otimes \chi_{W_1}, \chi_{V_2} \otimes \chi_{W_2} \rangle_{G_1 \times G_2} = \langle \chi_{V_1}, \chi_{W_1} \rangle_{G_1} \langle \chi_{V_2}, \chi_{W_2} \rangle_{G_2}.$$  

4.1.4. Induction and restriction. Given a subgroup $H < G$ and $CH$-module $U$, one can use the fact that $CG$ is a $(CG, CH)$-bimodule to form the induced $CG$-module

$$\text{Ind}^G_H U := CG \otimes_{CH} U.$$

The fact that $CG$ is free as a $(right-)CH$-module\(^{145}\) on basis elements $\{t_g\}_{gH \in G/H}$ makes this tensor product easy to analyze. For example one can compute its character

$$\chi_{\text{Ind}^G_H U}(g) = \frac{1}{|H|} \sum_{k \in G; k g k^{-1} \in H} \chi_U(k g k^{-1}).$$ \(^{146}\)

One can also recognize when a $CG$-module $V$ is isomorphic to $\text{Ind}^G_H U$ for some $CH$-module $U$: this happens if and only if there is an $H$-stable subspace $V \subset U$ having the property that $V = \bigoplus_{gH \in G/H} gU$.

The above construction of a $CG$-module $\text{Ind}^G_H U$ corresponding to any $CH$-module $U$ is part of a functor $\text{Ind}^G_H$ from the category of $CH$-modules to the category of $CG$-modules\(^{147}\); this functor is called induction.

Besides induction on $CH$-modules, one can define induction on class functions of $H$:

**Exercise 4.1.** Let $G$ be a finite group, and $H$ a subgroup of $G$. Let $f \in R_C(H)$ be a class function. We define the *induction* $\text{Ind}^G_H f$ of $f$ to be the function $G \to \mathbb{C}$ given by

$$\left( \text{Ind}^G_H f \right)(g) = \frac{1}{|H|} \sum_{k \in G; k g k^{-1} \in H} f(k g k^{-1}) \quad \text{for all } g \in G.$$  

(a) Prove that this induction $\text{Ind}^G_H f$ is a class function on $G$, hence belongs to $R_C(G)$.

(b) Let $J$ be a system of right coset representatives for $H \backslash G$, so that $G = \bigsqcup_{j \in J} H j$. Prove that

$$\left( \text{Ind}^G_H f \right)(g) = \sum_{j \in J; j g j^{-1} \in H} f(j g j^{-1}) \quad \text{for all } g \in G.$$  

\(^{144}\)This is proven in [146, §3.2, Thm. 10]. The fact that $\mathbb{C}$ is algebraically closed is essential for this!

\(^{145}\)... which also has a beautiful generalization to finite-dimensional Hopf algebras due to Nichols and Zoeller; see [117, §3.1].

\(^{146}\)See [146, §7.2, Prop. 20(ii)] for the proof of this equality. (Another proof is given in [49, Remark 5.9.2 (the Remark after Theorem 3.32 in the arXiv version)], but [49] uses a different definition of $\text{Ind}^G_H U$; see Remark 4.5 for why it is equivalent to ours. Yet another proof of (4.3) is given in Exercise 4.14(k).)

\(^{147}\)On morphisms, it sends any $f : U \to U'$ to $\text{id}_{CG} \otimes_{CH} f : CG \otimes_{CH} U \to CG \otimes_{CH} U'$. 


The induction \( \text{Ind}_H^G \) defined in Exercise 4.1 is a \( C \)-linear map \( R_C(H) \to R_C(G) \). Since every \( CH \)-module \( U \) satisfies
\[
\chi_{\text{Ind}_H^G U} = \text{Ind}_H^G(\chi_U)
\]
(4.5), this \( C \)-linear map \( \text{Ind}_H^G \) restricts to a \( Z \)-linear map \( R(H) \to R(G) \) (also denoted \( \text{Ind}_H^G \)) which sends the character \( \chi_U \) of any \( CH \)-module \( U \) to the character \( \chi_{\text{Ind}_H^G U} \) of the induced \( CG \)-module \( \text{Ind}_H^G U \).

Exercise 4.2. Let \( G, H \) and \( I \) be three finite groups such that \( I < H < G \). Let \( U \) be a \( CI \)-module. Prove that \( \text{Ind}_H^G \text{Ind}_I^H U \cong \text{Ind}_I^G U \). (This fact is often referred to as the transitivity of induction.)

Exercise 4.3. Let \( G_1 \) and \( G_2 \) be two groups. Let \( H_1 < G_1 \) and \( H_2 < G_2 \) be two subgroups. Let \( U_1 \) be a \( CH_1 \)-module, and \( U_2 \) be a \( CH_2 \)-module. Show that
\[
\text{Ind}_{H_1 \times H_2}^{G_1 \times G_2} (U_1 \otimes U_2) \cong \left( \text{Ind}_{H_1}^{G_1} U_1 \right) \otimes \left( \text{Ind}_{H_2}^{G_2} U_2 \right)
\]
as \( C[G_1 \times G_2] \)-modules.

The restriction operation \( V \mapsto \text{Res}_H^G V \) restricts a \( CG \)-module \( V \) to a \( CH \)-module. Frobenius reciprocity asserts the adjointness between \( \text{Ind}_H^G \) and \( \text{Res}_H^G \)
\[
\text{Hom}_{CG}(\text{Ind}_H^G U, V) \cong \text{Hom}_{CH}(U, \text{Res}_H^G V),
\]
as a special case (\( S = A = CG, R = CH, B = U, C = V \)) of the general adjoint associativity
\[
\text{Hom}_S(A \otimes_R B, C) \cong \text{Hom}_R(B, \text{Hom}_S(A, C))
\]
for \( S, R \) two rings, \( A \) an \((S, R)\)-bimodule, \( B \) a left \( R \)-module, \( C \) a left \( S \)-module.

We can define not just the restriction of a \( CG \)-module, but also the restriction of a class function \( f \in R_C(G) \). When \( H \) is a subgroup of \( G \), the restriction \( \text{Res}_H^G f \) of an \( f \in R_C(G) \) is defined as the result of restricting the map \( f : G \to C \) to \( H \). This \( \text{Res}_H^G f \) is easily seen to belong to \( R_C(H) \), and so \( \text{Res}_H^G \) is a \( C \)-linear map \( R_C(G) \to R_C(H) \). This map restricts to a \( Z \)-linear map \( R(G) \to R(H) \), since we have \( \text{Res}_H^G \chi_V = \chi_{\text{Res}_H^G V} \) for any \( CG \)-module \( V \). Taking characters in (4.7) (and recalling \( \text{Res}_H^G \chi_V = \chi_{\text{Res}_H^G V} \) and (4.5)), we obtain
\[
(\text{Ind}_H^G \chi_U, \chi_V)_G = (\chi_U, \text{Res}_H^G \chi_V)_H.
\]
By bilinearity, this yields the equality
\[
\left( \text{Ind}_H^G \alpha, \beta \right)_G = (\alpha, \text{Res}_H^G \beta)_H
\]
for any class functions \( \alpha \in R_C(H) \) and \( \beta \in R_C(G) \) (since \( R(G) \) spans \( R_C(G) \) as a \( C \)-vector space).

Exercise 4.4. Let \( G \) be a finite group, and let \( H < G \). Let \( U \) be a \( CH \)-module. If \( A \) and \( B \) are two algebras, \( P \) is a \((B, A)\)-bimodule and \( Q \) is a left \( B \)-module, then \( \text{Hom}_B(P, Q) \) is a left \( A \)-module (since \( CG \) is a \((CH, CG)\)-bimodule). As a consequence, \( \text{Hom}_{CG}(CG, U) \) is a \( CG \)-module. Prove that this \( CG \)-module is isomorphic to \( \text{Ind}_H^G U \).

Remark 4.5. Some texts define the induction \( \text{Ind}_H^G U \) of a \( CH \)-module \( U \) to be \( \text{Hom}_{CG}(CG, U) \) (rather than to be \( CG \otimes_{CH} U \), as we did).\(^{149}\) As Exercise 4.4 shows, this definition is equivalent to ours as long as \( G \) is finite (but not otherwise).

Exercise 4.4 yields the following “wrong-way” version of Frobenius reciprocity:

Exercise 4.6. Let \( G \) be a finite group; let \( H < G \). Let \( U \) be a \( CG \)-module, and let \( V \) be a \( CH \)-module. Prove that \( \text{Hom}_{CG}(U, \text{Ind}_H^G V) \cong \text{Hom}_{CH}(\text{Res}_H^G U, V) \).

\(^{148}\)This follows by comparing the value of \( \chi_{\text{Ind}_H^G U}(g) \) obtained from (4.3) with the value of \( \left( \text{Ind}_H^G(\chi_U) \right)(g) \) found using (4.4).

\(^{149}\)Or they define it as a set of morphisms of \( H \)-sets from \( G \) to \( U \) (this is how [49, Def. 5.8.1 (Def. 4.28 in the arXiv version)] defines it); this is easily seen to be equivalent to \( \text{Hom}_{CG}(CG, U) \).
4.1.5. Mackey’s formula. Mackey gave an alternate description of a module which has been induced and then restricted. To state it, for a subgroup $H < G$ and $g$ in $G$, let $g^H := g^{-1}Hg$ and $H^g := gHg^{-1}$. Given a $\mathbb{C} H$-module $U$, say defined by a homomorphism $H \xymatrix@1{\ar[r] & GL(U)}$, let $U^g$ denote the $\mathbb{C}[gHg^{-1}]$-module on the same $\mathbb{C}$-vector space $U$ defined by the composite homomorphism

$$g^H : \quad \begin{array}{ccc} H & \longrightarrow & H \\ h \mapsto & g^{-1}hg \end{array} \quad \xymatrix@1{\ar[r] & GL(U)}.$$

**Theorem 4.7.** (Mackey’s formula) Consider subgroups $H, K < G$, and any $\mathbb{C} H$-module $U$. If $\{g_1, \ldots, g_t\}$ are double coset representatives for $K \backslash G / H$, then

$$\text{Res}^G_H \text{Ind}^G_K U \cong \bigoplus_{i=1}^t \text{Ind}^K_{H \cap K} \left( (\text{Res}^H_{H \cap K} U)^{g_i} \right).$$

**Proof.** In this proof, all tensor product symbols $\otimes$ should be interpreted as $\otimes_{\mathbb{C} H}$. Recall $\mathbb{C} G$ has $\mathbb{C}$-basis $\{t_g\}_{g \in G}$. For subsets $S \subset G$, let $\mathbb{C}[S]$ denote the $\mathbb{C}$-span of $\{t_g\}_{g \in S}$ in $\mathbb{C} G$.

Note that each double coset $KgH$ gives rise to a sub-$\mathbb{C}[K, H]$-bimodule $\mathbb{C}[KgH]$ within $\mathbb{C} G$, and one has a $\mathbb{C} K$-module direct sum decomposition

$$\text{Ind}^G_K U = \mathbb{C} G \otimes U = \bigoplus_{i=1}^t \mathbb{C}[Kg_i H] \otimes U.$$

Hence it suffices to check for any element $g$ in $G$ that

$$\mathbb{C}[KgH] \otimes U \cong \text{Ind}^K_{H \cap K} \left( (\text{Res}^H_{H \cap K} U)^g \right).$$

Note that $g^H \cap K$ is the subgroup of $K$ consisting of the elements $k$ in $K$ for which $kgH = gH$. Hence by picking $\{k_1, \ldots, k_s\}$ to be coset representatives for $K / (g^H \cap K)$, one disjointly decomposes the double coset

$$KgH = \bigsqcup_{j=1}^s k_j (g^H \cap K)gH,$$

giving a $\mathbb{C}$-vector space direct sum decomposition

$$\mathbb{C}[KgH] \otimes U = \bigoplus_{j=1}^s \mathbb{C}[k_j (g^H \cap K)gH] \otimes U \cong \text{Ind}^K_{H \cap K} \left( \mathbb{C}[\{g^H \cap K \} gH] \otimes U \right).$$

So it remains to check that one has a $\mathbb{C}[g^H \cap K]$-module isomorphism

$$\mathbb{C}[\{g^H \cap K \} gH] \otimes U \cong (\text{Res}^H_{H \cap K} U)^g.$$

Bearing in mind that, for each $k$ in $g^H \cap K$ and $h$ in $H$, one has $g^{-1}kg$ in $H$ and hence

$$t_{kgh} \otimes u = t_g \cdot t_{g^{-1}kg} \otimes u = t_g \otimes g^{-1}kg \cdot u,$$

one sees that this isomorphism can be defined by mapping

$$t_{kgh} \otimes u \longmapsto g^{-1}kg \cdot u.$$ 

\hfill $\square$

4.1.6. Inflation and fixed points. There are two (adjoint) constructions on representations that apply when one has a normal subgroup $K \trianglelefteq G$. Given a $\mathbb{C}[G/K]$-module $U$, say defined by the homomorphism $G / K \xymatrix@1{\ar[r] & GL(U)}$, the inflation of $U$ to a $\mathbb{C} G$-module $\text{Infl}^G_{G/K} U$ has the same underlying space $U$, and is defined by the composite homomorphism $G \to G / K \xymatrix@1{\ar[r] & GL(U)}$. We will later use the easily-checked fact that when $H < G$ is any other subgroup, one has

$$(\text{4.10}) \quad \text{Res}^G_H \text{Infl}^G_{G/K} U = \text{Infl}^H_{H \cap K} \text{Res}^G_{H / H \cap K} U.$$

(We regard $H / H \cap K$ as a subgroup of $G / K$, since the canonical homomorphism $H / H \cap K \to G / K$ is injective.)
Inflation turns out to be adjoint to the $K$-fixed space construction sending a $\mathbb{C}G$-module $V$ to the $\mathbb{C}[G/K]$-module

$$V^K := \{ v \in V : k(v) = v \text{ for } k \in K \}$$

Note that $V^K$ is indeed a $G$-stable subspace: for any $v$ in $V^K$ and $g$ in $G$, one has that $g(v)$ lies in $V^K$ since an element $k$ in $K$ satisfies $kg(v) = g \cdot g^{-1}kg(v) = g(v)$ as $g^{-1}kg$ lies in $K$. One has this adjointness

$$\text{Hom}_{\mathbb{C}G}(\text{Infl}^G_{G/K} U, V) = \text{Hom}_{\mathbb{C}[G/K]}(U, V^K)$$

because any $\mathbb{C}G$-module homomorphism $\varphi$ on the left must have the property that $k\varphi(u) = \varphi(k(u)) = \varphi(u)$ for all $k$ in $K$, so that $\varphi$ actually lies on the right.

We will also need the following formula for the character $\chi_{V^K}$ in terms of the character $\chi_V$:

$$\chi_{V^K}(gK) = \frac{1}{|K|} \sum_{k \in K} \chi_V(gk). \tag{4.12}$$

To see this, note that when one has a $\mathbb{C}$-linear endomorphism $\varphi$ on a space $V$ that preserves some $\mathbb{C}$-subspace $W \subset V$, if $V \xrightarrow{\pi} W$ is any idempotent projection onto $W$, then the trace of the restriction $\varphi|_{W}$ equals the trace of $\varphi \circ \pi$ on $V$. Applying this to $W = V^K$ and $\varphi = g$, with $\pi = \frac{1}{|K|} \sum_{k \in K} k$, gives (4.12).\footnote{For another proof of (4.12), see Exercise 4.14(l).}

Another way to restate (4.12) is:

$$\chi_{V^K}(gK) = \frac{1}{|K|} \sum_{h \in gK} \chi_V(h). \tag{4.13}$$

Inflation and $K$-fixed space construction can also be defined on class functions. For inflation, this is particularly easy: Inflation $\text{Infl}^G_{G/K}$ $f$ of an $f \in R_C(G/K)$ is defined as the composition $G \longrightarrow G/K \xrightarrow{f} \mathbb{C}$. This is a class function of $G$ and thus lies in $R_C(G)$. Thus, inflation $\text{Infl}^G_{G/K}$ is a $\mathbb{C}$-linear map $R_C(G/K) \rightarrow R_C(G)$. It restricts to a $\mathbb{Z}$-linear map $R(G/K) \rightarrow R(G)$, since it is clear that every $\mathbb{C}(G/K)$-module $U$ satisfies $\text{Infl}^G_{G/K} \chi_U = \chi_{\text{Infl}^G_{G/K} U}$.

We can also use (4.12) (or (4.13)) as inspiration for defining a “$K$-fixed space construction” on class functions. Explicitly, for every class function $f \in R_C(G)$, we define a class function $f^K \in R_C(G/K)$ by

$$f^K(gK) = \frac{1}{|K|} \sum_{k \in K} f(gk) = \frac{1}{|K|} \sum_{h \in gK} f(h).$$

The map $(\cdot)^K : R_C(G) \rightarrow R_C(G/K)$, $f \mapsto f^K$ is $\mathbb{C}$-linear, and restricts to a $\mathbb{Z}$-linear map $R(G) \rightarrow R(G/K)$. Again, we have a compatibility with the $K$-fixed point construction on modules: We have $\chi_{V^K} = (\chi_V)^K$ for every $\mathbb{C}G$-module $V$.

Taking characters in (4.11), we obtain

$$\langle \text{Infl}^G_{G/K} \chi_U, \chi_V \rangle_G = \langle \chi_U, \chi_V^K \rangle_{G/K} \tag{4.14}$$

for any $\mathbb{C}[G/K]$-module $U$ and any $\mathbb{C}G$-module $V$ (since $\text{Infl}^G_{G/K} U = \text{Infl}^G_{G/K} \chi_U$ and $\chi_{V^K} = (\chi_V)^K$). By $\mathbb{Z}$-linearity, this implies that

$$\left( \text{Infl}^G_{G/K} \alpha, \beta \right)_G = \left( \alpha, \beta^K \right)_{G/K}$$

for any class functions $\alpha \in R_C(G/K)$ and $\beta \in R_C(G)$.

There is also an analogue of (4.6):

**Lemma 4.8.** Let $G_1$ and $G_2$ be two groups, and $K_1 < G_1$ and $K_2 < G_2$ be two respective subgroups. Let $U_i$ be a $\mathbb{C}G_i$-module for each $i \in \{1, 2\}$. Then,

$$U_1 \otimes U_2)^{K_1 \times K_2} = U_1^{K_1} \otimes U_2^{K_2} \tag{4.15}$$

(as subspaces of $U_1 \otimes U_2$).
Proof. The subgroup $K_1 = K_1 \times 1$ of $G_1 \times G_2$ acts on $U_1 \otimes U_2$, and its fixed points are $(U_1 \otimes U_2)^{K_1} = U_1^{K_1} \otimes U_2$ (because for a $\mathbb{C}K_1$-module, tensoring with $U_2$ is the same as taking a direct power, which clearly commutes with taking fixed points). Similarly, $(U_1 \otimes U_2)^{K_2} = U_1 \otimes U_2^{K_2}$. Now,

$$(U_1 \otimes U_2)^{K_1 \times K_2} = (U_1 \otimes U_2)^{K_1} \cap (U_1 \otimes U_2)^{K_2} = \left(U_1^{K_1} \otimes U_2\right) \cap \left(U_1 \otimes U_2^{K_2}\right) = U_1^{K_1} \otimes U_2^{K_2}$$

according to the known linear-algebraic fact stating that if $P$ and $Q$ are subspaces of two vector spaces $U$ and $V$, respectively, then $(P \otimes V) \cap (U \otimes Q) = P \otimes Q$.

Exercise 4.9. (a) Let $G_1$ and $G_2$ be two groups. Let $V_i$ and $W_i$ be finite-dimensional $\mathbb{C}G_i$-modules for every $i \in \{1, 2\}$. Prove that the $\mathbb{C}$-linear map

$\text{Hom}_{\mathbb{C}G_1}(V_1, W_1) \otimes \text{Hom}_{\mathbb{C}G_2}(V_2, W_2) \to \text{Hom}_{\mathbb{C}[G_1 \times G_2]}(V_1 \otimes V_2, W_1 \otimes W_2)$

sending each tensor $f \otimes g$ to the tensor product $f \otimes g$ of homomorphisms is a vector space isomorphism. (b) Use part (a) to give a new proof of (4.2).

As an aside, (4.10) has a “dual” analogue:

Exercise 4.10. Let $G$ be a finite group, and let $K \triangleleft H < G$. Let $U$ be a $\mathbb{C}H$-module. As usual, regard $H/(H \cap K)$ as a subgroup of $G/K$. Show that

$$\left(\text{Ind}^G_H U\right)^K \cong \text{Ind}^G_{H/(H \cap K)} \left(U^{H/(H \cap K)}\right)$$

as $\mathbb{C}[G/K]$-modules.

Inflation also “commutes” with induction:

Exercise 4.11. Let $G$ be a finite group, and let $K < H < G$ be such that $K \triangleleft G$. Thus, automatically, $K \triangleleft H$, and we regard the quotient $H/K$ as a subgroup of $G/K$. Let $V$ be a $\mathbb{C}[H/K]$-module. Show that

$$\text{Infl}^G_{H/K} \text{Ind}^G_H V \cong \text{Ind}^G_{H/(H \cap K)} \text{Infl}^G_{H/(H \cap K)} V$$

as $\mathbb{C}$-modules.

Exercise 4.12. Let $G$ be a finite group, and let $K \triangleleft G$. Let $V$ be a $\mathbb{C}G$-module. Let $I_{V,K}$ denote the $\mathbb{C}$-vector subspace of $V$ spanned by all elements of the form $v - kv$ for $k \in K$ and $v \in V$.

(a) Show that $I_{V,K}$ is a $\mathbb{C}G$-submodule of $V$.

(b) Let $V_K$ denote the quotient $\mathbb{C}G$-module $V/I_{V,K}$. (This module is occasionally called the $K$-coinvariant module of $V$, a name it sadly shares with at least two other non-equivalent constructions in algebra.) Show that $V_K \cong \text{Infl}^G_{H/K} (V^K)$ as $\mathbb{C}G$-modules. (Use char $\mathbb{C} = 0$.)

In the remainder of this subsection, we shall briefly survey generalized notions of induction and restriction, defined in terms of a group homomorphism $\rho$ rather than in terms of a group $G$ and a subgroup $H$. These generalized notions (defined by van Leeuwen in [94, §2.2]) will not be used in the rest of these notes, but they shed some new light on the facts about induction, restriction, inflation and fixed point construction discussed above. (In particular, they reveal that some of said facts have common generalizations.)

The reader might have noticed that the definitions of inflation and of restriction (both for characters and for modules) are similar. In fact, they both are particular cases of the following construction:

Remark 4.13. Let $G$ and $H$ be two finite groups, and let $\rho : H \to G$ be a group homomorphism.

- If $f \in R_C(G)$, then the $\rho$-restriction $\text{Res}_\rho f$ of $f$ is defined as the map $f \circ \rho : H \to \mathbb{C}$. This map is easily seen to belong to $R_C(H)$.

- If $V$ is a $\mathbb{C}G$-module, then the $\rho$-restriction $\text{Res}_\rho V$ of $V$ is the $\mathbb{C}H$-module with ground space $V$ and action given by

$$h \cdot v = \rho(h) \cdot v$$

for every $h \in H$ and $v \in V$.

This construction generalizes both inflation and restriction: If $H$ is a subgroup of $G$, and if $\rho : H \to G$ is the inclusion map, then $\text{Res}_\rho f = \text{Res}_H^G f$ (for any $f \in R_C(G)$) and $\text{Res}_\rho V = \text{Res}_H^G V$ (for any $\mathbb{C}G$-module $V$). If, instead, we have $G = H/K$ for a normal subgroup $K$ of $H$, and if $\rho : H \to G$ is the projection map, then $\text{Res}_\rho f = \text{Infl}_H^K f$ (for any $f \in R_C(H/K)$) and $\text{Res}_\rho V = \text{Infl}_H^K V$ (for any $\mathbb{C}[H/K]$-module $V$).

A subtler observation is that induction and fixed point construction can be generalized by a common notion. This is the subject of Exercise 4.14 below.

Exercise 4.14. Let $G$ and $H$ be two finite groups, and let $\rho : H \to G$ be a group homomorphism. We introduce the following notations:
- If \( f \in R_C(H) \), then the \( \rho \)-induction \( \text{Ind}_\rho f \) of \( f \) is a map \( G \to \mathbb{C} \) which is defined as follows:

\[
(\text{Ind}_\rho f) (g) = \frac{1}{|H|} \sum_{(h,k) \in H \times G; \quad k p(h) k^{-1} = g} f(h)
\]

for every \( g \in G \).

- If \( U \) is a \( CH \)-module, then the \( \rho \)-induction \( \text{Ind}_\rho U \) of \( U \) is defined as the \( CG \)-module \( CG \otimes_{CH} U \), where \( CG \) is regarded as a \((CG, CH)\)-bimodule according to the following rule: The left \( CG \)-module structure on \( CG \) is plain multiplication inside \( CG \); the right \( CH \)-module structure on \( CG \) is induced by the \( C \)-algebra homomorphism \( \mathbb{C} [\rho] : CH \to CG \) (thus, it is explicitly given by \( \gamma \eta = \gamma \cdot (\mathbb{C} [\rho]) \eta \) for all \( \gamma \in CG \) and \( \eta \in CH \)).

Prove the following properties of this construction:

(a) For every \( f \in R_C(H) \), we have \( \text{Ind}_\rho f \in R_C(G) \).
(b) For any finite-dimensional \( CH \)-module \( U \), we have \( \chi_{\text{Ind}_\rho U} = \text{Ind}_\rho \chi_U \).
(c) If \( H \) is a subgroup of \( G \), and if \( \rho : H \to G \) is the inclusion map, then \( \text{Ind}_\rho f = \text{Ind}_H^G f \) for every \( f \in R_C(H) \).
(d) If \( H \) is a subgroup of \( G \), and if \( \rho : H \to G \) is the inclusion map, then \( \text{Ind}_\rho U = \text{Ind}_H^G U \) for every \( CH \)-module \( U \).
(e) If \( G = H/K \) for some normal subgroup \( K \) of \( H \), and if \( \rho : H \to G \) is the projection map, then \( \text{Ind}_\rho f = f^K \) for every \( f \in R_C(H) \).
(f) If \( G = H/K \) for some normal subgroup \( K \) of \( H \), and if \( \rho : H \to G \) is the projection map, then \( \text{Ind}_\rho U \cong U^K \) for every \( CH \)-module \( U \).
(g) Any class functions \( \alpha \in R_C(H) \) and \( \beta \in R_C(G) \) satisfy

\[
(\text{Ind}_\rho \alpha, \beta)_G = (\alpha, \text{Res}_\rho \beta)_H
\]

and

\[
(\text{Ind}_\rho \alpha, \beta)_G = (\alpha, \text{Res}_\rho \beta)_H.
\]

(See Remark 4.13 for the definition of \( \text{Res}_\rho \beta \).)

(h) We have \( \text{Hom}_{CG}(\text{Ind}_\rho U, V) \cong \text{Hom}_{CH}(U, \text{Res}_\rho V) \) for every \( CH \)-module \( U \) and every \( CG \)-module \( V \). (See Remark 4.13 for the definition of \( \text{Res}_\rho V \).)

(i) Similarly to how we made \( CG \) into a \((CG, CH)\)-bimodule, let us make \( CG \) into a \((CH, CG)\)-bimodule (so the right \( CG \)-module structure is plain multiplication inside \( CG \), whereas the left \( CH \)-module structure is induced by the \( C \)-algebra homomorphism \( \mathbb{C} [\rho] : CH \to CG \)). If \( U \) is any \( CH \)-module, then the \( CG \)-module \( \text{Hom}_{CH}(CG, U) \) (defined as in Exercise 4.4 using the \((CH, CG)\)-bimodule structure on \( CG \)) is isomorphic to \( \text{Ind}_\rho U \).

(j) We have \( \text{Hom}_{CG}(U, \text{Ind}_\rho V) \cong \text{Hom}_{CH}(\text{Res}_\rho U, V) \) for every \( CG \)-module \( U \) and every \( CH \)-module \( V \). (See Remark 4.13 for the definition of \( \text{Res}_\rho V \).)

Furthermore:

(k) Use the above to prove the formula (4.3).
(l) Use the above to prove the formula (4.12).

[Hint: Part (b) of this exercise is hard. To solve it, it is useful to have a way of computing the trace of a linear operator without knowing a basis of the vector space it is acting on. There is a way to do this using a “finite dual generating system”, which is a somewhat less restricted notion than that of a basis\(^{151}\). Try to create a finite dual generating system for \( \text{Ind}_\rho U \) from one for \( U \) (and from the group \( G \)), and then use it to compute \( \chi_{\text{Ind}_\rho U} \).]

\(^{151}\)More precisely: Let \( K \) be a field, and \( V \) be a \( K \)-vector space. A \textit{finite dual generating system} for \( V \) means a triple \((I, (a_i)_{i \in I}, (f_i)_{i \in I})\), where

- \( I \) is a finite set;
- \((a_i)_{i \in I}\) is a family of elements of \( V \);
- \((f_i)_{i \in I}\) is a family of elements of \( V^* \) (where \( V^* \) means \( \text{Hom}_K(V, K) \))

such that every \( v \in V \) satisfies \( v = \sum_{i \in I} f_i(v) a_i \). For example, if \((e_j)_{j \in J}\) is a finite basis of the vector space \( V \), and if \((e'^*_j)_{j \in J}\) is the basis of \( V^* \) dual to this basis \((e_j)_{j \in J}\), then \((J, (e_j)_{j \in J}, (e'^*_j)_{j \in J})\) is a finite dual generating system for \( V \); however, most finite dual generating systems are not obtained this way.
The solution of part (i) is a modification of the solution of Exercise 4.4, but complicated by the fact that $H$ is no longer (necessarily) a subgroup of $G$. Part (f) can be solved by similar arguments, or using part (i), or using Exercise 4.12(b).

The result of Exercise 4.14(h) generalizes (4.7) (because of Exercise 4.14(d)), but also generalizes (4.11) (due to Exercise 4.14(f)). Similarly, Exercise 4.14(g) generalizes both (4.9) and (4.14). Similarly, Exercise 4.14(i) generalizes Exercise 4.4, and Exercise 4.14(j) generalizes Exercise 4.6.

Similarly, Exercise 4.3 is generalized by the following exercise:

**Exercise 4.15.** Let $G_1$, $G_2$, $H_1$ and $H_2$ be four finite groups. Let $\rho_1 : H_1 \to G_1$ and $\rho_2 : H_2 \to G_2$ be two group homomorphisms. These two homomorphisms clearly induce a group homomorphism $\rho_1 \times \rho_2 : H_1 \times H_2 \to G_1 \times G_2$. Let $U_1$ be a $\mathbb{C}H_1$-module, and $U_2$ be a $\mathbb{C}H_2$-module. Show that

$$\text{Ind}_{\rho_1 \times \rho_2} (U_1 \otimes U_2) \cong (\text{Ind}_{\rho_1} U_1) \otimes (\text{Ind}_{\rho_2} U_2)$$

as $\mathbb{C}[G_1 \times G_2]$-modules.

The Ind$_\rho$ and Res$_\rho$ operators behave “functorially” with respect to composition. Here is what this means:

**Exercise 4.16.** Let $G$, $H$ and $I$ be three finite groups. Let $\rho : H \to G$ and $\tau : I \to H$ be two group homomorphisms.

(a) We have Ind$_\rho$ Ind$_\tau$ $U \cong$ Ind$_{\rho \circ \tau}$ $U$ for every $\mathbb{C}I$-module $U$.

(b) We have Ind$_\rho$ Ind$_\tau$ $f =$ Ind$_{\rho \circ \tau}$ $f$ for every $f \in R_C(I)$.

(c) We have Res$_\rho$, Res$_\tau$ $V =$ Res$_{\rho \circ \tau}$ $V$ for every $\mathbb{C}G$-module $V$.

(d) We have Res$_\rho$, Res$_\tau$ $f =$ Res$_{\rho \circ \tau}$ $f$ for every $f \in R_C(G)$.

Exercise 4.16(a), of course, generalizes Exercise 4.2.

4.1.7. **Semidirect products.** Recall that a semidirect product is a group $G \ltimes K$ having two subgroups $G, K$ with

- $K \triangleleft (G \ltimes K)$ is a normal subgroup,
- $G \ltimes K = GK = KG$, and
- $G \cap K = \{e\}$.

In this setting one has two interesting adjoint constructions, applied in Section 4.5.

**Proposition 4.17.** Fix a $\mathbb{C}[G \ltimes K]$-module $V$.

(i) For any $\mathbb{C}G$-module $U$, one has $\mathbb{C}[G \ltimes K]$-module structure

$$\Phi(U) := U \otimes V,$$

determined via

$$k(u \otimes v) = u \otimes k(v),$$

$$g(u \otimes v) = g(u) \otimes g(v).$$

(ii) For any $\mathbb{C}[G \ltimes K]$-module $W$, one has $\mathbb{C}G$-module structure

$$\Psi(W) := \text{Hom}_{\mathbb{C}K}(\text{Res}_K^{G \ltimes K} V, \text{Res}_K^{G \ltimes K} W),$$

determined via $g(\varphi) = g \circ \varphi \circ g^{-1}$.

(iii) The maps

$$\mathbb{C}G - \text{mods} \xrightarrow{\Phi} \mathbb{C}[G \ltimes K] - \text{mods}$$

The crucial observation is now that if $(I_i, (a_i)_{i \in I_i}, (f_i)_{i \in I_i})$ is a finite dual generating system for a vector space $V$, and if $T$ is an endomorphism of $V$, then

$$\text{trace } T = \sum_{i \in I} f_i (T a_i).$$

Prove this!
are adjoint in the sense that one has an isomorphism
\[
\Hom_{CG}(U, \Psi(W)) \xrightarrow{\sim} \Hom_{C[G\times K]}(\Phi(U), W)
\]
\[
\Hom_{CG}(U, \Hom_{CK}(\Res_{K}^{G\times K} V, \Res_{K}^{G\times K} W)) \xrightarrow{\sim} \Hom_{C[G\times K]}(U \otimes V, W)
\]
\[\varphi \mapsto \varphi(u \otimes v) := \varphi(u)(v)\]

(iv) One has a \(CG\)-module isomorphism
\[
(\Psi \circ \Phi)(U) \cong U \otimes \End_{CK}(\Res_{K}^{G\times K} V).
\]

In particular, if \(\Res_{K}^{G\times K} V\) is a simple \(CK\)-module, then \((\Psi \circ \Phi)(U) \cong U\).

Proof. These are mostly straightforward exercises in the definitions. To check assertion (iv), for example, note that \(K\) acts only in the right tensor factor in \(\Res_{K}^{G\times K} (U \otimes V)\), and hence as \(CG\)-modules one has
\[
(\Psi \circ \Phi)(U) = \Hom_{CK}(\Res_{K}^{G\times K} V, \Res_{K}^{G\times K} U \otimes V)
\]
\[
= \Hom_{CK}(\Res_{K}^{G\times K} V, U \otimes \Res_{K}^{G\times K} V)
\]
\[
= U \otimes \Hom_{CK}(\Res_{K}^{G\times K} V, \Res_{K}^{G\times K} V)
\]
\[
= U \otimes \End_{CK}(\Res_{K}^{G\times K} V)
\]

\[\square\]

4.2. Three towers of groups. Here we consider three towers of groups
\[
G_{\ast} = (G_{0} < G_{1} < G_{2} < G_{3} < \cdots)
\]

where either
- \(G_{n} = \mathfrak{S}_{n}\), the symmetric group\(^{152}\), or
- \(G_{n} = \mathfrak{S}_{n}[\Gamma]\), the wreath product of the symmetric group with some arbitrary finite group \(\Gamma\), or
- \(G_{n} = GL_n(\mathbb{F}_q)\), the finite general linear group\(^{153}\).

Here the wreath product \(\mathfrak{S}_n[\Gamma]\) can be thought of informally as the group of *monomial* \(n \times n\) matrices whose nonzero entries lie in \(\Gamma\), that is, \(n \times n\) matrices having exactly one nonzero entry in each row and column, and that entry is an element of \(\Gamma\). E.g.,
\[
\begin{pmatrix}
0 & g_2 & 0 \\
g_1 & 0 & 0 \\
0 & g_3 & 0 \\
\end{pmatrix}
\begin{pmatrix}
0 & 0 & g_6 \\
0 & g_5 & 0 \\
g_4 & 0 & 0 \\
\end{pmatrix}
= \begin{pmatrix}
0 & g_2 g_5 & 0 \\
0 & 0 & g_1 g_6 \\
g_3 g_4 & 0 & 0 \\
\end{pmatrix}.
\]

More formally, \(\mathfrak{S}_n[\Gamma]\) is the semidirect product \(\mathfrak{S}_n \rtimes \Gamma^n\) in which \(\mathfrak{S}_n\) acts on \(\Gamma^n\) via \(\sigma(\gamma_1, \ldots, \gamma_n) = (\gamma_{\sigma^{-1}(1)}, \ldots, \gamma_{\sigma^{-1}(n)})\).

For each of the three towers \(G_{\ast}\), there are embeddings \(G_i \times G_j \hookrightarrow G_{i+j}\) and we introduce maps \(\ind^{i+j}_{i,j}\) taking \(\mathbb{C}[G_i \times G_j]\)-modules to \(\mathbb{C}G_{i+j}\)-modules, as well as maps \(\res^{i+j}_{i,j}\) carrying modules in the reverse direction which are adjoint:

\[
\Hom_{\mathbb{C}G_{i+j}}(\ind^{i+j}_{i,j} U, V) = \Hom_{\mathbb{C}[G_i \times G_j]}(U, \res^{i+j}_{i,j} V).
\]

**Definition 4.18.** For \(G_n = \mathfrak{S}_n\), one embeds \(\mathfrak{S}_i \times \mathfrak{S}_j\) into \(\mathfrak{S}_{i+j}\) as the permutations that permute \(\{1, 2, \ldots, i\}\) and \(\{i+1, i+2, \ldots, i+j\}\) separately. Here one defines
\[
\ind^{i+j}_{i,j} := \Ind_{\mathfrak{S}_i \times \mathfrak{S}_j}^{\mathfrak{S}_{i+j}},
\]
\[
\res^{i+j}_{i,j} := \Res_{\mathfrak{S}_i \times \mathfrak{S}_j}^{\mathfrak{S}_{i+j}}.
\]

\(^{152}\)The symmetric group \(\mathfrak{S}_0\) is the group of all permutations of the empty set \(\{1, 2, \ldots, 0\} = \emptyset\). It is a trivial group. (Note that \(\mathfrak{S}_1\) is also a trivial group.)

\(^{153}\)The group \(GL_0(\mathbb{F}_q)\) is a trivial group, consisting of the empty \(0 \times 0\) matrix.
For $G_n = \mathfrak{S}_n[\Gamma]$, similarly embed $\mathfrak{S}_i[\Gamma] \times \mathfrak{S}_j[\Gamma]$ into $\mathfrak{S}_{i+j}[\Gamma]$ as block monomial matrices whose two diagonal blocks have sizes $i, j$ respectively, and define

\[
\text{ind}_{i,j}^{i+j} := \text{Ind}_{\mathfrak{S}_i[\Gamma] \times \mathfrak{S}_j[\Gamma]}^{\mathfrak{S}_{i+j}[\Gamma]},
\]

\[
\text{res}_{i,j}^{i+j} := \text{Res}_{\mathfrak{S}_i[\Gamma] \times \mathfrak{S}_j[\Gamma]}^{\mathfrak{S}_{i+j}[\Gamma]}.
\]

For $G_n = GL_n(\mathbb{F}_q)$, which we will denote just $GL_n$, similarly embed $GL_i \times GL_j$ into $GL_{i+j}$ as block diagonal matrices whose two diagonal blocks have sizes $i, j$ respectively. However, one also introduces as an intermediate the parabolic subgroup $P_{i,j}$ consisting of the block upper-triangular matrices of the form

\[
\begin{pmatrix}
g_i & \ell \\
0 & g_j
\end{pmatrix}
\]

where $g_i, g_j$ lie in $GL_i, GL_j$, respectively, and $\ell$ in $\mathbb{F}_q^{i \times j}$ is arbitrary. One has a quotient map $P_{i,j} \rightarrow GL_i \times GL_j$

whose kernel $K_{i,j}$ is the set of matrices of the form

\[
\begin{pmatrix}
I_i & \ell \\
0 & I_j
\end{pmatrix}
\]

with $\ell$ again arbitrary. Here one defines

\[
\text{ind}_{i,j}^{i+j} := \text{Ind}_{P_{i,j}}^{GL_{i+j}} \text{Ind}_{GL_i \times GL_j}^{P_{i,j}},
\]

\[
\text{res}_{i,j}^{i+j} := \left(\text{Res}_{P_{i,j}}^{GL_{i+j}}(-)^{K_{i,j}}\right).
\]

In the case $G_n = GL_n$, the operation $\text{ind}_{i,j}^{i+j}$ is sometimes called parabolic induction or Harish-Chandra induction. The operation $\text{res}_{i,j}^{i+j}$ is essentially just the $K_{i,j}$-fixed point construction $V \mapsto V^{K_{i,j}}$. However writing it as the above two-step composite makes it more obvious, (via (4.7) and (4.11)) that $\text{res}_{i,j}^{i+j}$ is again adjoint to $\text{ind}_{i,j}^{i+j}$.

**Definition 4.19.** For each of the three towers $G_*$, define a graded $\mathbb{Z}$-module

\[
A := A(G_*) := \bigoplus_{n \geq 0} R(G_n)
\]

with a bilinear form $(\cdot, \cdot)_A$ whose restriction to $A_n := R(G_n)$ is the usual form $(\cdot, \cdot)_{G_n}$, and such that $\Sigma := \bigsqcup_{n \geq 0} \text{Irr}(G_n)$ gives an orthonormal $\mathbb{Z}$-basis. Notice that $A_0 = \mathbb{Z}$ has its basis element 1 equal to the unique irreducible character of the trivial group $G_0$.

Bearing in mind that $A_n = R(G_n)$ and

\[
A_i \otimes A_j = R(G_i) \otimes R(G_j) \cong R(G_i \times G_j)
\]

one then has candidates for product and coproduct defined by

\[
m := \text{ind}_{i,j}^{i+j} : A_i \otimes A_j \rightarrow A_{i+j}
\]

\[
\Delta := \bigoplus_{i+j = n} \text{res}_{i,j}^{i+j} : A_n \rightarrow \bigoplus_{i+j = n} A_i \otimes A_j.
\]

The coassociativity of $\Delta$ is an easy consequence of transitivity of the constructions of restriction and fixed points\(^{154}\). We could derive the associativity of $m$ from the transitivity of induction and inflation, but this would be more complicated\(^{155}\); we will instead prove it differently.

We first show that the maps $m$ and $\Delta$ are adjoint with respect to the forms $(\cdot, \cdot)_A$ and $(\cdot, \cdot)_{A \otimes A}$. In fact, if $U, V, W$ are modules over $CG_i, CG_j, CG_{i+j}$, respectively, then we can write the $\mathbb{C}[G_i \times G_j]$-module $\text{res}_{i,j}^{i+j} W$ as a direct sum $\bigoplus_k X_k \otimes Y_k$ with $X_k$ being $CG_i$-modules and $Y_k$ being $CG_j$-modules; we then have

\[
\text{res}_{i,j}^{i+j} \chi W = \sum_k \chi X_k \otimes \chi Y_k.
\]

\(^{154}\)More precisely, using this transitivity, it is easily reduced to proving that $K_{i+j,k} \cdot (K_{i,j} \times \{I_k\}) = K_{i+j+k} \cdot (\{I_i\} \times K_{j,k})$ (an equality between subgroups of $GL_{i+j+k}$). This equality can be proven by realizing that both of its sides equal the set of all block matrices of the form

\[
\begin{pmatrix}
I_i & \ell & m \\
0 & I_j & n \\
0 & 0 & I_k
\end{pmatrix}
\]

with $\ell, m$ and $n$ being matrices of sizes $i \times j, i \times k$ and $j \times k$, respectively.

\(^{155}\)See Exercise 4.31(c) for such a derivation.
and
\[(m(\chi_U \otimes \chi_V), \chi_W)_A = (\text{ind}_{i+j}^{i} (\chi_U \otimes \chi_V), \chi_W)_A = (\text{ind}_{i+j}^{i} (\chi_U \otimes \chi_V), \chi_W)_{G_{i+j}} = (\chi_U \otimes \chi_V, \sum_k \chi x_k \otimes \chi y_k)_{G_{i+j}} = \sum_k (\chi_U, \chi x_k)_{G_i} (\chi V, \chi y_k)_{G_j} = (\chi_U \otimes \chi V, \Delta (\chi_W))_{A \otimes A} \]

(the third equality sign follows by taking dimensions in (4.18) and recalling (4.1); the fourth equality sign follows from (4.19); the sixth one follows from (4.2)) and
\[(\chi_U \otimes \chi_V, \Delta (\chi_W))_{A \otimes A} = (\chi_U \otimes \chi_V, \text{res}_{i+j}^{i} \chi_W)_{A \otimes A} = (\chi_U \otimes \chi_V, \sum_k \chi x_k \otimes \chi y_k)_{A \otimes A} = \sum_k (\chi_U, \chi x_k)_{G_i} (\chi V, \chi y_k)_{G_j} ;\]

(the first equality sign follows by removing all terms in \(\Delta (\chi_W)\) whose scalar product with \(\chi_U \otimes \chi_V\) vanishes for reasons of gradedness; the second equality sign follows from (4.19)), which in comparison yield
\[(m(\chi_U \otimes \chi_V), \chi_W)_A = (\chi_U \otimes \chi_V, \Delta (\chi_W))_{A \otimes A}, \]

thus showing that \(m\) and \(\Delta\) are adjoint maps. Therefore, \(m\) is associative (since \(\Delta\) is coassociative).

Endowing \(A = \bigoplus_{n \geq 0} R(G_n)\) with the obvious unit and counit maps, it thus becomes a graded, finite-type \(\mathbb{Z}\)-algebra and \(\mathbb{Z}\)-coalgebra.

The next section addresses the issue of why they form a bialgebra. However, assuming this for the moment, it should be clear that each of these algebras \(A\) is a PSH having \(\Sigma = \bigcup_{n \geq 0} \text{Irr}(G_n)\) as its PSH-basis. \(\Sigma\) is self-dual because \(m\), \(\Delta\) are defined by adjoint maps, and it is positive because \(m\), \(\Delta\) take irreducible representations to genuine representations not just virtual ones, and hence have characters which are nonnegative sums of irreducible characters.

**Exercise 4.20.** Let \(i, j\) and \(k\) be three nonnegative integers. Let \(U\) be a \(\mathbb{C}G_i\)-module, let \(V\) be a \(\mathbb{C}G_j\)-module, and let \(W\) be a \(\mathbb{C}G_k\)-module. Show that there are canonical \(\mathbb{C}[G_i \times G_j \times G_k]\)-module isomorphisms
\[\text{Ind}_{G_1 \times G_2 \times G_3}^{G_1 \times G_2} (U \otimes V) \otimes W) \cong \text{Ind}_{G_1 \times G_2 \times G_3}^{G_1 \times G_2 \times G_3} (U \otimes V \otimes W) \cong \text{Ind}_{G_1 \times G_2 \times G_3}^{G_1 \times G_2 \times G_3} (U \otimes \text{Ind}_{G_1 \times G_2 \times G_3}^{G_1 \times G_2 \times G_3} (V \otimes W)).\]

(Similar statements hold for the other two towers of groups and their respective ind functors, although the one for the \(GL_n\) tower is harder to prove. See Exercise 4.31(a) for a more general result.)

### 4.3. Bialgebra and double cosets.
To show that the algebra and coalgebras \(A = A(G_n)\) are bialgebras, the central issue is checking the pentagonal diagram in (1.8), that is, as maps \(A \otimes A \to A \otimes A\), one has
\[\Delta \circ m = (m \otimes m) \circ (\text{id} \otimes T \otimes \text{id}) \circ (\Delta \otimes \Delta).\]

In checking this, it is convenient to have a lighter notation for various subgroups of the groups \(G_n\) corresponding to compositions \(\alpha\).

**Definition 4.21.**
(a) An almost-composition is a (finite) tuple \(\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_l)\) of nonnegative integers. Its length is defined to be \(l\) and denoted by \(l(\alpha)\); its size is defined to be \(\alpha_1 + \alpha_2 + \cdots + \alpha_l\) and denoted by \(|\alpha|\); its parts are its entries \(\alpha_1, \alpha_2, \ldots, \alpha_l\). The almost-compositions of size \(n\) are called the almost-compositions of \(n\).

(b) A composition is a finite tuple of positive integers. Of course, any composition is an almost-composition, and so all notions defined for almost-compositions (like size and length) make sense for compositions.

Note that any partition of \(n\) (written without trailing zeroes) is a composition of \(n\). We write \(\emptyset\) (and sometimes, sloppily, \((0)\) when there is no danger of mistaking it for the almost-composition \((0)\)) for the empty composition \((\emptyset)\).
Definition 4.22. Given an almost-composition \( \alpha = (\alpha_1, \ldots, \alpha_\ell) \) of \( n \), define a subgroup

\[
G_\alpha \cong G_{\alpha_1} \times \cdots \times G_{\alpha_\ell} < G_n
\]

via the block-diagonal embedding with diagonal blocks of sizes \( (\alpha_1, \ldots, \alpha_\ell) \). This \( G_\alpha \) is called a Young subgroup \( \mathfrak{S}_\alpha \) when \( G_n = \mathfrak{S}_n \) and a Levi subgroup when \( G_n = GL_n \). In the case when \( G_n = \mathfrak{S}_n[\Gamma] \), we also denote \( G_\alpha \) by \( \mathfrak{S}_\alpha[\Gamma] \). In the case where \( G_n = GL_n \), also define the parabolic subgroup \( P_\alpha \) to be the subgroup of \( G_n \) consisting of block-upper triangular matrices whose diagonal blocks have sizes \( (\alpha_1, \ldots, \alpha_\ell) \), and let \( K_\alpha \) be the kernel of the obvious surjection \( P_\alpha \to G_\alpha \) which sends a block upper-triangular matrix to the tuple of its diagonal blocks whose sizes are \( \alpha_1, \alpha_2, \ldots, \alpha_\ell \). Notice that \( P_{(i,j)} = P_{i,j} \) for any \( i \) and \( j \) with \( i + j = n \); similarly, \( K_{(i,j)} = K_{i,j} \) for any \( i \) and \( j \) with \( i + j = n \). We will also abbreviate \( G_{(i,j)} = G_i \times G_j \) by \( G_{i,j} \).

When \( (\alpha_1, \alpha_2, \ldots, \alpha_\ell) \) is an almost-composition, we abbreviate \( G_{(\alpha_1, \alpha_2, \ldots, \alpha_\ell)} \) (and similarly for the \( P's \)).

Definition 4.23. Let \( K \) and \( H \) be two groups, \( \tau : K \to H \) a group homomorphism, and \( U \) a \( CH \)-module. Then, \( U^\tau \) is defined as the \( CK \)-module with ground space \( U \) and action given by \( k \cdot u = \tau(k) \cdot u \) for all \( k \in K \) and \( u \in U \). \footnote{We have already met this \( CK \)-module \( U^\tau \) in Remark 4.13, where it was called \( Res_u U \).} This very simple construction generalizes the definition of \( U^g \) for an element \( g \in G \), where \( G \) is a group containing \( H \) as a subgroup; in fact, in this situation we have \( U^g = U^\tau \), where \( K = gH \) and \( \tau : K \to H \) is the map \( k \mapsto g^{-1}kg \).

Using homogeneity, checking the bialgebra condition (4.20) in the homogeneous component \( (A \otimes A)_n \) amounts to the following: for each pair of representations \( U_1, U_2 \) of \( G_{r_1}, G_{r_2} \) with \( r_1 + r_2 = n \), and for each \( (c_1, c_2) \) with \( c_1 + c_2 = n \), one must verify that

\[
\text{res}_{c_1,c_2}^n \left( \text{ind}_{r_1,r_2}^n (U_1 \otimes U_2) \right)
\]

\[
\cong \bigoplus_A \left( \text{ind}_{\alpha_{11},\alpha_{12}}^{r_1} \otimes \text{ind}_{\alpha_{21},\alpha_{22}}^{r_2} \right) \left( \text{res}_{\alpha_{11},\alpha_{12}}^{r_1,r_2} U_1 \otimes \text{res}_{\alpha_{21},\alpha_{22}}^{r_2} U_2 \right) \tau^{-1}_A
\]

(4.21)

where the direct sum is over all matrices \( A = \begin{bmatrix} \alpha_{11} & \alpha_{12} \\ \alpha_{21} & \alpha_{22} \end{bmatrix} \) in \( \mathbb{N}^{2 \times 2} \) with row sums \( (r_1, r_2) \) and column sums \( (c_1, c_2) \), and where \( \tau_A \) is the obvious isomorphism between the subgroups

\[
G_{\alpha_{11},\alpha_{12},\alpha_{21},\alpha_{22}} \cong G_{r_1,r_2}
\]

(4.22)

\[
G_{\alpha_{11},\alpha_{12},\alpha_{21},\alpha_{22}} \cong G_{c_1,c_2}
\]

(we are using the inverse \( \tau^{-1}_A \) of this isomorphism \( \tau_A \) to identify modules for the first subgroup with modules for the second subgroup, according to Definition 4.23).

As one might guess, (4.21) comes from the Mackey formula (Theorem 4.7), once one identifies the appropriate double coset representatives. This is just as easy to do in a slightly more general setting.

Definition 4.24. Given almost-compositions \( \alpha, \beta \) of \( n \) having lengths \( \ell, m \) and a matrix \( A \in \mathbb{N}^{\ell \times m} \) with row sums \( \alpha \) and column sums \( \beta \), define a permutation \( w_A \) in \( \mathfrak{S}_n \) as follows. Disjointly decompose \( [n] = \{1, 2, \ldots, n\} \) into consecutive intervals of numbers

\[
[n] = I_1 \sqcup \cdots \sqcup I_\ell
\]

\[
[n] = J_1 \sqcup \cdots \sqcup J_m
\]

such that \( |I_i| = \alpha_i, \ |J_j| = \beta_j \). For every \( j \in [m] \), disjointly decompose \( J_j \) into consecutive intervals of numbers \( J_{j,1} = J_{j,1} \sqcup J_{j,2} \sqcup \cdots \sqcup J_{j,\ell} \) such that every \( i \in [\ell] \) satisfies \( |J_{j,i}| = \alpha_{ij} \). For every \( i \in [\ell] \), disjointly decompose \( I_i \) into consecutive intervals of numbers \( I_{i,1} = I_{i,1} \sqcup I_{i,2} \sqcup \cdots \sqcup I_{i,m} \) such that every \( j \in [m] \) satisfies \( |I_{i,j}| = \alpha_{ij} \). Now, for every \( i \in [\ell] \) and \( j \in [m] \), let \( \pi_{i,j} \) be the increasing bijection from \( J_{j,i} \) to \( I_{i,j} \) (this is well-defined since these two sets both have cardinality \( \alpha_{ij} \)). The disjoint union of these bijections \( \pi_{i,j} \) over all \( i \) and \( j \) is a bijection \( [n] \to [n] \) (since the disjoint union of the sets \( J_{j,i} \) over all \( i \) and \( j \) is \( [n] \), and so is the disjoint union of the sets \( I_{i,j} \), that is, a permutation of \( [n] \); this permutation is what we call \( w_A \).

Example 4.25. Taking \( n = 9 \) and \( \alpha = (4, 5), \beta = (3, 4, 2) \), one has

\[
I_1 = \{1, 2, 3, 4\}, \quad I_2 = \{5, 6, 7, 8, 9\}
\]

\[
J_1 = \{1, 2, 3\}, \quad J_2 = \{4, 5, 6, 7\}, \quad J_3 = \{8, 9\}.
\]
Then one possible matrix $A$ having row and column sums $\alpha, \beta$ is $A = \begin{bmatrix} 2 & 2 & 0 \\ 1 & 2 & 2 \end{bmatrix}$, and its associated permutation $w_A$ written in two-line notation is
\[
\begin{pmatrix}
1 & 2 & 3 \\
1 & 2 & 5
\end{pmatrix}
\begin{pmatrix}
4 & 5 & 6 & 7 \\
3 & 4 & 6 & 7
\end{pmatrix}
\begin{pmatrix}
8 & 9 \\
8 & 9
\end{pmatrix}
\]
with vertical lines dividing the sets $J_j$ on top, and with elements of $I_i$ underlined $i$ times on the bottom.

Remark 4.26. Given almost-compositions $\alpha$ and $\beta$ of $n$ having lengths $\ell$ and $m$, and a permutation $w \in S_n$. It is easy to see that there exists a matrix $A \in \mathbb{N}^{\ell \times m}$ satisfying $w_A = w$ if and only if the restriction of $w$ to each $J_j$ and the restriction of $w^{-1}$ to each $I_i$ are increasing. In this case, the matrix $A$ is determined by $a_{ij} = |w(J_j) \cap I_i|$. 

Among our three towers $G_n$ of groups, the symmetric group tower $(G_n = S_n)$ is the simplest one. We will now see that it also embeds into the two others, in the sense that $S_n$ embeds into $S_n[\Gamma]$ for every $\Gamma$ and into $GL_n(F_q)$ for every $q$.

First, for every $n \in \mathbb{N}$ and any group $\Gamma$, we embed the group $S_n$ into $S_n[\Gamma]$ by means of the canonical embedding $S_n \to S_n \times \Gamma^n = S_n[\Gamma]$. If we regard elements of $S_n[\Gamma]$ as $n \times n$ monomial matrices with nonzero entries in $\Gamma$, then this boils down to identifying every $\pi \in S_n$ with the permutation matrix of $\pi$ (in which the 1’s are read as the neutral element of $\Gamma$). If $\alpha$ is an almost-composition of $n$, then this embedding $S_n \to S_n[\Gamma]$ makes the subgroup $S_\alpha$ of $S_n$ become a subgroup of $S_n[\Gamma]$, more precisely a subgroup of $S_\alpha[\Gamma] < S_n[\Gamma]$.

For every $n \in \mathbb{N}$ and every $q$, we embed the group $S_n$ into $GL_n(F_q)$ by identifying every permutation $\pi \in S_n$ with its permutation matrix in $GL_n(F_q)$. If $\alpha$ is an almost-composition of $n$, then this embedding makes the subgroup $S_\alpha$ of $S_n$ become a subgroup of $GL_n(F_q)$. If we let $G_n = GL_n(F_q)$, then $S_\alpha < G_\alpha < P_\alpha$.

The embeddings we have just defined commute with the group embeddings $G_n < G_{n+1}$ on both sides.

Proposition 4.27. The permutations $\{w_A\}$ as $A$ runs over all matrices in $\mathbb{N}^{\ell \times m}$ having row, column sums $\alpha, \beta$ give a system of double coset representatives for $S_\alpha \backslash S_n/S_\beta$.

\[
\begin{align*}
S_\alpha \backslash S_n[\Gamma] &\backslash S_\beta[\Gamma] \\
\alpha \backslash GL_n/\beta
\end{align*}
\]

Proof. First note that double coset representatives for $S_\alpha \backslash S_n/S_\beta$ should also provide double coset representatives for $S_\alpha[\Gamma] \backslash S_n[\Gamma]/S_\beta[\Gamma]$, since $S_\alpha[\Gamma] = S_\alpha \Gamma^n = \Gamma^n S_\alpha$.

We give an algorithm to show that every double coset $S_\alpha w S_\beta$ contains some $w_A$. Start by altering $w$ within its coset $wS_\beta$, that is, by permuting the positions within each set $J_j$, to obtain a representative $w'$ for $wS_\beta$ in which each set $w'(J_j)$ appears in increasing order in the second line of the two-line notation for $w'$. Then alter $w'$ within its coset $S_\alpha w'$, that is, by permuting the values within each set $I_i$, to obtain a representative $w_A$ having the elements of each set $I_i$ appearing in increasing order in the second line; because the values within each set $I_i$ are consecutive, this alteration will not ruin the property that one had each set $w'(J_j)$ appearing in increasing order. For example, one might have
\[
\begin{align*}
w &= \begin{pmatrix} 1 & 2 & 3 \\ 4 & 5 & 6 \\ 7 & 8 & 9 \end{pmatrix} \\
w' &= \begin{pmatrix} 1 & 2 & 3 \\ 2 & 4 & 8 \\ 1 & 3 & 5 & 9 \\ 6 & 7 & \end{pmatrix} \in wS_\beta \\
w_A &= \begin{pmatrix} 1 & 2 & 3 \\ 1 & 2 & 5 & 7 \\ 3 & 4 & 6 & 7 \\ 8 & 9 & \end{pmatrix} \in S_\alpha w' \subset S_\alpha w' S_\beta = S_\alpha w S_\beta
\end{align*}
\]

Next note that $S_\alpha w_A S_\beta = S_\alpha w S_\beta$ implies $A = B$, since the quantities $a_{ij}(w) := |w(J_j) \cap I_i|$ are easily seen to be constant on double cosets $S_\alpha w S_\beta$. 

A similar argument shows that $P_A P_B = P_A P_B$ implies $A = B$: for $g$ in $GL_n$, the rank $r_{ij}(g)$ of the matrix obtained by restricting $g$ to rows $I_i \sqcup I_{i+1} \sqcup \cdots \sqcup I_\ell$ and columns $J_1 \sqcup J_2 \sqcup \cdots \sqcup J_j$ is constant on double cosets $P_A P_B$, and for a permutation matrix $w$ one can recover $a_{i,j}(w)$ from the formula

$$a_{i,j}(w) = r_{i,j}(w) - r_{i,j-1}(w) - r_{i+1,j}(w) + r_{i+1,j-1}(w).$$

Thus it only remains to show that every double coset $P_A P_B$ contains some $w_A$. Since $S_n < P_n$, and we have seen already that every double coset $S_n w S_n$ contains some $w_A$, it suffices to show that every double coset $P_A P_B$ contains some permutation $w$. However, we claim that this is already true for the smaller double cosets $B g B$ where $B = P_i$ is the Borel subgroup of upper triangular invertible matrices, that is, one has the usual Bruhat decomposition

$$GL_n = \bigcup_{w \in S_n} B w B.$$

To prove this decomposition, we show how to find a permutation $w$ in each double coset $B g B$. The freedom to alter $g$ within its coset $g B$ allows one to scale columns and add scalar multiples of earlier columns to later columns. We claim that using such column operations, one can always find a representative $g'$ for coset $g B$ in which

- the bottommost nonzero entry of each column is 1 (call this entry a pivot),
- the entries to right of each pivot within its row are all 0, and
- there is one pivot in each row and each column, so that their positions are the positions of the 1's in some permutation matrix $w$.

In fact, we will see below that $B g B = B w B$ in this case. The algorithm which produces $g'$ from $g$ is simple: starting with the leftmost column, find its bottommost nonzero entry, and scale the column to make this entry a 1, creating the pivot in this column. Now use this pivot to clear out all entries in its row to its right, using column operations that subtract multiples of this column from later columns. Having done this, move on to the next column to the right, and repeat, scaling to create a pivot, and using it to eliminate entries to its right.\footnote{To see that this works, we need to check three facts:

(a) We will find a nonzero entry in every column during our algorithm.

(b) Our column operations preserve the zeroes lying to the right of already existing pivots.

(c) Every row contains exactly one pivot at the end of the algorithm.

But fact (a) simply says that our matrix can never have an all-zero column during the algorithm; this is clear (since the rank of the matrix remains constant during the algorithm and was $n$ at its beginning). Fact (b) holds because all our operations either scale columns (which clearly preserves zero entries) or subtract a multiple of the column $c$ containing the current pivot from a later column $d$ (which will preserve every zero lying to the right of an already existing pivot, because any already existing pivot must lie in a column $b < c$ and therefore both columns $c$ and $d$ have zeroes in its row). Fact (c) follows from noticing that there are $n$ pivots altogether at the end of the algorithm, but no row can contain two of them (since the entries to the right of a pivot in its row are 0).}

For example, the typical matrix $g$ lying in the double coset $B w B$ where

$$w = \begin{pmatrix} 1 & 2 & 3 & & & & 4 & 5 & 6 & 7 & & 8 & 9 \\ 4 & 8 & 2 & & & & 5 & 3 & 9 & 1 & & 7 & 6 \end{pmatrix}$$

from before is one that can be altered within its coset $g B$ to look like this:

$$g' = \begin{pmatrix} * & * & * & * & * & * & 1 & 0 & 0 \\ * & * & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\ * & * & 0 & * & 1 & 0 & 0 & 0 & 0 \\ 0 & * & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\ 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & * & 0 & 0 & 0 & * & 0 & 1 & 0 \\ 0 & * & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \end{pmatrix} \in g B.$$

Having found this $g'$ in $g B$, a similar algorithm using left multiplication by $B$ shows that $w$ lies in $B g' B = B w B$. This time no scalings are required to create the pivot entries: starting with the bottom row, one uses its pivot to eliminate all the entries above it in the same column (shown by stars * above) by adding
multiples of the bottom row to higher rows. Then do the same using the pivot in the next-to-bottom row, etc. The result is the permutation matrix for \( w \).

**Remark 4.28.** The Bruhat decomposition \( GL_n = \bigsqcup_{w \in \mathfrak{S}_n} BwB \) is related to the so-called LPU factorization – one of a myriad of matrix factorizations appearing in linear algebra.\(^{158}\) It is actually a fairly general phenomenon, and requires neither the finiteness of \( \mathbb{F} \), nor the invertibility, nor even the squareness of the matrices (see Exercise 4.29(b) for an analogue holding in a more general setup).

**Exercise 4.29.** Let \( \mathbb{F} \) be any field.

(a) For any \( n \in \mathbb{N} \) and any \( A \in GL_n(\mathbb{F}) \), prove that there exist a lower-triangular matrix \( L \in GL_n(\mathbb{F}) \), an upper-triangular matrix \( U \in GL_n(\mathbb{F}) \) and a permutation matrix \( P \in \mathfrak{S}_n \subset GL_n(\mathbb{F}) \) (here, we identify permutations with the corresponding permutation matrices) such that \( A = LPU \). This can be derived from the all of these bookkeeping details to the reader to check.\(^{160}\)

(b) Let \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \). Let \( F_{n,m} \) denote the set of all \( n \times m \)-matrices \( B \in \{0,1\}^{n \times m} \) such that each row of \( B \) contains at most one 1 and each column of \( B \) contains at most one 1. We regard \( F_{n,m} \) as a subset of \( \mathbb{F}^{n \times m} \) by means of regarding \( \{0,1\} \) as a subset of \( \mathbb{F} \).

For every \( k \in \mathbb{N} \), we let \( B_k \) denote the subgroup of \( GL_k(\mathbb{F}) \) consisting of all upper-triangular matrices.

Prove that

\[
\mathbb{F}^{n \times m} = \bigsqcup_{f \in F_{n,m}} B_nfB_m.
\]

**Corollary 4.30.** For each of the three towers of groups \( G_s \), the product and coproduct structures on \( A = A(G_s) \) endow it with a bialgebra structure, and hence they form PSH’s.

**Proof.** The first two towers \( G_n = \mathfrak{S}_n \) and \( G_n = \mathfrak{S}_n[\Gamma] \) have product, coproduct defined by induction, restriction along embeddings \( G_i \times G_j < G_{i+j} \). Hence the desired bialgebra equality (4.21) follows from Mackey’s Theorem 4.7, taking \( G = G_n, H = G_{(r_1,r_2)}, K = G_{(c_1,c_2)}, U = U_1 \otimes U_2 \) with double coset representatives\(^{159}\)

\[
\{g_1, \ldots, g_t\} = \{w_{A^t} : A \in \mathbb{N}^{2 \times 2}, A \text{ has row sums } (r_1, r_2) \text{ and column sums } (c_1, c_2)\}
\]

and checking for a given double coset

\[
\begin{align*}
KgH &= (G_{c_1,c_2})w_{A^t}(G_{r_1,r_2}), \\
H \cap K &= w_{A^t}(G_{r_1,r_2}) \cap G_{c_1,c_2},
\end{align*}
\]

indexed by a matrix \( A \in \mathbb{N}^{2 \times 2} \) with row sums \( (r_1, r_2) \) and column sums \( (c_1, c_2) \), that the two subgroups appearing on the left in (4.22) are exactly

\[
H \cap K^{w_{A^t}} = G_{r_1,r_2} \cap (G_{c_1,c_2})^{w_{A^t}},
\]

\[
w_{A^t}H \cap K = w_{A^t}(G_{r_1,r_2}) \cap G_{c_1,c_2},
\]

respectively. One should also apply (4.6) and check that the isomorphism \( \tau_A \) between the two subgroups in (4.22) is the conjugation isomorphism by \( w_{A^t} \) (that is, \( \tau_A(g) = w_{A^t}gw_{A^t}^{-1} \) for every \( g \in H \cap K^{w_{A^t}} \)). We leave all of these bookkeeping details to the reader to check.\(^{160}\)

\(^{158}\)Specifically, an LPU factorization of a matrix \( A \in GL_n(\mathbb{F}) \) (for an arbitrary field \( \mathbb{F} \)) means a way to write \( A \) as a product \( A = LPU \) with \( L \in GL_m(\mathbb{F}) \) being lower-triangular, \( U \in GL_n(\mathbb{F}) \) being upper-triangular, and \( P \in \mathfrak{S}_n \subset GL_n(\mathbb{F}) \) being a permutation matrix. Such a factorization always exists (although it is generally not unique). This can be derived from the Bruhat decomposition (see Exercise 4.29(a) for a proof). See also [159] for related discussion.

\(^{159}\)Proposition 4.27 gives as a system of double coset representatives for \( G_{(c_1,c_2)}/G_n/G_{(r_1,r_2)} \) the elements

\[
\{w_A : A \in \mathbb{N}^{2 \times 2}, A \text{ has row sums } (c_1, c_2) \text{ and column sums } (r_1, r_2)\} = \{w_{A^t} : A \in \mathbb{N}^{2 \times 2}, A \text{ has row sums } (r_1, r_2) \text{ and column sums } (c_1, c_2)\}
\]

where \( A^t \) denotes the transpose matrix of \( A \).

\(^{160}\)It helps to recognize \( w_{A^t} \) as the permutation written in two-line notation as

\[
\begin{bmatrix}
1 & 2 & \ldots & a_{11} & a_{11}+1 & a_{11}+2 & \ldots & r_1 & r_1+1 & r_1+2 & \ldots & a_{22}' & a_{22}'+1 & a_{22}'+2 & \ldots & n \\
1 & 2 & \ldots & a_{11} & c_1+1 & c_1+2 & \ldots & a_{22}' & a_{11}+1 & a_{11}+2 & \ldots & c_1 & a_{22}'+1 & a_{22}'+2 & \ldots & n
\end{bmatrix},
\]

where \( a_{22}' = r_1 + a_{21} = c_1 + a_{12} = n - a_{22} \). In matrix form, \( w_{A^t} \) is the block matrix

\[
\begin{bmatrix}
I_{a_{11}} & 0 & 0 & 0 \\
0 & I_{a_{21}} & 0 & 0 \\
0 & 0 & I_{a_{12}} & 0 \\
0 & 0 & 0 & I_{a_{22}}
\end{bmatrix}.
\]
For the tower with $G_n = GL_n$, there is slightly more work to be done to check the equality (4.21). Via Mackey’s Theorem 4.7 and Proposition 4.27, the left side is

$\text{res}^{n}_{c_1,c_2} \left( \text{ind}^{n}_{r_1,r_2} (U_1 \otimes U_2) \right)
\begin{equation}
= \left( \text{Res}^{G_n}_{P_{r_1,r_2}} \text{Ind}^{G_n}_{P_{r_1,r_2}} \text{Infl}^{P_{r_1,r_2}}_{G_{r_1,r_2}} (U_1 \otimes U_2) \right)_{K_{c_1,c_2}}
\end{equation}
\begin{equation}
= \bigoplus_{A} \left( \text{Ind}^{P_{r_1,r_2}}_{A} \text{Ind}^{P_{r_1,r_2}}_{P_{r_1,r_2} \cap P_{c_1,c_2}} \left( \left( \text{Res}^{P_{r_1,r_2}}_{P_{r_1,r_2} \cap P_{c_1,c_2}} \text{Infl}^{P_{r_1,r_2}}_{G_{r_1,r_2}} (U_1 \otimes U_2) \right)^{\tau_A^{-1}} \right) \right)_{K_{c_1,c_2}}
\end{equation}
(4.23)

where $A$ runs over the usual $2 \times 2$ matrices. The right side is a direct sum over this same set of matrices $A$:

$\bigoplus_{A} \left( \text{ind}^{G_{c_1,c_2}}_{a_{11},a_{21}} \otimes \text{ind}^{G_{c_2}}_{a_{12},a_{22}} \right) \left( \left( \text{res}^{a_{11},a_{12}}_{a_{11},a_{12}} U_1 \otimes \text{res}^{a_{21},a_{22}}_{a_{21},a_{22}} U_2 \right)^{\tau_A^{-1}} \right)
\begin{equation}
= \bigoplus_{A} \left( \text{Ind}^{G_{c_1}}_{a_{11},a_{21}} \otimes \text{Ind}^{G_{c_2}}_{a_{12},a_{22}} \circ \left( \text{Ind}^{P_{a_{11},a_{21}}}_{G_{a_{11},a_{21}}} \otimes \text{Ind}^{P_{a_{12},a_{22}}}_{G_{a_{12},a_{22}}} \right) \left( \left( \text{Res}^{G_{c_1}}_{P_{a_{11},a_{21}}} U_1 \right)_{K_{a_{11},a_{12}}} \otimes \left( \text{Res}^{G_{c_2}}_{P_{a_{21},a_{22}}} U_2 \right)_{K_{a_{21},a_{22}}} \right)^{\tau_A^{-1}} \right)
\end{equation}
\begin{equation}
= \bigoplus_{A} \text{Ind}^{G_{c_1,c_2}}_{P_{a_{11},a_{21}} \times P_{a_{12},a_{22}}} \left( \left( \text{Res}^{G_{c_1,c_2}}_{P_{a_{11},a_{12}} \times P_{a_{21},a_{22}}} (U_1 \otimes U_2) \right)_{K_{a_{11},a_{12}} \times K_{a_{21},a_{22}}} \right)^{\tau_A^{-1}}
\end{equation}
(4.24)

(by (4.6), (4.15) and their obvious analogues for restriction and inflation). Thus it suffices to check for each $2 \times 2$ matrix $A$ that any $CG_{c_1,c_2}$-module of the form $V_1 \otimes V_2$ has the same inner product with the $A$-summands of (4.23) and (4.24).

Abbreviate $w := w_{A^t}$ and $\tau := \tau_A^{-1}$.

Notice that $wP_{r_1,r_2}$ is the group of all matrices having the block form

$\begin{bmatrix}
g_{11} & h & i & j \\
0 & g_{21} & 0 & k \\
d & e & g_{12} & \ell \\
0 & f & 0 & g_{22}
\end{bmatrix}$

in which the diagonal blocks $g_{ij}$ for $i,j = 1,2$ are invertible of size $a_{ij} \times a_{ij}$, while the blocks $h, i, j, k, \ell, d, e, f$ are all arbitrary matrices$^{161}$ of the appropriate (rectangular) block sizes. Hence, $wP_{r_1,r_2} \cap P_{c_1,c_2}$ is the group of all matrices having the block form

$\begin{bmatrix}
g_{11} & h & i & j \\
0 & g_{21} & 0 & k \\
0 & 0 & g_{12} & \ell \\
0 & 0 & 0 & g_{22}
\end{bmatrix}$

in which the diagonal blocks $g_{ij}$ for $i,j = 1,2$ are invertible of size $a_{ij} \times a_{ij}$, while the blocks $h, i, j, k, \ell$ are all arbitrary matrices of the appropriate (rectangular) block sizes; then $wP_{r_1,r_2} \cap G_{c_1,c_2}$ is the subgroup where the blocks $i,j,k$ all vanish. The canonical projection $wP_{r_1,r_2} \cap P_{c_1,c_2} \rightarrow wP_{r_1,r_2} \cap G_{c_1,c_2}$ (obtained by restricting the projection $P_{c_1,c_2} \rightarrow G_{c_1,c_2}$) has kernel $wP_{r_1,r_2} \cap P_{c_1,c_2} \cap K_{c_1,c_2}$. Consequently,

$\left( wP_{r_1,r_2} \cap P_{c_1,c_2} \right) / \left( wP_{r_1,r_2} \cap P_{c_1,c_2} \cap K_{c_1,c_2} \right) = wP_{r_1,r_2} \cap G_{c_1,c_2}.$

Similarly,

$\left( P_{r_1,r_2} \cap P_{c_1,c_2} \right) / \left( P_{r_1,r_2} \cap P_{c_1,c_2} \cap K_{c_1,c_2} \right) = G_{r_1,r_2} \cap P_{c_1,c_2}.$

\(^{161}\)The blocks $i$ and $j$ have nothing to do with the indices $i, j$ in $g_{ij}$.
Computing first the inner product of $V_1 \otimes V_2$ with the $A$-summand of (4.23), and using adjointness properties, one gets

\[
\left( \left( \text{Res}_{P_{r_1}, r_2}^{P_{c_1}, c_2} \right. (U_1 \otimes U_2) \right)^\tau, \\
\text{Res}_{P_{r_1}, r_2}^{P_{c_1}, c_2} \left( V_1 \otimes V_2 \right) \right)_{P_{r_1}, r_2}^{P_{c_1}, c_2}
\]

(by (4.28) and (4.27)). One can compute this inner product by first recalling that $w P_{r_1, r_2}$ is the group of matrices having the block form (4.26) in which the diagonal blocks $g_{ij}$ for $i, j = 1, 2$ are invertible of size $a_{ij} \times a_{ij}$, while the blocks $h, i, j, k, \ell$ are all arbitrary matrices of the appropriate (rectangular) block sizes; then $w P_{r_1, r_2}$ is the subgroup where the blocks $i, j, k$ all vanish. The inner product above then becomes

\[
\frac{1}{|w P_{r_1, r_2} \cap P_{c_1, c_2}|} \sum_{(g_{ij})} \chi_U \left( \begin{array}{c} g_{11} \\ g_{12} \end{array} \right) \chi_U \left( \begin{array}{c} g_{21} \\ g_{22} \end{array} \right) \chi_V \left( \begin{array}{c} g_{11} \\ g_{21} \end{array} \right) \chi_V \left( \begin{array}{c} 0 \\ g_{22} \end{array} \right).
\]

If one instead computes the inner product of $V_1 \otimes V_2$ with the $A$-summand of (4.24), using adjointness properties and (4.13) one gets

\[
\left( \left( \text{Res}_{P_{a_{11}, a_{21}}^{a_{12}, a_{22}}}^{A_{a_{11}, a_{12}} \times A_{a_{21}, a_{22}}} (U_1 \otimes U_2) \right)^K_{a_{11}, a_{12}} \times K_{a_{21}, a_{22}} \right) \left( \text{Res}_{P_{a_{11}, a_{21}}^{a_{12}, a_{22}}}^{A_{a_{11}, a_{12}} \times A_{a_{21}, a_{22}}} (V_1 \otimes V_2) \right)^K_{a_{11}, a_{12}} \times K_{a_{21}, a_{22}} \right)_{G_{a_{11}, a_{12}} \times K_{a_{21}, a_{22}}}
\]

\[
= \frac{1}{|G_{a_{11}, a_{21}, a_{12}, a_{22}|} \sum_{(g_{ij})} \chi_U \left( \begin{array}{c} g_{11} \\ g_{12} \end{array} \right) \chi_U \left( \begin{array}{c} g_{21} \\ g_{22} \end{array} \right) \chi_V \left( \begin{array}{c} g_{11} \\ g_{21} \end{array} \right) \chi_V \left( \begin{array}{c} 0 \\ g_{22} \end{array} \right) \sum_{(h, \ell)} \chi_V \left( \begin{array}{c} 0 \\ g_{22} \end{array} \right) \chi_V \left( \begin{array}{c} 0 \\ g_{22} \end{array} \right).
\]

But this right hand side can be seen to equal (4.29), after one notes that

\[
|w P_{r_1, r_2} \cap P_{c_1, c_2}| = |G_{a_{11}, a_{21}, a_{12}, a_{22}|} \cdot |K_{a_{11}, a_{12}} \times K_{a_{21}, a_{22}}| \cdot |K_{a_{11}, a_{12}} \times K_{a_{21}, a_{22}}| \cdot \# \{ j \in \mathbb{F}_q^{a_{11} \times a_{22}} \}
\]

and that the summands in (4.29) are independent of the matrix $j$ in the summation.

We can also define a $\mathbb{C}$-vector space $A_{C_1}$ as the direct sum $\bigoplus_{n \geq 0} R_{C}(G_n)$. In the same way as we have made $A = \bigoplus_{n \geq 0} R(G_n)$ into a $\mathbb{Z}$-bialgebra, we can turn $A_{C_1} = \bigoplus_{n \geq 0} R_{C}(G_n)$ into a $\mathbb{C}$-bialgebra.$^{162}$ There is a $\mathbb{C}$-bilinear form $(\cdot, \cdot)_C$ on $A_{C_1}$ which can be defined either as the $\mathbb{C}$-bilinear extension of the $\mathbb{Z}$-bilinear form $(\cdot, \cdot)_A : A \times A \rightarrow \mathbb{Z}$ to $A_{C_1}$, or (equivalently) as the $\mathbb{C}$-bilinear form on $A_{C_1}$ which restricts to $(\cdot, \cdot)_{C_n}$ on every homogeneous component $R_{C}(G_n)$ and makes different homogeneous components mutually orthogonal. The obvious embedding of $A$ into the $\mathbb{C}$-bialgebra $A_{C_1}$ (obtained from the embeddings $R(G_n) \rightarrow R_{C}(G_n)$ for

$^{162}$The definitions of $m$ and $\Delta$ for this $\mathbb{C}$-bialgebra look the same as for $A$: For instance, $m$ is still defined to be $m^{i,j}_{i',j'}(a, b)$ on $(A_{C_1})_i \otimes (A_{C_1})_j$, where $m^{i,j}_{i',j'}(a, b)$ is defined by the same formulas as in Definition 4.18. However, the operators of induction, restriction, inflation and $K$-fixed space construction appearing in these formulas now act on class functions as opposed to modules.

The fact that these maps $m$ and $\Delta$ satisfy the axioms of a $\mathbb{C}$-bialgebra is easy to check: they are merely the $\mathbb{C}$-linear extensions of the maps $m$ and $\Delta$ of the $\mathbb{Z}$-bialgebra $A$ (this is because, for instance, induction of class functions and induction of modules are related by the identity (4.5)), and thus satisfy the same axioms as the latter.
all $n$) respects the bialgebra operations, and the $\mathbb{C}$-bialgebra $A_\mathbb{C}$ can be identified with $A \otimes \mathbb{C}$ (the result of extending scalars to $\mathbb{C}$ in $A$), because every finite group $G$ satisfies $R_C(G) \cong R(G) \otimes \mathbb{C}$. The embedding of $A$ into $A_\mathbb{C}$ also respects the bilinear forms.

**Exercise 4.31.** Let $G_*$ be one of the three towers.

For every almost-composition $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ of $n \in \mathbb{N}$, let us define a map $\mathrm{ind}^n_\alpha$ which takes $\mathbb{C}G_\alpha$-modules to $\mathbb{C}G_n$-modules as follows: If $G_* = \mathfrak{S}_n$ or $G_* = \mathfrak{S}_n [\Gamma]$, we set

$$\mathrm{ind}^n_\alpha := \mathrm{Ind}^{G_n}_{G_\alpha}.$$

If $G_* = GL_n$, then we set

$$\mathrm{ind}^n_\alpha := \mathrm{Ind}^{P_n}_{P_\alpha} \mathrm{Infl}^*_{G_n}.$$

(Note that $\mathrm{ind}^n_\alpha = \mathrm{ind}^n_{i,j}$ if $\alpha$ has the form $(i, j)$.)

Similarly, for every almost-composition $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ of $n \in \mathbb{N}$, let us define a map $\mathrm{res}^n_\alpha$ which takes $\mathbb{C}G_n$-modules to $\mathbb{C}G_\alpha$-modules as follows: If $G_* = \mathfrak{S}_n$ or $G_* = \mathfrak{S}_n [\Gamma]$, we set

$$\mathrm{res}^n_\alpha := \mathrm{Res}^{G_\alpha}_{G_n}.$$

If $G_* = GL_n$, then we set

$$\mathrm{res}^n_\alpha := \left( \mathrm{Res}^{P_\alpha}_{P_n} (-) \right)^{K_n}.$$

(Note that $\mathrm{res}^n_\alpha = \mathrm{res}^n_{i,j}$ if $\alpha$ has the form $(i, j)$.)

(a) If $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ is an almost-composition of an integer $n \in \mathbb{N}$ satisfying $\ell \geq 1$, and if $V_i$ is a $\mathbb{C}G_{\alpha_i}$-module for every $i \in \{1, 2, \ldots, \ell\}$, then show that

$$\mathrm{ind}^n_{\alpha_1+\alpha_2+\cdots+\alpha_{\ell-1}, \alpha_\ell} \left( \mathrm{ind}^{\alpha_1+\alpha_2+\cdots+\alpha_{\ell-1}}_{(\alpha_1, \alpha_2, \ldots, \alpha_{\ell-1})} (V_1 \otimes V_2 \otimes \cdots \otimes V_{\ell-1}) \otimes V_\ell \right)$$

$$\cong \mathrm{ind}^n_\alpha (V_1 \otimes V_2 \otimes \cdots \otimes V_\ell)$$

$$\cong \mathrm{ind}^n_{\alpha_1+\alpha_2+\cdots+\alpha_{\ell-1}} (V_1 \otimes \mathrm{ind}^{\alpha_2+\alpha_3+\cdots+\alpha_\ell}_{(\alpha_2, \alpha_3, \ldots, \alpha_\ell)} (V_2 \otimes V_3 \otimes \cdots \otimes V_\ell)).$$

(b) Solve Exercise 4.20 again using Exercise 4.31(a).

(c) We proved above that the map $m : A \otimes A \to A$ (where $A = A(G_*)$) is associative, by using the adjointness of $m$ and $\Delta$. Give a new proof of this fact, which makes no use of $\Delta$.

(d) If $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ is an almost-composition of an $n \in \mathbb{N}$, and if $\chi_i \in R(G_{\alpha_i})$ for every $i \in \{1, 2, \ldots, \ell\}$, then show that

$$\chi_1 \chi_2 \cdots \chi_\ell = \mathrm{ind}^n_\alpha (\chi_1 \otimes \chi_2 \otimes \cdots \otimes \chi_\ell)$$

in $A = A(G_*)$.

(e) If $n \in \mathbb{N}$, $\ell \in \mathbb{N}$ and $\chi \in R(G_n)$, then show that

$$\Delta^{(\ell-1)} \chi = \sum \mathrm{res}^n_\alpha \chi$$

in $A^{\otimes \ell}$, where $A = A(G_*)$. Here, the sum on the right hand side runs over all almost-compositions $\alpha$ of $n$ having length $\ell$.

4.4 Symmetric groups. Finally, some payoff. Consider the tower of symmetric groups $G_n = \mathfrak{S}_n$, and $A = A(G_*) =: A(\mathfrak{S})$. Denote by $1_{\mathfrak{S}_n}, \sgm_{\mathfrak{S}_n}$, the *trivial* and *sign* characters on $\mathfrak{S}_n$. For a partition $\lambda$ of $n$, denote by $1_{\mathfrak{S}_\lambda}, \sgm_{\mathfrak{S}_\lambda}$, the trivial and sign characters restricted to the Young subgroup $\mathfrak{S}_\lambda = \mathfrak{S}_{\lambda_1} \times \mathfrak{S}_{\lambda_2} \times \cdots$, and denote by $1_{\lambda}$ the class function which is the characteristic function for the $\mathfrak{S}_n$-conjugacy class of permutations of cycle type $\lambda$.

**Theorem 4.32.** (a) Irreducible complex characters $\{\chi^\lambda\}$ of $\mathfrak{S}_n$ are indexed by partitions $\lambda$ in $\text{Par}_n$, and one has a PSH-isomorphism, the **Frobenius characteristic map**,

$$A = A(\mathfrak{S}) \xrightarrow{\text{ch}} \Lambda$$

---

\footnote{This is because, for example, induction of class functions harmonizes with induction of modules (i.e., the equality (4.5) holds).}

\footnote{It is unrelated to the Frobenius endomorphisms from Exercise 2.85.}
that for \( n \geq 0 \) and \( \lambda \in \text{Par}_n \) sends
\[
\begin{align*}
1_{E_n} & \mapsto h_n \\
\text{sgn}_{E_n} & \mapsto e_n \\
\chi & \mapsto s_\lambda \\
\text{Ind}_{E_\lambda}^{E_n} 1_{E_n} & \mapsto h_\lambda \\
\text{Ind}_{E_\lambda}^{E_n} \text{sgn}_{E_\lambda} & \mapsto e_\lambda \\
1_\lambda & \mapsto \frac{\rho_\lambda}{z_\lambda}
\end{align*}
\]
(where \( \text{ch} \) is extended to a \( \mathbb{C} \)-linear map \( A_C \to \Lambda_C \) and for \( n \geq 1 \) sends
\[
1_{(n)} \mapsto \frac{\rho_n}{n}
\]
Here, \( z_\lambda \) is defined as in Proposition 2.43.

(b) For each \( n \geq 0 \), the involution on class functions \( f : \mathfrak{S}_n \to \mathbb{C} \) sending \( f \mapsto \text{sgn}_{E_n} \cdot f \) where
\[
(\text{sgn}_{E_n} \cdot f)(g) := \text{sgn}(g) f(g)
\]
preserves the \( \mathbb{Z} \)-sublattice \( R(\mathfrak{S}_n) \) of genuine characters. The direct sum of these involutions induces an involution on \( A = A(\mathfrak{S}) = \bigoplus_{n \geq 0} R(\mathfrak{S}_n) \) that corresponds under \( \text{ch} \) to the involution \( \omega \) on \( \Lambda \).

Proof. (a) Corollary 4.30 implies that the set \( \Sigma = \bigsqcup_{n \geq 0} \text{Irr}(\mathfrak{S}_n) \) gives a PSH-basis for \( A \). Since a character \( \chi \) of \( \mathfrak{S}_n \) has
\[
(4.30) \quad \Delta(\chi) = \bigoplus_{i+j=n} \text{Res}_{\mathfrak{S}_i \times \mathfrak{S}_j} \chi,
\]
such an element \( \chi \in \Sigma \cap A_n \) is never primitive for \( n \geq 2 \). Hence the unique irreducible character \( \rho = 1_{E_1} \) of \( \mathfrak{S}_1 \) is the only element of \( \mathcal{C} = \Sigma \cap p \).

Thus Theorem 3.30(g) tells us that there are two PSH-isomorphisms \( A \to \Lambda \), each of which sends \( \Sigma \) to the PSH-basis of Schur functions \( \{s_\lambda\} \) for \( \Lambda \). It also tells us that we can pin down one of the two isomorphisms to call \( \text{ch} \), by insisting that it map the two characters \( 1_{E_2}, \text{sgn}_{E_2} \) in \( \text{Irr}(\mathfrak{S}_2) \) to \( h_2, e_2 \) (and not \( e_2, h_2 \)).

Bearing in mind the coproduct formula \((4.30)\), and the fact that \( 1_{E_n}, \text{sgn}_{E_n} \), restrict, respectively, to trivial and sign characters of \( \mathfrak{S}_i \times \mathfrak{S}_j \) for \( i+j=n \), one finds that for \( n \geq 2 \) one has \( \text{sgn}_{E_2} \) annihilating \( 1_{E_n} \), and \( 1_{E_2} \) annihilating \( \text{sgn}_{E_n} \). Therefore Theorem 3.18(b) (applied to \( \Lambda \)) implies \( 1_{E_n}, \text{sgn}_{E_n} \) are sent under \( \text{ch} \) to \( h_n, e_n \). Then the fact that \( \text{Ind}_{E_\lambda}^{E_n} 1_{E_n}, \text{Ind}_{E_\lambda}^{E_n} \text{sgn}_{E_\lambda} \) are sent to \( h_\lambda, e_\lambda \) follows via induction products.

Recall that the \( \mathbb{C} \)-vector space \( A_C = \bigoplus_{n \geq 0} R_C(\mathfrak{S}_n) \) is a \( \mathbb{C} \)-bialgebra, and can be identified with \( A \otimes_{\mathbb{Z}} \mathbb{C} \). The multiplicity and the comultiplication of \( A_C \) are \( \mathbb{C} \)-linear extensions of those of \( A \), and are still given by the same formulas \( m = \text{ind}^{i+j}_{i,j} \) and \( \Delta = \bigoplus (-1)^{i+j=n} \text{res}^{i+j}_{i,j} \) as those of \( A \) (but now, induction and restriction are defined for class functions, not just for representations). The \( \mathbb{C} \)-bilinear form \( \langle \cdot, \cdot \rangle_A \) on \( A_C \) extends both the \( \mathbb{Z} \)-bilinear form \( \langle \cdot, \cdot \rangle_A \) on \( A \) and the \( \mathbb{C} \)-bilinear forms \( \langle \cdot, \cdot \rangle_{E_n} \) on all \( R_C(\mathfrak{S}_n) \).

For the assertion about \( 1_{(n)} \), note that it is primitive in \( A_C \) for \( n \geq 1 \), because as a class function, the indicator function of \( n \)-cycles vanishes upon restriction to \( \mathfrak{S}_i \times \mathfrak{S}_j \) for \( i+j=n \) if both \( i,j \geq 1 \); these subgroups contain no \( n \)-cycles. Hence Corollary 3.9 implies that \( \text{ch}(1_{(n)}) \) is a scalar multiple of \( p_n \). To pin down the scalar, note \( p_n = m(n) \) so \( (h_n, p_n)_\Lambda = (h_n, m_n)_\Lambda = 1 \), while \( \text{ch}^{-1}(h_n) = 1_{E_n} \) has
\[
1_{E_n} \cdot 1_{(n)} = \frac{1}{n!} \cdot (n-1)! = \frac{1}{n}. \tag{4.35}
\]
Thus \( \text{ch}(1_{(n)}) = \frac{\rho_n}{n} \). The fact that \( \text{ch}(1_\lambda) = \frac{\rho_\lambda}{z_\lambda} \) then follows via induction product calculations\(^{166}\). Part (b) follows from Exercise 4.35 below. \( \square \)

Remark 4.33. The paper of Lulevičius [98] gives a very elegant alternate approach to the Frobenius map as a Hopf isomorphism \( A(\mathfrak{S}) \xrightarrow{\text{ch}} \Lambda \), inspired by equivariant K-theory and vector bundles over spaces which are finite sets of points!

\(^{165}\)The first equality sign in this computation uses the fact that the number of all \( n \)-cycles in \( \mathfrak{S}_n \) is \((n-1)!\). This is because any \( n \)-cycle in \( \mathfrak{S}_n \) can be uniquely written in the form \( (i_1, i_2, \ldots, i_{n-1}, n) \) (in cycle notation) with \( (i_1, i_2, \ldots, i_{n-1}) \) being a permutation in \( \mathfrak{S}_{n-1} \) (written in one-line notation).

\(^{166}\)For instance, one can use (4.3) to show that \( z_\lambda \lambda = \lambda_1 \lambda_2 \cdots \lambda_\ell \cdot 1_{(\lambda_1)} 1_{(\lambda_2)} \cdots 1_{(\lambda_\ell)} \) if \( \lambda = (\lambda_1, \lambda_2, \ldots, \lambda_\ell) \) with \( \ell = \ell(\lambda) \). See Exercise 4.34(d) for the details.
Exercise 4.34. If $P$ is a subset of a group $G$, we denote by $1_P$ the map $G \to \mathbb{C}$ which sends every element of $P$ to 1 and all remaining elements of $G$ to 0. For any finite group $G$ and any $h \in G$, we introduce the following notations:

- Let $Z_G(h)$ denote the centralizer of $h$ in $G$.
- Let $\text{Conj}_G(h)$ denote the conjugacy class of $h$ in $G$.
- Define a map $\alpha_{G,h} : G \to \mathbb{C}$ by $\alpha_{G,h} = |Z_G(h)|^{-1}$. This map $\alpha_{G,h}$ is a class function.

(a) Prove that $\alpha_{G,h}(g) = \sum_{k \in G} \left[ khk^{-1} = g \right]$ for every finite group $G$ and any $h \in G$ and $g \in G$. Here, we are using the Iverson bracket notation (that is, for any statement $A$, we define $|A|$ to be the integer 1 if $A$ is true, and 0 otherwise).

(b) Prove that if $H$ is a subgroup of a finite group $G$, and if $h \in H$, then $\text{Ind}_H^G \alpha_{H,h} = \alpha_{G,h}$.

(c) Prove that if $G_1$ and $G_2$ are finite groups, and if $h_1 \in G_1$ and $h_2 \in G_2$, then the canonical isomorphism $R_C(G_1) \otimes R_C(G_2) \to R_C(G_1 \times G_2)$ sends $\alpha_{G_1,h_1} \otimes \alpha_{G_2,h_2}$ to $\alpha_{G_1 \times G_2,(h_1,h_2)}$.

(d) Let $\text{Conj}_G(h)$ denote the conjugacy class of $h$ in $G$.

Exercise 4.35. If $G$ is a group and $U_1$ and $U_2$ are two $CG$-modules, then the tensor product $U_1 \otimes U_2$ is a $\mathbb{C}[G \times G]$-module, which can be made into a $CG$-module by letting $g \in G$ act as $(g,g) \in G \times G$. This $CG$-module $U_1 \otimes U_2$ is called the inner tensor product of $U_1$ and $U_2$, and is a restriction of the outer tensor product $U_1 \otimes U_2$ using the inclusion map $G \to G \times G$, $g \mapsto (g,g)$.

Let $n \geq 0$, and let $\text{sgn} e_n$ be the $1$-dimensional $\mathbb{C}S_n$-module $\mathbb{C}$ on which every $g \in S_n$ acts as multiplication by $\text{sgn}(g)$. If $V$ is a $\mathbb{C}S_n$-module, show that the involution on $A(\mathbb{S}) = \bigoplus_{n \geq 0} R(\mathbb{S}_n)$ defined in Theorem 4.32(b) sends $\chi_V \mapsto \chi_{\text{sgn} e_n \otimes V}$ where $\text{sgn} e_n \otimes V$ is the inner tensor product of $\text{sgn} e_n$ and $V$. Use this to show that this involution is a nontrivial PSH-automorphism of $A(\mathbb{S})$, and deduce Theorem 4.32(b).

Exercise 4.36. (a) Show that for every $n \geq 0$, every $g \in S_n$ and every finite-dimensional $\mathbb{C}S_n$-module $V$, we have $\chi_V(g) \in \mathbb{Z}$.

(b) Show that for every $n \geq 0$ and every finite-dimensional $\mathbb{C}S_n$-module $V$, there exists a $\mathbb{Q}S_n$-module $W$ such that $V \cong \mathbb{C} \otimes_{\mathbb{Q}} W$. (In the representation theorists’ parlance, this says that all representations of $S_n$ are defined over $\mathbb{Q}$. This part of the exercise requires some familiarity with representation theory.)

Remark 4.37. Parts (a) and (b) of Exercise 4.36 both follow from an even stronger result: For every $n \geq 0$ and every finite-dimensional $\mathbb{C}S_n$-module $V$, there exists a $\mathbb{Z}S_n$-module $W$ which is finitely generated and free as a $\mathbb{Z}$-module and satisfies $V \cong \mathbb{C} \otimes_{\mathbb{Z}} W$ as $\mathbb{C}S_n$-modules. This follows from the combinatorial approach to the representation theory of $S_n$, in which the irreducible representations of $S_n$ (the Specht modules) are constructed using Young tableaux and tabloids. See the literature on the symmetric group, e.g., [138], [53, §7] or [167] for this approach.

The connection between $\Lambda$ and $A(\mathbb{S})$ as established in Theorem 4.32 benefits both the study of $\Lambda$ and that of $A(\mathbb{S})$. The following two exercises show some applications to $\Lambda$:

Exercise 4.38. If $G$ is a group and $U_1$ and $U_2$ are two $CG$-modules, then let $U_1 \boxtimes U_2$ denote the inner tensor product of $U_1$ and $U_2$ (as defined in Exercise 4.35). Consider also the binary operation $* \circ \lambda_\mathbb{Q}$ defined in Exercise 2.80(h).

(a) Show that $\chi_{U_1 \boxtimes U_2} = \chi_{U_1} * \chi_{U_2}$ for any $n \in \mathbb{N}$ and any two $\mathbb{C}S_n$-modules $U_1$ and $U_2$.

(b) Use this to obtain a new solution for Exercise 2.80(h).

(c) Show that $s_\mu * s_\nu = \sum_{\lambda \in \text{Par}} \mathbb{N}_{\lambda\lambda}$ for any two partitions $\mu$ and $\nu$.

[Hint: For any group $G$, introduce a binary operation $* \circ \lambda_\mathbb{Q}$ on $R_C(G)$ which satisfies $\chi_{U_1 \boxtimes U_2} = \chi_{U_1} * \chi_{U_2}$ for any two $CG$-modules $U_1$ and $U_2$.]
Exercise 4.39. Define a \( \mathbb{Q} \)-bilinear map \( \boxtimes : \Lambda_{\mathbb{Q}} \times \Lambda_{\mathbb{Q}} \to \Lambda_{\mathbb{Q}} \), which will be written in infix notation (that is, we will write \( a \boxtimes b \) instead of \( \boxtimes (a, b) \)), by setting
\[
p_{\lambda} \boxtimes p_{\mu} = \prod_{i,j} \prod_{1}^{\ell(\lambda)} \prod_{1}^{\ell(\mu)} p_{\gcd(\lambda_i, \mu_j)} \quad \text{for any partitions } \lambda \text{ and } \mu.
\]

(a) Show that \( \Lambda_{\mathbb{Q}} \), equipped with the binary operation \( \boxtimes \), becomes a commutative \( \mathbb{Q} \)-algebra with unity \( p_1 \).

(b) For every \( r \in \mathbb{Z} \), define the \( \mathbb{Q} \)-algebra homomorphism \( \epsilon_r : \Lambda_{\mathbb{Q}} \to \mathbb{Q} \) as in Exercise 2.80(c). Show that \( 1 \boxtimes f = \epsilon_1(f) \) for every \( f \in \Lambda_{\mathbb{Q}} \) (where 1 denotes the unity of \( \Lambda \)).

(c) Show that \( s_\mu \boxtimes s_\nu = \sum_{\lambda \in \text{Par}_n} N_s_{\lambda} \) for any two partitions \( \mu \) and \( \nu \).

(d) Show that \( f \boxtimes g \in \Lambda \) for any \( f \in \Lambda \) and \( g \in \Lambda \).

[Hint: For every set \( X \), let \( \mathfrak{S}_X \) denote the group of all permutations of \( X \). For two sets \( X \) and \( Y \), there is a canonical group homomorphism \( \mathfrak{S}_X \times \mathfrak{S}_Y \to \mathfrak{S}_{X \times Y} \), which is injective if \( X \) and \( Y \) are nonempty. For positive integers \( n \) and \( m \), this yields an embedding \( \mathfrak{S}_n \times \mathfrak{S}_m \to \mathfrak{S}_{\{1, \ldots, n\} \times \{1, \ldots, m\}} \), which, once \( \mathfrak{S}_{\{1, \ldots, n\} \times \{1, \ldots, m\}} \) is identified with \( \mathfrak{S}_{nm} \) (using an arbitrary but fixed bijection \( \{1, 2, \ldots, n\} \times \{1, 2, \ldots, m\} \to \{1, 2, \ldots, nm\} \)), can be regarded as an embedding \( \mathfrak{S}_n \times \mathfrak{S}_m \to \mathfrak{S}_{nm} \) and thus allows defining a \( \mathbb{C} \mathfrak{S}_{nm} \)-module \( \text{Ind}_{\mathfrak{S}_n \times \mathfrak{S}_m}^{\mathfrak{S}_{nm}}(U \otimes V) \) for any \( \mathbb{C} \mathfrak{S}_{nm} \)-module \( U \) and any \( \mathbb{C} \mathfrak{S}_n \)-module \( V \). This gives a binary operation on \( A(\mathfrak{S}) \). Show that this operation corresponds to \( \boxtimes \) under the PSH-isomorphism \( \text{ch} : A(\mathfrak{S}) \to \Lambda \).]

Remark 4.40. The statements (and the idea of the solution) of Exercise 4.39 are due to Manuel Maia and Miguel Méndez (see [108] and, more explicitly, [115]), who call the operation \( \boxtimes \) the arithmetic product. Li [97, Thm. 3.5] denotes it by \( \boxtimes \) and relates it to the enumeration of unlabelled graphs.

4.5. Wreath products. Next consider the tower of groups \( G_n = \mathfrak{S}_n[\Gamma] \) for a finite group \( \Gamma \), and the Hopf algebra \( A = A(G_n) = A(\mathfrak{S}[\Gamma]) \). Recall (from Theorem 4.32) that irreducible complex representations \( \chi^\lambda \) of \( \mathfrak{S}_n \) are indexed by partitions \( \lambda \) in \( \text{Par}_n \). Index the irreducible complex representations of \( \Gamma = \{\rho_1, \ldots, \rho_d\} \).

Definition 4.41. Define for a partition \( \lambda \in \text{Par}_n \) and \( \rho_i \in \text{Irr}(\Gamma) \) a representation \( \chi^{\lambda, \rho} \) of \( \mathfrak{S}_n[\Gamma] \) in which \( \sigma \) in \( \mathfrak{S}_n \) and \( \gamma = (\gamma_1, \ldots, \gamma_n) \) in \( \Gamma^n \) act on the space \( \chi^\lambda \otimes (\rho^\sigma) \) as follows
\[
\begin{align*}
\sigma(u \otimes v_1 \otimes \cdots \otimes v_n) &= \sigma(u) \otimes (v_{\sigma^{-1}(1)} \otimes \cdots \otimes v_{\sigma^{-1}(n)}) \\
\gamma(u \otimes v_1 \otimes \cdots \otimes v_n) &= \gamma_1 \cdot v_1 \otimes \cdots \otimes \gamma_n \cdot v_n
\end{align*}
\]

(4.31)

Theorem 4.42. The irreducible \( \mathbb{C} \mathfrak{S}_n[\Gamma] \)-modules are the induced characters
\[
\chi^\lambda := \text{Ind}_{\mathfrak{S}_n[\Gamma]}^{\mathfrak{S}_n[\Gamma]} \left( \chi^{\lambda_{(1)}, \rho_1} \otimes \cdots \otimes \chi^{\lambda_{(d)}, \rho_d} \right)
\]
as \( \lambda \) runs through all functions
\[
\text{Irr}(\Gamma) \xrightarrow{\lambda \mapsto \text{Par}} \lambda^{(i)}
\]
with the property that \( \sum_{i=1}^d \lambda_{(i)} = n \). Here, \( \text{deg}(\lambda) \) denotes the \( d \)-tuple \( (|\lambda^{(1)}|, |\lambda^{(2)}|, \ldots, |\lambda^{(d)}|) \in \mathbb{N}^d \), and \( \mathfrak{S}_{\text{deg}(\lambda)} \) is defined as the subgroup \( \mathfrak{S}_{|\lambda^{(1)}|} \times \mathfrak{S}_{|\lambda^{(2)}|} \times \cdots \times \mathfrak{S}_{|\lambda^{(d)}|} \) of \( \mathfrak{S}_n \).

Furthermore, one has a PSH-isomorphism
\[
A(\mathfrak{S}[\Gamma]) \xrightarrow{\chi^\lambda \mapsto s_{\lambda^{(1)}} \otimes \cdots \otimes s_{\lambda^{(d)}}} \Lambda^\otimes d
\]

Proof. We know from Corollary 4.30 that \( A(\mathfrak{S}[\Gamma]) \) is a PSH, with PSH-basis \( \Sigma \) given by the union of all irreducible characters of all groups \( \mathfrak{S}_n[\Gamma] \). Therefore Theorem 3.12 tells us that \( A(\mathfrak{S}[\Gamma]) \cong \bigotimes_{\rho \in \Sigma} A(\mathfrak{S}[\Gamma])(\rho) \) where \( \mathcal{C} \) is the set of irreducible characters which are also primitive. Just as in the case of \( \mathfrak{S}_n \), it is clear from the definition of the coproduct that an irreducible character \( \rho \) of \( \mathfrak{S}_n[\Gamma] \) is primitive if and only if \( n = 1 \), that in this case \( \mathfrak{S}_n[\Gamma] = \Gamma \), and \( \rho \) lies in \( \text{Irr}(\Gamma) = \{\rho_1, \ldots, \rho_d\} \).

\[\text{This is well-defined, since } (p_{\lambda})_{\lambda \in \text{Par}} \text{ is a } \mathbb{Q}-\text{module basis of } \Lambda_{\mathbb{Q}}.\]
The remaining assertions of the theorem will then follow from the definition of the induction product algebra structure on \( A(\mathfrak{S}[\Gamma]) \), once we have shown that, for every \( \rho \in \text{Irr}(\Gamma) \), there is a PSH-isomorphism sending

\[
A(\mathfrak{S}) \quad \xrightarrow{\chi} \quad A(\mathfrak{S}[\Gamma])(\rho)
\]

such an isomorphism comes from applying Proposition 4.17 to the semidirect product \( \mathfrak{S}_n[\Gamma] = \mathfrak{S}_n \rtimes \Gamma^n \), so that \( K = \Gamma^n, G = \mathfrak{S}_n \), and fixing \( V = \rho^{\otimes n} \) as \( \mathbb{C}\mathfrak{S}_n[\Gamma] \)-module with structure as defined in \((4.31)\) (but with \( \lambda \) set to \( (n) \)), so that \( \chi^\lambda \) is the trivial 1-dimensional \( G \)-module and therefore \( \chi^\lambda \) is the trivial 1-dimensional \( \mathfrak{S}_n \)-module. One obtains for each \( n \), maps

\[
R(\mathfrak{S}_n) \xrightarrow{\Phi} R(\mathfrak{S}_n[\Gamma])
\]

where

\[
\chi \xrightarrow{\Phi} \chi \otimes (\rho^{\otimes n})
\]

\[
\alpha \xrightarrow{\Psi} \text{Hom}_{\mathbb{C}\Gamma^n}(\rho^{\otimes n}, \alpha).
\]

Taking the direct sum of these maps for all \( n \) gives maps \( A(\mathfrak{S}) \xrightarrow{\Phi} A(\mathfrak{S}[\Gamma]) \).

These maps are coalgebra morphisms because of their interaction with restriction to \( \mathfrak{S}_i \times \mathfrak{S}_j \). Since Proposition 4.17(iii) gives the adjointness property that

\[
(\chi, \Psi(\alpha))_{A(\mathfrak{S})} = (\Phi(\chi), \alpha)_{A(\mathfrak{S}[\Gamma])},
\]

one concludes from the self-duality of \( A(\mathfrak{S}), A(\mathfrak{S}[\Gamma]) \) that \( \Phi, \Psi \) are also algebra morphisms. Since they take genuine characters to genuine characters, they are PSH-morphisms. Since \( \rho \) being a simple \( \mathbb{C}\Gamma \)-module implies that \( V = \rho^{\otimes n} \) is a simple \( \mathbb{C}\Gamma^n \)-module, Proposition 4.17(iv) shows that

\[
(\Psi \circ \Phi)(\chi) = \chi
\]

for all \( \mathfrak{S}_n \)-characters \( \chi \). Hence \( \Phi \) is an injective PSH-morphism. Using adjointness, \((4.33)\) also shows that \( \Phi \) sends \( \mathbb{C}\mathfrak{S}_n \)-simples \( \chi \) to \( \mathbb{C}[\mathfrak{S}_n[\Gamma]] \)-simples \( \Phi(\chi) \):

\[
(\Phi(\chi), \Phi(\chi))_{A(\mathfrak{S}[\Gamma])} = ((\Psi \circ \Phi)(\chi), \chi)_{A(\mathfrak{S})} = (\chi, \chi)_{A(\mathfrak{S})} = 1.
\]

Since \( \Phi(\chi) = \chi \otimes (\rho^{\otimes n}) \) has \( V = \rho^{\otimes n} \) as a constituent upon restriction to \( \Gamma^n \), Frobenius Reciprocity shows that the irreducible character \( \Phi(\chi) \) is a constituent of \( \text{Ind}_{\mathfrak{S}_n[\Gamma]}^{\mathfrak{S}_n[\Gamma]}(\rho^{\otimes n}) = \rho^n \). Hence the entire image of \( \Phi \) lies in \( A(\mathfrak{S}[\Gamma])(\rho) \) (due to how we defined \( A(\rho) \) in the proof of Theorem 3.12), and so \( \Phi \) must restrict to an isomorphism as desired in \((4.32)\).

One of Zelevinsky's sample applications of the theorem is this branching rule.

**Corollary 4.43.** Given \( \underline{\lambda} = (\lambda^{(1)}, \ldots, \lambda^{(d)}) \) with \( \sum_{i=1}^d |\lambda^{(i)}| = n \), one has

\[
\text{Res}_{\mathfrak{S}_n[\Gamma] \times \Gamma}^{\mathfrak{S}_n[\Gamma]} \chi(\underline{\lambda}) = \sum_{i=1}^d \sum_{|\lambda^{(i)}| \leq \lambda^{(i)}} \chi^{(\lambda^{(1)}, \ldots, \lambda^{(i)}, \ldots, \lambda^{(d)})} \otimes \rho_i.
\]

(We are identifying functions \( \lambda : \text{Irr}(\Gamma) \to \text{Par} \) with the corresponding \( d \)-tuples \( (\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(d)}) \) here.)

**Example 4.44.** For \( \Gamma \) a two-element group, so \( \text{Irr}(\Gamma) = \{ \rho_1, \rho_2 \} \) and \( d = 2 \), then

\[
\text{Res}_{\mathfrak{S}_n[\Gamma] \times \Gamma}^{\mathfrak{S}_n[\Gamma]} \chi^{((3,3),(1,1))} = \chi^{((3,3),(1,1))} \otimes \rho_1 + \chi^{((2,1),(1,1))} \otimes \rho_1 + \chi^{((3,1),(1))} \otimes \rho_2.
\]

**Proof of Corollary 4.43.** By Theorem 4.42, this is equivalent to computing in the Hopf algebra \( A := \Lambda^{\otimes d} \) the component of the coproduct of \( s_{\lambda^{(1)}} \otimes \cdots \otimes s_{\lambda^{(d)}} \) that lies in \( A_{n-1} \otimes A_1 \). Working within each tensor factor \( \Lambda \), the Pieri formula implies that the \( \Lambda_{|\lambda|-1} \otimes A_1 \)-component of \( \Delta(s_{\lambda}) \) is

\[
\sum_{\lambda \subseteq \lambda; |\lambda|<|\lambda'|} s_{\lambda'} \otimes \rho.
\]

One must apply this in each of the \( d \) tensor factors of \( A = \Lambda^{\otimes d} \), then sum on \( i \). \( \square \)
4.6. General linear groups. We now consider the tower of finite general linear groups $G_n = GL_n = GL_n(\mathbb{F}_q)$ and $A = A(G_n) := A(GL)$. Corollary 4.30 tells us that $A(GL)$ is a PSH, with PSH-basis $\Sigma$ given by the union of all irreducible characters of all groups $GL_n$. Therefore Theorem 3.12 tells us that

$$A(GL) \cong \bigotimes_{\rho \in \mathcal{C}} A(GL)(\rho)$$

where $\mathcal{C} = \Sigma \cap \mathfrak{p}$ is the set of primitive irreducible characters.

**Definition 4.45.** Call an irreducible representation $\rho$ of $GL_n$ cuspidal for $n \geq 1$ if it lies in $\mathcal{C}$, that is, its restriction to proper parabolic subgroups $P_{i,j}$ with $i+j = n$ and $i, j > 0$ contain no nonzero vectors which are $K_{i,j}$-invariant. Given an irreducible character $\sigma$ of $GL_n$, say that $d(\sigma) = n$, and let $\mathcal{C}_n := \{ \rho \in \mathcal{C} : d(\rho) = n \}$ for $n \geq 1$ denote the subset of cuspidal characters of $GL_n$.

Just as was the case for $\mathfrak{S}_1$ and $\mathfrak{S}_1[\Gamma] = \Gamma$, every irreducible character $\rho$ of $GL_1(\mathbb{F}_q) = \mathbb{F}_q^\times$ is cuspidal. However, this does not exhaust the cuspidal characters. In fact, one can predict the number of cuspidal characters in $\mathcal{C}_n$, using knowledge of the number of conjugacy classes in $GL_n$. Let $\mathcal{F}$ denote the set of all nonconstant monic irreducible polynomials $f(x) \neq x$ in $\mathbb{F}_q[x]$. Let $\mathcal{F}_n := \{ f \in \mathcal{F} : \deg(f) = n \}$ for $n \geq 1$.

**Proposition 4.46.** The number $|\mathcal{C}_n|$ of cuspidal characters of $GL_n(\mathbb{F}_q)$ is the number of $|\mathcal{F}_n|$ of irreducible monic degree $n$ polynomials $f(x) \neq x$ in $\mathbb{F}_q[x]$ with nonzero constant term.

**Proof.** We show $|\mathcal{C}_n| = |\mathcal{F}_n|$ for $n \geq 1$ by strong induction on $n$. For the base case$^{171}$ $n = 1$, just as with the families $G_n = \mathfrak{S}_n$ and $G_n = \mathfrak{S}_n[\Gamma]$, when $n = 1$ any irreducible character $\chi$ of $G_1 = GL_1(\mathbb{F}_q)$ gives a primitive element of $A = A(GL)$, and hence is cuspidal. Since $GL_1(\mathbb{F}_q) = \mathbb{F}_q^\times$ is abelian, there are $|\mathbb{F}_q^\times| = q-1$ such cuspidal characters in $\mathcal{C}_1$, which agrees with the fact that there are $q-1$ monic (irreducible) linear polynomials $f(x) \neq x$ in $\mathbb{F}_q[x]$, namely $\mathcal{F}_1 := \{ f(x) = x-c : c \in \mathbb{F}_q^\times \}$.

In the inductive step, use the fact that the number $|\Sigma_n|$ of irreducible complex characters $\chi$ of $GL_n(\mathbb{F}_q)$ equals its number of conjugacy classes. These conjugacy classes are uniquely represented by rational canonical forms, which are parametrized by functions $\lambda : \mathcal{F} \to \Par$ with the property that $\sum_{f \in \mathcal{F}} \deg(f)|\lambda(f)| = n$. On the other hand, (4.34) tells us that $|\Sigma_n|$ is similarly parametrized by the functions $\lambda : \mathcal{C} \to \Par$ having the property that $\sum_{\rho \in \mathcal{C}} d(\rho)|\lambda(\rho)| = n$. Thus we have parallel disjoint decompositions

$$\mathcal{C} = \bigsqcup_{n \geq 1} \mathcal{C}_n$$
where $\mathcal{C}_n := \{ \rho \in \mathcal{C} : d(\rho) = n \}$

$$\mathcal{F} = \bigsqcup_{n \geq 1} \mathcal{F}_n$$
where $\mathcal{F}_n := \{ f \in \mathcal{F} : \deg(f) = n \}$

and hence an equality for all $n \geq 1$

$$\left| \left\{ \mathcal{C} \xrightarrow{\lambda} \Par : \sum_{\rho \in \mathcal{C}} d(\rho)|\lambda(\rho)| = n \right\} \right| = |\Sigma_n| = \left| \left\{ \mathcal{F} \xrightarrow{\lambda} \Par : \sum_{f \in \mathcal{F}} \deg(f)|\lambda(f)| = n \right\} \right|.$$

Since there is only one partition $\lambda$ having $|\lambda| = 1$ (namely, $\lambda = (1)$), this leads to parallel recursions

$$|\mathcal{C}_n| = |\Sigma_n| - \left| \bigsqcup_{i=1}^{n-1} \left\{ \mathcal{C}_i \xrightarrow{\lambda} \Par : \sum_{\rho \in \mathcal{C}} d(\rho)|\lambda(\rho)| = n \right\} \right|$$

$$|\mathcal{F}_n| = |\Sigma_n| - \left| \bigsqcup_{i=1}^{n-1} \left\{ \mathcal{F}_i \xrightarrow{\lambda} \Par : \sum_{f \in \mathcal{F}} \deg(f)|\lambda(f)| = n \right\} \right|$$

and induction implies that $|\mathcal{C}_n| = |\mathcal{F}_n|$.

$\square$

We shall use the notation $\mathbb{1}_H$ for the trivial character of a group $H$ whenever $H$ is a finite group. This generalizes the notations $\mathbb{1}_{\mathfrak{S}_n}$ and $\mathbb{1}_{\mathfrak{S}_\lambda}$ introduced above.

$^{171}$Actually, we don’t need any base case for our strong induction. We nevertheless handle the case $n = 1$ as a warmup.
Example 4.47. Taking \( q = 2 \), let us list the sets \( \mathcal{F}_n \) of monic irreducible polynomials \( f(x) \neq x \) in \( \mathbb{F}_2[x] \) of degree \( n \) for \( n \leq 3 \), so that we know how many cuspidal characters of \( GL_n(\mathbb{F}_q) \) in \( \mathcal{C}_n \) to expect:

\[
\begin{align*}
\mathcal{F}_1 &= \{x + 1\} \\
\mathcal{F}_2 &= \{x^2 + x + 1\} \\
\mathcal{F}_3 &= \{x^3 + x + 1, x^3 + x^2 + 1\}
\end{align*}
\]

Thus we expect

- one cuspidal character of \( GL_1(\mathbb{F}_2) \), namely \( \rho_1(= 1_{GL_1(\mathbb{F}_2)}) \),
- one cuspidal character \( \rho_2 \) of \( GL_2(\mathbb{F}_2) \), and
- two cuspidal characters \( \rho_3, \rho_3' \) of \( GL_3(\mathbb{F}_2) \).

We will say more about \( \rho_2, \rho_3, \rho_3' \) in the next section.

Exercise 4.48. (a) Show that for \( n \geq 2 \),

\[
|\mathcal{C}_n| = \frac{1}{n} \sum_{d|n} \mu \left( \frac{n}{d} \right) q^d
\]

where \( \mu(m) \) is the number-theoretic Möbius function of \( m \), that is \( \mu(m) = (-1)^d \) if \( m = p_1 \cdots p_d \) for \( d \) distinct primes, and \( \mu(m) = 0 \) if \( m \) is not squarefree. (Here, the summation sign \( \sum_{d|n} \) means a sum over all positive divisors \( d \) of \( n \).)

(b) Show that (4.35) also counts the necklaces with \( n \) beads of \( q \) colors (= equivalence classes under the \( \mathbb{Z}/n\mathbb{Z} \)-action of cyclic rotation on sequences \( (a_1, \ldots, a_n) \) in \( \mathbb{F}_q^n \)) which are primitive in the sense that no nontrivial rotation fixes any of the sequences within the equivalence class. For example, when \( q = 2 \), here are representatives of these primitive necklaces for \( n = 2, 3, 4 \):

- \( n = 2 \): \( \{(0, 1)\} \)
- \( n = 3 \): \( \{(0, 0, 1), (0, 1, 1)\} \)
- \( n = 4 \): \( \{(0, 0, 0, 1), (0, 0, 1, 1), (0, 1, 1, 1)\} \).

4.7. Steinberg’s unipotent characters. Not surprisingly, the (cuspidal) character \( \iota := 1_{GL_1} \) of \( GL_1(\mathbb{F}_q) \) plays a distinguished role. The parabolic subgroup \( P(1^n) \) of \( GL_n(\mathbb{F}_q) \) is the Borel subgroup \( B \) of upper triangular matrices, and we have \( \iota^n = \text{Ind}^{GL_n}_{B} 1_B = \mathbb{C}[GL_n/B] \) (identifying representations with their characters as usual)\(^{172}\). The subalgebra \( A(GL)(\iota) \) of \( A(GL) \) is the \( \mathbb{Z} \)-span of the irreducible characters \( \sigma \) that appear as constituents of \( \iota^n = \text{Ind}^{GL_n}_{B} 1_B = \mathbb{C}[GL_n/B] \) for some \( n \).

Definition 4.49. An irreducible character \( \sigma \) of \( GL_n \) appearing as a constituent of \( \text{Ind}^{GL_n}_{B} 1_B = \mathbb{C}[GL_n/B] \) is called a unipotent character. Equivalently, by Frobenius reciprocity, \( \sigma \) is unipotent if it contains a nonzero \( B \)-invariant vector.

In particular, \( 1_{GL_n} \) is a unipotent character of \( GL_n \) for each \( n \).

Proposition 4.50. One can choose \( \Lambda \cong A(GL)(\iota) \) in Theorem 3.20(g) so that \( h_n \mapsto 1_{GL_n} \).

Proof. Theorem 3.18(a) tells us \( \iota^2 = \text{Ind}^{GL_2}_{B} 1_B \) must have exactly two irreducible constituents, one of which is \( 1_{GL_2} \); call the other one \( \text{St}_2 \). Choose the isomorphism so as to send \( h_2 \mapsto 1_{GL_2} \). Then \( h_n \mapsto 1_{GL_n} \) follows from the claim that \( \text{St}_2^2 (1_{GL_2}) = 0 \) for \( n \geq 2 \): one has

\[
\Delta(1_{GL_n}) = \sum_{i+j=n} \left( \text{Res}_{P_{i,j}}^{G_n} 1_{GL_n} \right)^{K_{i,j}} = \sum_{i+j=n} 1_{GL_i} \otimes 1_{GL_j}
\]

\(^{172}\) Proof. Exercise 4.31(d) (applied to \( G_* = GL_* \), \( \ell = n \), \( \alpha = (1^n) = \left( \begin{array}{c} 1, 1, \ldots, 1 \\ n \text{ times} \end{array} \right) \) and \( \chi_i = \iota \)) gives

\[
i^n = \text{Ind}^{(1^n)}_{B} \iota^n = \text{Ind}^{P_{(1^n)}}_{B} \iota^{P_{(1^n)}} = \text{Ind}^{GL_n}_{B} 1_B = \mathbb{C}[GL_n/B],
\]

where the last equality follows from the general fact that if \( G \) is a finite group and \( H \) is a subgroup of \( G \), then \( \text{Ind}^{H}_{B} 1_{H} \cong \mathbb{C}[G/H] \) as \( CG \)-modules.
so that $\text{St}_2^+ (1_{GL_n}) = (\text{St}_2, 1_{GL_2}) 1_{GL_{n-2}} = 0$ since $\text{St}_2 \neq 1_{GL_2}$.

This subalgebra $A(GL)(i)$, and the unipotent characters $\chi^\lambda_q$ corresponding under this isomorphism to the Schur functions $s_\lambda$, were introduced by Steinberg [156]. He wrote down $\chi^\lambda_q$ as a virtual sum of induced characters $\text{Ind}_{P_\alpha}^{G_n} 1_{P_\alpha} = 1_{G_n} \cdots 1_{G_{\alpha_i}}$, modelled on the Jacobi-Trudi determinantal expression for $s_\lambda = \det(h_{\lambda_{i+j}})$. Note that $\text{Ind}_{P_\alpha}^{G_n} 1_{P_\alpha}$ is the transitive permutation representation $\mathbb{C}[G/P_\alpha]$ for $GL_n$ permuting the finite partial flag variety $G/P_\alpha$, that is, the set of $\alpha$-flags of subspaces

$$\{0\} \subset V_{\alpha_1} \subset V_{\alpha_1 + \alpha_2} \subset \ldots \subset V_{\alpha_1 + \alpha_2 + \ldots + \alpha_{i-1}} \subset \mathbb{F}^n_q$$

where $\dim_q V_d = d$ in each case. This character has dimension equal to $|G/P_\alpha|$, with formula given by the $q$-multinomial coefficient (see e.g. Stanley [154, §1.7]):

$$\begin{bmatrix} n \\ \alpha \end{bmatrix}_q = \frac{[n]_q!}{[\alpha_1]_q! \cdots [\alpha_\ell_q]_q!}$$

where $[n]_q! = [n]_q[n-1]_q \cdots [2]_q[1]_q$ and $[n]_q := 1 + q + \cdots + q^{n-1} = \frac{q^n - 1}{q-1}$.

Our terminology $\text{St}_2$ is motivated by the $n = 2$ special case of the Steinberg character $\text{St}_n$, which is the unipotent character corresponding under the isomorphism in Proposition 4.50 to $e_n = s_{(1^n)}$. It can be defined by the virtual sum

$$\text{St}_n := \chi^q_{(1^n)} = \sum_{\alpha} (-1)^n - \ell(\alpha) \text{Ind}_{P_\alpha}^{G_n} 1_{P_\alpha}$$

in which the sum runs through all compositions $\alpha$ of $n$. This turns out to be the genuine character for $GL_n(\mathbb{F}_q)$ acting on the top homology group of its Tits building: the simplicial complex whose vertices are nonzero proper subspaces $V$ of $\mathbb{F}^n_q$, and whose simplices correspond to flags of nested subspaces. One needs to know that this Tits building has only top homology, so that one can deduce the above character formula from the Hopf trace formula; see Björner [21].

4.8. Examples: $GL_2(\mathbb{F}_2)$ and $GL_3(\mathbb{F}_2)$. Let’s get our hands dirty.

Example 4.51. For $n = 2$, there are two unipotent characters, $\chi^q_{(2)} = 1_{GL_2}$ and

$$(4.36) \quad \text{St}_2 := \chi^{(1,1)}_q = \frac{1}{1_{GL_1}} - \frac{1}{1_{GL_2}} = \text{Ind}_{B}^{GL_2} 1_B - 1_{GL_2}$$

since the Jacobi-Trudi formula (2.22) gives $s_{(1,1)} = \det \begin{bmatrix} h_1 & h_2 \\ 1 & h_1 \end{bmatrix} = h_1^2 - h_2$. The description (4.36) for this Steinberg character $\text{St}_2$ shows that it has dimension

$$|GL_2/B| - 1 = (q + 1) - 1 = q$$

and that one can think of it as follows: consider the permutation action of $GL_2$ on the $q+1$ lines $\{\ell_0, \ell_1, \ldots, \ell_q\}$ in the projective space $\mathbb{P}^1_{\mathbb{F}_q} = GL_2(\mathbb{F}_q)/B$, and take the invariant subspace perpendicular to the sum of basis elements $e_{\ell_0} + \cdots + e_{\ell_q}$.

Example 4.52. Continuing the previous example, but taking $q = 2$, we find that we have constructed two unipotent characters: $\chi^q_{(2)} = 1_{GL_2}$ of dimension 1, and $\text{St}_2 = \chi^{(1,1)}_q$ of dimension $q = 2$. This lets us identify the unique cuspidal character $\rho_2$ of $GL_2(\mathbb{F}_2)$, using knowledge of the character table of $GL_2(\mathbb{F}_2) \cong S_3$:

<table>
<thead>
<tr>
<th>$1_{GL_2} = \chi^q_{(2)}$</th>
<th>unipotent</th>
<th>1</th>
<th>1</th>
<th>1</th>
</tr>
</thead>
<tbody>
<tr>
<td>$1_{GL_2}$</td>
<td></td>
<td>1</td>
<td></td>
<td></td>
</tr>
<tr>
<td>$\text{St}_2 = \chi^{(1,1)}_q$</td>
<td>unipotent</td>
<td>2</td>
<td>0</td>
<td>-1</td>
</tr>
<tr>
<td>$\rho_2$</td>
<td></td>
<td>1</td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

In other words, the cuspidal character $\rho_2$ of $GL_2(\mathbb{F}_2)$ corresponds under the isomorphism $GL_2(\mathbb{F}_2) \cong S_3$ to the sign character $\text{sgn}_{S_3}$. 

\[ \square \]
Example 4.53. Continuing the previous example to \( q = 2 \) and \( n = 3 \) lets us analyze the irreducible characters of \( GL_3(\mathbb{F}_2) \). Recalling our labelling \( \rho_1, \rho_2, \rho_3, \rho_3' \) from Example 4.47 of the cuspidal characters of \( GL_n(\mathbb{F}_2) \) for \( n = 1, 2, 3 \), Zelevinsky’s Theorem 3.12 tells us that the \( GL_3(\mathbb{F}_2) \)-irreducible characters should be labelled by functions \( \{ \rho_1, \rho_2, \rho_3, \rho_3' \} \xrightarrow{\Delta} \text{Par} \) for which

\[
1 \cdot |\Delta(\rho_1)| + 2 \cdot |\Delta(\rho_2)| + 3 \cdot |\Delta(\rho_3)| + 3 \cdot |\Delta(\rho_3')| = 3
\]

We will label such an irreducible character \( \chi^{\Delta} = \chi^{(\Delta(\rho_1), \Delta(\rho_2), \Delta(\rho_3), \Delta(\rho_3))} \).

Three of these irreducibles will be the unipotent characters, mapping under the isomorphism from Proposition 4.50 as follows:

- \( s(3) = h_3 \mapsto \chi^{((3), \varnothing, \varnothing, \varnothing)} = 1_{GL_3} \) of dimension 1.
- \( s_{(2,1)} = \det \begin{bmatrix} h_2 & h_3 \\ 1 & h_1 \end{bmatrix} = h_2 h_3 - h_3 \mapsto \chi^{((2,1), \varnothing, \varnothing, \varnothing)} = \text{Ind}_{P_{2,1}}^{GL_3} 1_{P_{2,1}} - 1_{GL_3} \), of dimension \( \left[ \begin{array}{c} 3 \\ 2, 1 \end{array} \right]_q = [3]_q - 1 = q^2 + q \xrightarrow{q=2} 6. \)
- Lastly,
  \[
  s_{(1,1,1)} = \det \begin{bmatrix} h_1 & h_2 & h_3 \\ 1 & h_1 & h_2 \\ 0 & 1 & h_1 \end{bmatrix} = h_1^8 - h_2 h_1 - h_1 h_2 + h_3
  \]
  \( \mapsto \text{St}_3 = \chi^{((1,1,1), \varnothing, \varnothing, \varnothing)} = \text{Ind}_{B}^{GL_3} 1_B - \text{Ind}_{P_{2,1}}^{GL_3} 1_{P_{2,1}} - \text{Ind}_{P_{1,2}}^{GL_3} 1_{P_{1,2}} + 1_{GL_3} \),
  of dimension \( \left[ \begin{array}{c} 3 \\ 1, 1, 1 \end{array} \right]_q - \left[ \begin{array}{c} 3 \\ 2, 1 \end{array} \right]_q - \left[ \begin{array}{c} 3 \\ 1, 2 \end{array} \right]_q + \left[ \begin{array}{c} 3 \\ 3 \end{array} \right]_q \)
  \( = [3]_q! - [3]_q - [3]_q + 1 = q^3 \xrightarrow{q=2} 8. \)

There should also be one non-unipotent, non-cuspidal character, namely

\( \chi^{((1), (1), \varnothing, \varnothing)} = \rho_1 \rho_2 = \text{Ind}_{P_{1,2}}^{GL_3} \text{Ind}_{GL_1 \times GL_2}^{P_{1,2}} (1_{GL_1} \otimes \rho_2) \)

having dimension \( \left[ \begin{array}{c} 3 \\ 1, 2 \end{array} \right]_q \cdot 1 \cdot 1 = [3]_q \xrightarrow{q=2} 7. \)

Finally, we expect cuspidal characters \( \rho_3 = \chi^{(\varnothing, \varnothing, (1), \varnothing)} \), \( \rho_3' = \chi^{(\varnothing, \varnothing, (1), \varnothing)} \), whose dimensions \( d_3, d_3' \) can be deduced from the equation

\[
1^2 + 6^2 + 8^2 + 7^2 + d_3^2 + (d_3')^2 = |GL_3(\mathbb{F}_2)| = \left( (q^3 - q^0)(q^3 - q^1)(q^3 - q^2) \right)_{q=2} = 168.
\]

This forces \( d_3^2 + (d_3')^2 = 18 \), whose only solution in positive integers is \( d_3 = d_3' = 3 \).

We can check our predictions of the dimensions for the various \( GL_3(\mathbb{F}_2) \)-irreducible characters since \( GL_3(\mathbb{F}_2) \) is the finite simple group of order 168 (also isomorphic to \( PSL_3(\mathbb{F}_7) \)), with known character table (see James and Liebeck [78, p. 318]):

<table>
<thead>
<tr>
<th>( 1_{GL_3} )</th>
<th>( \chi^{((1), \varnothing, \varnothing, \varnothing)} )</th>
<th>( \chi^{((2,1), \varnothing, \varnothing, \varnothing)} )</th>
<th>( \chi^{((1,1,1), \varnothing, \varnothing, \varnothing)} )</th>
<th>( \chi^{((1), (1), \varnothing, \varnothing)} )</th>
<th>( \chi^{(\varnothing, \varnothing, (1), \varnothing)} )</th>
<th>( \chi^{(\varnothing, \varnothing, (1), \varnothing)} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>Dimension</td>
<td>168</td>
<td>8</td>
<td>4</td>
<td>3</td>
<td>7</td>
<td>7</td>
</tr>
<tr>
<td>Unipotent?/Cuspidal?</td>
<td>Yes</td>
<td>No</td>
<td>Yes</td>
<td>Yes</td>
<td>Yes</td>
<td>Yes</td>
</tr>
</tbody>
</table>

Here \( \alpha := -1/2 + i\sqrt{7}/2 \).

Remark 4.54. It is known (see e.g. Bump [24, Cor. 7.4]) that, for \( n \geq 2 \), the dimension of any cuspidal irreducible character \( \rho \) of \( GL_n(\mathbb{F}_q) \) is

\[
(q^{n-1} - 1)(q^{n-2} - 1) \cdots (q^2 - 1)(q - 1).
\]
Note that when \( q = 2 \),
- for \( n = 2 \) this gives \( 2^1 - 1 = 1 \) for the dimension of \( \rho_2 \), and
- for \( n = 3 \) it gives \( (2^2 - 1)(2 - 1) = 3 \) for the dimensions of \( \rho_3, \rho'_3 \),
agreeing with our calculations above. Much more is known about the character table of \( GL_n(\mathbb{F}_q) \); see Remark 4.68 below, Zelevinsky [170, Chap. 11], and Macdonald [106, Chap. IV].

4.9. The Hall algebra. There is another interesting Hopf subalgebra (and quotient Hopf algebra) of \( A(GL) \), related to unipotent conjugacy classes in \( GL_n(\mathbb{F}_q) \).

Definition 4.55. Say that an element \( g \) in \( GL_n(\mathbb{F}_q) \) is \textit{unipotent} if its eigenvalues are all equal to 1. Equivalently, \( g \in GL_n(\mathbb{F}_q) \) is unipotent if and only if \( g - \text{id}_{\mathbb{F}_q^n} \) is nilpotent. A conjugacy class in \( GL_n(\mathbb{F}_q) \) is \textit{unipotent} if its elements are unipotent.

Denote by \( H_n \) the \( \mathbb{C} \)-subspace of \( R_C(GL_n) \) consisting of those class functions which are supported only on unipotent conjugacy classes, and let \( H = \bigoplus_{n \geq 0} H_n \) as a \( \mathbb{C} \)-subspace of \( A_C(GL) = \bigoplus_{n \geq 0} R_C(GL_n) \).

Proposition 4.56. The subspace \( H \) is a Hopf subalgebra of \( A_C(GL) \), which is graded, connected, and of finite type, and self-dual with respect to the inner product on class functions inherited from \( A_C(GL) \). It is also a quotient Hopf algebra of \( A_C(GL) \), as the \( C \)-linear surjection \( A_C(GL) \to H \) restricting class functions to unipotent classes is a Hopf algebra homomorphism. This surjection has kernel \( H^\perp \), which is both an ideal and a two-sided coideal.

Proof. It is immediately clear that \( H^\perp \) is a graded \( \mathbb{C} \)-vector subspace of \( A_C(GL) \), whose \( n \)-th homogeneous component consists of those class functions on \( GL_n \) whose values on all unipotent classes are 0. (This holds no matter whether the perpendicular space is taken with respect to the Hermitian form \( \langle \cdot, \cdot \rangle_G \) or with respect to the bilinear form \( \langle \cdot, \cdot \rangle_G \).) In other words, \( H^\perp \) is the kernel of the surjection \( A_C(GL) \to H \) defined in the proposition.

Given two class functions \( \chi_i, \chi_j \) on \( GL_i, GL_j \) and \( g \) in \( GL_{i+j} \), one has
\[
(\chi_i \cdot \chi_j)(g) = \frac{1}{|P_{i,j}|} \sum_{h \in GL_{i+j}; \ h^{-1}gh = \begin{bmatrix} g_i & * \\ 0 & g_j \end{bmatrix} \in P_{i,j}} \chi_i(g_i)\chi_j(g_j).
\]
Since \( g \) is unipotent if and only if \( h^{-1}gh \) is unipotent and if and only if both \( g_i, g_j \) are unipotent, the formula (4.37) shows both that \( H \) is a subalgebra\(^{173}\) and that \( H^\perp \) is a two-sided ideal\(^{174}\). It also shows that the surjection \( A_C(GL) \to H \) restricting every class function to unipotent classes is an algebra homomorphism\(^{175}\).

Similarly, for class functions \( \chi \) on \( GL_n \) and \( (g_i, g_j) \) in \( GL_{i,j} = GL_i \times GL_j \), one has
\[
\Delta(\chi)(g_i, g_j) = \frac{1}{q^{ij}} \sum_{k \in \mathbb{F}_q^{i+j}} \chi \begin{bmatrix} g_i & k \\ 0 & g_j \end{bmatrix}
\]
using (4.13). This shows both that \( H \) is a sub-coalgebra of \( A = A_C(GL) \) (that is, it satisfies \( \Delta H \subset H \otimes H \)) and that \( H^\perp \) is a two-sided coideal (that is, we have \( \Delta(H^\perp) \subset H^\perp \otimes A + A \otimes H^\perp \)), since it shows that if \( \chi \) is supported only on unipotent classes, then \( \Delta(\chi) \) vanishes on \( (g_1, g_2) \) that have either \( g_1 \) or \( g_2 \) non-unipotent. It also shows that the surjection \( A_C(GL) \to H \) restricting every class function to unipotent classes is a coalgebra homomorphism. The rest follows. \( \square \)

The subspace \( H \) is called the \textit{Hall algebra}. It has an obvious orthogonal \( \mathbb{C} \)-basis, with interesting structure constants.

Definition 4.57. Given a partition \( \lambda \) of \( n \), let \( J_\lambda \) denote the \( GL_n \)-conjugacy class of unipotent matrices whose \textit{Jordan type} (that is, the list of the sizes of the Jordan blocks, in decreasing order) is given by \( \lambda \). Furthermore, let \( z_\lambda(q) \) denote the size of the centralizer of any element of this conjugacy class \( J_\lambda \).

\(^{173}\)Indeed, if \( \chi_i \) and \( \chi_j \) are both supported only on unipotent classes, then the same holds for \( \chi_i \cdot \chi_j \).
\(^{174}\)In fact, if one of \( \chi_i \) and \( \chi_j \) annihilates all unipotent classes, then so does \( \chi_i \cdot \chi_j \).
\(^{175}\)Because if \( g \) is unipotent, then the only values of \( \chi_i \) and \( \chi_j \) appearing on the right hand side of (4.37) are those on unipotent elements
The indicator class functions\(^{176}\) \(\{1_{I_\lambda}\}_{\lambda \in \text{Par}}\) form a \(\mathbb{C}\)-basis for \(\mathcal{H}\) whose multiplicative structure constants are called the Hall coefficients \(g_{\mu,\nu}^\lambda(q)\):

\[
1_{I_\lambda}1_{I_\nu} = \sum_{\lambda} g_{\mu,\nu}^\lambda(q) 1_{I_\lambda}.
\]

Because the dual basis to \(\{1_{I_\lambda}\}\) is \(\{z_\lambda(q)1_{I_\lambda}\}\), self-duality of \(\mathcal{H}\) shows that the Hall coefficients are (essentially) also structure constants for the comultiplication:

\[
\Delta 1_{I_\lambda} = \sum_{\mu,\nu} g_{\mu,\nu}^\lambda(q) \frac{z_\mu(q)z_\nu(q)}{z_\lambda(q)} \cdot 1_{I_\mu} \otimes 1_{I_\nu}.
\]

The Hall coefficient \(g_{\mu,\nu}^\lambda(q)\) has the following interpretation.

**Proposition 4.58.** Fix any \(g\) in \(GL_n(\mathbb{F}_q)\) acting unipotently on \(\mathbb{F}_q^n\) with Jordan type \(\lambda\). Then \(g_{\mu,\nu}^\lambda(q)\) counts the \(g\)-stable \(\mathbb{F}_q\)-subspaces \(V \subset \mathbb{F}_q^n\) for which the restriction \(g|V\) acts with Jordan type \(\mu\), and the induced map \(\bar{g}\) on the quotient space \(\mathbb{F}_q^n/V\) has Jordan type \(\nu\).

**Proof.** Given \(\mu,\nu\) partitions of \(i, j\) with \(i + j = n\), taking \(\chi_i, \chi_j\) equal to \(1_{J_\mu}, 1_{J_\nu}\) in (4.37) shows that for any \(g\) in \(GL_n\), the value of \(\left(\frac{1}{|P_{i,j}|}\left\{h \in GL_n : h^{-1}gh = \begin{bmatrix} g_i & * \\ 0 & g_j \end{bmatrix} \text{ with } g_i \in J_\mu, g_j \in J_\nu\right\}\right)\).

Let \(S\) denote the set appearing in (4.38), and let \(\mathbb{F}_q^i\) denote the \(i\)-dimensional subspace of \(\mathbb{F}_q^n\) spanned by the first \(i\) standard basis vectors. Note that the condition on an element \(h\) in \(S\) saying that \(h^{-1}gh\) is in block upper-triangular form can be re-expressed by saying that the subspace \(V := h(\mathbb{F}_q^i)\) is \(g\)-stable. One then sees that the map \(h \overset{\varphi}{\mapsto} V = h(\mathbb{F}_q^i)\) surjects \(S\) onto the set of \(i\)-dimensional \(g\)-stable subspaces \(V\) of \(\mathbb{F}_q^n\) for which \(g|V\) and \(\bar{g}\) are unipotent of types \(\mu, \nu\), respectively. Furthermore, for any particular such \(V\), its fiber \(\varphi^{-1}(V)\) in \(S\) is a coset of the stabilizer within \(GL_n\) of \(V\), which is conjugate to \(P_{i,j}\), and hence has cardinality \(|\varphi^{-1}(V)| = |P_{i,j}|\). This proves the assertion of the proposition. \(\square\)

The Hall algebra \(\mathcal{H}\) will turn out to be isomorphic to the ring \(\Lambda_C\) of symmetric functions with \(C\) coefficients, via a composite \(\varphi\) of three maps

\[
\Lambda_C \longrightarrow A(GL(i)_C) \longrightarrow A(GL)_C \longrightarrow \mathcal{H}
\]

in which the first map is the isomorphism from Proposition 4.50, the second is inclusion, and the third is the quotient map from Proposition 4.56.

**Theorem 4.59.** The above composite \(\varphi\) is a Hopf algebra isomorphism, sending

\[
\begin{align*}
\quad h_n & \longrightarrow \sum_{\lambda \in \text{Par}_n} 1_{I_\lambda}, \\
\quad e_n & \longrightarrow q^\binom{n}{2} 1_{I_n}, \\
\quad p_n & \longrightarrow \sum_{\lambda \in \text{Par}_n} (q; q)_\ell(\lambda) - 1 1_{I_\lambda} \quad \text{for } n > 0,
\end{align*}
\]

where

\[
(x; q)_m := (1 - x)(1 - qx)(1 - q^2 x) \cdots (1 - q^{m-1} x) \quad \text{for all } m \in \mathbb{N}.
\]

**Proof.** That \(\varphi\) is a graded Hopf morphism follows because it is a composite of three such morphisms. We claim that once one shows the formula for the (nonzero) image \(\varphi(p_n)\) given above is correct, then this will already show \(\varphi\) is an isomorphism, by the following argument. Note first that \(\Lambda_C\) and \(\mathcal{H}\) both have dimension \(|\text{Par}_n|\) for their \(n^{th}\) homogeneous components, so it suffices to show that the graded map \(\varphi\) is injective. On the other hand, both \(\Lambda_C\) and \(\mathcal{H}\) are (graded, connected, finite type) self-dual Hopf algebras (although with respect to a sesquilinear form), so Theorem 3.8 says that each is the symmetric algebra on its space of primitive elements. Thus it suffices to check that \(\varphi\) is injective when restricted to their subspaces.

\(^{176}\) Here we use the following notation: Whenever \(P\) is a subset of a group \(G\), we denote by \(1_P\) the map \(G \to C\) which sends every element of \(P\) to 1 and all remaining elements of \(G\) to 0. This is not in conflict with the notation \(1_G\) for the trivial character of \(G\), since \(1_P = 1_G\) for \(P = G\). Note that \(1_P\) is a class function when \(P\) is a union of conjugacy classes of \(G\).
of primitives.\(^{177}\) For \(\Lambda_C\), by Corollary 3.9 the primitives are spanned by \(\{p_1, p_2, \ldots\}\), with only one basis element in each degree \(n \geq 1\). Hence \(\varphi\) is injective on the subspace of primitives if and only if it does not annihilate any \(p_n\).

Thus it only remains to show the above formulas for the images of \(h_n, e_n, p_n\) under \(\varphi\). This is clear for \(h_n\), since Proposition 4.50 shows that it maps under the first two composites to the indicator function \(\mathbb{1}_{GL_n}\) which then restricts to the sum of indicators \(\sum_{\lambda \in \text{Par}_n} \mathbb{1}_{J_\lambda}\) in \(\mathcal{H}\). For \(e_n, p_n\), we resort to generating functions. Let \(\hat{h}_n, \hat{e}_n, \hat{p}_n\) denote the three putative images in \(\mathcal{H}\) of \(h_n, e_n, p_n\), appearing on the right side in the theorem, and define generating functions
\[
\hat{H}(t) := \sum_{n \geq 0} \hat{h}_n t^n, \quad \hat{E}(t) := \sum_{n \geq 0} \hat{e}_n t^n, \quad \hat{P}(t) := \sum_{n \geq 0} \hat{p}_n t^n \quad \text{in } \mathcal{H}[[t]].
\]

We wish to show that the map \(\varphi[[t]] : \Lambda_C[[t]] \to \mathcal{H}[[t]]\) (induced by \(\varphi\)) maps \(H(t), E(t), P(t)\) in \(\Lambda[[t]]\) to these three generating functions\(^{178}\). Since we have already shown this is correct for \(H(t)\), by (2.16), (2.30), it suffices to check that in \(\mathcal{H}[[t]]\) one has

\[
\hat{H}(t) \hat{E}(-t) = 1, \quad \text{or equivalently,} \quad \sum_{k=0}^{n} (-1)^k \hat{e}_k \hat{h}_{n-k} = \delta_{0,n};
\]

\[
\hat{H}'(t) \hat{E}(-t) = \hat{P}(t), \quad \text{or equivalently,} \quad \sum_{k=0}^{n} (-1)^k (n-k) \hat{e}_k \hat{h}_{n-k} = \hat{p}_n.
\]

Thus it would be helpful to evaluate the class function \(\hat{e}_k \hat{h}_{n-k}\). Note that a unipotent \(g\) in \(GL_n\) having \(\ell\) Jordan blocks has an \(\ell\)-dimensional 1-eigenspace, so that the number of \(k\)-dimensional \(q\)-stable \(\mathbb{F}_q\)-subspaces of \(\mathbb{F}_q^n\) on which \(g\) has Jordan type \((1^k)\) (that is, on which \(g\) acts as the identity) is the \(q\)-binomial coefficient
\[
\binom{\ell}{k}_q = \frac{(q; q)_\ell}{(q; q)_k (q; q)_{\ell-k}},
\]
counting \(k\)-dimensional \(\mathbb{F}_q\)-subspaces \(V\) of an \(\ell\)-dimensional \(\mathbb{F}_q\)-vector space; see, e.g., [154, §1.7]. Hence, for a unipotent \(g\) in \(GL_n\) having \(\ell\) Jordan blocks, we have
\[
(\hat{e}_k \hat{h}_{n-k})(g) = q^{\binom{k}{2}} \cdot \left( \frac{1}{J_{\ell, k}} ; \hat{h}_{n-k} \right)(g) = q^{\binom{k}{2}} \cdot \sum_{\nu \in \text{Par}_{n-k}} \left( \frac{1}{J_{\ell, k}} ; \frac{1}{J_{\nu}} \right)(g) = q^{\binom{k}{2}} \left[ \frac{\ell}{k} \right]_q
\]
(by Proposition 4.58). Thus one needs for \(\ell \geq 1\) that
\[
\sum_{k=0}^{\ell} (-1)^k q^{\binom{k}{2}} \left[ \frac{\ell}{k} \right]_q = 0,
\]
\[
\sum_{k=0}^{\ell} (-1)^k (n-k) q^{\binom{k}{2}} \left[ \frac{\ell}{k} \right]_q = (q; q)_{\ell-1}.
\]

Identity (4.39) comes from setting \(x = 1\) in the \(q\)-binomial theorem [154, Exer. 3.119]:
\[
\sum_{k=0}^{\ell} (-1)^k q^{\binom{k}{2}} \left[ \frac{\ell}{k} \right]_q x^{\ell-k} = (x-1)(x-q)(x-q^2) \cdots (x-q^{\ell-1}).
\]

Identity (4.40) comes from taking \(\frac{d}{dx}\) in the \(q\)-binomial theorem, then setting \(x = 1\), and finally adding \((n-\ell)\) times (4.39).

\[\text{Exercise 4.60.} \text{ Fix a prime power } q. \text{ For any } k \in \mathbb{N}, \text{ and any } k \text{ partitions } \lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)}, \text{ we define a family } \left( g_{\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)}}^\lambda (q) \right)_{\lambda \in \text{Par}} \text{ of elements of } \mathbb{C} \text{ by the equation}
\]
\[
\frac{1}{J_{\lambda^{(1)}}} \frac{1}{J_{\lambda^{(2)}}} \cdots \frac{1}{J_{\lambda^{(k)}}} = \sum_{\lambda \in \text{Par}} g_{\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)}}^\lambda (q) \frac{1}{J_{\lambda}}
\]
in \(\mathcal{H}\). This notation generalizes the notation \(g_{\lambda, \nu}^\mu (q)\) we introduced in Definition 4.57. Note that \(g_{\lambda, \emptyset}^\mu (q) = \delta_{\lambda, \emptyset}\) for any two partitions \(\lambda\) and \(\mu\), and that \(g^\lambda (q) = \delta_{\lambda, \emptyset}\) for any partition \(\lambda\) (where \(g^\lambda (q)\) is to be understood as \(g_{\lambda, \emptyset}^\emptyset (q)\) for \(k = 0\)).

\(^{177}\)An alternative way to see that it suffices to check this is by recalling Exercise 1.54(c).

\(^{178}\)See (2.14), (2.15), (2.30) for the definitions of \(H(t), E(t), P(t)\).
(a) Let \( \lambda \in \text{Par} \), and let \( n = |\lambda| \). Let \( V \) be an \( n \)-dimensional \( \mathbb{F}_q \)-vector space, and let \( g \) be a unipotent endomorphism of \( V \) having Jordan type \( \lambda \). Let \( k \in \mathbb{N} \), and let \( \lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)} \) be \( k \) partitions. A \((\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)})\)-compatible \( g \)-flag will mean a sequence \( 0 = V_0 \subset V_1 \subset V_2 \subset \cdots \subset V_k = V \) of \( g \)-invariant \( \mathbb{F}_q \)-vector subspaces \( V_i \) of \( V \) such that for every \( i \in \{1, 2, \ldots, k\} \), the endomorphism of \( V_i/V_{i-1} \) induced by \( g \) has Jordan type \( \lambda^{(i)} \).

Show that \( g^\lambda_{(\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)})}(q) \) is the number of \((\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)})\)-compatible \( g \)-flags.\(^{179}\)

(b) Let \( \lambda \in \text{Par} \). Let \( k \in \mathbb{N} \), and let \( \lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)} \) be \( k \) partitions. Show that \( g^\lambda_{(\lambda^{(1)}, \lambda^{(2)}, \ldots, \lambda^{(k)})}(q) = 0 \) unless \( |\lambda^{(1)}| + |\lambda^{(2)}| + \cdots + |\lambda^{(k)}| = |\lambda| \) and \( \lambda^{(1)} + \lambda^{(2)} + \cdots + \lambda^{(k)} \sim \lambda \). (Here and in the following, we are using the notations of Exercise 2.93).

(c) Let \( \lambda \in \text{Par} \), and let us write the transpose partition \( \lambda' \) as \( \lambda' = ((\lambda')_1, (\lambda')_2, \ldots, (\lambda')_q) \). Show that \( g^\lambda_{((\lambda')_1), (\lambda')_2, \ldots, (\lambda')_q}(q) \neq 0 \).

(d) Let \( n \in \mathbb{N} \) and \( \lambda \in \text{Par}_n \). Show that

\[
\varphi(e_\lambda) = \sum_{\mu \in \text{Par}_n : \lambda' \sim \mu} \alpha_{\lambda, \mu} 1_{J_\mu}
\]

for some coefficients \( \alpha_{\lambda, \mu} \in \mathbb{C} \) satisfying \( \alpha_{\lambda, \lambda'} \neq 0 \).

(e) Give another proof of the fact that the map \( \varphi \) is injective.

[Hint: For (b), use Exercise 2.96(b).]

We next indicate, without proof, how \( \mathcal{H} \) relates to the classical Hall algebra.

**Definition 4.61.** Let \( p \) be a prime. The usual Hall algebra, or what Schiffmann [140, §2.3] calls Steinitz’s classical Hall algebra (see also Macdonald [106, Chap. II]), has \( \mathbb{Z} \)-basis elements \( \{u_\lambda\}_{\lambda \in \text{Par}} \), with the multiplicative structure constants \( g^\lambda_{\mu, \nu}(p) \) in

\[
u \mu = \sum_{\lambda} g^\lambda_{\mu, \nu}(p) u_\lambda
\]

defined as follows: fix a finite abelian \( p \)-group \( L \) of type \( \lambda \), meaning that

\[
\lambda \in \bigoplus_{i=1}^{\ell(\lambda)} \mathbb{Z}/p^{i} \mathbb{Z},
\]

and let \( g^\lambda_{\mu, \nu}(p) \) be the number of subgroups \( M \) of \( L \) of type \( \mu \), for which the quotient \( N := L/M \) is of type \( \nu \). In other words, \( g^\lambda_{\mu, \nu}(p) \) counts, for a fixed abelian \( p \)-group \( L \) of type \( \lambda \), the number of short exact sequences \( 0 \rightarrow M \rightarrow L \rightarrow N \rightarrow 0 \) in which \( M, N \) have types \( \mu, \nu \), respectively (modulo isomorphism of short exact sequences restricting to the identity on \( L \)).

We claim that when one takes the finite field \( \mathbb{F}_q \) of order \( q = p \) a prime, the \( \mathbb{Z} \)-linear map

\[
u \mu \mapsto 1_{J_\lambda}
\]

gives an isomorphism from this classical Hall algebra to the \( \mathbb{Z} \)-algebra \( \mathcal{H}_\mathbb{Z} \subset \mathcal{H} \). The key point is Hall’s Theorem, a non-obvious statement for which Macdonald includes two proofs in [106, Chap. II], one of them due to Zelevinsky.\(^{181}\) To state it, we first recall some notions about discrete valuation rings.

**Definition 4.62.** A discrete valuation ring (DVR) \( \mathfrak{o} \) is a principal ideal domain having only one maximal ideal \( \mathfrak{m} \neq 0 \), with quotient \( k = \mathfrak{o}/\mathfrak{m} \) called its residue field.

\(^{179}\)This is well-defined. In fact, both \( V_1 \) and \( V_i - 1 \) are \( g \)-invariant, so that \( g \) restricts to an endomorphism of \( V_i \), which further restricts to an endomorphism of \( V_i - 1 \), and thus gives rise to an endomorphism of \( V_i/V_{i-1} \).

\(^{180}\)This can be seen as a generalization of Proposition 4.58. In fact, if \( \mu \) and \( \nu \) are two partitions, then a \((\mu, \nu)\)-compatible \( g \)-flag is a sequence \( 0 = V_0 \subset V_1 \subset V_2 = V \) of \( g \)-invariant \( \mathbb{F}_q \)-vector subspaces \( V_i \) of \( V \) such that the endomorphism of \( V_i/V_0 \) is induced by \( g \) has Jordan type \( \mu \), and the endomorphism of \( V_2/V_0 \) is induced by \( g \) has Jordan type \( \nu \). Choosing such a sequence amounts to choosing \( V_0 \) (since there is only one choice for each of \( V_0 \) and \( V_2 \)), and the conditions on this \( V_i \) are precisely the conditions on \( V \) in Proposition 4.58.

\(^{181}\)See also [140, Thm. 2.6, Prop. 2.7] for quick proofs of part of it, similar to Zelevinsky’s. Another proof, based on a recent category-theoretical paradigm, can be found in [45, Theorem 3.53].

Theorem 4.65. Assume for all partitions $\mu$ with respect to its basis $(\lambda_{i})_{i=1}^{\ell}$ of type $\lambda$ in this situation.

Here are the two crucial examples for us.

**Example 4.63.** For any field $F$, the power series ring $\mathfrak{o} = F[[t]]$ is a DVR with maximal ideal $m = (t)$ and residue field $k = \mathfrak{o}/m = F[[t]]/(t) \cong F$. An $\mathfrak{o}$-module $L$ of type $\lambda$ is an $F$-vector space together with an $F$-linear transformation $T \in \text{End} L$ that acts on $L$ nilpotently (so that $g := T + 1$ acts unipotently, where $1 = \text{id}_L$) with Jordan blocks of sizes given by $\lambda$: each summand $\mathfrak{o}/m^{\lambda_i} = F[[t]]/(t^{\lambda_i})$ of $L$ has an $F$-basis \{1, t, t^2, \ldots, t^{\lambda_i-1}\} on which the map $T$ that multiplies by $t$ acts as a nilpotent Jordan block of size $\lambda_i$. Note also that, in this setting, $\mathfrak{o}$-submodules are the same as $T$-stable (or $g$-stable) $F$-subspaces.

**Example 4.64.** The ring of $p$-adic integers $\mathfrak{o} = \mathbb{Z}_p$ is a DVR with maximal ideal $m = (p)$ and residue field $k = \mathfrak{o}/m = \mathbb{Z}_p/p\mathbb{Z}_p \cong \mathbb{Z}/p\mathbb{Z}$. An $\mathfrak{o}$-module $L$ of type $\lambda$ is an abelian $p$-group of type $\lambda$: for each summand, $\mathfrak{o}/m^{\lambda_i} = \mathbb{Z}_p/p^{\lambda_i}\mathbb{Z}_p \cong \mathbb{Z}/p^{\lambda_i}\mathbb{Z}$. Note also that, in this setting, $\mathfrak{o}$-submodules are the same as subgroups.

One last notation: $n(\lambda) := \sum_{i \geq 1} (i - 1)\lambda_i$, for $\lambda$ in Par. Hall’s Theorem is as follows.

**Theorem 4.65.** Assume $\mathfrak{o}$ is a DVR with maximal ideal $m$, and that its residue field $k = \mathfrak{o}/m$ is finite of cardinality $q$. Fix an $\mathfrak{o}$-module $L$ of type $\lambda$. Then the number of $\mathfrak{o}$-submodules $M$ of type $\mu$ for which the quotient $N = L/M$ is of type $\nu$ can be written as the specialization

$$[g_{\mu,\nu}^\lambda(t)]_{t=q}$$

of a polynomial $g_{\mu,\nu}^\lambda(t)$ in $\mathbb{Z}[t]$, called the **Hall polynomial**.

Furthermore, the Hall polynomial $g_{\mu,\nu}^\lambda(t)$ has degree at most $n(\lambda) - (n(\mu) + n(\nu))$, and its coefficient of $t^{n(\lambda) - (n(\mu) + n(\nu))}$ is the Littlewood-Richardson coefficient $c_{\mu,\nu}^\lambda$.

Comparing what Hall’s Theorem says in Examples 4.63 and 4.64, shows that the map (4.41) gives the desired isomorphism from the classical Hall algebra to $H_{\mathbb{Z}}$.

We close this section with some remarks on the vast literature on Hall algebras that we will not discuss here.

**Remark 4.66.** Macdonald’s version of Hall’s Theorem [106, (4.3)] is stronger than Theorem 4.65, and useful for certain applications: he shows that $g_{\mu,\nu}^\lambda(t)$ is the zero polynomial whenever the Littlewood-Richardson coefficient $c_{\mu,\nu}^\lambda$ is zero.

**Remark 4.67.** In general, not all coefficients of the Hall polynomials $g_{\mu,\nu}^\lambda(t)$ are nonnegative (see Butler/Hales [26] for a study of when they are): it often happens that $g_{\mu,\nu}^\lambda(1) = 0$ despite $g_{\mu,\nu}^\lambda(t)$ not being the zero polynomial\(^{182}\). However, in [80, Thm. 4.2], Klein showed that the polynomial values $g_{\mu,\nu}^\lambda(p)$ for $p$ prime are always positive when $c_{\mu,\nu}^\lambda \neq 0$. (This easily yields the same result for $p$ a prime power.)

**Remark 4.68.** Zelevinsky in [170, Chaps 10, 11] uses the isomorphism $\Lambda_\mathbb{C} \rightarrow H$ to derive J. Green’s formula for the value of any irreducible character $\chi$ of $GL_n$ on any unipotent class $J_\chi$. The answer involves values of irreducible characters of $S_n$ along with Green’s polynomials $Q^\lambda(g)$ (see Macdonald [106, §III.7]; they are denoted $Q^\lambda(g)$ by Zelevinsky), which express the images under the isomorphism of Theorem 4.59 of the symmetric function basis $\{p_\lambda\}$ in terms of the basis $\{\Delta_\lambda\}$.

**Remark 4.69.** The Hall polynomials $g_{\mu,\nu}^\lambda(t)$ also essentially give the multiplicative structure constants for $\Lambda(\mathfrak{o})[[t]]$ with respect to its basis of Hall-Littlewood symmetric functions $P_\lambda = P_\lambda(x; t)$:

$$P_\mu P_\nu = \sum_\lambda t^{n(\lambda) - (n(\mu) + n(\nu))} g_{\mu,\nu}^\lambda(t^{-1}) P_\lambda.$$  

See Macdonald [106, §III.3].

\(^{182}\)Actually, Butler/Hales show in [26, proof of Prop. 2.4] that the values $g_{\mu,\nu}^\lambda(1)$ are the structure constants of the ring $\Lambda$ with respect to its basis $(m_\lambda)_{\lambda \in Par^t}$, we have

$$m_\mu m_\nu = \sum_{\lambda \in Par^t} g_{\mu,\nu}^\lambda(1)m_\lambda$$

for all partitions $\mu$ and $\nu$.  

Remark 4.70. Schiffmann [140] discusses self-dual Hopf algebras which vastly generalize the classical Hall algebra called Ringel-Hall algebras, associated to abelian categories which are hereditary. Examples come from categories of nilpotent representations of quivers; the quiver having exactly one node and one arc recovers the classical Hall algebra $\mathcal{H}_Z$ discussed above.

Remark 4.71. The general linear groups $GL_n(F_q)$ are one of four families of so-called classical groups. Progress has been made on extending Zelevinsky’s PSH theory to the other families:

(a) Work of Thiem and Vinroot [164] shows that the tower $\{G_\star\}$ of finite unitary groups $U_n(F_{q^2})$ give rise to another positive self-dual Hopf algebra $A = \bigoplus_{n \geq 0} R(U_n(F_{q^2}))$, in which the role of Harish-Chandra induction is played by Deligne-Lusztig induction. In this theory, character and degree formulas for $U_n(F_{q^2})$ are related to those of $GL_n(F_q)$ by substituting $q \mapsto -q$, along with appropriate scalings by $\pm 1$, a phenomenon sometimes called Ennola duality. See also [155, §4].

(b) van Leeuwen [94] has studied $\bigoplus_{n \geq 0} R(Sp_{2n}(F_q))$, $\bigoplus_{n \geq 0} R(O_{2n}(F_q))$ and $\bigoplus_{n \geq 0} R(U_n(F_{q^2}))$ not as Hopf algebras, but rather as so-called twisted PSH-modules over the PSH $A(GL)$ (a “deformed” version of the older notion of Hopf modules). He classified these PSH-modules axiomatically similarly to Zelevinsky’s above classification of PSH’s.

(c) In a recent honors thesis [149], Shelley-Abrahamson defined yet another variation of the concept of Hopf modules, named 2-compatible Hopf modules, and identified $\bigoplus_{n \geq 0} R(Sp_{2n}(F_q))$ and $\bigoplus_{n \geq 0} R(O_{2n+1}(F_q))$ as such modules over $A(GL)$. 
5. Quasisymmetric functions and $P$-partitions

We discuss here our next important example of a Hopf algebra arising in combinatorics: the quasisymmetric functions of Gessel [59], with roots in work of Stanley [151] on $P$-partitions.

5.1. Definitions, and Hopf structure. The definitions of quasisymmetric functions require a totally ordered variable set. Usually we will use a variable set denoted $x = (x_1, x_2, \ldots)$ with the usual ordering $x_1 < x_2 < \ldots$. However, it is good to have some flexibility in changing the ordering, which is why we make the following definition.

Definition 5.1. Given any totally ordered set $I$, create a totally ordered variable set $\{x_i\}_{i \in I}$, and then let $R(\{x_i\}_{i \in I})$ denote the power series of bounded degree in $\{x_i\}_{i \in I}$ having coefficients in $k$.

The quasisymmetric functions $QSym(\{x_i\}_{i \in I})$ over the alphabet $(\{x_i\}_{i \in I})$ will be the $k$-submodule consisting of the elements $f$ in $R(\{x_i\}_{i \in I})$ that have the same coefficient on the monomials $x_{i_1}^{\alpha_1} \cdots x_{i_\ell}^{\alpha_\ell}$ and $x_{j_1}^{\beta_1} \cdots x_{j_\ell}^{\beta_\ell}$ whenever both $i_1 < \cdots < i_\ell$ and $j_1 < \cdots < j_\ell$ in the total order on $I$. We write $QSym_k(\{x_i\}_{i \in I})$ instead of $QSym(\{x_i\}_{i \in I})$ to stress the choice of base ring $k$.

It immediately follows from this definition that $QSym(\{x_i\}_{i \in I})$ is a free $k$-submodule of $R(\{x_i\}_{i \in I})$, having as $k$-basis elements the monomial quasisymmetric functions

$$M_\alpha(\{x_i\}_{i \in I}) := \sum_{i_1 < \cdots < i_\ell \in I} x_{i_1}^{\alpha_1} \cdots x_{i_\ell}^{\alpha_\ell}$$

for all compositions $\alpha$ satisfying $\ell(\alpha) \leq |I|$. When $I$ is infinite, this means that the $M_\alpha$ for all compositions $\alpha$ form a basis of $QSym(\{x_i\}_{i \in I})$.

Note that $QSym(\{x_i\}_{i \in I}) = \bigoplus_{n \geq 0} QSym_n(\{x_i\}_{i \in I})$ is a graded $k$-module of finite type, where $QSym_n(\{x_i\}_{i \in I})$ is the $k$-submodule of quasisymmetric functions which are homogeneous of degree $n$. Letting Comp denote the set of all compositions $\alpha$, and $Comp_n$ the compositions $\alpha$ of $n$ (that is, compositions whose parts sum to $n$), the subset $\{M_\alpha \mid \alpha \in Comp_n, \ell(\alpha) \leq |I|\}$ gives a $k$-basis for $QSym_n(\{x_i\}_{i \in I})$.

Example 5.2. Taking the variable set $x = (x_1 < x_2 < \cdots)$ to define $QSym(x)$, for $n = 0, 1, 2, 3$, one has these basis elements in $QSym_n(x)$:

$$M_{()} = M_{\emptyset} = 1$$

$$M_{(1)} = x_1 + x_2 + x_3 + \cdots$$

$$M_{(1)} = x_1^2 + x_2^2 + x_3^2 + \cdots$$

$$M_{(1)} = x_1 x_2 + x_1 x_3 + x_2 x_3 + \cdots$$

$$M_{(2)} = x_1^3 + x_2^3 + x_3^3 + \cdots$$

$$M_{(2)} = x_1^2 x_2 + x_1^2 x_3 + x_2^2 x_3 + \cdots$$

$$M_{(2)} = x_1 x_2^2 + x_1 x_3^2 + x_2 x_3^2 + \cdots$$

$$M_{(3)} = x_1 x_2 x_3 + x_1 x_2 x_4 + x_1 x_3 x_4 + \cdots$$

It is not obvious that $QSym(x)$ is a subalgebra of $R(x)$, but we will show this momentarily. For example,

$$M_\alpha M_{(b,c)} = (x_1^a + x_2^a + x_3^a + \cdots)(x_1^b x_2 + x_1^b x_3 + x_2^b x_3 + \cdots)$$

$$= x_1^{a+b} x_2 + x_1^{a} x_3 x_2 + x_1^{a} x_3 x_2 + x_2^{a} x_3 x_2 + \cdots$$

$$= M_{(a+b,c)} + M_{(a,b+c)} + M_{(a,b,c)} + M_{(a,b,c)} + M_{(b,a,c)}$$

Proposition 5.3. For any infinite totally ordered set $I$, one has that $QSym(\{x_i\}_{i \in I})$ is a $k$-subalgebra of $R(\{x_i\}_{i \in I})$, with multiplication in the $\{M_\alpha\}$-basis as follows: Fix three pairwise disjoint chain posets $(i_1 < \cdots < i_\ell)$, $(j_1 < \cdots < j_m)$ and $(k_1 < k_2 < \cdots)$. Now, if $\alpha = (\alpha_1, \ldots, \alpha_\ell)$, $\beta = (\beta_1, \ldots, \beta_m)$ then

$$(5.1) \quad M_\alpha M_\beta = \sum_f M_{\alpha(b), \beta(f)}$$

in which the sum is over all $p \in \mathbb{N}$ and all maps $f$ from the disjoint union of two chains to a chain $i_1 < \cdots < i_\ell \cup j_1 < \cdots < j_m \to k_1 < \cdots < k_p$. 

\[M_{(i_1 < \cdots < i_\ell)} \cup (j_1 < \cdots < j_m) \to (k_1 < \cdots < k_p)\]
which are both surjective and strictly order-preserving \((x < y \implies f(x) < f(y))\), and where the composition \(w(f) := (w_1(f), \ldots, w_p(f))\) is defined by \(w_s(f) := \sum_{i\in f^{-1}(k_s)} \alpha_i + \sum_{j\in f^{-1}(k_s)} \beta_j\). 

**Example 5.4.** For this example, set \(\alpha = (2, 1)\) and \(\beta = (3, 4, 2)\). Let us compute \(M_2M_3\) using (5.1). Indeed, the length of \(\alpha\) is \(\ell = 2\), and the length of \(\beta\) is \(m = 3\), so the sum on the right hand side of (5.1) is a sum over all \(p \in \mathbb{N}\) and all surjective strictly order-preserving maps \(f\) from the disjoint union \((i_1 < i_2) \cup (j_1 < j_2 < j_3)\) of two chains to the chain \((k_1 < k_2 < \cdots < k_p)\). Such maps can exist only when \(p \leq 5\) (due to having to be surjective) and only for \(p \geq 3\) (since, being strictly order-preserving, they have to be injective when restricted to \((j_1 < j_2 < j_3)\)). Hence, enumerating them is a finite problem. The reader can check that the value obtained for \(M_2M_3\) is 

\[
M_{(2,1,3,4,2)} + M_{(2,3,1,4,2)} + M_{(2,3,4,1,2)} + M_{(2,3,4,2,1)} + M_{(3,2,1,4,2)} \\
+ M_{(3,2,4,1,2)} + M_{(3,2,4,2,1)} + M_{(3,4,2,1,2)} + M_{(3,4,2,2,1)} + M_{(3,4,4,2,1)} \\
+ M_{(3,2,4,3)} + M_{(3,3,5,2)} + M_{(3,4,4,2)} + M_{(3,2,4,3)} + M_{(3,2,5,3)} + M_{(3,4,2,3)} \\
+ M_{(3,4,4,1)} + M_{(3,6,1,2)} + M_{(3,6,2,1)} + M_{(5,1,4,2)} + M_{(5,4,1,2)} + M_{(5,4,2,1)} \\
+ M_{(5,4,3)} + M_{(5,5,2)} + M_{(3,6,3)}.
\]

Here, we have listed the addends corresponding to \(p = 5\) on the first two rows, the addends corresponding to \(p = 4\) on the next two rows, and those corresponding to \(p = 3\) on the third row. The reader might notice that the first two rows (i.e., the addends with \(p = 5\)) are basically a list of shuffles of \(\alpha\) and \(\beta\): In general, the maps (5.2) for \(p = \ell + m\) are in bijection with the elements of \(\text{Sh}_{\ell,m}\) \(^{183}\), and the corresponding compositions \(wt\) are the shuffles of \(\alpha\) and \(\beta\). Therefore the name “overlapping shuffle product”.

**Proof of Proposition 5.3.** Formula (5.1) comes from considering a typical product of two monomials in the expansion of Formula (5.1) comes from considering a typical product of two monomials in the expansion of

\[
(x_1^{\alpha_1} \cdots x_i^{\alpha_i}) (x_j^{\beta_1} \cdots x_{j_m}^{\beta_m}) = x_{k_1}^{\gamma_1} \cdots x_{k_p}^{\gamma_p}
\]

for subscript sequences \(i_1 < \cdots < i_\ell\) and \(j_1 < \cdots < j_m\) and \(k_1 < \cdots < k_p\) with

\[
\{k_1, \ldots, k_p\} = \{i_1, \ldots, i_\ell\} \cup \{j_1, \ldots, j_m\}.
\]

Thinking of \(\{i_s\}\) and \(\{j_t\}\) as disjoint sets, multiplication gives a surjective map \(f\) as in (5.2), with \(\gamma_s = wt_s(f)\).

The multiplication rule (5.1) shows that the \(k\)-algebra \(\text{QSym}(\{x_i\}_{i \in I})\) does not depend much on \(I\), as long as \(I\) is infinite. More precisely, all such \(k\)-algebras are mutually isomorphic. We can use this to define a \(k\)-algebra of quasisymmetric functions without any reference to \(I\):

**Definition 5.5.** Let \(\text{QSym}\) be the \(k\)-algebra defined as having \(k\)-basis \(\{M_{\alpha}\}_{\alpha \in \text{Comp}}\) and with multiplication defined \(k\)-linearly by (5.1). This is called the \(k\)-algebra of quasisymmetric functions. We write \(\text{QSym}_k\) instead of \(\text{QSym}\) to stress the choice of base ring \(k\).

The \(k\)-algebra \(\text{QSym}\) is graded, and its \(n\)-th graded component \(\text{QSym}_n\) has \(k\)-basis \(\{M_{\alpha}\}_{\alpha \in \text{Comp}_n}\).

For every infinite totally ordered set \(I\), the \(k\)-algebra \(\text{QSym}\) is isomorphic to the \(k\)-algebra \(\text{QSym}(\{x_i\}_{i \in I})\).

The isomorphism sends \(M_{\alpha} \mapsto M_{\alpha}(\{x_i\}_{i \in I})\).

In particular, we obtain the isomorphism \(\text{QSym} \cong \text{QSym}(x)\) for \(x\) being the infinite chain \((x_1 < x_2 < x_3 < \cdots)\).

We will identify \(\text{QSym}\) with \(\text{QSym}(x)\) along this isomorphism. This allows us to regard quasisymmetric functions either as power series in a specific set of variables (“alphabet”), or as formal linear combinations of \(M_{\alpha}\)’s, whatever is more convenient.

For any infinite alphabet \(\{x_i\}_{i \in I}\) and any \(f \in \text{QSym}\), we denote by \(f(\{x_i\}_{i \in I})\) the image of \(f\) under the algebra isomorphism \(\text{QSym} \to \text{QSym}(\{x_i\}_{i \in I})\) defined in Definition 5.5.

The comultiplication of \(\text{QSym}\) will extend the one that we defined for \(\Lambda\), but we need to take care about the order of the variables this time. We consider the linear order from (2.12) on two sets of variables \((x, y) = (x_1 < x_2 < \cdots < y_1 < y_2 < \cdots)\), and we embed the \(k\)-algebra \(\text{QSym}(x) \otimes \text{QSym}(y)\) into the \(k\)-algebra \(R(x, y)\) by identifying every \(f \otimes g \in \text{QSym}(x) \otimes \text{QSym}(y)\) with \(fg \in R(x, y)\) (this embedding is

\(^{183}\) The bijection takes a map \(f\) to the inverse of the permutation \(\sigma \in \mathfrak{S}_p\) which sends every \(x \in \{1, 2, \ldots, \ell\}\) to the index \(y\) satisfying \(f(x) = k_y\), and sends every \(x \in \{\ell + 1, \ell + 2, \ldots, \ell + m\}\) to the index \(y\) satisfying \(f(x) = k_y\).
indeed injective, as can easily be verified on the basis elements $M_\beta(x) \otimes M_\gamma(y)$ of the former $k$-algebra. It can then be seen that

$$QSym(x, y) \subset QSym(x) \otimes QSym(y)$$

(where the right hand side is viewed as $k$-subalgebra of $R(x, y)$ via said embedding)\(^{184}\), so that one can define $QSym \overset{\Delta}{\to} QSym \otimes QSym$ as the composite of the maps in the bottom row here:

$$R(x, y) \cup f \mapsto f(x, y) = f(x_1, x_2, \ldots, y_1, y_2, \ldots).$$

(Recall that $f(x, y)$ is formally defined as the image of $f$ under the algebra isomorphism $QSym \to QSym(x, y)$ defined in Definition 5.5.)

**Example 5.6.** For example,

$$\Delta M_{(a, b, c)} = M_{(a, b, c)}(x_1, x_2, \ldots, y_1, y_2, \ldots)$$

$$= x_1^a x_2^b x_3^c + x_1^b x_2^a x_3^c + \cdots$$

$$+ x_1^a x_2^c y_1^c + x_1^b x_2^c y_1^c + \cdots$$

$$+ x_1^a y_1^c y_2^c + x_1^b y_1^c y_2^c + \cdots$$

$$+ y_1^a y_2^b y_3^c + y_1^b y_2^b y_3^c + \cdots$$

$$= M_{(a, b, c)}(x) + M_{(a, b)}(x) M_{(c)}(y) + M_{(a)}(x) M_{(b, c)}(y) + M_{(a, b, c)}(y)$$

Defining the *concatenation* $\beta \cdot \gamma$ of two compositions $\beta = (\beta_1, \ldots, \beta_r), \gamma = (\gamma_1, \ldots, \gamma_s)$ to be the composition $(\beta_1, \ldots, \beta_r, \gamma_1, \ldots, \gamma_s)$, one has the following description of the coproduct in the $\{M_\alpha\}$ basis.

**Proposition 5.7.** For a composition $\alpha = (\alpha_1, \ldots, \alpha_\ell)$, one has

$$\Delta M_\alpha = \sum_{k=0}^\ell M_{(\alpha_1, \ldots, \alpha_k)} \otimes M_{(\alpha_{k+1}, \ldots, \alpha_\ell)} = \sum_{\beta, \gamma : \beta \cdot \gamma = \alpha} M_\beta \otimes M_\gamma$$

*Proof.* This comes from expressing a monomial in $\Delta M_\alpha = M_\alpha(x, y)$ uniquely in the form $x_1^{i_1} \cdots x_k^{i_k} \cdot y_{j_1}^{\alpha_{k+1}} \cdots y_{j_{\ell-k}}^{\alpha_\ell}$ for some $k \in \{0, 1, \ldots, n\}$ and some subscripts $i_1 < \cdots < i_k$ and $j_1 < \cdots < j_{\ell-k}$. \(\Box\)

**Proposition 5.8.** The quasisymmetric functions $QSym$ form a graded connected Hopf algebra of finite type, which is commutative, and contains the symmetric functions $\Lambda$ as a Hopf subalgebra.

*Proof.* To prove coassociativity of $\Delta$, we need to be slightly careful. It seems reasonable to argue by $(\Delta \otimes \text{id}) \circ \Delta f = f(x, y, z) = (\text{id} \otimes \Delta) \circ \Delta f$ as in the case of $\Lambda$, but this would now require further justification, as terms like $f(x, y)$ and $f(x, y, z)$ are no longer directly defined as evaluations of $f$ on some sequences (but rather are defined as images of $f$ under certain homomorphisms). However, it is very easy to see that $\Delta$ is coassociative by checking $(\Delta \otimes \text{id}) \circ \Delta = (\text{id} \otimes \Delta) \circ \Delta$ on the $\{M_\alpha\}$ basis: Proposition 5.7 yields

$$((\Delta \otimes \text{id}) \circ \Delta) M_\alpha = \sum_{k=0}^\ell \Delta(M_{(\alpha_1, \ldots, \alpha_k)}) \otimes M_{(\alpha_{k+1}, \ldots, \alpha_\ell)}$$

$$= \sum_{k=0}^\ell \left( \sum_{i=0}^k M_{(\alpha_1, \ldots, \alpha_i)} \otimes M_{(\alpha_{i+1}, \ldots, \alpha_k)} \right) \otimes M_{(\alpha_{k+1}, \ldots, \alpha_\ell)}$$

$$= \sum_{k=0}^\ell \sum_{i=0}^k M_{(\alpha_1, \ldots, \alpha_i)} \otimes M_{(\alpha_{i+1}, \ldots, \alpha_k)} \otimes M_{(\alpha_{k+1}, \ldots, \alpha_\ell)}$$

\(^{184}\)This is not completely obvious, but can be easily checked by verifying that $M_\alpha(x, y) = \sum_{\beta, \gamma : \beta \cdot \gamma = \alpha} M_\beta(x) \otimes M_\gamma(y)$ for every composition $\alpha$ (see the proof of Proposition 5.7 for why this holds).
and the same expression for \((\text{id} \otimes \Delta) \circ \Delta\) \(M_\alpha\).

The coproduct \(\Delta\) of \(\text{QSym}\) is an algebra morphism because it is defined as a composite of algebra morphisms in the bottom row of (5.3). To prove that the restriction of \(\Delta\) to the subring \(\Lambda\) of \(\text{QSym}\) is the comultiplication of \(\Lambda\), it thus is enough to check that it sends the elementary symmetric function \(e_n\) to \(\sum_{i=0}^n e_i \otimes e_{n-i}\) for every \(n \in \mathbb{N}\). This again follows from Proposition 5.7, since \(e_n = M_{(1,1,\ldots,1)}\) (with \(n\) times 1).

The counit is as usual for a graded connected coalgebra, and just as in the case of \(\Lambda\), sends a quasisymmetric function \(f(\mathbf{x})\) to its constant term \(f(0,0,\ldots)\). This is an evaluation, and hence an algebra map. Hence \(\text{QSym}\) forms a bialgebra, and as it is graded and connected, also a Hopf algebra by Proposition 1.36. It is clearly of finite type and contains \(\Lambda\) as a Hopf subalgebra. 

We will identify the antipode in \(\text{QSym}\) shortly, but we first deal with another slightly subtle issue. In addition to the counit evaluation \(\epsilon(f) = f(0,0,\ldots)\), starting in Section 7.1, we will want to specialize elements in \(\text{QSym}(\mathbf{x})\) by making other variable substitutions, in which all but a finite list of variables are set to zero. We justify this here.

**Proposition 5.9.** Fix a totally ordered set \(I\), a commutative \(k\)-algebra \(A\), a finite list of variables \(x_{i_1},\ldots,x_{i_m}\), say with \(i_1 < \ldots < i_m\) in \(I\), and an ordered list of elements \((a_1,\ldots,a_m) \in A^m\).

Then there is a well-defined evaluation homomorphism

\[
\text{QSym}(\{x_i\}_{i \in I}) \longrightarrow A
f \longmapsto [f]_{x_{i_1}=a_1,\ldots,x_{i_m}=a_m, x_j=0\text{ for }j \notin \{i_1,\ldots,i_m\}}.
\]

Furthermore, this homomorphism depends only upon the list \((a_1,\ldots,a_m)\), as it coincides with the following:

\[
\text{QSym}(\{x_i\}_{i \in I}) \equiv \text{QSym}(x_1,x_2,\ldots) \longrightarrow A
f(x_1,x_2,\ldots) \longmapsto f(a_1,\ldots,a_m,0,0\ldots).
\]

(This latter statement is stated for the case when \(I\) is infinite; otherwise, read “\(x_1,x_2,\ldots,x_{|I|}\)” for “\(x_1,x_2,\ldots\),” and interpret \((a_1,\ldots,a_m,0,0\ldots)\) as an \(|I|\)-tuple.)

**Proof.** One already can make sense of evaluating \(x_{i_1} = a_1,\ldots,x_{i_m} = a_m\) and \(x_j = 0\) for \(j \notin \{i_1,\ldots,i_m\}\) in the ambient ring \(R(\{x_i\}_{i \in I})\) containing \(\text{QSym}(\{x_i\}_{i \in I})\), since a power series \(f\) of bounded degree will have finitely many monomials that only involve the variables \(x_{i_1},\ldots,x_{i_m}\). The last assertion follows from quasisymmetry of \(f\), and is perhaps checked most easily when \(f = M_\alpha(\{x_i\}_{i \in I})\) for some \(\alpha\). \(\square\)

The antipode in \(\text{QSym}\) has a reasonably simple expression in the \(\{M_\alpha\}\) basis, but requiring a definition.

**Definition 5.10.** For \(\alpha,\beta\) in \(\text{Comp}_n\), say that \(\alpha\) refines \(\beta\) or \(\beta\) coarsens \(\alpha\) if, informally, one can obtain \(\beta\) from \(\alpha\) by combining some of its adjacent parts. Alternatively, one has a bijection \(\text{Comp}_n \rightarrow 2^{[n-1]}\) where \([n-1] := \{1,2,\ldots,n-1\}\) which sends \(\alpha = (\alpha_1,\ldots,\alpha_\ell)\) having length \(\ell(\alpha) = \ell\) to its subset of partial sums

\[
D(\alpha) := \{\alpha_1,\alpha_1 + \alpha_2,\ldots,\alpha_1 + \cdots + \alpha_{\ell-1}\},
\]

and this sends the refinement ordering to the inclusion ordering on the Boolean algebra \(2^{[n-1]}\) (to be more precise: a composition \(\alpha\) refines a composition \(\beta\) if and only if \(D(\alpha) \supset D(\beta)\)). There is also a bijection sending \(\alpha\) to its ribbon diagram: the skew diagram \(\lambda/\mu\) having rows of sizes \((\alpha_1,\ldots,\alpha_\ell)\) read from bottom to top with exactly one column of overlap between adjacent rows. These bijections and the refinement partial
Proof. Proposition 5.7, so Proposition 1.37 shows when \( \ell \) this involution switches the sign of \((-1)^{\ell(\alpha)}\) triples \((i,\beta,f)\) that we should be assigning types to, not terms.

For example, \( S(M_{(a,b,c)}) = - (M_{(c,b,a)} + M_{(b+c,a)} + M_{(c,a+b)} + M_{(a+b+c)}) \).

**Theorem 5.11.** For any composition \( \alpha \) in \( \text{Comp} \),

\[
S(M_\alpha) = (-1)^{\ell(\alpha)} \sum_{\gamma \in \text{Comp}: \gamma \text{ coarsens } \text{rev}(\alpha)} M_\gamma
\]

For example,

\[
S(M_{(a,b,c)}) = - (M_{(c,b,a)} + M_{(b+c,a)} + M_{(c,a+b)} + M_{(a+b+c)})
\]

**Proof.** We give Ehrenborg’s proof\(^{185}\) [47, Prop. 3.4] via induction on \( \ell = \ell(\alpha) \). One has easy base cases when \( \ell(\alpha) = 0 \), where \( S(M_\alpha) = S(1) = 1 = (-1)^0 M_{\text{rev}(\emptyset)} \), and when \( \ell(\alpha) = 1 \), where \( M_\alpha \) is primitive by Proposition 5.7, so Proposition 1.37 shows \( S(M_\alpha) = -M_\alpha = (-1)^1 M_{\text{rev}(\alpha)} \).

For the inductive step, apply the inductive definition of \( S \) from the proof of Proposition 1.36:

\[
S(M_{(a_1,\ldots,a_\ell)}) = - \sum_{i=0}^{\ell-1} S(M_{(a_1,\ldots,a_i)}) M_{(a_{i+1},\ldots,a_\ell)}
\]

\[
= \sum_{i=0}^{\ell-1} \sum_{\beta \text{ coarsening } (a_1,a_{i+1},\ldots,a_\ell)} (-1)^{i+1} M_\beta M_{(a_{i+1},\ldots,a_\ell)}
\]

The idea will be to cancel terms of opposite sign that appear in the expansions of the products \( M_\beta M_{(a_{i+1},\ldots,a_\ell)} \).

Note that each composition \( \beta \) appearing above has first part \( \beta_1 \) of the form \( a_1 + a_{i+1} + \cdots + a_k \) for some \( h \leq i \) (unless \( \beta = \emptyset \)), and hence each term \( M_\gamma \) in the expansion of the product \( M_\beta M_{(a_{i+1},\ldots,a_\ell)} \) has \( \gamma_1 \) (that is, the first entry of \( \gamma \)) a sum that can take one of these three forms:

- \( a_i + a_i + \cdots + a_h \),
- \( a_{i+1} + (a_i + a_{i+1} + \cdots + a_h) \),
- \( a_{i+1} \).

Say that the type of \( \gamma \) is \( i \) in the first case, and \( i + 1 \) in the second two cases\(^{186}\); in other words, the type is the largest subscript \( k \) on a part \( a_k \) which was combined in the sum \( \gamma_1 \). It is not hard to see that a given \( \gamma \) for which the type \( k \) is strictly smaller than \( \ell \) arises from exactly two pairs \((\beta,\gamma),(\beta',\gamma)\), having opposite signs \((-1)^h \) and \((-1)^{k+1} \) in the above sum\(^{187}\). For example, if \( \alpha = (a_1,\ldots,a_8) \), then the composition

\[^{185}\]A different proof was given by Malvenuto and Reutenauer [110, Cor. 2.3], and is sketched in Remark 5.36 below.

\[^{186}\]We imagine that we label the terms obtained by expanding \( M_\beta M_{(a_{i+1},\ldots,a_\ell)} \) by distinct labels, so that each term knows how exactly it was created (i.e., which \( i \), which \( \beta \) and which map \( f \) as in (5.2) gave rise to it). Strictly speaking, it is these triples \((i,\beta,f)\) that we should be assigning types to, not terms.

\[^{187}\]Strictly speaking, this means that we have an involution on the set of our \((i,\beta,f)\) triples having type smaller than \( \ell \), and this involution switches the sign of \((-1)^f M_{wt,f} \).
Thus one can cancel almost all the terms, excepting those with
\[ \gamma = (\alpha_6 + \alpha_5 + \alpha_4, \alpha_3, \alpha_7, \alpha_8 + \alpha_2 + \alpha_1) \] of type 6 can arise from either of
\[ \beta = (\alpha_6 + \alpha_5 + \alpha_4, \alpha_3, \alpha_2 + \alpha_1) \] with \( i = 6 \) and sign \(-1\)^7.
\[ \beta' = (\alpha_5 + \alpha_4, \alpha_2 + \alpha_1) \] with \( i = 5 \) and sign \(-1\)^6.

Similarly, \( \gamma = (\alpha_6, \alpha_5 + \alpha_4, \alpha_3, \alpha_7, \alpha_8 + \alpha_2 + \alpha_1) \) can arise from either of
\[ \beta = (\alpha_6, \alpha_5 + \alpha_4, \alpha_3, \alpha_2 + \alpha_1) \] with \( i = 6 \) and sign \(-1\)^7.
\[ \beta' = (\alpha_5 + \alpha_4, \alpha_2 + \alpha_1) \] with \( i = 5 \) and sign \(-1\)^6.

Thus one can cancel almost all the terms, excepting those with \( \gamma \) of type \( \ell \) among the terms \( M_\gamma \) in the expansion of the last \( (i = \ell - 1) \) summand \( M_\beta M_{(\alpha_1)} \). A bit of thought shows that these are the \( \gamma \) coarsening \( \rev(\alpha) \), and all have sign \(-1\)^\( \ell \). \( \square \)

5.2. The fundamental basis and \( P \)-partitions. There is a second important basis for \( \text{QSym} \) which arose originally in Stanley’s \( P \)-partition theory [151].

**Definition 5.12.** A labelled poset will here mean a partially ordered set \( P \) whose underlying set is some finite subset of the integers. A \( P \)-partition is a function \( P \xrightarrow{f} \{1, 2, \ldots \} \) with the property that
- \( i <_P j \) and \( i <_\mathbb{Z} j \) implies \( f(i) \leq f(j) \), and
- \( i <_P j \) and \( i >_\mathbb{Z} j \) implies \( f(i) < f(j) \).

Denote by \( \mathcal{A}(P) \) the set of all \( P \)-partitions \( f \), and let \( F_P(x) := \sum_{f \in \mathcal{A}(P)} x_f \) where \( x_f := \prod_{i \in P} x_{f(i)} \). This \( F_P(x) \) is an element of \( \mathbb{k}[[x]] := \mathbb{k}[x_1, x_2, \ldots] \).

**Example 5.13.** Depicted is a labelled poset \( P \), along with the relations among the four values \( f = (f(1), f(2), f(3), f(4)) \) that define its \( P \)-partitions \( f \):

\[
\begin{array}{cccc}
4 & 2 & \leq & f(1) \\
\downarrow & & \downarrow & \\
3 & 1 & \leq & f(2) \\
\downarrow & & \downarrow & \\
 & f(3) & < & f(4)
\end{array}
\]

**Remark 5.14.** Stanley’s treatment of \( P \)-partitions in [154, §3.15 and §7.19] uses a language different from ours. First, Stanley works not with labelled posets \( P \), but with pairs \((P, \omega)\) of a poset \( P \) and a bijective labelling \( \omega : P \to [n] \). Thus, the relation \( \lesssim \) is not given on \( P \) a priori, but has to be pulled back from \([n]\) using \( \omega \) (and it depends on \( \omega \), whence Stanley speaks of “\((P, \omega)\)-partitions”). Furthermore, Stanley uses the notations \( F_P \) and \( F_{P, \omega} \) for something different from what we denote by \( F_P \), whereas what we call \( F_P \) is dubbed \( K_{P, \omega} \) in [154, §7.19].

The so-called fundamental quasisymmetric functions are an important special case of the \( F_P(x) \). We shall first define them directly and then see how they are obtained as \( P \)-partition enumerators \( F_P(x) \) for some special labelled posets \( P \).

**Definition 5.15.** Let \( n \in \mathbb{N} \) and \( \alpha \in \text{Comp}_n \). We define the fundamental quasisymmetric function \( L_\alpha = L_\alpha(x) \in \text{QSym} \) by
\[
L_\alpha := \sum_{\beta \in \text{Comp}_n: \beta \text{ refines } \alpha} M_\beta.
\]

**Example 5.16.** The extreme cases for \( \alpha \) in \( \text{Comp}_n \) give quasisymmetric functions \( L_\alpha \) which are symmetric:
\[
L_{(1^n)} = M_{(1^n)} = e_n, \\
L_{(n)} = \sum_{\alpha \in \text{Comp}_n} M_\alpha = h_n
\]

---

\[ ^{188} \text{See [60] for a history of } P \text{-partitions; our notations, however, strongly differ from those in [60].} \]
One sees that the \( \{ L_\alpha \}_{\alpha \in \text{Comp}} \) are a \( k \)-basis for \( \text{QSym} \), as inclusion-exclusion applied to (5.4) gives

\[
M_\alpha = \sum_{\beta \in \text{Comp}_n: \beta \text{ refines } \alpha} (-1)^{\ell(\beta) - \ell(\alpha)} L_\beta.
\]

**Proposition 5.17.** Let \( n \in \mathbb{N} \). Let \( \alpha \) be a composition of \( n \). Let \( I \) be an infinite totally ordered set. Then,

\[
L_\alpha \left( \{ x_i \}_{i \in I} \right) = \sum_{\substack{i_1 \leq i_2 \leq \ldots \leq i_n \text{ in } I; \ \ i_j < i_{j+1} \text{ if } j \in D(\alpha)}} x_{i_1} x_{i_2} \cdots x_{i_n},
\]

where \( L_\alpha \left( \{ x_i \}_{i \in I} \right) \) is defined as the image of \( L_\alpha \) under the isomorphism \( \text{QSym} \to \text{QSym} \left( \{ x_i \}_{i \in I} \right) \) obtained in Definition 5.5. In particular, for the standard (totally ordered) variable set \( x = (x_1 < x_2 < \ldots) \), we obtain (5.5)

\[
L_\alpha = L_\alpha (x) = \sum_{\substack{(1 \leq) \ i_1 \leq i_2 \leq \ldots \leq i_n; \ \ i_j < i_{j+1} \text{ if } j \in D(\alpha)}} x_{i_1} x_{i_2} \cdots x_{i_n}.
\]

**Proof.** Every composition \( \beta = (\beta_1, \ldots, \beta_t) \) of \( n \) satisfies

\[
M_\beta \left( \{ x_i \}_{i \in I} \right) = \sum_{\substack{k_1 < \cdots < k_t \in I}} x^{\beta_1}_{k_1} \cdots x^{\beta_t}_{k_t} = \sum_{\substack{i_1 \leq i_2 \leq \ldots \leq i_n \text{ in } I; \ \ i_j < i_{j+1} \text{ if and only if } j \in D(\beta)}} x_{i_1} x_{i_2} \cdots x_{i_n}.
\]

Applying the ring homomorphism \( \text{QSym} \to \text{QSym} \left( \{ x_i \}_{i \in I} \right) \) to (5.4), we obtain

\[
L_\alpha \left( \{ x_i \}_{i \in I} \right) = \sum_{\beta \in \text{Comp}_n: \beta \text{ refines } \alpha} \sum_{\beta \in \text{Comp}_n: D(\alpha) \subset D(\beta)} \sum_{i_1 \leq \ldots \leq i_n \text{ in } I; \ \ i_j < i_{j+1} \text{ if and only if } j \in D(\beta)} x_{i_1} x_{i_2} \cdots x_{i_n}.
\]

\[
= \sum_{Z \subseteq [n-1]: \ D(\alpha) \subset Z} \sum_{i_1 \leq \ldots \leq i_n \text{ in } I; \ \ i_j < i_{j+1} \text{ if and only if } j \in Z} x_{i_1} x_{i_2} \cdots x_{i_n}.
\]

\[\square\]

**Proposition 5.18.** When the labelled poset \( P \) is a total or linear order \( w = (w_1 < \ldots < w_n) \), the generating function \( F_w(x) \) depends only upon the descent set

\[
\text{Des}(w) := \{ i : w_i \geq w_{i+1} \} \subset \{ 1, 2, \ldots, n-1 \}
\]

and its associated composition \( \alpha \in \text{Comp}_n \) having partial sums \( D(\alpha) = \text{Des}(w) \): one has that \( F_w(x) \) equals the fundamental quasi-symmetric function \( L_\alpha \).

E.g., total order \( w = 31542 \) has \( \text{Des}(w) = \{2, 4\} \) and composition \( \alpha = (2, 2, 1) \), so

\[
F_{35142}(x) = \sum_{f(3) \leq f(5) < f(1) \leq f(4) < f(2)} x_{f(3)} x_{f(5)} x_{f(1)} x_{f(4)} x_{f(2)}
\]

\[
= \sum_{i_3 \leq i_5 < i_4 \leq i_1} x_{i_1} x_{i_2} x_{i_3} x_{i_4} x_{i_5}
\]

\[
= L_{(2,2,1)} = M_{(2,2,1)} + M_{(1,1,2,1)} + M_{(1,1,1,2,1)} + M_{(1,1,1,1,1,1)}.
\]

**Proof of Proposition 5.18.** Write \( F_w(x) \) as a sum of monomials \( x_{f(w_1)} \cdots x_{f(w_n)} \) over all \( P \)-partitions \( f \). These \( P \)-partitions are exactly the sequences \( f(w_1) \leq \cdots \leq f(w_n) \) having strict inequalities \( f(w_i) < f(w_{i+1}) \) whenever \( i \) is in \( \text{Des}(w) \). Hence, comparison with (5.5) reveals that this sum equals \( L_\alpha \). \[\square\]

The next proposition ([154, Cor. 7.19.5], [104, Cor. 3.3.24]) is an algebraic shadow of Stanley’s main lemma [154, Thm. 7.19.4] in \( P \)-partition theory. It expands any \( F_P(x) \) in the \( \{ L_\alpha \} \) basis, as a sum over the set \( \mathcal{L}(P) \) of all linear extensions \( w \) of \( P \) \footnote{Let us explain what we mean by linear extensions and how we represent them.}. E.g., the poset \( P \) from Example 5.13 has \( \mathcal{L}(P) = \{3124, 3142, 3412\} \).
Theorem 5.19. For any labelled poset \( P \),
\[
F_P(x) = \sum_{w \in \mathcal{L}(P)} F_w(x).
\]

Proof. We give Gessel’s proof [59, Thm. 1], via induction on the number of pairs \( i, j \) which are incomparable in \( P \). When this quantity is 0, then \( P \) is itself a linear order \( w \), so that \( \mathcal{L}(P) = \{ w \} \) and there is nothing to prove.

In the inductive step, let \( i, j \) be incomparable elements. Consider the two posets \( P_{i<j} \) and \( P_{j<i} \) which are obtained from \( P \) by adding in an order relation between \( i \) and \( j \), and then taking the transitive closure; it is not hard to see that these transitive closures cannot contain a cycle, so that these really do define two posets. The result then follows by induction applied to \( P_{i<j}, P_{j<i} \), once one notices that \( \mathcal{L}(P) = \mathcal{L}(P_{i<j}) \cup \mathcal{L}(P_{j<i}) \) since every linear extension \( w \) of \( P \) either has \( i \) before \( j \) or vice-versa, and \( \mathcal{A}(P) = \mathcal{A}(P_{i<j}) \cup \mathcal{A}(P_{j<i}) \) since, assuming that \( i \leq_P j \) without loss of generality, every \( f \) in \( \mathcal{A}(P) \) either satisfies \( f(i) \leq f(j) \) or \( f(i) \geq f(j) \).

Example 5.20. To illustrate the induction in the above proof, consider the poset \( P \) from Example 5.13, having \( \mathcal{L}(P) = \{3124, 3142, 3412\} \). Then choosing as incomparable pair \((i, j) = (1, 4)\), one has

\[
\begin{align*}
P_{i<j} &= 4 \quad \downarrow \quad 2 \quad \downarrow \quad f(4) \quad \downarrow \quad \leq \quad f(2) \\
1 &\quad \downarrow \quad \leq \quad f(1) \\
3 &\quad \downarrow \quad < \quad f(3) \\
\mathcal{L}(P_{i<j}) &= \{3124, 3142\}
\end{align*}
\]

\[
\begin{align*}
P_{j<i} &= 2 \quad \downarrow \quad f(2) \\
1 &\quad \downarrow \quad \leq \quad f(1) \\
3 &\quad \downarrow \quad < \quad \downarrow \quad \leq \quad f(4) \\
\mathcal{L}(P_{j<i}) &= \{3412\}
\end{align*}
\]


[Hint: For every \( f : P \to \{1, 2, 3, \ldots\} \), we can define a binary relation \( \prec_f \) on the set \( P \) by letting \( i \prec_f j \) hold if and only if
\[
(f(i) < f(j)) \text{ or } (f(i) = f(j) \text{ and } i <_P j).
\]
Show that this binary relation \( \prec_f \) is (the smaller relation of) a total order. When \( f \) is a \( P \)-partition, then endowing the set \( P \) with this total order yields a linear extension of \( P \). Use this to show that the set \( \mathcal{A}(P) \) is the union of its pairwise disjoint subsets \( \mathcal{A}(w) \) with \( w \in \mathcal{L}(P) \).

We next wish to describe the structure maps for the Hopf algebra \( \text{QSym} \) in the basis \( \{ L_\alpha \} \) of fundamental quasisymmetric functions. For this purpose, two more definitions are useful.

If \( P \) is a finite poset, then a linear extension of \( P \) denotes a total order \( w \) on the set \( P \) having the property that \( i <_P j \) (for two elements \( i \) and \( j \) of \( P \)) implies \( i <_w j \). (In other words, it is a linear order on the ground set \( P \) which extends \( P \) as a poset; therefore the name.) We identify such a total order \( w \) with the list \( (p_1, p_2, \ldots, p_n) \) containing all elements of \( P \) in \( w \)-increasing order (that is, \( p_1 <_w p_2 <_w \cdots <_w p_n \)).

(Stanley, in [154, §3.5], defines linear extensions in a slightly different way: For him, a linear extension of a finite poset \( P \) is an order-preserving bijection from \( P \) to the subposet \( \{1, 2, \ldots, |P|\} \) of \( \mathbb{Z} \). But this is equivalent to our definition, since a bijection like this can be used to transport the order relation of \( \{1, 2, \ldots, |P|\} \) back to \( P \), thus resulting in a total order on \( P \) which is a linear extension of \( P \) in our sense.)
Definition 5.22. Given two nonempty compositions $\alpha = (\alpha_1, \ldots, \alpha_\ell)$, $\beta = (\beta_1, \ldots, \beta_m)$, their near-concatenation is

$$\alpha \circ \beta := (\alpha_1, \ldots, \alpha_{\ell-1}, \alpha_\ell + \beta_1, \beta_2, \ldots, \beta_m)$$

For example, the figure below depicts for $\alpha = (1, 3, 3)$ (black squares) and $\beta = (4, 2)$ (white squares) the concatenation and near-concatenation as ribbons:

Lastly, given $\alpha$ in $\text{Comp}_n$, let $\omega(\alpha)$ be the unique composition in $\text{Comp}_n$ whose partial sums $D(\omega(\alpha))$ form the complementary set within $[n-1]$ to the partial sums $D(\text{rev}(\alpha))$; alternatively, one can check this means that the ribbon for $\omega(\alpha)$ is obtained from that of $\alpha$ by conjugation or transposing, that is, if $\alpha = \lambda/\mu$ then $\omega(\alpha) = \lambda'/\mu'$. E.g. if $\alpha = (4, 2, 2)$ so that $n = 8$, then $\text{rev}(\alpha) = (2, 2, 4)$ has $D(\text{rev}(\alpha)) = \{2, 4\} \subset [7]$, complementary to the set $\{1, 3, 5, 6, 7\}$ which are the partial sums for $\omega(\alpha) = (1, 2, 2, 1, 1, 1)$, and the ribbon diagrams of $\alpha, \omega(\alpha)$ are

Proposition 5.23. The structure maps for the Hopf algebra $\text{QSym}$ in the basis $\{L_\alpha\}$ of fundamental quasisymmetric functions are as follows:

$$\Delta L_\alpha = \sum_{(\beta, \gamma): \beta \cdot \gamma = \alpha} L_\beta \otimes L_\gamma$$  \hfill (5.7)

$$L_\alpha L_\beta = \sum_{w \in w_\alpha \sqcup w_\beta} L_{\gamma(w)}$$  \hfill (5.8)

$$S(L_\alpha) = (-1)^{|\alpha|} L_{\omega(\alpha)}.$$  \hfill (5.9)

Here we are making use of the following notations in (5.8) (recall also Definition 1.68):

- A labelled linear order will mean a labelled poset $P$ whose order $<_P$ is a total order. We will identify any labelled linear order $P$ with the word (over the alphabet $\mathbb{Z}$) obtained by writing down the elements of $P$ in increasing order (with respect to the total order $<_P$). This way, every word (over the alphabet $\mathbb{Z}$) which has no two equal letters becomes identified with a labelled linear order.
- $w_\alpha$ is any labelled linear order with underlying set $\{1, 2, \ldots, |\alpha|\}$ such that $\text{Des}(w_\alpha) = D(\alpha)$.
- $w_\beta$ is any labelled linear order with underlying set $\{|\alpha| + 1, |\alpha| + 2, \ldots, |\alpha| + |\beta|\}$ such that $\text{Des}(w_\beta) = D(\beta)$.
- $\gamma(w)$ is the unique composition of $|\alpha| + |\beta|$ with $D(\gamma(w)) = \text{Des}(w)$.

At first glance the formula (5.7) for $\Delta L_\alpha$ might seem more complicated than the formula of Proposition 5.7 for $\Delta M_\alpha$. However, it is equally simple when viewed in terms of ribbon diagrams: it cuts the ribbon diagram $\alpha$ into two smaller ribbons $\beta$ and $\gamma$, in all $|\alpha| + 1$ possible ways, via horizontal cuts ($\beta \cdot \gamma = \alpha$) or vertical cuts ($\beta \circ \gamma = \alpha$). For example,
Example 5.24. To multiply $L_{(1,1)} L_{(2)}$, one could pick $w_\alpha = 21$ and $w_\beta = 34$, and then

$$L_{(1,1)} L_{(2)} = \sum_{w \in L_{(1,1)}} L_{(2)} w = \gamma(w) = \gamma_{(2314)} + \gamma_{(3214)} + \gamma_{(2341)} + \gamma_{(3241)} + \gamma_{(3421)}.$$

Proof of Proposition 5.23. To prove formula (5.7) for $\alpha$ in $\text{Comp}_n$, note that

$$\Delta L_\alpha = L_\alpha(x, y) = \sum_{k=0}^{n} \sum_{1 \leq i_1 \leq \ldots \leq i_k \leq n: \ i_r < i_{r+1} \ for \ r \in \{D(\alpha) \setminus \{k\}\} \ } x_{i_1} \cdots x_{i_k} \cdot y_{i_k+1} \cdots y_n$$

by Proposition 5.17. One then realizes that the inner sums corresponding to values of $k$ that lie (resp. do not lie) in $D(\alpha) \cup \{0, n\}$ correspond to the terms $L_\beta(x) L_\gamma(y)$ for pairs $(\beta, \gamma)$ in which $\beta \cdot \gamma = \alpha$ (resp. $\beta \odot \gamma = \alpha$).

For formula (5.8), let $P$ be the labelled poset which is the disjoint union of linear orders $w_\alpha, w_\beta$. Then

$$L_\alpha L_\beta = F_{w_\alpha}(x) F_{w_\beta}(x) = F_P(x) = \sum_{w \in L \in \{P\}} \sum_{w \in w_\alpha \cup w_\beta} L_\gamma(w)$$

where the first equality used Proposition 5.18, the second equality comes from the definition of a $P$-partition, the third equality from Theorem 5.19, and the fourth from the equality $L(P) = w_\alpha \cup w_\beta$.

To prove formula (5.9), compute using Theorem 5.11 that

$$S(L_\alpha) = \sum_{\beta \text{ refining } \alpha} S(M_\beta) = \sum_{(\beta, \gamma): \ \beta \text{ refines } \alpha, \ \gamma \text{ coarsens } \text{rev}(\beta)} (-1)^{\ell(\beta)} M_\gamma = \sum_{\gamma} M_\gamma \sum_{\beta} (-1)^{\ell(\beta)}$$

in which the last inner sum is over $\beta$ for which

$$D(\beta) \supset D(\alpha) \cup D(\text{rev}(\gamma)).$$

The alternating signs make such inner sums vanish unless they have only the single term where $D(\beta) = [n-1]$ (that is, $\beta = (1^n)$). This happens exactly when $D(\text{rev}(\gamma)) \cup D(\alpha) = [n-1]$ or equivalently, when $D(\text{rev}(\gamma))$ contains the complement of $D(\alpha)$, that is, when $D(\gamma)$ contains the complement of $D(\text{rev}(\alpha))$, that is, when $\gamma$ refines $\omega(\alpha)$. Thus

$$S(L_\alpha) = \sum_{\gamma \in \text{Comp}_n} M_\gamma \cdot (-1)^n = (-1)^{|\alpha|} L_{\omega(\alpha)}.$$

The antipode formula (5.9) for $L_\alpha$ leads to a general interpretation for the antipode of QSym acting on $P$-partition enumerators $F_P(x)$.

Definition 5.25. Given a labelled poset $P$ on $\{1, 2, \ldots, n\}$, let the opposite or dual labelled poset $P^{opp}$ have $i < P^{opp} j$ if and only if $j < P i$.

For example,

$$P = \begin{array}{ccc}
4 & 2 & \text{opp} \\
3 & 1 &
\end{array} \quad P^{opp} = \begin{array}{ccc}
4 & 3 & \text{opp} \\
2 & 1 &
\end{array}$$

The following observation is straightforward.

Proposition 5.26. When $P$ is a linear order corresponding to some permutation $w = (w_1, \ldots, w_n)$ in $S_n$, then $w^{opp} = w\omega_0$ where $\omega_0 \in S_n$ is the permutation that swaps $i \leftrightarrow n+1-i$ (this is the so-called longest permutation, thus named due to it having the highest “Coxeter length” among all permutations in $S_n$). Furthermore, in this situation one has $F_w(x) = L_\alpha$, that is, $\text{Des}(w) = D(\alpha)$ if and only if $\text{Des}(w^{opp}) = D(\omega(\alpha))$, that is $F_w^{opp}(x) = L_{\omega(\alpha)}$. Thus,

$$S(F_w(x)) = (-1)^n F_w^{opp}(x).$$
For example, given the compositions considered earlier

\[ \alpha = (4, 2, 2) = \begin{array}{ccc}
\square & \square \\
\square & \square
\end{array} \quad \omega(\alpha) = (1, 2, 2, 1, 1) = \begin{array}{ccc}
\square & \square & \square \\
\square & \square & \square
\end{array} \]

if one picks \( w = 1235 \cdot 47 \cdot 68 \) (with descent positions marked by dots) having \( \text{Des}(w) = \{4, 6\} = D(\alpha) \), then \( w_{\text{opp}} = w w_0 = 8 \cdot 67 \cdot 45 \cdot 3 \cdot 2 \cdot 1 \) has \( \text{Des}(w_{\text{opp}}) = \{1, 3, 5, 6, 7\} = D(\omega(\alpha)) \).

**Corollary 5.27.** For any labelled poset \( P \) on \( \{1, 2, \ldots, n\} \), one has

\[ S(F_P(x)) = (-1)^n F_{P_{\text{opp}}}(x). \]

**Proof.** Since \( S \) is linear, one can apply Theorem 5.19 and Proposition 5.26

\[ S(F_P(x)) = \sum_{w \in \mathcal{L}(P)} S(F_w(x)) = \sum_{w \in \mathcal{L}(P)} (-1)^n F_{w_{\text{opp}}}(x) = (-1)^n F_{P_{\text{opp}}}(x). \]

as \( \mathcal{L}(P_{\text{opp}}) = \{w_{\text{opp}} : w \in \mathcal{L}(P)\} \).

**Remark 5.28.** Malvenuto and Reutenauer, in [111, Theorem 3.1], prove an even more general antipode formula, which encompasses our Corollary 5.27, Proposition 5.26, Theorem 5.11 and (5.9).

We remark on a special case of Corollary 5.27 to which we alluded earlier, related to skew Schur functions.

**Corollary 5.29.** In \( \Lambda \), the action of \( \omega \) and the antipode \( S \) on skew Schur functions \( s_{\lambda/\mu} \) are as follows:

\[ \omega(s_{\lambda/\mu}) = s_{\lambda'/\mu'} \]

\[ S(s_{\lambda/\mu}) = (-1)^{|\lambda/\mu|} s_{\lambda'/\mu'}. \]

**Proof.** Given a skew shape \( \lambda/\mu \), one can always create a labelled poset \( P \) which is its skew Ferrers poset, together with one of many column-strict labellings, in such a way that \( F_P(x) = s_{\lambda/\mu}(x) \). An example is shown here for \( \lambda/\mu = (4, 4, 2)/(1, 1, 0) \):

\[ \lambda/\mu = \begin{array}{ccc}
\square & \square & \square \\
\square & \square & \square \\
\square & \square & \square \\
\square & \square & \square
\end{array} \quad P = \begin{array}{ccc}
5 & 4 & 2 \\
8 & 3 & 1 \\
7 & 6 & \\
\end{array} \quad w = \begin{array}{ccc}
f(5) & f(4) & f(2) \\
f(8) & f(7) & f(3) & f(1) \\
\end{array} \quad w_{\text{opp}} = \begin{array}{ccc}
f(6) & f(3) & f(1) \\
f(7) & f(4) & f(2) \\
\end{array} \quad F_{P_{\text{opp}}}(x) = \sum_T x^{\text{cont}(T)} \]

as \( T \) ranges over all reverse semistandard tableaux or column-strict plane partitions of \( \lambda'/\mu' \):

\[ \lambda'/\mu' = \begin{array}{ccc}
\square & \square & \square \\
\square & \square & \square \\
\square & \square & \square \\
\square & \square & \square
\end{array} \quad P_{\text{opp}} = \begin{array}{ccc}
6 & 3 & 1 \\
8 & 4 & 2 \\
5 & 7 & \\
\end{array} \quad w = \begin{array}{ccc}
f(6) & f(3) & f(1) \\
f(7) & f(4) & f(2) \\
\end{array} \quad w_{\text{opp}} = \begin{array}{ccc}
f(5) & f(3) & f(1) \\
f(7) & f(4) & f(2) \\
\end{array} \quad F_{P_{\text{opp}}}(x) = \sum_T x^{\text{cont}(T)} \]

as \( T \) ranges over all reverse semistandard tableaux or column-strict plane partitions of \( \lambda'/\mu' \).
But this means that $F_{P_{opp}}(x) = s_{\lambda'/\mu'}(x)$, since the fact that skew Schur functions lie in $\Lambda$ implies that they can be defined either as generating functions for column-strict tableaux or reverse semistandard tableaux; see Remark 2.12 above, or [154, Prop. 7.10.4].

Thus we have

$$F_P(x) = s_{\lambda/\mu}(x)$$

$$F_{P_{opp}}(x) = s_{\lambda'/\mu'}(x).$$

Proposition 1.46(c) tell us that the antipode for QSym must specialize to the antipode for $\Lambda$ (see also Remark 5.43 below), so (5.11) is a special case of Corollary 5.27. Then (5.10) follows from the relation (2.20) that $S(f) = (-1)^n\omega(f)$ for $f$ in $\Lambda_n$.

\begin{remark}

Before leaving $P$-partitions temporarily, we mention two open questions about them.

The first is a conjecture of Stanley from his thesis [151]. As mentioned in the proof of Corollary 5.29, each skew Schur function $s_{\lambda'/\mu}$ is a special instance of $P$-partition enumerator $F_P(x)$.

Conjecture 5.31. A labelled poset $P$ has $F_P(x)$ symmetric, and not just quasisymmetric, if and only if $P$ is a column-strict labelling of some skew Ferrers poset $\lambda'/\mu$.

A somewhat weaker result in this direction was proven by Malvenuto in her thesis [109, Thm. 6.4], showing that if a labelled poset $P$ has the stronger property that its set of linear extensions $\mathcal{L}(P)$ is a union of plactic or Knuth equivalence classes, then $P$ must be a column-strict labelling of a skew Ferrers poset.

The next question is due to P. McNamara, and is suggested by the obvious factorizations of $P$-partition enumerators $F_{P_{1\cup P_2}}(x) = F_{P_1}(x)F_{P_2}(x)$.

Question 5.32. If $k$ is a field, does a connected labelled poset $P$ always have $F_P(x)$ irreducible within the ring $QSym$?

The phrasing of this question requires further comment. It is assumed here that $x = (x_1, x_2, \ldots)$ is infinite; for example when $P$ is a 2-element chain labelled “against the grain” (i.e., the bigger element of the chain has the smaller label), then $F_P(x) = e_2(x)$ is irreducible, but its specialization to two variables $x = (x_1, x_2)$ is $e_2(x_1, x_2) = x_1x_2$, which is reducible. If one wishes to work in finitely many variables $x = (x_1, \ldots, x_m)$ one can perhaps assume that $m$ is at least $|P| + 1$.

When working in $QSym = QSym(x)$ in infinitely many variables, it is perhaps not so clear where factorizations occur. For example, if $f$ lies in $QSym$ and factors $f = g \cdot h$ with $g, h$ in $R(x)$, does this imply that $g, h$ also lie in $QSym$? The answer is “Yes” (for $k = \mathbb{Z}$), but this is not obvious, and was proven by P. Pylyavskyy in [130, Chap. 11].

One also might wonder whether $QSym_{\mathbb{Z}}$ is a unique factorization domain, but this follows from the result of M. Hazewinkel ([64] and [68, Thm. 6.7.5], and Theorem 6.74 further below) who proved a conjecture of Ditters that $QSym_{\mathbb{Z}}$ is a polynomial algebra; earlier Malvenuto and Reutenauer [110, Cor. 2.2] had shown that $QSym_{\mathbb{Q}}$ is a polynomial algebra. In fact, one can find polynomial generators $\{P_\alpha\}$ for $QSym_{\mathbb{Q}}$ as a subset of the dual basis to the $Q$-basis $\{\xi_\alpha\}$ for $NSym_{\mathbb{Q}}$ which comes from taking products $\xi_\alpha := \xi_{\alpha_1} \cdots \xi_{\alpha_k}$ of the elements $\{\xi_\alpha\}$ defined in Remark 5.36 below. Specifically, one takes those $P_\alpha$ for which the composition $\alpha$ is a Lyndon composition; see the first proof of Proposition 6.75 for a mild variation on this construction.

Hazewinkel’s proof [68, Thm. 6.7.5] of the polynomiality of $QSym_{\mathbb{Z}}$ also shows that $QSym$ is a polynomial ring over $\Lambda$ (see Corollary 6.119); in particular, this yields that $QSym$ is a free $\Lambda$-module.\footnote{The latter statement has an analogue in finitely many indeterminates, proven by Lauve and Mason in [91, Corollary 13]: The quasisymmetric functions $QSym(\{x_i\}_{i \in I})$ are free as a $\Lambda(\{x_i\}_{i \in I})$-module for any totally ordered set $I$, infinite or not. In the case of finite $I$, this cannot be derived by Hazewinkel’s arguments, as the ring $QSym(\{x_i\}_{i \in I})$ is not in general a polynomial ring (e.g., when $k = \mathbb{Q}$ and $I = \{1, 2\}$, this ring is not even a UFD, as witnessed by $(x_1^2x_2) \cdot (x_1x_2^2) = (x_1x_2)^3$).}

An affirmative answer to Question 5.32 is known at least in the special case where $P$ is a connected column-strict labelling of a skew Ferrers diagram, that is, when $F_P(x) = s_{\lambda'/\mu}(x)$ for some connected skew diagram $\lambda'/\mu$; see [12].

5.3. The Hopf algebra $NSym$ dual to $QSym$. We introduce here the (graded) dual Hopf algebra to $QSym$. This is well-defined, as $QSym$ is connected graded of finite type.
**Definition 5.33.** Let \( \text{NSym} := \text{QSym}^o \), with dual pairing \( \text{NSym} \otimes \text{QSym} \to k \). Let \( \{H_\alpha\} \) be the \( k \)-basis of \( \text{NSym} \) dual to the \( k \)-basis \( \{M_\alpha\} \) of \( \text{QSym} \), so that
\[
(H_\alpha, M_\beta) = \delta_{\alpha, \beta}.
\]
When the base ring \( k \) is not clear from the context, we write \( \text{NSym}_k \) in lieu of \( \text{NSym} \).

The Hopf algebra \( \text{NSym} \) is known as the Hopf algebra of noncommutative symmetric functions. Its study goes back to \([57]\).

**Theorem 5.34.** Letting \( H_n := H_n \) for \( n = 0, 1, 2, \ldots \), with \( H_0 = 1 \), one has that
(5.12) \( \text{NSym} \cong k\langle H_1, H_2, \ldots \rangle \),
the free associative (but not commutative) algebra on generators \( \{H_1, H_2, \ldots \} \) with coproduct determined by
(5.13) \( \Delta H_n = \sum_{i+j=n} H_i \otimes H_j. \)

**Proof.** Since Proposition 5.7 asserts that \( \Delta M_\alpha = \sum_{(\beta, \gamma): \beta, \gamma = \alpha} M_\beta \otimes M_\gamma \), and since \( \{H_\alpha\} \) are dual to \( \{M_\alpha\} \), one concludes that for any compositions \( \beta, \gamma \), one has
\[ H_\beta H_\gamma = H_{\beta \cdot \gamma} \]
Iterating this gives
(5.14) \( H_\alpha = H_{(\alpha_1, \ldots, \alpha_\ell)} = H_{\alpha_1} \cdots H_{\alpha_\ell}. \)
Since the \( H_\alpha \) are a \( k \)-basis for \( \text{NSym} \), this shows \( \text{NSym} \cong k\langle H_1, H_2, \ldots \rangle \).

Note that \( H_n = H_n \) is dual to \( M_n \), so to understand \( \Delta H_n \), one should understand how \( M_n \) can appear as a term in the product \( M_\alpha M_\beta \). By (5.1) this occurs only if \( \alpha = (i) \), \( \beta = (j) \) where \( i + j = n \), where
\[ M_{(i)} M_{(j)} = M_{(i+j)} + M_{(i,j)} + M_{(j,i)} \]
(\text{where the } M_{(i,j)} \text{ and } M_{(j,i)} \text{ addends have to be disregarded if one of } i \text{ and } j \text{ is 0}). By duality, this implies the formula (5.13). \( \square \)

**Corollary 5.35.** The algebra homomorphism defined by
\[
\begin{align*}
\text{NSym} & \xrightarrow{\pi} \Lambda \\
H_n & \mapsto h_n
\end{align*}
\]
is a Hopf algebra surjection, and adjoint to the inclusion \( \Lambda \hookrightarrow \text{QSym} \) (with respect to the dual pairing \( \text{NSym} \otimes \text{QSym} \to k \)).

**Proof.** As an algebra map \( \pi \) may be identified with the surjection \( T(V) \to \text{Sym}(V) \) from the tensor algebra on a graded free \( k \)-module \( V \) with basis \( \{H_1, H_2, \ldots \} \) to the symmetric algebra on \( V \), since
\[
\text{NSym} \cong k\langle H_1, H_2, \ldots \rangle
\]
\[
\Lambda \cong k[h_1, h_2, \ldots]
\]
As (5.13) and Proposition 2.24(iii) assert that
\[
\begin{align*}
\Delta H_n & = \sum_{i+j=n} H_i \otimes H_j \\
\Delta h_n & = \sum_{i+j=n} h_i \otimes h_j
\end{align*}
\]
this map \( \pi \) is also a bialgebra morphism, and hence a Hopf morphism by Proposition 1.46(c).
To check \( \pi \) is adjoint to \( i \), let \( \lambda(\alpha) \) denote the partition which is the weakly decreasing rearrangement of the composition \( \alpha \), and note that the bases \( \{H_\alpha\} \) of NSym and \( \{m_\lambda\} \) of \( \Lambda \) satisfy
\[
(\pi(H_\alpha), m_\lambda) = (h_{\lambda(\alpha)}, m_\lambda) = \begin{cases} 
1 & \text{if } \lambda(\alpha) = \lambda \\
0 & \text{otherwise}
\end{cases} = \begin{pmatrix} \sum_{\lambda, \beta: \lambda(\beta) = \lambda} M_\beta \\ H_\alpha \end{pmatrix} = (H_\alpha, i(m_\lambda)).
\]

\[\square\]

**Remark 5.36.** For those who prefer generating functions to sign-reversing involutions, we sketch here Malvenuto and Reutenauer’s elegant proof [110, Cor. 2.3] of the antipode formula (Theorem 5.11). One needs to know that when \( \mathbb{Q} \) is a subring of \( \mathbb{k} \), and \( A \) is a \( \mathbb{k} \)-algebra (possibly noncommutative), in the ring of power series \( A[[t]] \) where \( t \) commutes with all of \( A \), one has familiar facts, such as
\[
a(t) = \log b(t) \quad \text{if and only if} \quad b(t) = \exp a(t)
\]
and whenever \( a(t), b(t) \) commute in \( A[[t]] \), one has
\[
\begin{align*}
\exp (a(t) + b(t)) &= \exp a(t) \exp b(t) \\
\log (a(t)b(t)) &= \log a(t) + \log b(t)
\end{align*}
\]
Start by assuming WLOG that \( \mathbb{k} = \mathbb{Z} \) (as \( \text{NSym}_\mathbb{k} = \text{NSym}_2 \otimes_2 \mathbb{k} \) in the general case). Now, define \( \text{NSym}_\mathbb{Q} = \text{NSym} \otimes_\mathbb{Q} \mathbb{Q} \) the elements \( \{\xi_1, \xi_2, \ldots\} \) via generating functions in \( \text{NSym}_\mathbb{Q}[[t]] \):
\[
\tilde{H}(t) := \sum_{n \geq 0} H_n t^n,
\]
\[
\xi(t) := \sum_{n \geq 1} \xi_n t^n = \log \tilde{H}(t)
\]
One first checks that this makes each \( \xi_n \) primitive, via a computation in the ring \( (\text{NSym}_\mathbb{Q} \otimes \text{NSym}_\mathbb{Q})[[t]] \) (into which we “embed” the ring \( (\text{NSym}_\mathbb{Q}[[t]]) \otimes_{\mathbb{Q}[[t]]} (\text{NSym}_\mathbb{Q}[[t]]) \) via the canonical ring homomorphism from the latter into the former \footnote{This ring homomorphism might fail to be injective, whence the “embed” stands in quotation marks. This does not need to worry us, since we will not draw any conclusions in \( (\text{NSym}_\mathbb{Q}[[t]]) \otimes_{\mathbb{Q}[[t]]} (\text{NSym}_\mathbb{Q}[[t]]) \) from our computation.}
\[
\Delta \xi(t) = \Delta \left( \log \sum_{n \geq 0} H_n t^n \right) = \log \sum_{n \geq 0} \Delta(H_n) t^n = \log \sum_{n \geq 0} \left( \sum_{i+j=n} H_i \otimes H_j \right) t^n
\]
\[
= \log \left( \left( \sum_{i \geq 0} H_i t^i \right) \otimes \left( \sum_{j \geq 0} H_j t^j \right) \right) = \log \left( \sum_{i \geq 0} H_i t^i \otimes 1 + 1 \otimes \sum_{j \geq 0} H_j t^j \right)
\]
\[
= \log \tilde{H}(t) \otimes 1 + 1 \otimes \log \tilde{H}(t) = \xi(t) \otimes 1 + 1 \otimes \xi(t).
\]
Comparing coefficients in this equality yields \( \Delta(\xi_n) = \xi_n \otimes 1 + 1 \otimes \xi_n \). Thus \( S(\xi_n) = -\xi_n \), by Proposition 1.37. This allows one to determine \( S(H_\alpha) \) and \( S(H_\alpha) \), after one first inverts the relation (5.17) to get that \( \tilde{H}(t) = \exp \xi(t) \), and hence
\[
S(\tilde{H}(t)) = S(\exp \xi(t)) = \exp S(\xi(t)) = \exp (-\xi(t)) = (\exp(\xi(t)))^{-1}
\]
\[
= \tilde{H}(t)^{-1} = (1 + H_1 t + H_2 t^2 + \cdots)^{-1}.
\]
Upon expanding the right side, and comparing coefficients of \( t^n \), this gives
\[
S(H_n) = \sum_{\beta \in \text{Comp}_n} (-1)^{\ell(\beta)} H_\beta
\]
and hence
\[
S(H_\alpha) = S(H_{\alpha_1}) \cdots S(H_{\alpha_2}) S(H_{\alpha_1}) = \sum_{\gamma \text{ refines } \text{rev}(\alpha)} (-1)^{\ell(\gamma)} H_\gamma.
\]
As \( S_{\text{NSym}}, S_{\text{QSym}} \) are adjoint, and \( \{ H_\alpha \}, \{ M_\alpha \} \) are dual bases, this is equivalent to Theorem 5.11:

\[
S(M_\alpha) = (-1)^{\ell(\alpha)} \sum_{\gamma \text{ coarses } \text{rev}(\alpha)} M_\gamma
\]

(because if \( \mu \) and \( \nu \) are two compositions, then \( \mu \) coarsens \( \nu \) if and only if \( \text{rev}(\mu) \) coarsens \( \text{rev}(\nu) \)). Thus, Theorem 5.11 is proven once again.

Let us say a bit more about the elements \( \xi_n \) defined in (5.17) above. The elements \( n \xi_n \) are noncommutative analogues of the power sum symmetric functions \( p_n \) (and, indeed, are lifts of the latter to \( \text{NSym} \), as Exercise 5.37 below shows). They are called the cospecial power sums of the second kind and, indeed, are lifts of the latter to \( \text{NSym} \). They are furthermore useful in studying the so-called Eulerian idempotent of a cocommutative Hopf algebra, as shown in Exercise 5.38 below.

**Exercise 5.37.** Assume that \( \mathbb{Q} \) is a subring of \( \mathbb{k} \). Define a sequence of elements \( \xi_1, \xi_2, \xi_3, \ldots \) of \( \text{NSym} = \text{NSym}_k \) by (5.17).

(a) For every \( n \geq 1 \), show that \( \xi_n \) is a primitive homogeneous element of \( \text{NSym} \) of degree \( n \).

(b) For every \( n \geq 1 \), show that \( \pi (n \xi_n) \) is the \( n \)-th power sum symmetric function \( p_n \in \Lambda \).

(c) For every \( n \geq 1 \), show that

\[
\xi_n = \sum_{\alpha \in \text{Comp}_n} (-1)^{\ell(\alpha)} (\tfrac{\ell(\alpha)}{\ell(\alpha)})! H_\alpha.
\]

(d) For every composition \( \alpha \), define an element \( \xi_\alpha \) of \( \text{NSym} \) by \( \xi_\alpha = \xi_{\alpha_1} \xi_{\alpha_2} \cdots \xi_{\alpha_\ell} \), where \( \alpha \) is written in the form \( \alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell) \) with \( \ell = \ell(\alpha) \). Show that

\[
H_n = \sum_{\alpha \in \text{Comp}_n} \frac{1}{\ell(\alpha)!} \xi_\alpha
\]

for every \( n \in \mathbb{N} \).

Use this to prove that \((\xi_\alpha)_{\alpha \in \text{Comp}_n}\) is a \( k \)-basis of \( \text{NSym}_n \) for every \( n \in \mathbb{N} \).

**Exercise 5.38.** Assume that \( \mathbb{Q} \) is a subring of \( \mathbb{k} \). Let \( A \) be a cocommutative connected graded \( \mathbb{k} \)-bialgebra. Let \( A = \bigoplus_{n \geq 0} A_n \) be the decomposition of \( A \) into homogeneous components. If \( f \) is any \( \mathbb{k} \)-linear map \( A \to A \) annihilating \( A_0 \), then \( f \) is locally \(*\)-nilpotent\(^{194}\), and so the sum \( \log^* (f + u e) := \sum_{n \geq 1} (-1)^{n-1} \frac{1}{n} f^* e^n \) is a well-defined endomorphism of \( A \). Let \( \epsilon \) denote the endomorphism \( \log^* (\text{id}_A) \) of \( A \) (obtained by setting \( f = \text{id}_A - u e : A \to A \)). Show that \( \epsilon \) is a projection from \( A \) to the \( \mathbb{k} \)-submodule \( p \) of all primitive elements of \( A \) (and thus, in particular, is idempotent).

**Hint:** For every \( n \geq 0 \), let \( \pi_n : A \to A \) be the projection onto the \( n \)-th homogeneous component \( A_n \). Since \( \text{NSym} \) is the free \( \mathbb{k} \)-algebra with generators \( H_1, H_2, H_3, \ldots \), we can define a \( \mathbb{k} \)-algebra homomorphism \( \mathcal{W} : \text{NSym} \to (\text{End} A, \ast) \) by sending \( H_n \) to \( \pi_n \). Show that:

(a) The map \( \epsilon : A \to A \) is graded. For every \( n \geq 0 \), we will denote the map \( \pi_n \circ \epsilon = \epsilon \circ \pi_n : A \to A \) by \( \epsilon_n \).

(b) We have \( \mathcal{W}(\xi_n) = \epsilon_n \) for all \( n \geq 1 \), where \( \xi_n \) is defined as in Exercise 5.37.

(c) If \( w \) is an element of \( \text{NSym} \), and if we write \( \Delta(w) = \sum (w) w_1 \otimes w_2 \) using Sweedler’s notation, then

\[
\Delta \circ \mathcal{W}(w) = \left( \sum_{(w)} \mathcal{W}(w_1) \otimes \mathcal{W}(w_2) \right) \circ \Delta.
\]

(d) We have \( \epsilon_n(A) \subset p \) for every \( n \geq 0 \).

(e) We have \( \epsilon(A) \subset p \).

(f) The map \( \epsilon \) fixes any element of \( p \).

**Remark 5.39.** The endomorphism \( \epsilon \) of Exercise 5.38 is known as the Eulerian idempotent of \( A \), and can be contrasted with the Dynkin idempotent of Remark 1.66. It has been studied in [123], [126], [25] and [44], and relates to the Hochschild cohomology of commutative algebras [99, §4.5.2].

**Exercise 5.40.** Assume that \( \mathbb{Q} \) is a subring of \( \mathbb{k} \). Let \( A_n \) and \( \epsilon \) be as in Exercise 5.38.

(a) Show that \( \epsilon^* n \circ \epsilon^* m = n! \delta_{n,m} \epsilon^* n \) for all \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \).
E.g., if \( \alpha \) ribbon because column-strict tableaux \( T \) bottom to top. We claim

The second assertion follows from the first by inclusion-exclusion.

**Proof.** (a) For the first assertion, note that

\[
H_\alpha = \sum_{\beta \text{ coarsens } \alpha} R_\beta
\]

(5.20)

\[
R_\alpha = \sum_{\beta \text{ coarsens } \alpha} (-1)^{l(\beta) - l(\alpha)} H_\beta
\]

(5.21)

(b) The surjection \( \text{NSym} \twoheadrightarrow \Lambda \) sends \( R_\alpha \mapsto s_\alpha \), the skew Schur function associated to the ribbon \( \alpha \).

(c) Furthermore,

\[
R_\alpha R_\beta = R_{\alpha \cdot \beta} + R_{\alpha \bigcirc \beta} \quad \text{if } \alpha \text{ and } \beta \text{ are nonempty}
\]

(5.22)

\[
S(R_\alpha) = (-1)^{|\alpha|} R_{\omega(\alpha)}
\]

(5.23)

Finally, \( R_\beta \) is the multiplicative identity of \( \text{NSym} \).

**Proof.** (a) For the first assertion, note that

\[
H_\alpha = \sum_{\beta} (H_\alpha, L_\beta) R_\beta = \sum_{\beta} \left( H_\alpha, \sum_{\gamma \text{ refines } \beta} M_\gamma \right) R_\beta = \sum_{\beta \text{ coarsens } \alpha} R_\beta.
\]

The second assertion follows from the first by inclusion-exclusion.

(b) Write \( \alpha \) as \( (\alpha_1, \ldots, \alpha_\ell) \). To show that \( \pi(R_\alpha) = s_\alpha \), we instead examine \( \pi(H_\alpha) \):

\[
\pi(H_\alpha) = \pi(H_{\alpha_1} \cdots H_{\alpha_\ell}) = h_{\alpha_1} \cdots h_{\alpha_\ell} = s_{(\alpha_1)} \cdots s_{(\alpha_\ell)} = s_{\alpha_1 \oplus \cdots \oplus \alpha_\ell}
\]

where \( (\alpha_1) \oplus \cdots \oplus (\alpha_\ell) \) is some skew shape which is a horizontal strip having rows of lengths \( \alpha_1, \ldots, \alpha_\ell \) from bottom to top. We claim

\[
s_{(\alpha_1) \oplus \cdots \oplus (\alpha_\ell)} = \sum_{\beta \text{ coarsens } \alpha} s_\beta
\]

because column-strict tableaux \( T \) of shape \( (\alpha_1) \oplus \cdots \oplus (\alpha_\ell) \) bijection to column-strict tableaux \( T' \) of some ribbon \( \beta \) coarsening \( \alpha \), as follows: let \( a_i, b_i \) denote the leftmost, rightmost entries of the \( i \)th row from the bottom in \( T \), of length \( \alpha_i \), and

- if \( b_i \leq a_i+1 \), merge parts \( \alpha_i, \alpha_i+1 \) in \( \beta \), and concatenate the rows of length \( \alpha_i, \alpha_i+1 \) in \( T' \), or
- if \( b_i > a_i+1 \), do not merge parts \( \alpha_i, \alpha_i+1 \) in \( \beta \), and let these two rows overlap in one column in \( T' \).

E.g., if \( \alpha = (3,3,2,3,2) \), this \( T \) of shape \( (3) \oplus \cdots \oplus (5) \) maps to this \( T' \) of shape \( \beta = (3,8,2) \):

\[
T = \begin{array}{ccc}
4 & 4 & 5 \\
2 & 2 & 3 \\
1 & 1 & 3
\end{array} \quad \mapsto \quad \begin{array}{cccccc}
3 & 4 & 4 & 4 & 5 \\
2 & 2 & 3 & 4 & 4 & 4 & 5
\end{array}
\]

The reverse bijection breaks the rows of \( T' \) into the rows of \( T \) of lengths dictated by the parts of \( \alpha \). Having shown \( \pi(H_\alpha) = \sum_{\beta \text{ coarsens } \alpha} s_\beta \), the relation (5.20) and inclusion-exclusion show \( \pi(R_\alpha) = s_\alpha \).

(c) Finally, (5.22) and (5.23) follow from (5.7) and (5.9) by duality. \( \square \)

**Remark 5.43.** Since the maps

\[
\begin{array}{ccc}
\text{NSym} & \xrightarrow{\pi} & \text{QSym}
\end{array}
\]

\[\alpha\]

\[\pi\]

\[i\]
are Hopf morphisms, they must respect the antipodes $S_A, S_{QSym}, S_{NSym}$, but it is interesting to compare them explicitly using the fundamental basis for $QSym$ and the ribbon basis for $NSym$.

On one hand, $(5.9)$ shows that $S_{QSym}(L_\alpha) = (-1)^{i_\alpha}L_{\omega(\alpha)}$ extends the map $S_A$ since $L_{(1^n)} = e_n$ and $L_{(n)} = h_n$, as observed in Example 5.16, and $\omega(n) = (1^n)$.

On the other hand, $(5.23)$ shows that $S_{NSym}(R_\alpha) = (-1)^{i\alpha}R_{\omega(\alpha)}$ lifts the map $S_A$ to $S_{NSym}$: Theorem 5.42(b) showed that $R_\alpha$ lifts the skew Schur function $s_\alpha$, while (2.21) asserted that $S(s_{\lambda/\mu}) = (-1)^{|\lambda/\mu|}s_{\lambda/\mu}$, and a ribbon $\alpha = \lambda/\mu$ has $\omega(\alpha) = \lambda^\ell/\mu^\ell$.

Exercise 5.44. (a) Show that any integers $n$ and $i$ with $0 \leq i < n$ satisfy

$$R_{(1^i,n-i)} = \sum_{j=0}^{i} (-1)^{i-j} R_{(1^j)} H_{n-j}.$$  

(Here, as usual, $1^i$ stands for the number 1 repeated $i$ times.)

(b) Show that any integers $n$ and $i$ with $0 \leq i < n$ satisfy

$$(-1)^i R_{(1^i,n-i)} = \sum_{j=0}^{i} S(H_j) H_{n-j}.$$  

(c) For every positive integer $n$, define an element $\Psi_n$ of $NSym$ by

$$\Psi_n = \sum_{i=0}^{n-1} (-1)^i R_{(1^i,n-i)}.$$  

Show that $\Psi_n = (S \ast E)(H_n)$, where the map $E : NSym \rightarrow NSym$ is defined as in Exercise 1.65 (for $A = NSym$). Conclude that $\Psi_n$ is primitive.

(d) Prove that

$$\sum_{k=0}^{n-1} H_k \Psi_{n-k} = nH_n$$  

for every $n \in \mathbb{N}$.

(e) Define two power series $\psi(t)$ and $\overline{H}(t)$ in $NSym[[t]]$ by

$$\psi(t) = \sum_{n \geq 1} \Psi_n t^{n-1};$$  

$$\overline{H}(t) = \sum_{n \geq 0} H_n t^n.$$  

Show that $\text{195} \frac{d}{dt} \overline{H}(t) = \overline{H}(t) \cdot \psi(t)$.

(The functions $\Psi_n$ are called noncommutative power sums of the first kind; they are studied in [57]. The power sums of the second kind are the $n\xi_n$ in Remark 5.36.)

(f) Show that $\pi(\Psi_n)$ equals the power sum symmetric function $p_n$ for every positive integer $n$.

(g) Show that every positive integer $n$ satisfies

$$p_n = \sum_{i=0}^{n-1} (-1)^i s_{(n-i,1^i)} \quad \text{in } \Lambda.$$  

(h) For every nonempty composition $\alpha$, define a positive integer $lp(\alpha)$ by $lp(\alpha) = \alpha_\ell$, where $\alpha$ is written in the form $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ with $\ell = \ell(\alpha)$. (Thus, $lp(\alpha)$ is the last part of $\alpha$.) Show that every positive integer $n$ satisfies

$$\Psi_n = \sum_{\alpha \in \text{Comp}_n} (-1)^{\ell(\alpha)-1} lp(\alpha) H_{\alpha}.$$  

\text{195} The derivative $\frac{d}{dt}Q(t)$ of a power series $Q(t) \in R[[t]]$ over a noncommutative ring $R$ is defined just as in the case of $R$ commutative: by setting $\frac{d}{dt}Q(t) = \sum_{i \geq 1} q_i t^{i-1}$, where $Q(t)$ is written in the form $Q(t) = \sum_{i \geq 0} q_i t^i$. 

(i) Assume that \( \mathbb{Q} \) is a subring of \( \mathbf{k} \). For every composition \( \alpha \), define an element \( \Psi_\alpha \) of \( \text{NSym} \) by \( \Psi_\alpha = \Psi_{\alpha_1} \Psi_{\alpha_2} \cdots \Psi_{\alpha_\ell} \), where \( \alpha \) is written in the form \( \alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell) \) with \( \ell = \ell(\alpha) \). For every composition \( \alpha \), define \( \pi_\alpha(\alpha) \) to be the positive integer \( \alpha_1 (\alpha_1 + \alpha_2) \cdots (\alpha_1 + \alpha_2 + \cdots + \alpha_\ell) \), where \( \alpha \) is written in the form \( \alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell) \) with \( \ell = \ell(\alpha) \). Show that

\[
H_\alpha = \sum_{\alpha \in \text{Comp}_n} \frac{1}{\pi_\alpha(\alpha)} \Psi_\alpha
\]

for every \( n \in \mathbb{N} \).

Use this to prove that \( (\Psi_\alpha)_{\alpha \in \text{Comp}_n} \) is a \( \mathbf{k} \)-basis of \( \text{NSym}_n \) for every \( n \in \mathbb{N} \).

(j) Assume that \( \mathbb{Q} \) is a subring of \( \mathbf{k} \). Let \( V \) be the free \( \mathbf{k} \)-module with basis \( (b_n)_{n \in \{1, 2, 3, \ldots\}} \). Define a \( \mathbf{k} \)-module homomorphism \( f : V \to \text{NSym} \) by requiring that \( f(b_n) = \Psi_n \) for every \( n \in \{1, 2, 3, \ldots\} \). Let \( F \) be the \( \mathbf{k} \)-algebra homomorphism \( T(V) \to \text{NSym} \) induced by this \( f \) (using the universal property of the tensor algebra \( T(V) \)). Show that \( F \) is a Hopf algebra isomorphism (where the Hopf algebra structure on \( T(V) \) is as in Example 1.38).

(k) Assume that \( \mathbb{Q} \) is a subring of \( \mathbf{k} \). Let \( V \) be as in Exercise 5.44(j). Show that \( \text{QSym} \) is isomorphic to the shuffle algebra \( \text{Sh}(V) \) (defined as in Proposition 1.73) as Hopf algebras.

(l) Solve parts (a) and (b) of Exercise 2.90 again using the ribbon basis functions \( R_\alpha \).

One might wonder whether the Frobenius endomorphisms of \( \Lambda \) (defined in Exercise 2.85) and the Verschiebung endomorphisms of \( \Lambda \) (defined in Exercise 2.86) generalize to analogous operators on either \( \text{QSym} \) or \( \text{NSym} \). The next two exercises (whose claims mostly come from [65, §13]) answer this question: The Frobenius endomorphisms extend to \( \text{QSym} \), and the Verschiebung ones lift to \( \text{NSym} \).

**Exercise 5.45.** For every \( n \in \{1, 2, 3, \ldots\} \), define a map \( \mathbf{F}_n : \text{QSym} \to \text{QSym} \) by setting \( \mathbf{F}_n(a) = a(x_1^n, x_2^n, x_3^n, \ldots) \) for every \( a \in \text{QSym} \).

(1) Show that \( \mathbf{F}_n : \text{QSym} \to \text{QSym} \) is a \( \mathbf{k} \)-algebra homomorphism for every \( n \in \{1, 2, 3, \ldots\} \).

(b) Show that \( \mathbf{F}_n \circ \mathbf{F}_m = \mathbf{F}_{nm} \) for any two positive integers \( n \) and \( m \).

(c) Show that \( \mathbf{F}_1 = \text{id} \).

(d) Prove that \( \mathbf{F}_n(M(\beta_1, \beta_2, \ldots, \beta_s)) = M(n\beta_1, n\beta_2, \ldots, n\beta_s) \) for every \( n \in \{1, 2, 3, \ldots\} \) and \( (\beta_1, \beta_2, \ldots, \beta_s) \in \text{Comp} \).

(e) Prove that \( \mathbf{F}_n : \text{QSym} \to \text{QSym} \) is a Hopf algebra homomorphism for every \( n \in \{1, 2, 3, \ldots\} \).

(f) Consider the maps \( f_n : \Lambda \to \Lambda \) defined in Exercise 2.85. Show that \( \mathbf{F}_n \mid_{\Lambda} = f_n \) for every \( n \in \{1, 2, 3, \ldots\} \).

(g) Assume that \( \mathbf{k} = \mathbb{Z} \). Prove that \( f_p(a) \equiv a^p \text{ mod } p \) \( \text{QSym} \) for every \( a \in \text{QSym} \) and every prime number \( p \).

(h) Give a new solution to Exercise 2.85(d).

**Exercise 5.46.** For every \( n \in \{1, 2, 3, \ldots\} \), define a \( \mathbf{k} \)-algebra homomorphism \( \mathbf{V}_n : \text{NSym} \to \text{NSym} \) by

\[
\mathbf{V}_n(H_m) = \begin{cases} 
H_{m/n}, & \text{if } n \mid m; \\
0, & \text{if } n \nmid m;
\end{cases}
\]

for every positive integer \( m \).

(a) Show that any positive integers \( n \) and \( m \) satisfy

\[
\mathbf{V}_n(\Psi_m) = \begin{cases} 
\Psi_{m/n}, & \text{if } n \mid m; \\
0, & \text{if } n \nmid m;
\end{cases}
\]

where the elements \( \Psi_m \) and \( \Psi_{m/n} \) of \( \text{NSym} \) are as defined in Exercise 5.44(c).

(b) Show that if \( \mathbb{Q} \) is a subring of \( \mathbf{k} \), then any positive integers \( n \) and \( m \) satisfy

\[
\mathbf{V}_n(\xi_m) = \begin{cases} 
\xi_{m/n}, & \text{if } n \mid m; \\
0, & \text{if } n \nmid m;
\end{cases}
\]

where the elements \( \xi_m \) and \( \xi_{m/n} \) of \( \text{NSym} \) are as defined in Exercise 5.37.

196 This is well-defined, since \( \text{NSym} \) is (isomorphic to) the free associative algebra with generators \( H_1, H_2, H_3, \ldots \) (according to (5.12)).
(c) Prove that $V_n \circ V_m = V_{nm}$ for any two positive integers $n$ and $m$.
(d) Prove that $V_1 = \text{id}$.
(e) Prove that $V_n : \text{NSym} \to \text{NSym}$ is a Hopf algebra homomorphism for every $n \in \{1, 2, 3, \ldots\}$.

Now, consider also the maps $F_n : \text{QSym} \to \text{QSym}$ defined in Exercise 2.85. Fix a positive integer $n$.
(f) Prove that the maps $F_n : \text{QSym} \to \text{QSym}$ and $V_n : \text{NSym} \to \text{NSym}$ are adjoint with respect to the dual pairing $\text{NSym} \otimes \text{QSym} \to \kappa$.
(g) Consider the maps $v_n : \Lambda \to \Lambda$ defined in Exercise 2.86. Show that the surjection $\pi : \text{NSym} \to \Lambda$ satisfies $v_n \circ \pi = \pi \circ V_n$ for every $n \in \{1, 2, 3, \ldots\}$.
(h) Give a new solution to Exercise 2.86(f).

6. POLYNOMIAL GENERATORS FOR QSYM AND LYNDON WORDS

In this chapter, we shall construct an algebraically independent generating set for $\text{QSym}$ as a $k$-algebra, thus showing that $\text{QSym}$ is a polynomial algebra over $k$. This has been done by Malvenuto [109, Cor. 4.19] when $k$ is a field of characteristic 0, and by Hazewinkel [64] in the general case. We will begin by introducing the notion of Lyndon words (Section 6.1), on which both of these constructions rely; we will then (Section 6.2) elucidate the connection of Lyndon words with shuffles, and afterwards (Section 6.3) apply it to prove Radford’s theorem stating that the shuffle algebra of a free $k$-module over a commutative $\mathbb{Q}$-algebra is a polynomial algebra (Theorem 6.63). The shuffle algebra is not yet $\text{QSym}$, but Radford’s theorem on the shuffle algebra serves as a natural stepping stone for the study of the more complicated algebra $\text{QSym}$. We will prove – in two ways – that $\text{QSym}$ is a polynomial algebra when $\mathbb{Q}$ is a subring of $k$ in Section 6.4, and then we will finally prove the general case in Section 6.4.

Strictly speaking, this whole Chapter 6 is a digression, as it involves almost no coalgebraic or Hopf-algebraic structures, and its results will not be used in further chapters (which means it can be skipped if so desired). However, it sheds additional light on $\text{QSym}$ and serves as an excuse to study Lyndon words, which are a combinatorial object of independent interest (and are involved in the study of free algebras and Hopf algebras, apart from $\text{QSym}$ – see [131] and [134] [197]).

We will take a scenic route to the proof of Hazewinkel’s theorem. A reader only interested in the proof proper can restrict themselves to reading only the following:

- from Section 6.1, everything up to Corollary 6.6, then from Definition 6.13 up to Proposition 6.18, then from Definition 6.25 up to Lemma 6.28, and finally Theorem 6.30. (Proposition 6.19 and Theorem 6.20 are also relevant if one wants to use a different definition of Lyndon words, as they prove the equivalence of most such definitions.)
- from Section 6.3, Definition 6.60, Lemma 6.65, and Lemma 6.68.
- from Section 6.4, Definition 6.72, Theorem 6.74, then from Proposition 6.76 up to Definition 6.80, and Lemma 6.82.
- all of Section 6.5.

6.1. Lyndon words. Lyndon words have been independently defined by Shirshov [150], Lyndon [105], Radford [131, §2] and de Bruijn/Klarner [23] (though using different and sometimes incompatible notations). They have since been surfacing in various places in noncommutative algebra (particularly the study of free Lie algebras); expositions of their theory can be found in [103, §5] and [134, §5.1]. We will follow our own approach to the properties of Lyndon words that we need.

Definition 6.1. We fix a totally ordered set $\mathfrak{A}$, which we call the alphabet. Throughout Section 6.1 and Section 6.2, we will understand “word” to mean a word over $\mathfrak{A}$.

We recall that a word is just a (finite) tuple of elements of $\mathfrak{A}$. In other words, a word is an element of the set $\bigsqcup_{n \geq 0} \mathfrak{A}^n$. We denote this set by $\mathfrak{A}^*$.

The empty word is the unique tuple with 0 elements. It is denoted by $\varnothing$. If $w \in \mathfrak{A}^n$ is a word and $i \in \{1, 2, \ldots, n\}$, then the $i$-th letter of $w$ means the $i$-th entry of the $n$-tuple $w$. This $i$-th letter will be denoted by $w_i$.

[197]They also are involved in indexing basis elements of combinatorial Hopf algebras other than $\text{QSym}$. See Bergeron/Zabrocki [17].
The length \( \ell(w) \) of a word \( w \in \bigcup_{n \geq 0} \mathfrak{A}^n \) is defined to be the \( n \in \mathbb{N} \) satisfying \( w \in \mathfrak{A}^n \). Thus, \( w = (w_1, w_2, \ldots, w_{\ell(w)}) \) for every word \( w \).

The concatenation of two words \( u \) and \( v \) is defined to be the word \( (u_1, u_2, \ldots, u_{\ell(u)}, v_1, v_2, \ldots, v_{\ell(v)}) \). This concatenation is denoted by \( uv \) or \( u \cdot v \). The set \( \mathfrak{A}^* \) of all words is a monoid with respect to concatenation, with neutral element \( \emptyset \). It is precisely the free monoid on generators \( \mathfrak{A} \). If \( u \) is a word and \( i \in \mathbb{N} \), we will understand \( u^i \) to mean the \( i \)-th power of \( u \) in this monoid (that is, the word \( uu \cdots u \)).

The elements of \( \mathfrak{A} \) are called letters, and will be identified with elements of \( \mathfrak{A}^1 \subset \bigcup_{n \geq 0} \mathfrak{A}^n = \mathfrak{A}^* \). This identification equates every letter \( u \in \mathfrak{A} \) with the one-letter word \( (u) \in \mathfrak{A}^1 \). Thus, every word \( (u_1, u_2, \ldots, u_n) \in \mathfrak{A}^* \) equals the concatenation \( u_1 u_2 \cdots u_n \) of letters, hence allowing us to use \( u_1 u_2 \cdots u_n \) as a brief notation for the word \( (u_1, u_2, \ldots, u_n) \).

If \( w \) is a word, then:

- a prefix of \( w \) means a word of the form \( (w_1, w_2, \ldots, w_i) \) for some \( i \in \{0, 1, \ldots, \ell(w)\} \);
- a suffix of \( w \) means a word of the form \( (w_{i+1}, w_{i+2}, \ldots, w_{\ell(w)}) \) for some \( i \in \{0, 1, \ldots, \ell(w)\} \);
- a proper suffix of \( w \) means a word of the form \( (w_{i+1}, w_{i+2}, \ldots, w_{\ell(w)}) \) for some \( i \in \{1, 2, \ldots, \ell(w)\} \).

In other words,

- a prefix of \( w \in \mathfrak{A}^* \) is a word \( u \in \mathfrak{A}^* \) such that there exists a \( v \in \mathfrak{A}^* \) satisfying \( w = uv \);

- a suffix of \( w \in \mathfrak{A}^* \) is a word \( v \in \mathfrak{A}^* \) such that there exists a \( u \in \mathfrak{A}^* \) satisfying \( w = uv \);

- a proper suffix of \( w \in \mathfrak{A}^* \) is a word \( v \in \mathfrak{A}^* \) such that there exists a nonempty \( u \in \mathfrak{A}^* \) satisfying \( w = uv \).

We define a relation \( \leq \) on the set \( \mathfrak{A}^* \) as follows: For two words \( u \in \mathfrak{A}^* \) and \( v \in \mathfrak{A}^* \), we set \( u \leq v \) to hold if and only if

- either there exists an \( i \in \{1, 2, \ldots, \min \{\ell(u), \ell(v)\}\} \) such that \( u_i < v_i \), and every \( j \in \{1, 2, \ldots, i - 1\} \) satisfies \( u_j = v_j \),

  or the word \( u \) is a prefix of \( v \).

This order relation (taken as the smaller-or-equal relation) makes \( \mathfrak{A}^* \) into a poset (by Proposition 6.2(a) below), and we will always be regarding \( \mathfrak{A}^* \) as endowed with this poset structure (thus, notations such as \( < \), \( \leq \), \( > \) and \( \geq \) will be referring to this poset structure). This poset is actually totally ordered (see Proposition 6.2(a)).

Here are some examples of words compared by the relation \( \leq \):

\[
\begin{align*}
113 & \leq 114, & 113 & \leq 132, & 19 & \leq 195, & 41 & \leq 412, \\
41 & \leq 421, & 539 & \leq 54, & \emptyset & \leq 21, & \emptyset & \leq \emptyset
\end{align*}
\]

(where \( \mathfrak{A} \) is the alphabet \( \{1 < 2 < 3 < \cdots \} \)).

Notice that if \( u \) and \( v \) are two words of the same length (i.e., we have \( u, v \in \mathfrak{A}^n \) for one and the same \( n \)), then \( u \leq v \) holds if and only if \( u \) is lexicographically smaller-or-equal to \( v \). In other words, the relation \( \leq \) is an extension of the lexicographic order on every \( \mathfrak{A}^n \) to \( \mathfrak{A}^* \). This is the reason why this relation \( \leq \) is usually called the lexicographic order on \( \mathfrak{A}^* \). In particular, we will be using this name.\(^{198}\) However, unlike the lexicographic order on \( \mathfrak{A}^n \), it does not always respect concatenation from the right: It can happen that \( u, v, w \in \mathfrak{A}^* \) satisfy \( u \leq v \) but not \( uw \leq vw \). (For example, \( u = 1, v = 13 \) and \( w = 4 \), again with \( \mathfrak{A} = \{1 < 2 < 3 < \cdots \} \).) We will see in Proposition 6.2 that this is rather an exception than the rule and the relation \( \leq \) still behaves mostly predictably with respect to concatenation.

Some basic properties of the order relation \( \leq \) just defined are collected in the following proposition:

**Proposition 6.2.**

(a) The order relation \( \leq \) is (the smaller-or-equal relation of) a total order on the set \( \mathfrak{A}^* \).

(b) If \( a, c, d \in \mathfrak{A}^* \) satisfy \( c \leq d \), then \( ac \leq ad \).

(c) If \( a, c, d \in \mathfrak{A}^* \) satisfy \( ac \leq ad \), then \( c \leq d \).

(d) If \( a, b, c \in \mathfrak{A}^* \) satisfy \( a \leq c \), then either we have \( ab \leq cd \) or the word \( a \) is a prefix of \( c \).

\(^{198}\)The relation \( \leq \) is also known as the dictionary order, due to the fact that it is the order in which words appear in a dictionary.
(e) If $a, b, c, d \in \mathbb{A}^*$ satisfy $ab \leq cd$, then either we have $a \leq c$ or the word $c$ is a prefix of $a$.

(f) If $a, b, c, d \in \mathbb{A}^*$ satisfy $ab \leq cd$ and $\ell(a) \leq \ell(c)$, then $a \leq c$.

(g) If $a, b, c \in \mathbb{A}^*$ satisfy $a \leq b \leq ac$, then $a$ is a prefix of $b$.

(h) If $a \in \mathbb{A}^*$ is a prefix of $b \in \mathbb{A}^*$, then $a \leq b$.

(i) If $a$ and $b$ are two prefixes of $c \in \mathbb{A}^*$, then either $a$ is a prefix of $b$, or $b$ is a prefix of $a$.

(j) If $a, b, c \in \mathbb{A}^*$ are such that $a \leq b$ and $\ell(a) \leq \ell(b)$, then $ac \leq bc$.

(k) If $a \in \mathbb{A}^*$ and $b \in \mathbb{A}^*$ are such that $b$ is nonempty, then $a < ab$.

**Exercise 6.3.** Prove Proposition 6.2.

**Hint:** No part of Proposition 6.2 requires more than straightforward case analysis. However, the proof of (a) can be simplified by identifying the order relation $\leq$ on $\mathbb{A}^*$ as a restriction of the lexicographic order on the set $\mathbb{B}^\infty$, where $\mathbb{B}$ is a suitable extension of the alphabet $\mathbb{A}$. What is this extension, and how to embed $\mathbb{A}^*$ into $\mathbb{B}^\infty$?

Proposition 6.2 provides a set of tools for working with the lexicographic order without having to refer to its definition; we shall use it extensively. Proposition 6.2(h) (and its equivalent form stating that $a \leq ac$ for every $a \in \mathbb{A}^*$ and $c \in \mathbb{A}^*$) and Proposition 6.2(k) will often be used without explicit mention.

Before we define Lyndon words, let us show two more facts about words which will be used later. First, when do words commute?

**Proposition 6.4.** Let $u, v \in \mathbb{A}^*$ satisfy $uw = vu$. Then, there exist a $t \in \mathbb{A}^*$ and two nonnegative integers $n$ and $m$ such that $u = t^n$ and $v = t^m$.

**Proof.** We prove this by strong induction on $\ell(u) + \ell(v)$. We assume WLOG that $\ell(u)$ and $\ell(v)$ are positive (because otherwise, one of $u$ and $v$ is the empty word, and everything is trivial). It is easy to see that either $u$ is a prefix of $v$, or $v$ is a prefix of $u$.\footnote{\textit{Proof.} The word $u$ is a prefix of $v$. But the word $v$ is also a prefix of $w$ (since $wv = vu$). Hence, Proposition 6.2(i) applied to $a = u, b = v$ and $c = w$ yields that either $u$ is a prefix of $w$ or $w = vu$.}

We assume WLOG that $u$ is a prefix of $v$ (since our situation is symmetric). Thus, we can write $v$ in the form $v = uw$ for some $w \in \mathbb{A}^*$. Consider this $w$. Clearly,

$$\ell(u) + \ell(w) = \ell\left(\underbrace{uw}_{=v}\right) = \ell(w) < \ell(u) + \ell(v)$$

(since $\ell(v)$ is positive). Since $v = uw$, the equality $uw = vu$ becomes $uwv = uuw$. C cancelling $u$ from this equality, we obtain $wv = uv$. Now, we can apply Proposition 6.4 to $w$ instead of $v$ (by the induction assumption, since $\ell(u) + \ell(w) < \ell(u) + \ell(v)$), and obtain that there exist a $t \in \mathbb{A}^*$ and two nonnegative integers $n$ and $m$ such that $u = t^n$ and $v = t^m$. Consider this $t$ and these $n$ and $m$. Of course, $u = t^n$ and $v = t^m$.

Proposition 6.4 is proven.

**Proposition 6.5.** Let $u, v, w \in \mathbb{A}^*$ be nonempty words satisfying $uv \geq vu$, $vw \geq wu$ and $wu \geq uw$. Then, there exist a $t \in \mathbb{A}^*$ and three nonnegative integers $n$, $m$, and $p$ such that $u = t^n$, $v = t^m$, and $w = t^p$.

**Proof.** We prove this by strong induction on $\ell(u) + \ell(v) + \ell(w)$. Clearly, $\ell(u)$, $\ell(v)$ and $\ell(w)$ are positive (since $u$, $v$, and $w$ are nonempty). We assume WLOG that $\ell(u) = \min\{\ell(u), \ell(v), \ell(w)\}$ (because there is a cyclic symmetry in our situation). Thus, $\ell(u) \leq \ell(v)$ and $\ell(u) \leq \ell(w)$. But $vu \leq uv$. Hence, Proposition 6.2(e) (applied to $a = u, b = u, c = u$ and $d = v$) yields that either we have $v \leq u$ or the word $u$ is a prefix of $v$. But Proposition 6.2(f) (applied to $a = u, b = w, c = w$ and $d = u$) yields $u \leq w$ (since $uv < uw$ and $\ell(u) \leq \ell(w)$). Furthermore, $wu \leq uv$. Hence, Proposition 6.2(e) (applied to $a = w, b = v, c = v$ and $d = w$) yields that either we have $w \leq v$ or the word $w$ is a prefix of $v$.

From what we have found so far, it is easy to see that $u$ is a prefix of $v$. In other words, there exists a $v' \in \mathbb{A}^*$ such that $v = uv'$. Consider this $v'$.

\footnote{\textit{Proof.} Assume the contrary. Then, $u$ is not a prefix of $v$. Hence, we must have $v \leq u$ (since either we have $v \leq u$ or the word $u$ is a prefix of $v$), and in fact $v < u$ (because $v = u$ would contradict to $u$ not being a prefix of $v$). Thus, $v < u \leq w$. But recall that either we have $w \leq v$ or the word $v$ is a prefix of $w$. Thus, $v$ must be a prefix of $w$ (because $v < w$ rules out $w \leq v$). In other words, there exists a $q \in \mathbb{A}^*$ such that $w = vq$. Consider this $q$. We have $v < u \leq w = vq$. Thus, Proposition 6.2(g) (applied to $a = v, b = u$ and $c = q$) yields that $v$ is a prefix of $u$. In light of $\ell(u) \leq \ell(v)$, this is only possible if $v = u$, but this contradicts $v < u$. This contradiction completes this proof.}
If the word \( v' \) is empty, then the statement of Proposition 6.5 can be easily deduced from Proposition 6.4\(^{201} \). Thus, we assume WLOG that this is not the case. Hence, \( v' \) is nonempty.

Using \( v = uv' \), we can rewrite \( wv \geq vu \) as \( uvw' \geq wuv' \). That is, \( uvw' \leq wuv' \) (by Proposition 6.2(c), applied to \( a = u, c = v'u \) and \( d = wv' \)). That is, \( uv' \geq v'u \). But \( \ell(wu) = \ell(u) + \ell(w) = \ell(w) + \ell(u) = \ell(wu) \geq \ell(wu) \). Hence, Proposition 6.2(i) (applied to \( a = wu, b = wu \) and \( c = v' \)) yields \( uuv' \leq wuv' \) (since \( uw \leq wu \)). Now, \( uuv' w = wv' w = wuv' \geq wuv' \) (since \( uvw' \leq wuv' \)), so that \( uuv' \leq wv' w \). Hence, \( wv' \leq wv' w \) (by Proposition 6.2(c), applied to \( a = u, c = wv' \) and \( d = v'w \)), so that \( v'w \geq wv' \).

Proof. Assume that the word \( v'w \) is empty, then everything is obvious. We thus WLOG assume that \( v'w \) is nonempty. Then, \( wv' \geq wv' w \). This contradicts \( wu < wu \), so we have Proposition 6.5.

Proposition 6.5 to \( v' \) instead of \( v \) (by the induction hypothesis, because \( \ell(u) + \ell(v') + \ell(w) = \ell(v) + \ell(w) < \ell(u) + \ell(v) + \ell(w) \)). As a result, we see that there exist a \( t \in \mathfrak{A}^* \) and three nonnegative integers \( n, m, p \) such that \( u = t^n, v' = t^m \) and \( w = t^p \). Clearly, this \( t \) and these \( n, m, p \) satisfy \( v = v' = t^n t^m = t^{n+m} \), and so the statement of Proposition 6.5 is satisfied. The induction step is thus complete.

**Corollary 6.6.** Let \( u, v, w \in \mathfrak{A}^* \) be words satisfying \( uw \geq vu \) and \( wv \geq wv \). Assume that \( v \) is nonempty. Then, \( uw \geq wu \).

Proof. Assume the contrary. Thus, \( uw < wu \), so that \( wu \geq uw \).

If \( u \) or \( w \) is empty, then everything is obvious. We thus WLOG assume that \( u \) and \( w \) are nonempty. Thus, Proposition 6.5 shows that there exist a \( t \in \mathfrak{A}^* \) and three nonnegative integers \( n, m, p \) such that \( u = t^n, v' = t^m \) and \( w = t^p \). But this yields \( uw = t^p t^n = t^{p+n} = t^{n+m} \), contradicting \( uw < wu \).

This contradiction finishes the proof.

**Exercise 6.7.** Find an alternative proof of Corollary 6.6 which does not use Proposition 6.5.

The above results have a curious consequence, which we are not going to use:

**Corollary 6.8.** We can define a preorder on the set \( \mathfrak{A}^* \setminus \{\emptyset\} \) of all nonempty words by defining a nonempty word \( u \) to be greater-or-equal-to a nonempty word \( v \) (with respect to this preorder) if and only if \( uv \geq vu \). Two nonempty words \( u, v \) are equivalent with respect to the equivalence relation induced by this preorder if and only if there exist a \( t \in \mathfrak{A}^* \) and two nonnegative integers \( n, m \) such that \( u = t^n \) and \( v = t^m \).

Proof. The alleged preorder is transitive (by Corollary 6.6) and reflexive (obviously), and hence is really a preorder. The claim in the second sentence follows from Proposition 6.4.

As another consequence of Proposition 6.5, we obtain a classical property of words [103, Proposition 1.3.1]:

**Exercise 6.9.** Let \( u \) and \( v \) be words and \( n \) and \( m \) be positive integers such that \( u^n = v^m \). Prove that there exists a word \( t \) and positive integers \( i \) and \( j \) such that \( u = t^i \) and \( v = t^j \).

Here is another application of Corollary 6.6:

**Exercise 6.10.** Let \( n \) and \( m \) be positive integers. Let \( u \in \mathfrak{A}^* \) and \( v \in \mathfrak{A}^* \) be two words. Prove that \( uv \geq vu \) holds if and only if \( u^n v^m \geq v^m u^n \) holds.

**Exercise 6.11.** Let \( n \) and \( m \) be positive integers. Let \( u \in \mathfrak{A}^* \) and \( v \in \mathfrak{A}^* \) be two words satisfying \( n \ell(u) = m \ell(v) \). Prove that \( uv \geq vu \) holds if and only if \( u^n \geq v^m \) holds.

We can also generalize Propositions 6.4 and 6.5:

**Exercise 6.12.** Let \( u_1, u_2, \ldots, u_k \) be nonempty words such that every \( i \in \{1, 2, \ldots, k\} \) satisfies \( u_i u_{i+1} \geq u_{i+1} u_i \), where \( u_{k+1} \) means \( u_1 \). Show that there exist a word \( t \) and nonnegative integers \( n_1, n_2, \ldots, n_k \) such that \( u_1 = t^{n_1}, u_2 = t^{n_2}, \ldots, u_k = t^{n_k} \).

\(^{201}\)Proof. Assume that the word \( v' \) is empty. Then, \( v = uv' \) becomes \( v = u \). Therefore, \( vu \geq uv \) becomes \( uv \geq wu \). Combined with \( wu \geq uv \), this yields \( uw = wu \). Hence, Proposition 6.4 (applied to \( u \) instead of \( v \)) yields that there exist a \( t \in \mathfrak{A}^* \) and two nonnegative integers \( n \) and \( m \) such that \( u = t^n \) and \( w = t^m \). Clearly, \( v = u = t^n \) as well, and so the statement of Proposition 6.5 is true.
Now, we define the notion of a Lyndon word. There are several definitions in literature, some of which will be proven equivalent in Theorem 6.20.

**Definition 6.13.** A word $w \in \mathcal{A}^*$ is said to be Lyndon if it is nonempty and satisfies the following property: Every nonempty proper suffix $v$ of $w$ satisfies $v > w$.

For example, the word 113 is Lyndon (because its nonempty proper suffixes are 13 and 3, and these are both $> 113$), and the word 242427 is Lyndon (its nonempty proper suffixes are 24247, 2427, 427, 27 and 7, and again these are each $> 242427$). The words 2424 and 35346 are not Lyndon (the word 2424 has a nonempty proper suffix 24 $\leq$ 2424, and the word 35346 has a nonempty proper suffix 346 $\leq$ 35346). Every word of length 1 is Lyndon (since it has no nonempty proper suffixes). A word $w = (w_1, w_2)$ with two letters is Lyndon if and only if $w_1 < w_2$. A word $w = (w_1, w_2, w_3)$ of length 3 is Lyndon if and only if $w_1 < w_3$ and $w_1 \leq w_2$. A four-letter word $w = (w_1, w_2, w_3, w_4)$ is Lyndon if and only if $w_1 < w_4$, $w_1 \leq w_3$, $w_1 \leq w_2$ and (if $w_1 = w_3$ then $w_2 < w_4$). (These rules only get more complicated as the words grow longer.)

We will show several properties of Lyndon words now. We begin with trivialities which will make some arguments a bit shorter:

**Proposition 6.14.** Let $w$ be a Lyndon word. Let $u$ and $v$ be words such that $w = uv$.

(a) If $v$ is nonempty, then $v \geq w$.

(b) If $v$ is nonempty, then $v > u$.

(c) If $u$ and $v$ are nonempty, then $vu > uv$.

(d) We have $vu \geq uw$.

**Proof.** (a) Assume that $v$ is nonempty. Clearly, $v$ is a suffix of $w$ (since $w = uv$). If $v$ is a proper suffix of $w$, then the definition of a Lyndon word yields that $v > w$ (since $w$ is a Lyndon word); otherwise, $v$ must be $w$ itself. In either case, we have $v \geq w$. Hence, Proposition 6.14(a) is proven.

(b) Assume that $v$ is nonempty. From Proposition 6.14(a), we obtain $v \geq w = uv > u$ (since $v$ is nonempty). This proves Proposition 6.14(b).

(c) Assume that $u$ and $v$ are nonempty. Since $u$ is nonempty, we have $vu > v \geq w$ (by Proposition 6.14(a)). Since $w = uv$, this becomes $vu > uv$. This proves Proposition 6.14(c).

(d) We need to prove that $vu \geq w$. If either $u$ or $v$ is empty, $vu$ and $uw$ are obviously equal, and thus $vu \geq w$ is true in this case. Hence, we can WLOG assume that $u$ and $v$ are nonempty. Assume this. Then, $vu \geq uw$ follows from Proposition 6.14(c). This proves Proposition 6.14(d). □

**Corollary 6.15.** Let $w$ be a Lyndon word. Let $v$ be a nonempty suffix of $w$. Then, $v \geq w$.

**Proof.** Since $v$ is a nonempty suffix of $w$, there exists $u \in \mathcal{A}^*$ such that $w = uv$. Thus, $v \geq w$ follows from Proposition 6.14(a). □

Our next proposition is [68, Lemma 6.5.4]; its part (a) is also [134, (5.1.2)]:

**Proposition 6.16.** Let $u$ and $v$ be two Lyndon words such that $u < v$. Then:

(a) The word $uv$ is Lyndon.

(b) We have $uw < v$.

**Proof.** (b) The word $u$ is Lyndon and thus nonempty. Hence, $uw \neq v$. If $uw \leq v$, then Proposition 6.16(b) easily follows. Hence, for the rest of this proof, we can WLOG assume that we don’t have $uw \leq v$. Assume this.

We have $u < v$. Hence, Proposition 6.2(d) (applied to $a = u$, $b = v$, $c = v$ and $d = \emptyset$) yields that either we have $uw \leq v\emptyset$ or the word $u$ is a prefix of $v$. Since we don’t have $uw \leq v\emptyset$, we thus see that the word $u$ is a prefix of $v$. In other words, there exists a $t \in \mathcal{A}^*$ satisfying $v = ut$. Consider this $t$. Then, $t$ is nonempty (else we would have $v = u\omega = u$ in contradiction to $u < v$).

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202 **Proof.** Assume the contrary. Then, $uw = v$. Thus, $uv = v = \omega v$. Cancelling $v$ from this equation, we obtain $u = \emptyset$. That is, $u$ is empty. This contradicts the fact that $u$ is nonempty. This contradiction proves that our assumption was wrong, qed.

203 **Proof.** Assume that $uw \leq v\emptyset$. Thus, $uv \leq v\emptyset = v$. Since $uv \neq v$, this becomes $uv < v$, so that Proposition 6.16(b) is proven.
Now, $v = ut$. Hence, $t$ is a proper suffix of $v$ (proper because $u$ is nonempty). Thus, $t$ is a nonempty proper suffix of $v$. Since every nonempty proper suffix of $v$ is $> v$ (because $v$ is Lyndon), this shows that $t > v$. Hence, $v < t$. Thus, Proposition 6.2(b) (applied to $a = u$, $c = v$ and $d = t$) yields $uv \leq ut = v$. Combined with $w \neq v$, this yields $uv < v$. Hence, Proposition 6.16(b) is proven.

(a) The word $v$ is nonempty (since it is Lyndon). Hence, $uv$ is nonempty. It thus remains to check that every nonempty proper suffix $p$ of $uv$ satisfies $p > uv$.

So let $p$ be a nonempty proper suffix of $uv$. We must show that $p > uv$. Since $p$ is a nonempty proper suffix of $uv$, we must be in one of the following two cases (depending on whether this suffix begins before the suffix $v$ of $uv$ begins or afterwards):

Case 1: The word $p$ is a nonempty suffix of $v$. (Note that $p = v$ is allowed.)

Case 2: The word $p$ has the form $qv$ where $q$ is a nonempty proper suffix of $u$.

Let us first handle Case 1. In this case, $p$ is a nonempty suffix of $v$. Since $v$ is Lyndon, this yields that $p \geq v$ (by Corollary 6.15, applied to $v$ and $p$ instead of $w$ and $v$). But Proposition 6.16(b) yields $uv < v$, thus $v > uv$. Hence, $p \geq v > uv$. We thus have proven $p > uv$ in Case 1.

Let us now consider Case 2. In this case, $p$ has the form $qv$ where $q$ is a nonempty proper suffix of $u$. Consider this $q$. Clearly, $q > u$ (since $u$ is Lyndon and since $q$ is a nonempty proper suffix of $u$), so that $u \leq q$. Thus, Proposition 6.2(d) (applied to $a = u$, $b = v$, $c = q$ and $d = v$) yields that either we have $uv \leq qv$ or the word $u$ is a prefix of $q$. Since $u$ being a prefix of $q$ is impossible (in fact, $q$ is a proper suffix of $u$, thus shorter than $u$), we thus must have $uv \leq qv$. Since $uv \neq qv$ (because otherwise we would have $uv = qv$, thus $u = q$ (because we can cancel $v$ from the equality $uv = qv$), contradicting $q > u$), this can be strengthened to $uv < qv = p$. Thus, $p > uv$ is proven in Case 2 as well.

Now that $p > uv$ is shown to hold in both cases, we conclude that $p > uv$ always holds.

Now, let us forget that we fixed $p$. We have thus shown that every nonempty proper suffix $p$ of $uv$ satisfies $p > uv$. Since $uv$ is nonempty, this yields that $uv$ is Lyndon (by the definition of a Lyndon word). Thus, the proof of Proposition 6.16(a) is complete.

Proposition 6.16(b), combined with Corollary 6.6, leads to a technical result which we will find good use for later:

Corollary 6.17. Let $u$ and $v$ be two Lyndon words such that $u < v$. Let $z$ be a word such that $zv \geq vz$ and $uz \geq zu$. Then, $z$ is the empty word.

Proof. Assume the contrary. Then, $z$ is nonempty. Thus, Corollary 6.6 (applied to $z$ and $v$ instead of $v$ and $w$) yields $uv \geq vu$. But Proposition 6.16(b) yields $uv < v \leq vu$, contradicting $uv \geq vu$. This contradiction completes our proof.

We notice that the preorder of Corollary 6.8 becomes particularly simple on Lyndon words:

Proposition 6.18. Let $u$ and $v$ be two Lyndon words. Then, $u \geq v$ if and only if $uv \geq vu$.

Proof. We distinguish between three cases:

Case 1: We have $u < v$.

Case 2: We have $u = v$.

Case 3: We have $u > v$.

Let us consider Case 1. In this case, we have $u < v$. Thus,

\[
uv < v \quad \text{ (by Proposition 6.16(b))}
\]

\[
\leq vu.
\]

Hence, we have neither $u \geq v$ nor $uv \geq vu$ (because we have $u < v$ and $uv < vu$). Thus, Proposition 6.18 is proven in Case 1.

In Case 2, we have $u = v$. Therefore, in Case 2, both inequalities $u \geq v$ and $uv \geq vu$ hold (and actually are equalities). Thus, Proposition 6.18 is proven in Case 2 as well.

Let us finally consider Case 3. In this case, we have $u > v$. In other words, $v < u$. Thus,

\[
vu < u \quad \text{ (by Proposition 6.16(b), applied to $v$ and $u$ instead of $u$ and $v$)}
\]

\[
\leq uv.
\]
Hence, we have both \( u \geq v \) and \( uv \geq vu \) (because we have \( v < u \) and \( vu < uv \)). Thus, Proposition 6.18 is proven in Case 3.

Proposition 6.18 is now proven in all three possible cases. \( \square \)

**Proposition 6.19.** Let \( w \) be a nonempty word. Let \( v \) be the (lexicographically) smallest nonempty suffix of \( w \). Then:

(a) The word \( v \) is a Lyndon word.

(b) Assume that \( w \) is not a Lyndon word. Then there exists a nonempty \( u \in \mathcal{A}^* \) such that \( w = uv \), \( u \geq v \) and \( uv \geq vu \).

**Proof.** (a) Every nonempty proper suffix of \( v \) is \( \geq v \) (since every nonempty proper suffix of \( v \) is a nonempty suffix of \( w \), but \( v \) is the smallest such suffix) and therefore \( > v \) (since a proper suffix of \( v \) cannot be \( = v \)). Combined with the fact that \( v \) is nonempty, this yields that \( v \) is Lyndon. Proposition 6.19(a) is proven.

(b) Assume that \( w \) is not a Lyndon word. Then, \( w \neq v \) (since \( v \) is Lyndon (by Proposition 6.19(a)) while \( w \) is not). Now, \( v \) is a suffix of \( w \). Thus, there exists an \( u \in \mathcal{A}^* \) such that \( w = uv \). Consider this \( u \). Clearly, \( u \) is nonempty (since \( uv = w \neq v \)). Assume (for the sake of contradiction) that \( u < v \). Let \( v' \) be the (lexicographically) smallest nonempty suffix of \( u \). Then, \( v' \) is a Lyndon word (by Proposition 6.19(a), applied to \( u \) and \( v' \) instead of \( w \) and \( v \)) and satisfies \( v' \leq u \) (since \( u \) is a nonempty suffix of \( u \), whereas \( v' \) is the smallest such suffix). Thus, \( v' \) and \( v \) are Lyndon words such that \( v' \leq u < v \). Proposition 6.16(a) (applied to \( v' \) instead of \( u \)) now yields that the word \( v'v \) is Lyndon. Hence, every nonempty proper suffix of \( v'v \) is \( > v'v \). Since \( v \) is a nonempty proper suffix of \( v'v \), this yields that \( v > v'v \).

But \( v'v \) is a nonempty suffix of \( u \), so that \( v'v \) is a nonempty suffix of \( uv = w \). Since \( v \) is the smallest such suffix, this yields that \( v'v \geq v \). This contradicts \( v > v'v \). Our assumption (that \( u < v \)) therefore fails. We conclude that \( u \geq v \).

It remains to prove that \( uv \geq vu \). Assume the contrary. Then, \( uv < vu \). Thus, there exists at least one suffix \( t \) of \( u \) such that \( tv < vt \) (namely, \( t = u \)). Let \( p \) be the minimum-length such suffix. Then, \( pv < vp \). Thus, \( p \) is nonempty.

Since \( p \) is a suffix of \( u \), it is clear that \( pv \) is a suffix of \( uv = w \). So we know that \( pv \) is a nonempty suffix of \( w \). Since \( v \) is the smallest such suffix, this yields that \( v \leq pv < vp \). Thus, Proposition 6.2(g) (applied to \( a = v \), \( b = pv \) and \( c = p \)) yields that \( v \) is a prefix of \( pv \). In other words, there exists a \( q \in \mathcal{A}^* \) such that \( pv = vq \). Consider this \( q \). This \( q \) is nonempty (because otherwise we would have \( pv = v q = v \), contradicting the fact that \( p \) is nonempty). From \( vq = pv < vp \), we obtain \( q \leq p \) (by Proposition 6.2(c), applied to \( a = v \), \( c = q \) and \( d = p \)).

We know that \( q \) is a suffix of \( pv \) (since \( vq = pv \)), whereas \( pv \) is a suffix of \( w \). Thus, \( q \) is a suffix of \( w \). So \( q \) is a nonempty suffix of \( w \). Since \( v \) is the smallest such suffix, this yields that \( v \leq q \). We now have \( v \leq q \leq p \leq pv < vp \). Hence, \( v \) is a prefix of \( p \) (by Proposition 6.2(g), applied to \( a = v \), \( b = p \) and \( c = p \)). In other words, there exists an \( r \in \mathcal{A}^* \) such that \( p = vr \). Consider this \( r \). Clearly, \( r \) is a suffix of \( p \), while \( r \) is a suffix of \( u \); therefore, \( r \) is a suffix of \( u \). Also, \( pv < vp \) rewrites as \( vr < vvr \) (because \( p = vr \)). Thus, Proposition 6.2(c) (applied to \( a = v \), \( c = rc \) and \( d = vr \)) yields \( rv \leq vr \). Since \( rv \neq vr \) (because otherwise, we would have \( rv = vr \), thus \( vrv = vvr \), contradicting \( vvr < vvr \)), this becomes \( rv < vr \).

Now, \( r \) is a suffix of \( u \) such that \( rv < vr \). Since \( p \) is the minimum-length such suffix, this yields \( \ell (r) \geq \ell (p) \).

But this contradicts the fact that \( \ell (p) = \ell (vr) = \ell (v) + \ell (r) > \ell (r) \). This contradiction proves our assumption wrong; thus, we have shown that \( uv \geq vu \). Proposition 6.19(b) is proven. \( \square \)

**Theorem 6.20.** Let \( w \) be a nonempty word. The following four assertions are equivalent:

- **Assertion A:** The word \( w \) is Lyndon.
- **Assertion B:** Any nonempty words \( u \) and \( v \) satisfying \( w = uv \) satisfy \( v > w \).
- **Assertion C:** Any nonempty words \( u \) and \( v \) satisfying \( w = uv \) satisfy \( v > u \).
- **Assertion D:** Any nonempty words \( u \) and \( v \) satisfying \( w = uv \) satisfy \( vu > w \).
Exercise 6.23. Let \(a\) and every positive integer \(n\) and Radford calls them “primes” and uses \(<\) instead of \(\leq\); also, de Bruijn/Klarner call Lyndon words “normal words”). Chen-Fox-Lyndon [30, §1], Reutenauer [134] and Radford [131] use our definition (but Chen-Fox-Lyndon call the Lyndon words “standard sequences”), and Radford calls them “primes” and uses \(<\) instead of \(\geq\).

Theorem 6.20 appears (with different notations) in Zhou-Lu [171, Proposition 1.4]. The equivalence \(D \iff A\) of our Theorem 6.20 is equivalent to [103, Proposition 5.12] and to [30, \(A'' = A'''\)].

The following exercise provides a different (laborious) approach to Theorem 6.20:

Exercise 6.21. (a) Prove that if \(u \in \mathcal{A}^*\) and \(v \in \mathcal{A}^*\) are two words satisfying \(uv < vu\), then there exists a nonempty suffix \(s\) of \(u\) satisfying \(sv < v\).

(b) Give a new proof of Theorem 6.20 (avoiding the use of Proposition 6.19).

[Hint: For (a), perform strong induction on \(\ell(u) + \ell(v)\), assume the contrary, and distinguish between the case when \(u \leq v\) and the case when \(v\) is a prefix of \(u\). For (b), use part (a) in proving the implication \(D \implies B\), and factor \(v\) as \(v = w^m v'\) with \(m\) maximal in the proof of the implication \(C \implies B\).]

The following two exercises are taken from [66].

Exercise 6.22. Let \(w\) be a nonempty word. Prove that \(w\) is Lyndon if and only if every nonempty word \(t\) and every positive integer \(n\) satisfy (if \(w \leq t^n\), then \(w \leq t\)).

Exercise 6.23. Let \(w_1, w_2, \ldots, w_n\) be \(n\) Lyndon words, where \(n\) is a positive integer. Assume that \(w_1 \leq w_2 \leq \cdots \leq w_n\) and \(w_1 < w_n\). Show that \(w_1 w_2 \cdots w_n\) is a Lyndon word.

The following exercise is a generalization (albeit not in an obvious way) of Exercise 6.23:

Exercise 6.24. Let \(w_1, w_2, \ldots, w_n\) be \(n\) Lyndon words, where \(n\) is a positive integer. Assume that \(w_i w_{i+1} \cdots w_n \geq w_1 w_2 \cdots w_n\) for every \(i \in \{1, 2, \ldots, n\}\). Show that \(w_1 w_2 \cdots w_n\) is a Lyndon word.

We are now ready to meet the one of the most important features of Lyndon words: a bijection between all words and multisets of Lyndon words; it is clear that such a bijection is vital for constructing polynomial generating sets of commutative algebras with bases indexed by words, such as \(\text{QS}ym\) or shuffle algebras. This is given by the Chen-Fox-Lyndon factorization:

Definition 6.25. Let \(w\) be a word. A Chen-Fox-Lyndon factorization (in short, CFL factorization) of \(w\) means a tuple \((a_1, a_2, \ldots, a_k)\) of Lyndon words satisfying \(w = a_1 a_2 \cdots a_k\) and \(a_1 \geq a_2 \geq \cdots \geq a_k\).
Example 6.26. The tuple \((23, 2, 14, 1323, 13, 12, 12, 1)\) is a CFL factorization of the word 2321413323132121 over the alphabet \(\{1, 2, 3, \ldots\}\) (ordered by \(1 < 2 < 3 < \cdots\)), since 23, 2, 14, 1323, 13, 12, 12 and 1 are Lyndon words satisfying 2321413323132121 = 23 · 2 · 14 · 1323 · 13 · 12 · 12 · 1 and 23 ≥ 2 ≥ 14 ≥ 1323 ≥ 13 ≥ 12 ≥ 12 ≥ 1.

The bijection is given by the following Chen-Fox-Lyndon theorem ([68, Theorem 6.5.5], [103, Thm. 5.1.5], [131, part of Thm. 2.1.4]):

**Theorem 6.27.** Let \(w\) be a word. Then, there exists a unique CFL factorization of \(w\).

Before we prove this, we need to state a lemma (which is [103, Proposition 5.1.6]):

**Lemma 6.28.** Let \((a_1, a_2, \ldots, a_k)\) be a CFL factorization of a nonempty word \(w\). Let \(p\) be a nonempty suffix of \(w\). Then, \(p \geq a_k\).

**Proof of Lemma 6.28.** We will prove Lemma 6.28 by induction over the (obviously) positive integer \(k\).

**Induction base:** Assume that \(k = 1\). Thus, \((a_1, a_2, \ldots, a_k) = (a_1)\) is a tuple of Lyndon words satisfying \(w = a_1 a_2 \cdots a_k = a_1\) (since \(k = 1\)), so that \(w\) is a Lyndon word (since \(a_1\) is a Lyndon word). Thus, Corollary 6.15 (applied to \(v = p\)) yields \(p \geq w = a_1 = a_k\) (since \(1 = k\)). Thus, Lemma 6.28 is proven in the case \(k = 1\). The induction base is complete.

**Induction step:** Let \(K\) be a positive integer. Assume (as the induction hypothesis) that Lemma 6.28 holds for \(k = K\). We now need to show that Lemma 6.28 holds for \(k = K + 1\).

So let \((a_1, a_2, \ldots, a_{K+1})\) be a CFL factorization of a nonempty word \(w\). Let \(p\) be a nonempty suffix of \(w\). We need to prove that \(p \geq a_{K+1}\).

By the definition of a CFL factorization, \((a_1, a_2, \ldots, a_{K+1})\) is a tuple of Lyndon words satisfying \(w = a_1 a_2 \cdots a_{K+1}\) and \(a_1 \geq a_2 \geq \cdots \geq a_{K+1}\). Let \(w' = a_2 a_3 \cdots a_{K+1}\); then, \(w = a_1 a_2 \cdots a_{K+1} = (a_2 a_3 \cdots a_{K+1}) = a_1 w'\). Hence, every nonempty suffix of \(w\) is either a nonempty suffix of \(w'\), or has the form \(qw'\) for a nonempty suffix \(q\) of \(a_1\). Since \(p\) is a nonempty suffix of \(w\), we thus must be in one of the following two cases:

**Case 1:** The word \(p\) is a nonempty suffix of \(w'\).

**Case 2:** The word \(p\) has the form \(qw'\) for a nonempty suffix \(q\) of \(a_1\).

Let us first consider Case 1. In this case, \(p\) is a nonempty suffix of \(w'\). The \(K\)-tuple \((a_2, a_3, \ldots, a_{K+1})\) of Lyndon words satisfies \(w' = a_2 a_3 \cdots a_{K+1}\) and \(a_2 \geq a_3 \geq \cdots \geq a_{K+1}\); therefore, \((a_2, a_3, \ldots, a_{K+1})\) is a CFL factorization of \(w'\). We can thus apply Lemma 6.28 to \(K\), \(w'\) and \((a_2, a_3, \ldots, a_{K+1})\) instead of \(k, w\) and \((a_1, a_2, \ldots, a_{K+1})\) (because we assumed that Lemma 6.28 is proven for \(k = K\)). As a result, we obtain that \(p \geq a_{K+1}\). Thus, \(p \geq a_{K+1}\) is proven in Case 1.

Let us now consider Case 2. In this case, \(p\) has the form \(qw'\) for a nonempty suffix \(q\) of \(a_1\). Consider this \(q\). Since \(a_1\) is a Lyndon word, we have \(q \geq a_1\) (by Corollary 6.15, applied to \(a_1\) and \(q\) instead of \(w\) and \(v\)). Thus, \(q \geq a_1 \geq a_2 \geq \cdots \geq a_{K+1}\), so that \(p = qw' \geq q \geq a_{K+1}\). Thus, \(p \geq a_{K+1}\) is proven in Case 2.

We have now proven \(p \geq a_{K+1}\) in all cases. This proves that Lemma 6.28 holds for \(k = K + 1\). The induction step is thus finished, and with it the proof of Lemma 6.28.

**Proof of Theorem 6.27.** Let us first prove that there exists a CFL factorization of \(w\).

Indeed, there clearly exists a tuple \((a_1, a_2, \ldots, a_k)\) of Lyndon words satisfying \(w = a_1 a_2 \cdots a_k\). Fix such a tuple with minimum \(k\). We claim that \(a_1 \geq a_2 \geq \cdots \geq a_k\).

Indeed, if some \(i \in \{1, 2, \ldots, k - 1\}\) would satisfy \(a_i < a_{i+1}\), then the word \(a_ia_{i+1}\) would be Lyndon (by Proposition 6.16(a), applied to \(u = a_i\) and \(v = a_{i+1}\)), whence \((a_1, a_2, \ldots, a_i, a_{i+1}, a_{i+2}, a_{i+3}, \ldots, a_k)\) would also be a tuple of Lyndon words satisfying \(w = a_1 a_2 \cdots a_i (a_ia_{i+1}) a_{i+2} a_{i+3} \cdots a_k\) but having length \(k - 1 < k\), contradicting the fact that \(k\) is the minimum length of such a tuple. Hence, no \(i \in \{1, 2, \ldots, k - 1\}\) can satisfy \(a_i < a_{i+1}\). In other words, every \(i \in \{1, 2, \ldots, k - 1\}\) satisfies \(a_i \geq a_{i+1}\). In other words, \(a_1 \geq a_2 \geq \cdots \geq a_k\). Thus, \((a_1, a_2, \ldots, a_k)\) is a CFL factorization of \(w\), so we have shown that such a CFL factorization exists.

It remains to show that there exists at most one CFL factorization of \(w\). We shall prove this by induction over \(\ell(w)\). Thus, we fix a word \(w\) and assume that

\[
(6.1) \quad \text{for every word } v \text{ with } \ell(v) < \ell(w), \text{ there exists at most one CFL factorization of } v.
\]

For instance, the tuple \((w_1, w_2, \ldots, w_{\ell(w)})\) of one-letter words is a valid example (recall that one-letter words are always Lyndon).
We now have to prove that there exists at most one CFL factorization of \( w \).

Indeed, let \((a_1, a_2, \ldots, a_k)\) and \((b_1, b_2, \ldots, b_m)\) be two CFL factorizations of \( w \). We need to prove that \((a_1, a_2, \ldots, a_k) = (b_1, b_2, \ldots, b_m)\). If \( w \) is empty, then this is obvious, so we WLOG assume that it is not; thus, \( k > 0 \) and \( m > 0 \).

Since \((b_1, b_2, \ldots, b_m)\) is a CFL factorization of \( w \), we have \( w = b_1 b_2 \cdots b_m \), and thus \( b_m \) is a nonempty suffix of \( w \). Thus, Lemma 6.28 (applied to \( p = b_m \)) yields \( b_m \geq a_k \). The same argument (but with the roles of \((a_1, a_2, \ldots, a_k)\) and \((b_1, b_2, \ldots, b_m)\) switched) shows that \( a_k \geq b_m \). Combined with \( b_m \geq a_k \), this yields \( a_k = b_m \). Now let \( v = a_1 a_2 \cdots a_{k-1} \). Then, \((a_1, a_2, \ldots, a_{k-1})\) is a CFL factorization of \( v \) (since \( a_1 \geq a_2 \geq \cdots \geq a_{k-1} \)).

Since \((a_1, a_2, \ldots, a_k)\) is a CFL factorization of \( w \), we have \( w = a_1 a_2 \cdots a_k = a_1 a_2 \cdots a_{k-1} a_k = v b_m \), so that

\[
v b_m = w = b_1 b_2 \cdots b_m = b_1 b_2 \cdots b_{m-1} b_m.
\]

Cancelling \( b_m \) yields \( v = b_1 b_2 \cdots b_{m-1} \). Thus, \((b_1, b_2, \ldots, b_{m-1})\) is a CFL factorization of \( v \) (since \( b_1 \geq b_2 \geq \cdots \geq b_{m-1} \)). Since \( \ell(v) < \ell(w) \) (because \( v = a_1 a_2 \cdots a_{k-1} \) is shorter than \( w = a_1 a_2 \cdots a_k \)), we can apply (6.1) to obtain that there exists at most one CFL factorization of \( v \). But we already know two such CFL factorizations: \((a_1, a_2, \ldots, a_{k-1})\) and \((b_1, b_2, \ldots, b_{m-1})\). Thus, \((a_1, a_2, \ldots, a_{k-1}) = (b_1, b_2, \ldots, b_{m-1})\), which, combined with \( a_k = b_m \), leads to \( (a_1, a_2, \ldots, a_k) = (b_1, b_2, \ldots, b_m) \). This is exactly what we needed to prove.

So we have shown (by induction) that there exists at most one CFL factorization of \( w \). This completes the proof of Theorem 6.27.

The CFL factorization allows us to count all Lyndon words of a given length if \( \mathfrak{A} \) is finite:

**Exercise 6.29.** Assume that the alphabet \( \mathfrak{A} \) is finite. Let \( q = |\mathfrak{A}| \). Show that the number of Lyndon words of length \( n \) equals \( \frac{1}{n} \sum_{d|n} \mu(d) q^{n/d} \) for every positive integer \( n \) (where \( \sum_{d|n} \) means a sum over all positive divisors of \( n \), and where \( \mu \) is the number-theoretic Möbius function).\(^{207}\)

Exercise 6.29 is a well-known result and appears, e.g., in [30, Theorem 1.5].

We will now study another kind of factorization: not of an arbitrary word into Lyndon words, but of a Lyndon word into two smaller Lyndon words. This factorization is called standard factorization ([103, §5.1]) or canonical factorization ([68, Lemma 6.5.33]); we only introduce it from the viewpoint we are interested in, namely its providing a way to do induction over Lyndon words.\(^{208}\) Here is what we need to know:

**Theorem 6.30.** Let \( w \) be a Lyndon word of length \( \geq 1 \). Let \( v \) be the (lexicographically) smallest nonempty proper suffix of \( w \). Since \( v \) is a proper suffix of \( w \), there exists a nonempty \( u \in \mathfrak{A}^* \) such that \( w = uv \). Consider this \( u \). Then:

(a) The words \( u \) and \( v \) are Lyndon.

(b) We have \( u < w < v \).

**Proof.** Every nonempty proper suffix of \( v \) is \( \geq v \) (since every nonempty proper suffix of \( v \) is a nonempty proper suffix of \( w \), but \( v \) is the smallest such suffix) and therefore \( > v \) (since a proper suffix of \( v \) cannot be \( = v \)). Combined with the fact that \( v \) is nonempty, this yields that \( v \) is Lyndon.

Since \( w \) is Lyndon, we know that every nonempty proper suffix of \( w \) is \( > w \). Applied to the nonempty proper suffix \( v \) of \( w \), this yields that \( v > w \). Hence, \( w < v \). Since \( v \) is nonempty, we have \( u < uv = w < v \). This proves Theorem 6.30(b).

Let \( p \) be a nonempty proper suffix of \( u \). Then, \( pv \) is a nonempty proper suffix of \( uv = w \). Thus, \( pv > w \) (since every nonempty proper suffix of \( w \) is \( > w \)). Thus, \( pv > w = uv \), so that \( uv < pv \). Thus, Proposition 6.2(c) (applied to \( a = u, b = v, c = p \) and \( d = v \)) yields that either we have \( u \leq p \) or the word \( p \) is a prefix of \( u \).

Let us assume (for the sake of contradiction) that \( p \leq u \). Then, \( p < u \) (because \( p \) is a proper suffix of \( u \), and therefore \( p \neq u \)). Hence, we cannot have \( u \leq p \). Thus, the word \( p \) is a prefix of \( u \) (since either we have

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\(^{207}\) In particular, \( \frac{1}{n} \sum_{d|n} \mu(d) q^{n/d} \) is an integer.

\(^{208}\) e.g., allowing to solve Exercise 6.24 in a simpler way
Moreover, \(qv\) is nonempty (since \(v\) is nonempty). Hence, \(qv\) is a nonempty proper suffix of \(w\). Since \(v\) is the smallest such suffix, this entails that \(v \leq qv\). Proposition 6.2(b) (applied to \(a = p, c = v\) and \(d = qv\)) thus yields \(pv \leq pqv\). Hence, \(pv \leq pqv = w\), which contradicts \(pv > w\). This contradiction shows that our assumption (that \(p \leq u\)) was false. We thus have \(p > u\).

We now have shown that \(p > u\) whenever \(p\) is a nonempty proper suffix of \(u\). Combined with the fact that \(u\) is nonempty, this shows that \(u\) is a Lyndon word. This completes the proof of Theorem 6.30(a). \(\square\)

Another approach to the standard factorization is given in the following exercise:

**Exercise 6.31.** Let \(w\) be a Lyndon word of length \(\geq 1\). Let \(v\) be the longest proper suffix of \(w\) such that \(v\) is Lyndon.\(^{209}\) Since \(v\) is a proper suffix of \(w\), there exists a nonempty \(u \in \mathcal{A}^*\) such that \(w = uv\). Consider this \(u\). Prove that:

(a) The words \(u\) and \(v\) are Lyndon.
(b) We have \(u < v < v\).
(c) The words \(u\) and \(v\) are precisely the words \(u\) and \(v\) constructed in Theorem 6.30.

Notice that a well-known recursive characterization of Lyndon words \([30, \mathcal{A}' = \mathcal{A}'']\) can be easily derived from Theorem 6.30 and Proposition 6.16(a). We will not dwell on it.

The following exercise surveys some variations on the characterizations of Lyndon words:

**Exercise 6.32.** Let \(w\) be a nonempty word. Consider the following nine assertions:

- **Assertion A’:** The word \(w\) is a power of a Lyndon word.
- **Assertion B’:** If \(u\) and \(v\) are nonempty words satisfying \(w = uv\), then either we have \(v \geq w\) or the word \(v\) is a prefix of \(w\).
- **Assertion C’:** If \(u\) and \(v\) are nonempty words satisfying \(w = uv\), then either we have \(v \geq u\) or the word \(v\) is a prefix of \(u\).
- **Assertion D’:** If \(u\) and \(v\) are nonempty words satisfying \(w = uv\), then we have \(vu \geq uv\).
- **Assertion E’:** If \(u\) and \(v\) are nonempty words satisfying \(w = uv\), then either we have \(v \geq u\) or the word \(v\) is a prefix of \(w\).
- **Assertion F’:** The word \(w\) is a prefix of a Lyndon word in \(\mathcal{A}^*\).
- **Assertion F’’:** Let \(m\) be an object not in the alphabet \(\mathcal{A}\). Let us equip the set \(\mathcal{A} \cup \{m\}\) with a total order which extends the total order on the alphabet \(\mathcal{A}\) and which satisfies \((a < m\) for every \(a \in \mathcal{A}\)). Then, the word \(wm \in (\mathcal{A} \cup \{m\})^*\) (the concatenation of the word \(w\) with the one-letter word \(m\)) is a Lyndon word.
- **Assertion G’:** There exists a Lyndon word \(t \in \mathcal{A}^*\), a positive integer \(\ell\) and a prefix \(p\) of \(t\) (possibly empty) such that \(w = tp\).
- **Assertion H’:** There exists a Lyndon word \(t \in \mathcal{A}^*\), a nonnegative integer \(\ell\) and a prefix \(p\) of \(t\) (possibly empty) such that \(w = tp\).

(a) Prove the equivalence \(A' \iff D'\).
(b) Prove the equivalence \(B' \iff C' \iff E' \iff F' \iff G' \iff H'\).
(c) Prove the implication \(F' \implies B'\).
(d) Prove the implication \(D' \implies B'\). (The implication \(B' \implies D'\) is false, as witnessed by the word \(11211\).)
(e) Prove that if there exists a letter \(\mu \in \mathcal{A}\) such that \((\mu > a\) for every letter \(a\) of \(w\)), then the equivalence \(F' \iff F''\) holds.
(f) Prove that if there exists a letter \(\mu \in \mathcal{A}\) such that \((\mu > a\) for some letter \(a\) of \(w\)), then the equivalence \(F' \iff F''\) holds.

The next exercise (based on work of Hazewinkel [67]) extends some of the above properties of Lyndon words (and words in general) to a more general setting, in which the alphabet \(\mathcal{A}\) is no longer required to be totally ordered, but only needs to be a poset:

\(^{209}\)This is well-defined, because there exists at least one proper suffix \(v\) of \(w\) such that \(v\) is Lyndon. (Indeed, the last letter of \(w\) forms such a suffix, because it is a proper suffix of \(w\) (since \(w\) has length \(\geq 1\)) and is Lyndon (since it is a one-letter word, and since every one-letter word is Lyndon).)
Exercise 6.33. In this exercise, we shall loosen the requirement that the alphabet $\mathcal{A}$ be a totally ordered set; instead, we will only require $\mathcal{A}$ to be a poset. The resulting more general setting will be called the partial-order setting, to distinguish it from the total-order setting in which $\mathcal{A}$ is required to be a totally ordered set. All results in Chapter 6 so far address the total-order setting. In this exercise, we will generalize some of them to the partial-order setting.

All notions that we have defined in the total-order setting (the notion of a word, the relation $\leq$, the notion of a Lyndon word, etc.) are defined in precisely the same way in the partial-order setting. However, the poset $\mathcal{A}^*$ is no longer totally ordered in the partial-order setting.

(a) Prove that Proposition 6.2 holds in the partial-order setting, as long as one replaces “a total order” by “a partial order” in part (a) of this Proposition.

(b) Prove (in the partial-order setting) that if $a, b, c, d \in \mathcal{A}^*$ are four words such that the words $ab$ and $cd$ are comparable (with respect to the partial order $\leq$), then the words $a$ and $c$ are comparable.


(d) Find a counterexample to Exercise 6.22 in the partial-order setting.

(e) Salvage Exercise 6.22 in the partial-order setting (i.e., find a statement which is easily equivalent to this exercise in the total-order setting, yet true in the partial-order setting).

(f) In the partial-order setting, a Hazewinkel-CFL factorization of a word $w$ will mean a tuple $(a_1, a_2, \ldots, a_k)$ of Lyndon words such that $w = a_1a_2\cdots a_k$ and such that no $i \in \{1, 2, \ldots, k-1\}$ satisfies $a_i < a_{i+1}$.

Prove that every word $w$ has a unique Hazewinkel-CFL factorization (in the partial-order setting).

(g) Prove that Exercise 6.32 still holds in the partial-order setting.

The reader is invited to try extending other results to the partial-order setting (it seems that no research has been done on this except for Hazewinkel’s [67]). We shall now, however, return to the total-order setting (which has the most known applications).

Lyndon words are related to various other objects in mathematics, such as free Lie algebras ([134]), shuffles and shuffle algebras (Sections 6.2 and 6.3 below), QSym (Sections 6.4 and 6.5), Markov chains on combinatorial Hopf algebras ([40]), de Bruijn sequences ([52], [118]), symmetric functions (specifically, the transition matrices between the bases $(h_\lambda)_{\lambda \in \operatorname{Par}}$, $(e_\lambda)_{\lambda \in \operatorname{Par}}$ and $(m_\lambda)_{\lambda \in \operatorname{Par}}$; see [83] for this), and the Burrows-Wheeler algorithm for data compression ([82]). They are also connected to necklaces (in the combinatorial sense) – a combinatorial object that also happens to be related to a lot of algebra ([137, Chapter 5], [38]).

We only survey the basics of this latter classical connection in an exercise:

Exercise 6.34. Let $\mathcal{A}$ be any set (not necessarily totally ordered). Let $C$ denote the infinite cyclic group, written multiplicatively. Fix a generator $c$ of $C$.

211 Fix a positive integer $n$. The group $C$ acts on $\mathcal{A}^n$ from the left according to the rule

$$c \cdot (a_1, a_2, \ldots, a_n) = (a_2, a_3, \ldots, a_n, a_1) \quad \text{for all } (a_1, a_2, \ldots, a_n) \in \mathcal{A}^n.$$  

212 The orbits of this $C$-action will be called $n$-necklaces; they form a set partition of the set $\mathcal{A}^n$.

The $n$-necklace containing a given $n$-tuple $w \in \mathcal{A}^n$ will be denoted by $[w]$.

(a) Prove that every $n$-necklace $N$ is a finite nonempty set and satisfies $|N| | n$. (Recall that $N$ is an orbit, thus a set; as usual, $|N|$ denotes the cardinality of this set.)

The period of an $n$-necklace $N$ is defined as the positive integer $|N|$. (This $|N|$ is indeed a positive integer, since $N$ is a finite nonempty set.)

An $n$-necklace is said to be aperiodic if its period is $n$.

210 This result, as well as the validity of Proposition 6.16 in the partial-order setting, are due to Hazewinkel [67].

211 So $C$ is a group isomorphic to $(\mathbb{Z}, +)$, and the isomorphism $(\mathbb{Z}, +) \to C$ sends every $n \in \mathbb{Z}$ to $c^n$. (Recall that we write the binary operation of $C$ as · instead of $+$.)

212 In other words, $c$ rotates any $n$-tuple of elements of $\mathcal{A}$ cyclically to the left. Thus, $c^n \in C$ acts trivially on $\mathcal{A}^n$, and so this action of $C$ on $\mathcal{A}^n$ factors through $C / \langle c^n \rangle$ (a cyclic group of order $n$).

213 Classically, one visualizes them as necklaces of $n$ beads of $|\mathcal{A}|$ colors. (The colors are the elements of $\mathcal{A}$.) The intuition behind this is that a necklace is an object that doesn’t really change when we rotate it in its plane. However, to make this intuition match the definition, we need to think of a necklace as being stuck in its (fixed) plane, so that we cannot lift it up and turn it around, dropping it back to its plane in a reflected state.
(b) Given any $n$-tuple $w = (w_1, w_2, \ldots, w_n) \in \mathfrak{A}^n$, prove that the $n$-necklace $[w]$ is aperiodic if and only if every $k \in \{1, 2, \ldots, n-1\}$ satisfies $(w_{k+1}, w_{k+2}, \ldots, w_n, w_1, w_2, \ldots, w_k) \neq w$.

From now on, we assume that the set $\mathfrak{A}$ is totally ordered. We use $\mathfrak{A}$ as our alphabet to define the notions of words, the lexicographic order, and Lyndon words. All notations that we introduced for words will thus be used for elements of $\mathfrak{A}^n$.

(c) Prove that every aperiodic $n$-necklace contains exactly one Lyndon word.

(d) If $N$ is an $n$-necklace which is not aperiodic, then prove that $N$ contains no Lyndon word.

(e) Show that the aperiodic $n$-necklaces are in bijection with Lyndon words of length $n$.

From now on, we assume that the set $\mathfrak{A}$ is finite. Define the number-theoretic Möbius function $\mu$ and the Euler totient function $\phi$ as in Exercise 2.82.

(f) Prove that the number of all aperiodic $n$-necklaces is

$$\frac{1}{n} \sum_{d|n} \mu(d) |\mathfrak{A}|^{n/d}.$$ 

(g) Prove that the number of all $n$-necklaces is

$$\frac{1}{n} \sum_{d|n} \phi(d) |\mathfrak{A}|^{n/d}.$$ 

(h) Solve Exercise 6.29 again.

(i) Forget that we fixed $\mathfrak{A}$. Show that every $q \in \mathbb{Z}$ satisfies $n \mid \sum_{d|n} \mu(d) q^{n/d}$ and $n \mid \sum_{d|n} \phi(d) q^{n/d}$.

[Hint: For (c), use Theorem 6.20. For (i), either use parts (f) and (g) and a trick to extend to $q$ negative; or recall Exercise 2.84.]

### 6.2. Shuffles and Lyndon words

We will now connect the theory of Lyndon words with the notion of shuffle products. We have already introduced the latter notion in Definition 1.68, but we will now study it more closely and introduce some more convenient notations (e.g., we will need a notation for single shuffles, not just the whole multiset).

**Definition 6.35.**  
(a) Let $n \in \mathbb{N}$ and $m \in \mathbb{N}$. Then, $\text{Sh}_{n,m}$ denotes the subset

$$\{ \sigma \in \mathfrak{S}_{n+m} : \sigma^{-1}(1) < \sigma^{-1}(2) < \cdots < \sigma^{-1}(n) ; \sigma^{-1}(n+1) < \cdots < \sigma^{-1}(n+m) \}$$

of the symmetric group $\mathfrak{S}_{n+m}$.

(b) Let $u = (u_1, u_2, \ldots, u_n)$ and $v = (v_1, v_2, \ldots, v_m)$ be two words. If $\sigma \in \text{Sh}_{n,m}$, then, $u \uparrow_{\sigma} v$ will denote the word $(w_{\sigma(1)}, w_{\sigma(2)}, \ldots, w_{\sigma(n+m)})$, where $(w_1, w_2, \ldots, w_{n+m})$ is the concatenation $u \cdot v = (u_1, u_2, \ldots, u_n, v_1, v_2, \ldots, v_m)$. We notice that the multiset of all letters of $u \uparrow_{\sigma} v$ is the disjoint union of the multiset of all letters of $u$ with the multiset of all letters of $v$. As a consequence, $\ell(u \uparrow_{\sigma} v) = \ell(u) + \ell(v)$.

(c) Let $u = (u_1, u_2, \ldots, u_n)$ and $v = (v_1, v_2, \ldots, v_m)$ be two words. The **multiset of shuffles of $u$ and $v$** is defined as the multiset $\{ (w_{\sigma(1)}, w_{\sigma(2)}, \ldots, w_{\sigma(n+m)}) : \sigma \in \text{Sh}_{n,m} \}_{\text{multiset}}$, where $(w_1, w_2, \ldots, w_{n+m})$ is the concatenation $u \cdot v = (u_1, u_2, \ldots, u_n, v_1, v_2, \ldots, v_m)$. In other words, the multiset of shuffles of $u$ and $v$ is the multiset

$$\{ u \uparrow_{\sigma} v : \sigma \in \text{Sh}_{n,m} \}_{\text{multiset}}.$$ 

It is denoted by $u \uparrow v$.

The next fact provides the main connection between Lyndon words and shuffles:

**Theorem 6.36.** Let $u$ and $v$ be two words.

Let $(a_1, a_2, \ldots, a_p)$ be the CFL factorization of $u$. Let $(b_1, b_2, \ldots, b_q)$ be the CFL factorization of $v$.

---

214Parts (a) and (c) of the below Definition 6.35 define notions which have already been introduced in Definition 1.68. Of course, the definitions of these notions are equivalent; however, the variables are differently labelled in the two definitions (for example, the variables $u$, $v$, $w$ and $\sigma$ of Definition 6.35(c) correspond to the variables $a$, $b$, $c$ and $w$ of Definition 1.68). The labels in Definition 6.35 have been chosen to match with the rest of Section 6.2.
(a) Let \((c_1, c_2, \ldots, c_{p+q})\) be the result of sorting the list \((a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q)\) in decreasing order. Then, the lexicographically highest element of the multiset \(w \sqcup v\) is \(c_1c_2\cdots c_{p+q}\) (and \((c_1, c_2, \ldots, c_{p+q})\) is the CFL factorization of this element).

(b) Let \(\mathcal{L}\) denote the set of all Lyndon words. If \(w\) is a Lyndon word and \(z\) is any word, let \(\text{mult}_w z\) denote the number of terms in the CFL factorization of \(z\) which are equal to \(w\). The multiplicity with which the lexicographically highest element of the multiset \(w \sqcup v\) appears in the multiset \(w \sqcup v\) is \(\prod_{w \in \mathcal{L}} \left(\frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u}\right)\). (This product is well-defined because almost all of its factors are 1.)

(c) If \(a_i \geq b_j\) for every \(i \in \{1, 2, \ldots, p\}\) and \(j \in \{1, 2, \ldots, q\}\), then the lexicographically highest element of the multiset \(w \sqcup v\) is \(uv\).

(d) If \(a_i > b_j\) for every \(i \in \{1, 2, \ldots, p\}\) and \(j \in \{1, 2, \ldots, q\}\), then the multiplicity with which the word \(uv\) appears in the multiset \(w \sqcup v\) is 1.

(e) Assume that \(w\) is a Lyndon word. Also, assume that \(u \geq v\) for every \(v \in \{1, 2, \ldots, q\}\). Then, the lexicographically highest element of the multiset \(w \sqcup v\) is \(uv\), and the multiplicity with which this word \(uv\) appears in the multiset \(w \sqcup v\) is \(\text{mult}_w v + 1\).

Example 6.37. For this example, let \(u\) and \(v\) be the words \(u = 23232\) and \(v = 323221\) over the alphabet \(\mathcal{A} = \{1, 2, 3, \ldots\}\) with total order given by \(1 < 2 < 3 < \cdots\). The CFL factorizations of \(u\) and \(v\) are \((23, 23, 2)\) and \((3, 23, 2, 2, 1)\), respectively. Thus, using the notations of Theorem 6.36, we have \(p = 3\), \((a_1, a_2, \ldots, a_p) = (23, 23, 2)\), \(q = 5\) and \((b_1, b_2, \ldots, b_q) = (3, 23, 2, 2, 1)\). Thus, Theorem 6.36(a) predicts that the lexicographically highest element of the multiset \(w \sqcup v\) equals \(23232323232322\), where \(c_1, c_2, c_3, c_4, c_5, c_6, c_7, c_8\) are the words \(23, 23, 2, 3, 23, 2, 2, 1\) listed in decreasing order (in other words, \((c_1, c_2, c_3, c_4, c_5, c_6, c_7, c_8) = (3, 23, 23, 2, 2, 2, 1)\)). In other words, Theorem 6.36(a) predicts that the lexicographically highest element of the multiset \(w \sqcup v\) is \(23232322\). We could verify this by brute force, but this would be laborious since the multiset \(w \sqcup v\) has \(\binom{5+6}{5} = 462\) elements (with multiplicities). Theorem 6.36(b) predicts that this lexicographically highest element \(23232322\) appears in the multiset \(w \sqcup v\) with a multiplicity of \(\prod_{w \in \mathcal{L}} \left(\frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u}\right)\). This product \(\prod_{w \in \mathcal{L}} \left(\frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u}\right)\) is infinite, but all but finitely many of its factors are 1 and therefore can be omitted; the only factors which are not 1 are those corresponding to Lyndon words \(w\) which appear both in the CFL factorization of \(u\) and in the CFL factorization of \(v\) (since for any other factor, at least one of the numbers \(\text{mult}_w u\) or \(\text{mult}_w v\) equals 0, and therefore the binomial coefficient \(\left(\frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u}\right)\) equals 1). Thus, in order to compute the product \(\prod_{w \in \mathcal{L}} \left(\frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u}\right)\), we only need to multiply these factors. In our example, these are the factors for \(w = 23\) and for \(w = 2\) (these are the only Lyndon words which appear both in the CFL factorization \((23, 23, 2)\) of \(u\) and in the CFL factorization \((3, 23, 2, 2, 1)\) of \(v\)). So we have

\[
\prod_{w \in \mathcal{L}} \left(\frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u}\right) = \left(\frac{\text{mult}_{23} u + \text{mult}_{23} v}{\text{mult}_{23} u}\right) = 3 \cdot 3 = 9.
\]

The word \(23232322\) must thus appear in the multiset \(w \sqcup v\) with a multiplicity of 9. This, too, could be checked by brute force.

Theorem 6.36 (and Theorem 6.56 further below, which describes more precisely how the lexicographically highest element of \(w \sqcup v\) emerges by shuffling \(u\) and \(v\)) is fairly close to [131, Theorem 2.2.2] (and will be used for the same purposes), the main difference being that we are talking about the shuffle product of two (not necessarily Lyndon) words, while Radford (and most other authors) study the shuffle product of many Lyndon words.

In order to prove Theorem 6.36, we will need to make some stronger statements, for which we first have to introduce some more notation:

\[215\text{with respect to the total order on } \mathcal{A}^+ \text{ whose greater-or-equal relation is } \geq\]
Definition 6.38.  
(a) If $p$ and $q$ are two integers, then $[p : q]^+$ denotes the interval $\{p + 1, p + 2, \ldots, q\}$ of $\mathbb{Z}$. Note that $[p : q]^+ = q - p$ if $q \geq p$.

(b) If $I$ and $J$ are two nonempty intervals of $\mathbb{Z}$, then we say that $I < J$ if and only if every $i \in I$ and $j \in J$ satisfy $i < j$. This defines a partial order on the set of nonempty intervals of $\mathbb{Z}$. (Roughly speaking, $I < J$ if the interval $I$ ends before $J$ begins.)

(c) If $w$ is a word with $n$ letters (for some $n \in \mathbb{N}$), and $I$ is an interval of $\mathbb{Z}$ such that $I \subset [0 : n]^+$, then $w[I]$ will denote the word $(w_{p+1}, w_{p+2}, \ldots, w_q)$, where $I$ is written in the form $I = [p : q]^+$ with $q \geq p$. Obviously, $\ell(w[I]) = |I| = q - p$. A word of the form $w[I]$ for an interval $I \subset [0 : n]^+$ (equivalently, a word which is a prefix of a suffix of $w$) is called a factor of $w$.

(d) Let $\alpha$ be a composition. Then, we define a tuple $\text{intsys} \alpha$ of intervals of $\mathbb{Z}$ as follows: Write $\alpha$ in the form $(\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ (so that $\ell = \ell(\alpha)$). Then, set $\text{intsys} \alpha = (I_1, I_2, \ldots, I_\ell)$, where

$$I_i = \left[ \sum_{k=1}^{i-1} \alpha_k : \sum_{k=1}^{i} \alpha_k \right]^+ \quad \text{for every } i \in \{1, 2, \ldots, \ell\}.$$

This $\ell$-tuple $\text{intsys} \alpha$ is a tuple of nonempty intervals of $\mathbb{Z}$. This tuple $\text{intsys} \alpha$ is called the interval system corresponding to $\alpha$. (This is precisely the $\ell$-tuple $(I_1, I_2, \ldots, I_\ell)$ constructed in Definition 4.24.) The length of the tuple $\text{intsys} \alpha$ is $\ell(\alpha)$.

Example 6.39.  
(a) We have $[2 : 4]^+ = \{3, 4\}$ and $[3 : 3]^+ = \emptyset$.


(c) If $w$ is the word 915352, then $w\left[[0 : 3]^+\right] = (w_1, w_2, w_3) = 915$ and $w\left[[2 : 4]^+\right] = (w_3, w_4) = 53$.

(d) If $\alpha$ is the composition $(4, 1, 4, 2, 3)$, then the interval system corresponding to $\alpha$ is

$$\text{intsys} \alpha = \left(\left[0 : 4\right]^+, \left[4 : 5\right]^+, \left[5 : 9\right]^+, \left[9 : 11\right]^+, \left[11 : 14\right]^+\right) = \{(1, 2, 3, 4), \{5\}, \{6, 7, 8, 9\}, \{10, 11\}, \{12, 13, 14\}\}.$$

The following properties of the notions introduced in the preceding definition are easy to check:

Remark 6.40.  
(a) If $I$ and $J$ are two nonempty intervals of $\mathbb{Z}$ satisfying $I < J$, then $I$ and $J$ are disjoint.

(b) If $I$ and $J$ are two disjoint nonempty intervals of $\mathbb{Z}$, then either $I < J$ or $J < I$.

(c) Let $\alpha$ be a composition. Write $\alpha$ in the form $(\alpha_1, \alpha_2, \ldots, \alpha_\ell)$ (so that $\ell = \ell(\alpha)$). The interval system $\text{intsys} \alpha$ can be described as the unique $\ell$-tuple $(I_1, I_2, \ldots, I_\ell)$ of nonempty intervals of $\mathbb{Z}$ satisfying the following three properties:

- The intervals $I_1, I_2, \ldots, I_\ell$ form a set partition of the set $[0 : n]^+$, where $n = |\alpha|$.
- We have $I_1 < I_2 < \cdots < I_\ell$.
- We have $|I_i| = \alpha_i$ for every $i \in \{1, 2, \ldots, \ell\}$.


The following two lemmas are collections of more or less trivial consequences of what it means to be an element of $\text{Sh}_{n,m}$ and what it means to be a shuffle:

Lemma 6.42. Let $n \in \mathbb{N}$ and $m \in \mathbb{N}$. Let $\sigma \in \text{Sh}_{n,m}$.

(a) If $I$ is an interval of $\mathbb{Z}$ such that $I \subset [0 : n + m]^+$, then $\sigma(I) \cap [0 : n]^+$ and $\sigma(I) \cap [n : n + m]^+$ are intervals.

(b) Let $K$ and $L$ be nonempty intervals of $\mathbb{Z}$ such that $K \subset [0 : n]^+$ and $L \subset [0 : n]^+$ and $K < L$ and such that $K \cup L$ is an interval. Assume that $\sigma^{-1}(K)$ and $\sigma^{-1}(L)$ are intervals, but $\sigma^{-1}(K) \cup \sigma^{-1}(L)$ is not an interval. Then, there exists a nonempty interval $P \subset [n : n + m]^+$ such that $\sigma^{-1}(P)$, $\sigma^{-1}(K) \cup \sigma^{-1}(P)$ and $\sigma^{-1}(P) \cup \sigma^{-1}(L)$ are intervals and such that $\sigma^{-1}(K) < \sigma^{-1}(P) < \sigma^{-1}(L)$.

(c) Lemma 6.42(b) remains valid if “$K \subset [0 : n]^+$ and $L \subset [0 : n]^+$” and “$P \subset [n : n + m]^+$” are replaced by “$K \subset [n : n + m]^+$ and $L \subset [n : n + m]^+$” and “$P \subset [0 : n]^+$”, respectively.

Exercise 6.43. Prove Lemma 6.42.

Lemma 6.44. Let $u$ and $v$ be two words. Let $n = \ell(u)$ and $m = \ell(v)$. Let $\sigma \in \text{Sh}_{n,m}$.
(a) If $I$ is an interval of $\mathbb{Z}$ satisfying either $I \subset [0 : n]^+$ or $I \subset [n : n + m]^+$, and if $\sigma^{-1}(I)$ is an interval, then

\begin{equation}
\left(u \cup_{\sigma} v\right) \left[\sigma^{-1}(I)\right] = (uv) [I].
\end{equation}

(b) Assume that $u \cup_{\sigma} v$ is the lexicographically highest element of the multiset $u \cup_{\tau} v$. Let $I \subset [0 : n]^+$ and $J \subset [n : n + m]^+$ be two nonempty intervals. Assume that $\sigma^{-1}(I)$ and $\sigma^{-1}(J)$ are also intervals, that $\sigma^{-1}(I) < \sigma^{-1}(J)$, and that $\sigma^{-1}(I) \cup \sigma^{-1}(J)$ is an interval as well. Then, $(uv)[I] \cdot (uv)[J] \geq (uv)[J : (uv)[I]].$

(c) Lemma 6.44(b) remains valid if “$I \subset [0 : n]^+$ and $J \subset [n : n + m]^+$” is replaced by “$I \subset [n : n + m]^+$ and $J \subset [0 : n]^+$.”

**Exercise 6.45.** Prove Lemma 6.44.

[Hint: For (b), show that there exists a $\tau \in \text{Sh}_{n,m}$ such that $u \cup_{\tau} v$ differs from $u \cup_{\sigma} v$ only in the order of the subwords $(uv)[I]$ and $(uv)[J].$]

We are still a few steps away from stating our results in a way that allows comfortably proving Theorem 6.36. For the latter aim, we introduce the notion of \emph{alpha-clumping permutations}, and characterize them in two ways:

**Definition 6.46.** Let $n \in \mathbb{N}$. Let $\alpha$ be a composition of $n$. Let $\ell = \ell(\alpha)$.

(a) For every set $S$ of positive integers, let $\overrightarrow{S}$ denote the list of all elements of $S$ in increasing order (with each element appearing exactly once). Notice that this list $\overrightarrow{S}$ is a word over the set of positive integers.

(b) For every $\tau \in \mathcal{S}_\ell$, we define a permutation $\text{iper} (\alpha, \tau) \in \mathcal{S}_n$ as follows:

The interval system corresponding to $\alpha$ is an $\ell$-tuple of intervals (since $\ell(\alpha) = \ell$); denote this $\ell$-tuple by $(I_1, I_2, \ldots, I_\ell)$. Now, define $\text{iper}(\alpha, \tau)$ to be the permutation in $\mathcal{S}_n$ which (in one-line notation) is the word $I_{\tau(1)} I_{\tau(2)} \cdots I_{\tau(\ell)}$ (a concatenation of $\ell$ words). This is well-defined\footnote{In fact, from the properties of interval systems, we know that the intervals $I_1, I_2, \ldots, I_\ell$ form a set partition of the set $[0 : n]^+$. Hence, the intervals $I_{\tau(1)}, I_{\tau(2)}, \ldots, I_{\tau(\ell)}$ form a set partition of the set $[0 : n]^+$. As a consequence, the word $I_{\tau(1)} I_{\tau(2)} \cdots I_{\tau(\ell)}$ is a permutation of the word $12\ldots n$, and so there exists a permutation in $\mathcal{S}_n$ (in one-line notation) this word, qed.}, hence, $\text{iper}(\alpha, \tau) \in \mathcal{S}_n$ is defined.

(c) The interval system corresponding to $\alpha$ is an $\ell$-tuple of intervals (since $\ell(\alpha) = \ell$); denote this $\ell$-tuple by $(I_1, I_2, \ldots, I_\ell)$.

A permutation $\sigma \in \mathcal{S}_n$ is said to be \emph{alpha-clumping} if every $i \in \{1, 2, \ldots, \ell\}$ has the two properties that:

- the set $\sigma^{-1}(I_i)$ is an interval;
- the restriction of the map $\sigma^{-1}$ to the interval $I_i$ is increasing.

**Example 6.47.** For this example, let $n = 7$ and $\alpha = (2, 1, 3, 1)$. Then, $\ell(\alpha) = 4$ and $(I_1, I_2, I_3, I_4) = ((1, 2), (3), (4, 5, 6), (7))$ (where we are using the notations of Definition 6.46). Hence, $I_1 = 12$, $I_2 = 3$, $I_3 = 456$ and $I_4 = 7$.

(a) If $\tau \in \mathcal{S}_\ell = \mathcal{S}_4$ is the permutation $(2, 3, 1, 4)$, then $\text{iper}(\alpha, \tau)$ is the permutation in $\mathcal{S}_7$ which (in one-line notation) is the word $I_{\tau(1)} I_{\tau(2)} I_{\tau(3)} I_{\tau(4)} = I_3 I_4 I_6 1 I_2 = 3456127$.

If $\tau \in \mathcal{S}_\ell = \mathcal{S}_4$ is the permutation $(3, 1, 4, 2)$, then $\text{iper}(\alpha, \tau)$ is the permutation in $\mathcal{S}_7$ which (in one-line notation) is the word $I_{\tau(1)} I_{\tau(2)} I_{\tau(3)} I_{\tau(4)} = I_3 I_4 I_2 I_6 = 4561237$.

(b) The permutation $\sigma = (3, 7, 4, 5, 6, 1, 2) \in \mathcal{S}_7$ (given here in one-line notation) is $\alpha$-clumping, because:

- every $i \in \{1, 2, \ldots, \ell\} = \{1, 2, 3, 4\}$ has the property that $\sigma^{-1}(I_i)$ is an interval (namely, $\sigma^{-1}(I_1) = \sigma^{-1}((1, 2)) = (6, 7)$, $\sigma^{-1}(I_2) = \sigma^{-1}((3)) = (1)$, $\sigma^{-1}(I_3) = \sigma^{-1}((4, 5, 6)) = (3, 4, 5)$ and $\sigma^{-1}(I_4) = \sigma^{-1}((7)) = (2)$), and
- the restrictions of the map $\sigma^{-1}$ to the intervals $I_i$ are increasing (this means that $\sigma^{-1}(1) < \sigma^{-1}(2)$ and $\sigma^{-1}(4) < \sigma^{-1}(5) < \sigma^{-1}(6)$, since the one-element intervals $I_2$ and $I_4$ do not contribute anything to this condition).
Here is a more or less trivial observation:

**Proposition 6.48.** Let \( n \in \mathbb{N} \). Let \( \alpha \) be a composition of \( n \). Let \( \ell = \ell (\alpha) \). Write \( \alpha \) in the form \((a_1, a_2, \ldots, a_p)\). The interval system corresponding to \( \alpha \) is an \( \ell \)-tuple of intervals (since \( \ell (\alpha) = \ell \)); denote this \( \ell \)-tuple by \((I_1, I_2, \ldots, I_\ell)\). Let \( \tau \in \mathcal{S}_\ell \). Set \( \sigma = \text{iper} (\alpha) \).

(a) We have \( \sigma^{-1} \left( I_{\tau(j)} \right) = \left[ \sum_{k=1}^{j-1} \alpha (\tau(k)) + \sum_{k=1}^{j} \alpha (\tau(k)) \right]^+ \) for every \( j \in \{1, 2, \ldots, \ell \} \).

(b) For every \( j \in \{1, 2, \ldots, \ell \} \), the restriction of the map \( \sigma^{-1} \) to the interval \( I_{\tau(j)} \) is increasing.

(c) The permutation \( \text{iper}(\alpha, \tau) \) is \( \alpha \)-clumping.

(d) Let \( i \in \{1, 2, \ldots, \ell - 1\} \). Then, the sets \( \sigma^{-1} \left( I_{\tau(i)} \right) \), \( \sigma^{-1} \left( I_{\tau(i+1)} \right) \), and \( \sigma^{-1} \left( I_{\tau(i)} \right) \cup \sigma^{-1} \left( I_{\tau(i+1)} \right) \) are nonempty intervals. Also, \( \sigma^{-1} \left( I_{\tau(i)} \right) < \sigma^{-1} \left( I_{\tau(i+1)} \right) \).

**Exercise 6.49.** Prove Proposition 6.48.

**Proposition 6.50.** Let \( n \in \mathbb{N} \). Let \( \alpha \) be a composition of \( n \). Let \( \ell = \ell (\alpha) \).

(a) Define a map

\[
\text{iper}_\alpha : \mathcal{S}_\ell \rightarrow \{ \omega \in \mathcal{S}_n | \omega \text{ is } \alpha\text{-clumping} \} ,
\]

\[
\tau \mapsto \text{iper}(\alpha, \tau)
\]

217 This map is well-defined because for every \( \tau \in \mathcal{S}_\ell \), the permutation \( \text{iper}(\alpha, \tau) \) is \( \alpha \)-clumping (according to Proposition 6.48(a)).

(b) Let \( \sigma \in \mathcal{S}_n \) be an \( \alpha \)-clumping permutation. Then, there exists a unique \( \tau \in \mathcal{S}_\ell \) satisfying \( \sigma = \text{iper}(\alpha, \tau) \).

**Exercise 6.51.** Prove Proposition 6.50.

Next, we recall that the concatenation \( \alpha \cdot \beta \) of two compositions \( \alpha \) and \( \beta \) is defined in the same way as the concatenation of two words; if we regard compositions as words over the alphabet \( \{1, 2, 3, \ldots\} \), then the concatenation \( \alpha \cdot \beta \) of two compositions \( \alpha \) and \( \beta \) is the concatenation \( \alpha \beta \) of the words \( \alpha \) and \( \beta \). Thus, we are going to write \( \alpha \beta \) for the concatenation \( \alpha \cdot \beta \) of two compositions \( \alpha \) and \( \beta \) from now on.

**Proposition 6.52.** Let \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \). Let \( \alpha \) be a composition of \( n \), and \( \beta \) be a composition of \( m \). Let \( p = \ell (\alpha) \) and \( q = \ell (\beta) \). Let \( \tau \in \mathcal{S}_{p+q} \). Notice that \( \text{iper}(\alpha \beta, \tau) \in \mathcal{S}_{n+m} \) (since \( \alpha \beta \) is a composition of \( n + m \) having length \( \ell (\alpha \beta) = \ell (\alpha) + \ell (\beta) = p + q \)). Then, \( \tau \in \text{Sh}_{p,q} \) if and only if \( \text{iper}(\alpha \beta, \tau) \in \text{Sh}_{n,m} \).

**Exercise 6.53.** Prove Proposition 6.52.

Here is one more simple fact:

**Lemma 6.54.** Let \( u \) and \( v \) be two words. Let \( n = \ell (u) \) and \( m = \ell (v) \). Let \( \alpha \) be a composition of \( n \), and let \( \beta \) be a composition of \( m \). Let \( p = \ell (\alpha) \) and \( q = \ell (\beta) \). The concatenation \( \alpha \beta \) is a composition of \( n + m \) having length \( \ell (\alpha \beta) = \ell (\alpha) + \ell (\beta) = p + q \). Thus, the interval system corresponding to \( \alpha \beta \) is a \( (p + q) \)-tuple of intervals which covers \([0 : n + m]^+\). Denote this \((p + q)\)-tuple by \((I_1, I_2, \ldots, I_{p+q})\).

Let \( \tau \in \text{Sh}_{p,q} \). Set \( \sigma = \text{iper}(\alpha \beta, \tau) \). Then,

\[
uu \tau = (uv) \left[ I_{\tau(1)} \right] \cdot (uv) \left[ I_{\tau(2)} \right] \cdots \cdot (uv) \left[ I_{\tau(p+q)} \right].
\]

**Exercise 6.55.** Prove Lemma 6.54.

Having these notations and trivialities in place, we can say a bit more about the lexicographically highest element of a shuffle product than what was said in Theorem 6.36:

**Theorem 6.56.** Let \( u \) and \( v \) be two words. Let \( n = \ell (u) \) and \( m = \ell (v) \).

Let \((a_1, a_2, \ldots, a_p)\) be the CFL factorization of \( u \). Let \((b_1, b_2, \ldots, b_q)\) be the CFL factorization of \( v \).

Let \( \alpha \) be the \( p \)-tuple \((\ell (a_1), \ell (a_2), \ldots, \ell (a_p)) \). Then, \( \alpha \) is a composition\(^{218} \) of length \( p \) and size \( \sum_{k=1}^{p} \ell (a_k) = \ell (u) = n \).
Let \( \beta \) be the \( q \)-tuple \( (\ell(b_1), \ell(b_2), \ldots, \ell(b_q)) \). Then, \( \beta \) is a composition of length \( q \) and size \( \sum_{k=1}^{q} \ell(b_k) = m \).

Now, \( \alpha \) is a composition of length \( p \) and size \( n \), and \( \beta \) is a composition of length \( q \) and size \( m \). Thus, the concatenation \( \alpha \beta \) of these two tuples is a composition of length \( p + q \) and size \( n + m \). The interval system corresponding to this composition \( \alpha \beta \) is a \((p + q)\)-tuple (since said composition has length \( p + q \)); denote this \((p + q)\)-tuple by \((I_1, I_2, \ldots, I_{p+q})\).

(a) If \( \tau \in Sh_{p,q} \) satisfies \( (uv)[I_{\tau(1)}] \geq (uv)[I_{\tau(2)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}] \), and if we set \( \sigma = \text{iper}(\alpha \beta, \tau) \), then \( \sigma \in Sh_{n,m} \) and the word \( uv \uplus v \) is the lexicographically highest element of the multiset \( u \uplus v \).

(b) Let \( \sigma \in Sh_{n,m} \) be a permutation such that \( u \uplus v \) is the lexicographically highest element of the multiset \( u \uplus v \). Then, there exists a unique permutation \( \tau \in Sh_{p,q} \) satisfying \( (uv)[I_{\tau(1)}] \geq (uv)[I_{\tau(2)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}] \) and \( \sigma = \text{iper}(\alpha \beta, \tau) \).

\[ (I_1, I_2, \ldots, I_{p+q}) = \text{intsys}(\alpha \beta). \]

But since \( \alpha = (\ell(a_1), \ell(a_2), \ldots, \ell(a_p)) \) and \( \beta = (\ell(b_1), \ell(b_2), \ldots, \ell(b_q)) \), we have

[\[ \alpha \beta = (\ell(a_1), \ell(a_2), \ldots, \ell(a_p), \ell(b_1), \ell(b_2), \ldots, \ell(b_q)). \]

Thus, \((I_1, I_2, \ldots, I_{p+q})\) rewrites as

\[ (I_1, I_2, \ldots, I_{p+q}) = \text{intsys}(\ell(a_1), \ell(a_2), \ldots, \ell(a_p), \ell(b_1), \ell(b_2), \ldots, \ell(b_q)). \]

By the definition of \( \text{intsys}(\ell(a_1), \ell(a_2), \ldots, \ell(a_p), \ell(b_1), \ell(b_2), \ldots, \ell(b_q)) \), we thus have

\[ I_i = \left[ \sum_{k=1}^{i-1} \ell(a_k) : \sum_{k=1}^{i} \ell(a_k) \right]^+ \]

for every \( i \in \{1, 2, \ldots, p\} \), and besides

\[ I_{p+j} = \left[ n + \sum_{k=1}^{j} \ell(b_k) : n + \sum_{k=1}^{j+1} \ell(b_k) \right]^+ \]

for every \( j \in \{1, 2, \ldots, q\} \).

Moreover, Remark 6.40(c) shows that \((I_1, I_2, \ldots, I_{p+q})\) is a \((p + q)\)-tuple of nonempty intervals of \( Z \) and satisfies the following three properties:

- The intervals \( I_1, I_2, \ldots, I_{p+q} \) form a set partition of the set \([0 : n + m]^+\).
- We have \( I_1 < I_2 < \cdots < I_{p+q} \).
- We have \(|I_i| = \ell(a_i)\) for every \( i \in \{1, 2, \ldots, p\} \) and \(|I_{p+j}| = \ell(b_j)\) for every \( j \in \{1, 2, \ldots, q\} \).

Of course, every \( i \in \{1, 2, \ldots, p\} \) satisfies

\[ I_i \subset [0 : n]^+ \]

and

\[ (uv)[I_i] = u[I_i] = a_i. \]

Meanwhile, every \( i \in \{p + 1, p + 2, \ldots, p + q\} \) satisfies

\[ I_i \subset [n : n + m]^+ \]

and

\[ (uv)[I_i] = v[I_i - n] = b_{i-p}. \]

(\( I_i - n \) denotes the interval \( \{k - n \mid k \in I_i\} \)). We thus see that

\[ (uv)[I_i] \]

is a Lyndon word for every \( i \in \{1, 2, \ldots, p + q\} \).

By the definition of a CFL factorization, we have \( a_1 \geq a_2 \geq \cdots \geq a_p \) and \( b_1 \geq b_2 \geq \cdots \geq b_q \).

We have \( \sigma \in Sh_{n,m} \), so that \( \sigma^{-1}(1) < \sigma^{-1}(2) < \cdots < \sigma^{-1}(n) \) and \( \sigma^{-1}(n + 1) < \sigma^{-1}(n + 2) < \cdots < \sigma^{-1}(n + m) \). In other words, the restriction of the map \( \sigma^{-1} \) to the interval \([0 : n]^+\) is strictly increasing, and so is the restriction of the map \( \sigma^{-1} \) to the interval \([n : n + m]^+\).

(b) We will first show that

\[ (uv)[J] \]

is an interval such that the word \((uv)[J]\) is Lyndon, then \( \sigma^{-1}(J) \) is an interval.

\[ \text{Proof.} \]

The proof of this is the same as the proof of the fact that \( \alpha \) is a composition of length \( p \) and size \( \sum_{k=1}^{p} \ell(a_k) = n \).

Indeed, when \( i \leq \rho \), this follows from (6.4) and the fact that \( a_i \) is Lyndon; whereas in the other case, this follows from (6.5) and the fact that \( b_{i-p} \) is Lyndon.
Proof of (6.7): We will prove (6.7) by strong induction over $|J|$.

So, fix some $N \in \mathbb{N}$. Assume (as the induction hypothesis) that (6.7) has been proven whenever $|J| < N$. We now need to prove (6.7) when $|J| = N$.

Let $J \subset [0 : n]^{+}$ be an interval such that the word $(uv) [J]$ is Lyndon and such that $|J| = N$. We have to prove that $\sigma^{-1} (J)$ is an interval. This is obvious if $|J| = 1$ (because in this case, $\sigma^{-1} (J)$ is a one-element set, thus trivially an interval). Hence, we WLOG assume that we don’t have $|J| = 1$. We also don’t have $|J| = 0$, because $(uv) [J]$ has to be Lyndon (and the empty word is not). So we have $|J| > 1$. Now, $\ell ((uv) [J]) = |J| > 1$, and thus $(uv) [J]$ is a Lyndon word of length $> 1$. Let $v'$ be the (lexicographically) smallest nonempty proper suffix of $(uv) [J]$. Since $v'$ is a proper suffix of $w$, there exists a nonempty $u' \in \mathcal{A}^{*}$ such that $(uv) [J] = u'v'$. Consider this $u'$.

Now, Theorem 6.30(a) (applied to $(uv) [J]$, $u'$ and $v'$ instead of $w$, $u$ and $v$) yields that the words $u'$ and $v'$ are Lyndon. Also, Theorem 6.30(b) (applied to $(uv) [J]$, $u'$ and $v'$ instead of $w$, $u$ and $v$) yields that $u' < u [J] < v'$.

But from the fact that $(uv) [J] = u'v'$ and $u'$ and $v'$ both being nonempty, it becomes immediately clear that we can write $J$ as a union of two disjoint nonempty intervals $K$ and $L$ such that $K < L$, $u' = (uv) [K]$ and $v' = (uv) [L]$. Consider these $K$ and $L$. The intervals $K$ and $L$ are nonempty and have their sizes add up to $N$ (since they are disjoint and their union is $J$), and hence both must have size smaller than $|J| = N$. So $K \subset [0 : n]^{+}$ is an interval of size $|K| < N$ having the property that $(uv) [K]$ is Lyndon (since $(uv) [K] = u'$ is Lyndon). Thus, we can apply (6.7) to $K$ instead of $J$ (because of the induction hypothesis). As a result, we conclude that $\sigma^{-1} (K)$ is an interval. Similarly, we can apply (6.7) to $L$ instead of $J$ (we know that $(uv) [L]$ is Lyndon since $(uv) [L] = v'$), and learn that $\sigma^{-1} (L)$ is an interval. The intervals $\sigma^{-1} (K)$ and $\sigma^{-1} (L)$ are both nonempty (since $K$ and $L$ are nonempty), and their union is $\sigma^{-1} (J)$ (because the union of $K$ and $L$ is $J$). The nonempty intervals $K$ and $L$ both are subsets of $[0 : n]^{+}$ (since their union is $J \subset [0 : n]^{+}$), and their union $K \cup L$ is an interval (since their union $K \cup L$ is $J$, and we know that $J$ is an interval).

Now, assume (for the sake of contradiction) that $\sigma^{-1} (J)$ is not an interval. Since $J$ is the union of $K$ and $L$, we have $J = K \cup L$ and thus $\sigma^{-1} (J) = \sigma^{-1} (K \cup L) = \sigma^{-1} (K) \cup \sigma^{-1} (L)$ (since $\sigma$ is a bijection). Therefore, $\sigma^{-1} (K) \cup \sigma^{-1} (L)$ is not an interval (since $\sigma^{-1} (J)$ is not an interval). Thus, Lemma 6.42(b) yields that there exists a nonempty interval $P \subset [n : n + m]^{+}$ such that $\sigma^{-1} (P)$, $\sigma^{-1} (K) \cup \sigma^{-1} (P)$ and $\sigma^{-1} (P) \cup \sigma^{-1} (L)$ are intervals and such that $\sigma^{-1} (K) < \sigma^{-1} (P) < \sigma^{-1} (L)$. Consider this $P$. Since $P$ is nonempty, we have $|P| \neq 0$.

Lemma 6.44(b) (applied to $K$ and $P$ instead of $I$ and $J$) yields

\[(uv) [K] \cdot (uv) [P] \geq (uv) [P] \cdot (uv) [K].\]

Since $(uv) [K] = u'$, this rewrites as

\[(u') \cdot (uv) [P] \geq (uv) [P] \cdot u'.\]  

But Lemma 6.44(c) (applied to $P$ and $L$ instead of $I$ and $J$) yields

\[(uv) [P] \cdot (uv) [L] \geq (uv) [L] \cdot (uv) [P].\]

Since $(uv) [L] = v'$, this rewrites as

\[(uv) [P] \cdot v' \geq v' \cdot (uv) [P].\]

Recall also that $u' < v'$, and that both words $u'$ and $v'$ are Lyndon. Now, Corollary 6.17 (applied to $u'$, $v'$ and $(uv) [P]$ instead of $u$, $v$ and $z$) yields that $(uv) [P]$ is the empty word (because of (6.9) and (6.11)), so that $\ell ((uv) [P]) = 0$. This contradicts $\ell ((uv) [P]) = |P| \neq 0$. This contradiction shows that our assumption (that $\sigma^{-1} (J)$ is not an interval) was wrong. Hence, $\sigma^{-1} (J)$ is an interval. This completes the induction step, and thus (6.7) is proven.

Similarly to (6.7), we can show that

\[(uv) [J] \subset [n : n + m]^{+}\]

is an interval such that the word $(uv) [J]$ is Lyndon, then $\sigma^{-1} (J)$ is an interval.

Now, let $i \in \{1, 2, \ldots, p + q\}$ be arbitrary. We are going to prove that

\[\sigma^{-1} (I_i) \text{ is an interval}.\]

Proof of (6.13): We must be in one of the following two cases:

Case 1: We have $i \in \{1, 2, \ldots, p\}$.
Case 2: We have $i \in \{p+1, p+2, \ldots, p+q\}$.

Let us first consider Case 1. In this case, we have $i \in \{1, 2, \ldots, p\}$. Thus, $I_i \subset [0 : n]^+$ (by (6.4)). Also, (6.4) yields that $(uv)[I_i] = a_i$ is a Lyndon word. Hence, (6.7) (applied to $J = I_i$) yields that $\sigma^{-1}(I_i)$ is an interval. Thus, (6.13) is proven in Case 1.

Similarly, we can prove (6.13) in Case 2, using (6.5) and (6.12) instead of (6.4) and (6.7), respectively. Hence, (6.13) is proven.

So we know that $\sigma^{-1}(I_i)$ is an interval. But we also know that either $I_i \subset [0 : n]^+$ or $I_i \subset [n : n+m]^+$ (depending on whether $i \leq a$ or $i > p$). As a consequence, the restriction of the map $\sigma^{-1}$ to the interval $I_i$ is increasing (because the restriction of the map $\sigma^{-1}$ to the interval $[0 : n]^+$ is strictly increasing, and so is the restriction of the map $\sigma^{-1}$ to the interval $[n : n+m]^+$).

Now, let us forget that we fixed $i$. We thus have shown that every $i \in \{1, 2, \ldots, p+q\}$ has the two properties:

- the set $\sigma^{-1}(I_i)$ is an interval;
- the restriction of the map $\sigma^{-1}$ to the interval $I_i$ is increasing.

In other words, the permutation $\sigma$ is $(\alpha\beta)$-clumping (since $(I_1, I_2, \ldots, I_{p+q})$ is the interval system corresponding to the composition $\alpha\beta$). Hence, Proposition 6.50(b) (applied to $n + m$, $\alpha\beta$ and $p+q$ instead of $n$, $\alpha$ and $\ell$) shows that there exists a unique $\tau \in \mathcal{S}_{p+q}$ satisfying $\sigma = \text{iper}(\alpha\beta, \tau)$. Thus, the uniqueness part of Theorem 6.56(b) (i.e., the claim that the $\tau$ in Theorem 6.56(b) is unique if it exists) is proven.

It now remains to prove the existence part of Theorem 6.56(b), i.e., to prove that there exists at least one permutation $\tau \in \mathcal{S}_{p,q}$ satisfying $\sigma = \text{iper}(\alpha\beta, \tau)$. We already know that there exists a unique $\tau \in \mathcal{S}_{p+q}$ satisfying $\sigma = \text{iper}(\alpha\beta, \tau)$. Consider this $\tau$. We will now prove that $(uv)[I_{\tau(i)}] \geq (uv)[I_{\tau(i+1)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}]$ and $\tau \in \mathcal{S}_{p,q}$. Once this is done, the existence part of Theorem 6.56(b) will be proven, and thus the proof of Theorem 6.56(b) will be complete.

Proposition 6.52 yields that $\tau \in \mathcal{S}_{p,q}$ if and only if $\text{iper}(\alpha\beta, \tau) \in \mathcal{S}_{n,m}$. Since we know that $\text{iper}(\alpha\beta, \tau) = \sigma \in \mathcal{S}_{n,m}$, we thus conclude that $\tau \in \mathcal{S}_{p,q}$. The only thing that remains to be proven now is that

$$ (uv)[I_{\tau(i)}] \geq (uv)[I_{\tau(i+1)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}]. $$

Proof of (6.14): We have $\tau \in \mathcal{S}_{p,q}$. In other words, $\tau^{-1}(1) < \tau^{-1}(2) < \cdots < \tau^{-1}(p)$ and $\tau^{-1}(p+1) < \tau^{-1}(p+2) < \cdots < \tau^{-1}(p+q)$. In other words, the restriction of the map $\tau^{-1}$ to the interval $[0 : p]^+$ is strictly increasing, and so is the restriction of the map $\tau^{-1}$ to the interval $[p : p+q]^+$.

Let $i \in \{1, 2, \ldots, p+q-1\}$. We will show that

$$ (uv)[I_{\tau(i)}] \geq (uv)[I_{\tau(i+1)}]. $$

Clearly, both $\tau(i)$ and $\tau(i+1)$ belong to $\{1, 2, \ldots, p+q\} = \{1, 2, \ldots, p\} \cup \{p+1, p+2, \ldots, p+q\}$. Thus, we must be in one of the following four cases:

Case 1: We have $\tau(i) \in \{1, 2, \ldots, p\}$ and $\tau(i+1) \in \{1, 2, \ldots, p\}$.

Case 2: We have $\tau(i) \in \{1, 2, \ldots, p\}$ and $\tau(i+1) \in \{p+1, p+2, \ldots, p+q\}$.

Case 3: We have $\tau(i) \in \{p+1, p+2, \ldots, p+q\}$ and $\tau(i+1) \in \{1, 2, \ldots, p\}$.

Case 4: We have $\tau(i) \in \{p+1, p+2, \ldots, p+q\}$ and $\tau(i+1) \in \{p+1, p+2, \ldots, p+q\}$.

Let us consider Case 1 first. In this case, we have $\tau(i) \in \{1, 2, \ldots, p\}$ and $\tau(i+1) \in \{1, 2, \ldots, p\}$. From the fact that the restriction of the map $\tau^{-1}$ to the interval $[0 : p]^+$ is strictly increasing, we can easily deduce $\tau(i) < \tau(i+1)$ \footnote{Proof: Assume the contrary. Then, $\tau(i) \geq \tau(i+1)$. Since both $\tau(i)$ and $\tau(i+1)$ belong to $\{1, 2, \ldots, p\} = [0 : p]^+$, this yields $\tau^{-1}(\tau(i)) \geq \tau^{-1}(\tau(i+1))$ (since the restriction of the map $\tau^{-1}$ to the interval $[0 : p]^+$ is strictly increasing), which contradicts $\tau^{-1}(\tau(i)) = i < i+1 = \tau^{-1}(\tau(i+1))$. This contradiction proves the assumption wrong, qed.}.

Therefore, $a_{\tau(i)} \geq a_{\tau(i+1)}$ (since $a_1 \geq a_2 \geq \cdots \geq a_p$).

But $(uv)[I_{\tau(i)}] = a_{\tau(i)}$ (by (6.4), applied to $\tau$ instead of $\alpha\beta$) and $(uv)[I_{\tau(i+1)}] = a_{\tau(i+1)}$ (similarly). In view of these equalities, the inequality $a_{\tau(i)} \geq a_{\tau(i+1)}$ rewrites as $(uv)[I_{\tau(i)}] \geq (uv)[I_{\tau(i+1)}]$. Thus, (6.15) is proven in Case 1.

Similarly, we can show (6.15) in Case 4 (observing that $(uv)[I_{\tau(i)}] = b_{\tau(i)-p}$ and $(uv)[I_{\tau(i+1)}] = b_{\tau(i+1)-p}$ in this case).

Let us now consider Case 2. In this case, we have $\tau(i) \in \{1, 2, \ldots, p\}$ and $\tau(i+1) \in \{p+1, p+2, \ldots, p+q\}$. From $\tau(i) \in \{1, 2, \ldots, p\}$, we conclude that $I_{\tau(i)} \subset [0 : n]^+$. From $\tau(i+1) \in \{p+1, p+2, \ldots, p+q\}$, we conclude that $I_{\tau(i+1)} \subset [n : n+m]^+$. The intervals $I_{\tau(i)}$ and $I_{\tau(i+1)}$ are clearly nonempty.
Proposition 6.48(d) (applied to \( n + m, \alpha, \beta, p + q \) and \((I_1, I_2, \ldots, I_{p+q})\) instead of \( n, \alpha, \ell \) and \((I_1, I_2, \ldots, I_{\ell})\)) yields that the sets \( \sigma^{-1}(I_{\tau(i)}), \sigma^{-1}(I_{\tau(i+1)}) \) and \( \sigma^{-1}(I_{\tau(i)}) \cup \sigma^{-1}(I_{\tau(i+1)}) \) are nonempty intervals, and that we have \( \sigma^{-1}(I_{\tau(i)}) < \sigma^{-1}(I_{\tau(i+1)}) \). Hence, Lemma 6.44(b) (applied to \( I = I_{\tau(i)} \) and \( J = I_{\tau(i+1)} \)) yields
\[
(uv) \left[ I_{\tau(i)} \right] \cdot (uv) \left[ I_{\tau(i+1)} \right] \geq (uv) \left[ I_{\tau(i)} \right] \cdot (uv) \left[ I_{\tau(i+1)} \right].
\]
But \((uv) \left[ I_{\tau(i)} \right] \) and \((uv) \left[ I_{\tau(i+1)} \right] \) are Lyndon words (as a consequence of (6.6)). Thus, Proposition 6.18 (applied to \((uv) \left[ I_{\tau(i)} \right] \) and \((uv) \left[ I_{\tau(i+1)} \right] \) instead of \( u \) and \( v \)) shows that \((uv) \left[ I_{\tau(i)} \right] \geq (uv) \left[ I_{\tau(i+1)} \right] \) if and only if \((uv) \left[ I_{\tau(i)} \right], (uv) \left[ I_{\tau(i+1)} \right] \) \((uv) \left[ I_{\tau(i)} \right] \geq (uv) \left[ I_{\tau(i+1)} \right] \cdot (uv) \left[ I_{\tau(i)} \right] \). Since we know that \((uv) \left[ I_{\tau(i)} \right], (uv) \left[ I_{\tau(i+1)} \right] \) \((uv) \left[ I_{\tau(i)} \right] \geq (uv) \left[ I_{\tau(i+1)} \right] \) holds, we thus conclude that \((uv) \left[ I_{\tau(i)} \right] \geq (uv) \left[ I_{\tau(i+1)} \right] \). Thus, (6.15) is proven in Case 2.

The proof of (6.15) in Case 3 is analogous to that in Case 2 (the main difference being that Lemma 6.44(c) is used in lieu of Lemma 6.44(b)).

Thus, (6.15) is proven in all possible cases. So we always have (6.15). In other words, \((uv) \left[ I_{\tau(i)} \right] \geq (uv) \left[ I_{\tau(i+1)} \right] \).

Now, forget that we fixed \( i \). We hence have shown that \((uv) \left[ I_{\tau(i)} \right] \geq (uv) \left[ I_{\tau(i+1)} \right] \) for all \( i \in \{1, 2, \ldots, p + q - 1\} \). This proves (6.14), and thus completes our proof of Theorem 6.56(b).

(a) Let \( \tau \in Sh_{p,q} \) be such that
\[
(uv) \left[ I_{\tau(1)} \right] \geq (uv) \left[ I_{\tau(2)} \right] \geq \cdots \geq (uv) \left[ I_{\tau(p+q)} \right].
\]
Set \( \sigma = \text{iper}(\alpha \beta, \tau) \). Then, Proposition 6.52 yields that \( \tau \in Sh_{p,q} \) if and only if \( \text{iper}(\alpha \beta, \sigma) \in Sh_{n,m} \). Since we know that \( \tau \in Sh_{p,q} \), we can deduce from this that \( \text{iper}(\alpha \beta, \sigma) \in Sh_{n,m} \), so that \( \sigma = \text{iper}(\alpha \beta, \tau) \in Sh_{n,m} \).

It remains to prove that the word \( u \uplus v \) is the lexicographically highest element of the multiset \( u \uplus v \).

It is clear that the multiset \( u \uplus v \) has some lexicographically highest element. This element has the form \( u \uplus v \) for some \( \bar{\sigma} \in Sh_{n,m} \) (because any element of this multiset has such a form). Consider this \( \bar{\sigma} \).

Theorem 6.56(b) (applied to \( \bar{\sigma} \) instead of \( \sigma \)) yields that there exists a unique permutation \( \bar{\tau} \in Sh_{p,q} \) satisfying \((uv) \left[ I_{\bar{\tau}(1)} \right] \geq (uv) \left[ I_{\bar{\tau}(2)} \right] \geq \cdots \geq (uv) \left[ I_{\bar{\tau}(p+q)} \right] \) and \( \bar{\sigma} = \text{iper}(\alpha \beta, \bar{\tau}) \). (What we call \( \bar{\tau} \) here is what has been called \( \tau \) in Theorem 6.56(b).)

Now, the chain of inequalities \((uv) \left[ I_{\bar{\tau}(1)} \right] \geq (uv) \left[ I_{\bar{\tau}(2)} \right] \geq \cdots \geq (uv) \left[ I_{\bar{\tau}(p+q)} \right] \) shows that the list \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) is the result of sorting the list \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) in decreasing order. But the chain of inequalities (6.16) shows that the list \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) is the result of sorting the same list \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) in decreasing order. So each of the two lists \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) and \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) is the result of sorting one and the same list \((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) in decreasing order. Since the result of sorting a given list in decreasing order is unique, this yields
\[
((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)) = ((uv) \left[ I_{\bar{\tau}(1)} \right], (uv) \left[ I_{\bar{\tau}(2)} \right], \ldots, (uv) \left[ I_{\bar{\tau}(p+q)} \right] \)).
\]
Hence,
\[
(uv) \left[ I_{\bar{\tau}(1)} \right] \cdot (uv) \left[ I_{\bar{\tau}(2)} \right] \cdots (uv) \left[ I_{\bar{\tau}(p+q)} \right] = (uv) \left[ I_{\bar{\tau}(1)} \right] \cdot (uv) \left[ I_{\bar{\tau}(2)} \right] \cdots (uv) \left[ I_{\bar{\tau}(p+q)} \right] \).
\]
But Lemma 6.54 yields
\[
(uv) \left[ I_{\bar{\tau}(1)} \right] = (uv) \left[ I_{\tau(1)} \right] \cdot (uv) \left[ I_{\tau(2)} \right] \cdots (uv) \left[ I_{\tau(p+q)} \right] .
\]
Meanwhile, Lemma 6.54 (applied to \( \bar{\tau} \) and \( \bar{\sigma} \) instead of \( \tau \) and \( \sigma \)) yields
\[
(uv) \left[ I_{\bar{\tau}(1)} \right] \cdot (uv) \left[ I_{\bar{\tau}(2)} \right] \cdots (uv) \left[ I_{\bar{\tau}(p+q)} \right] = (uv) \left[ I_{\tau(1)} \right] \cdot (uv) \left[ I_{\tau(2)} \right] \cdots (uv) \left[ I_{\tau(p+q)} \right] \quad \text{(by (6.17))}
\]
\[
= u \uplus v \quad \text{(by (6.18))}.
\]
Thus, \( u \uplus v \) is the lexicographically highest element of the multiset \( u \uplus v \) (since we know that \( u \uplus v \) is the lexicographically highest element of the multiset \( u \uplus v \)). This proves Theorem 6.56(a).

Now, in order to prove Theorem 6.36, we record a very simple fact about counting shuffles:
Proposition 6.57. Let \( p \in \mathbb{N} \) and \( q \in \mathbb{N} \). Let \( \mathfrak{W} \) be a totally ordered set, and let \( h : \{1, 2, \ldots, p+q\} \to \mathfrak{W} \) be a map. Assume that \( h(1) \geq h(2) \geq \cdots \geq h(p) \) and \( h(p+1) \geq h(p+2) \geq \cdots \geq h(p+q) \).

For every \( w \in \mathfrak{W} \), let \( a(w) \) denote the number of all \( i \in \{1, 2, \ldots, p\} \) satisfying \( h(i) = w \), and let \( b(w) \) denote the number of all \( i \in \{p+1, p+2, \ldots, p+q\} \) satisfying \( h(i) = w \).

Then, the number of \( \tau \in \text{Sh}_{p,q} \) satisfying \( h(\tau(1)) \geq h(\tau(2)) \geq \cdots \geq h(\tau(p+q)) \) is \( \prod_{w \in \mathfrak{W}} \left( a(w) + b(w) \right) \). For every \( i \in \{1, 2, \ldots, p+q\} \) satisfying \( h(i) = w \).

(Of course, all but finitely many factors of this product are 1.)


Proof of Theorem 6.36. Let \( n = \ell(u) \) and \( m = \ell(v) \). Define \( \alpha, \beta \) and \( (I_1, I_2, \ldots, I_{p+q}) \) as in Theorem 6.56.

Since \( (a_1, a_2, \ldots, a_p) \) is the CFL factorization of \( u \), we have \( a_1 \geq a_2 \geq \cdots \geq a_p \) and \( a_1 a_2 \cdots a_p = u \). Similarly, \( b_1 \geq b_2 \geq \cdots \geq b_q \) and \( b_1 b_2 \cdots b_q = v \).

From (6.4), we see that \( (uv)[I_i] = a_i \) for every \( i \in \{1, 2, \ldots, p\} \). From (6.5), we see that \( (uv)[I_i] = b_{i-p} \) for every \( i \in \{p+1, p+2, \ldots, p+q\} \). Combining these two equalities, we obtain

\[
(6.19) \quad (uv)[I_i] = \begin{cases} 
  a_i, & \text{if } i \leq p; \\
  b_{i-p}, & \text{if } i > p 
\end{cases} \quad \text{for every } i \in \{1, 2, \ldots, p+q\}.
\]

In other words,

\[
(6.20) \quad ((uv)[I_1], (uv)[I_2], \ldots, (uv)[I_{p+q}]) = (a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q).
\]

(a) Let \( z \) be the lexicographically highest element of the multiset \( u \uplus v \). We must prove that \( z = c_1 c_2 \cdots c_{p+q} \).

Since \( z \in u \uplus v \), we can write \( z \) in the form \( u \uplus v \) for some \( \sigma \in \text{Sh}_{n,m} \) (since we can write any element of \( u \uplus v \) in this form). Consider this \( \sigma \). Then, \( u \uplus v = z \) is the lexicographically highest element of the multiset \( u \uplus v \). Hence, Theorem 6.56(b) yields that there exists a unique permutation \( \tau \in \text{Sh}_{p,q} \) satisfying \( (uv)[I_{\tau(1)}] \geq (uv)[I_{\tau(2)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}] \) and \( \sigma = \text{iper}(\alpha \beta, \tau) \). Consider this \( \tau \).

Now, \( \tau \in \text{Sh}_{p,q} \subseteq \mathfrak{S}_{p+q} \) is a permutation, and thus the list \( ((uv)[I_{\tau(1)}], (uv)[I_{\tau(2)}], \ldots, (uv)[I_{\tau(p+q)}]) \) is a rearrangement of the list \( ((uv)[I_1], (uv)[I_2], \ldots, (uv)[I_{p+q}]) \). Due to (6.20), this rewrites as follows: The list \( ((uv)[I_{\tau(1)}], (uv)[I_{\tau(2)}], \ldots, (uv)[I_{\tau(p+q)}]) \) is a rearrangement of the list \( (a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q) \).

Hence, \( ((uv)[I_{\tau(1)}], (uv)[I_{\tau(2)}], \ldots, (uv)[I_{\tau(p+q)}]) \) is the result of sorting the list \( (a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q) \) in decreasing order (since \( (uv)[I_{\tau(1)}] \geq (uv)[I_{\tau(2)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}] \)). But since the result of sorting the list \( (a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q) \) in decreasing order is \( (c_1, c_2, \ldots, c_{p+q}) \), this becomes

\[
((uv)[I_{\tau(1)}], (uv)[I_{\tau(2)}], \ldots, (uv)[I_{\tau(p+q)}]) = (c_1, c_2, \ldots, c_{p+q}).
\]

Hence,

\[
(uv)[I_{\tau(1)}] \cdot (uv)[I_{\tau(2)}] \cdots (uv)[I_{\tau(p+q)}] = c_1 \cdot c_2 \cdots c_{p+q}.
\]

But Lemma 6.54 yields

\[
uu_{\sigma} = (uv)[I_{\tau(1)}] \cdot (uv)[I_{\tau(2)}] \cdots (uv)[I_{\tau(p+q)}].
\]

Altogether, we have

\[
z = uu_{\sigma} = (uv)[I_{\tau(1)}] \cdot (uv)[I_{\tau(2)}] \cdots (uv)[I_{\tau(p+q)}] = c_1 \cdot c_2 \cdots c_{p+q} = c_1 c_2 \cdots c_{p+q}.
\]

This proves Theorem 6.36(a).

(b) Recall that \( uu_{\sigma} = \left\{ uu_{\sigma} : \sigma \in \text{Sh}_{n,m} \right\}_{\text{multiset}} \). Hence,

\[
(\text{the multiplicity with which the lexicographically highest element of the multiset } \uu_{\sigma} \text{ appears in the multiset } uu_{\sigma})
\]

\[
= \left( \text{the number of all } \sigma \in \text{Sh}_{n,m} \text{ such that } uu_{\sigma} \text{ is the lexicographically highest element of the multiset } uu_{\sigma} \right).
\]
But for a given $\sigma \in Sh_{n,m}$, we know that $u \cup \sigma v$ is the lexicographically highest element of the multiset $u \cup v$ if and only if $\sigma$ can be written in the form $\sigma = \text{iper}(\alpha \beta, \tau)$ for some $\tau \in Sh_{p,q}$ satisfying $(uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}]$. Hence,

\[
\left( \text{the number of all } \sigma \in Sh_{n,m} \text{ such that } u \cup \sigma v \text{ is the lexicographically highest element of the multiset } u \cup v \right) = \left( \text{the number of all } \sigma \in Sh_{n,m} \text{ which can be written in the form } \sigma = \text{iper}(\alpha \beta, \tau) \text{ for some } \tau \in Sh_{p,q} \text{ satisfying } (uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}] \right) = \left. \left( \text{the number of all } \tau \in Sh_{p,q} \text{ satisfying } (uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}] \right) \right| \text{ (because if } \sigma \in Sh_{n,m} \text{ can be written in the form } \sigma = \text{iper}(\alpha \beta, \tau) \text{ for some } \tau \in Sh_{p,q} \text{ satisfying } (uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}], \text{ then } \sigma \text{ can be written uniquely in this form).} \right.
\]

Thus, we can apply Proposition 6.57 to $Sh_{p,q}$.

Proof. In fact, the “if” part of this assertion follows from Theorem 6.56(a), whereas its “only if” part follows from Theorem 6.56(b).

\[
\text{Proof. } \text{Let } \sigma \in Sh_{n,m} \text{ be such that } \sigma \text{ can be written in the form } \sigma = \text{iper}(\alpha \beta, \tau) \text{ for some } \tau \in Sh_{p,q} \text{ satisfying } (uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}]. \text{ Then, the word } u \cup v \text{ is the lexicographically highest element of the multiset } u \cup v \text{ (according to Theorem 6.56(a)). Hence, there exists a unique permutation } \tau \in Sh_{p,q} \text{ satisfying } (uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}] \text{ and } \sigma = \text{iper}(\alpha \beta, \tau) \text{ (according to Theorem 6.56(b)). In other words, } \sigma \text{ can be written uniquely in the form } \sigma = \text{iper}(\alpha \beta, \tau) \text{ for some } \tau \in Sh_{p,q} \text{ satisfying } (uv) [I_{\tau(1)}] \geq (uv) [I_{\tau(2)}] \geq \cdots \geq (uv) [I_{\tau(p+q)}], \text{ qed.}
\]
In other words,

\[(6.22) \quad \prod_{w \in \mathcal{L}} \left( \text{mult}_w u + \text{mult}_w v \right).\]

But for every \(i \in \{1, 2, \ldots, p+q\}\), we have

\[h(i) = \begin{cases} 
  a_i, & \text{if } i \leq p; \\
  b_{i-p}, & \text{if } i > p
\end{cases} = (uv)[I_i] \quad \text{(by (6.19)).}\]

Hence, for any \(\tau \in \text{Sh}_{p,q}\), the condition \(h(\tau(1)) \geq h(\tau(2)) \geq \cdots \geq h(\tau(p+q))\) is equivalent to \((uv)[I_{\tau(1)}] \geq (uv)[I_{\tau(2)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}]\). Thus,

\[
\text{the number of all } \tau \in \text{Sh}_{p,q} \text{ satisfying } h(\tau(1)) \geq h(\tau(2)) \geq \cdots \geq h(\tau(p+q))
\]

is equivalent to

\[
(\text{the number of all } \tau \in \text{Sh}_{p,q} \text{ satisfying } (uv)[I_{\tau(1)}] \geq (uv)[I_{\tau(2)}] \geq \cdots \geq (uv)[I_{\tau(p+q)}])
\]

and

\[
= \text{(the multiplicity with which the lexicographically highest element of the multiset } u \sqcup v \text{ appears in the multiset } u \sqcup v)
\]

(by (6.21)). Compared with (6.22), this yields

\[
\text{(the multiplicity with which the lexicographically highest element of the multiset } u \sqcup v \text{ appears in the multiset } u \sqcup v)
\]

\[= \prod_{w \in \mathcal{L}} \left( \text{mult}_w u + \text{mult}_w v \right).\]

This proves Theorem 6.36(b).

(c) We shall use the notations of Theorem 6.36(a) and Theorem 6.36(b).

Assume that \(a_i \geq b_j\) for every \(i \in \{1, 2, \ldots, p\}\) and \(j \in \{1, 2, \ldots, q\}\). This, combined with \(a_1 \geq a_2 \geq \cdots \geq a_p\) and \(b_1 \geq b_2 \geq \cdots \geq b_q\), yields that \(a_1 \geq a_2 \geq \cdots \geq a_p \geq b_1 \geq b_2 \geq \cdots \geq b_q\). Thus, the list \((a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q)\) is weakly decreasing. Thus, the result of sorting the list \((a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q)\) in decreasing order is the list \((a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q)\) itself. But since this result is \((c_1, c_2, \ldots, c_{p+q})\), this shows that \((c_1, c_2, \ldots, c_{p+q}) = (a_1, a_2, \ldots, a_p, b_1, b_2, \ldots, b_q)\). Hence, \(c_1 c_2 \cdots c_{p+q} = a_1 a_2 \cdots a_p b_1 b_2 \cdots b_q = uv\). Now, Theorem 6.36(a) yields that the lexicographically highest element of the multiset \(u \sqcup v\) is \(c_1 c_2 \cdots c_{p+q} = uv\). This proves Theorem 6.36(c).

(d) We shall use the notations of Theorem 6.36(a) and Theorem 6.36(b).

Assume that \(a_i > b_j\) for every \(i \in \{1, 2, \ldots, p\}\) and \(j \in \{1, 2, \ldots, q\}\). Thus, \(a_i > b_j\) for every \(i \in \{1, 2, \ldots, p\}\) and \(j \in \{1, 2, \ldots, q\}\). Hence, Theorem 6.36(c) yields that the lexicographically highest element of the multiset \(u \sqcup v\) is \(uv\). Therefore, Theorem 6.36(b) shows that the multiplicity with which this word \(uv\) appears in the multiset \(u \sqcup v\) is \(\prod_{w \in \mathcal{L}} \left( \text{mult}_w u + \text{mult}_w v \right)\).
Now, every \( w \in \mathcal{L} \) satisfies \( \left( \frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u} \right) = 1 \). Thus, as we know, the multiplicity with which this word \( uv \) appears in the multiset \( u \sqcup v \) is \( \prod_{w \in \mathcal{L}} \left( \frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u} \right) = \prod_{\mathcal{L}} 1 = 1 \). This proves Theorem 6.36(d).

(c) We shall use the notations of Theorem 6.36(a) and Theorem 6.36(b).

Since \( u \) is a Lyndon word, the 1-tuple \( (u) \) is the CFL factorization of \( u \). Hence, we can apply Theorem 6.36(c) to 1 and \( (u) \) instead of \( p \) and \( \{a_1, a_2, \ldots, a_p\} \). As a result, we conclude that the lexicographically highest element of the multiset \( u \sqcup v \) is \( uv \). It remains to prove that the multiplicity with which this word \( uv \) appears in the multiset \( u \sqcup v \) is \( \text{mult}_u v + 1 \).

For every \( w \in \mathcal{L} \) satisfying \( w \neq u \), we have

\[
\text{mult}_w u = 0
\]

Also, \( \text{mult}_u u = 1 \) (for a similar reason). But \( uv \) is the lexicographically highest element of the multiset \( u \sqcup v \). Hence, the multiplicity with which the word \( uv \) appears in the multiset \( u \sqcup v \) is the multiplicity with which the lexicographically highest element of the multiset \( u \sqcup v \) appears in the multiset \( u \sqcup v \). According to Theorem 6.36(b), the latter multiplicity is

\[
\prod_{w \in \mathcal{L}} \left( \frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u} \right) = \left( \frac{1 + \text{mult}_u v}{1} \right) \cdot \prod_{u \neq w} \left( \frac{0 + \text{mult}_w v}{\text{mult}_w u} \right) = (\text{mult}_u v + 1) \cdot \prod_{u \neq w} 1 = \text{mult}_u v + 1.
\]

This proves Theorem 6.36(c).

\[\square\]

As an application of our preceding results, we can prove a further necessary and sufficient criterion for a word to be Lyndon; this criterion is due to Chen/Fox/Lyndon [30, \( \mathfrak{A}' = \mathfrak{A}^{op} \)]:

**Exercise 6.59.** Let \( w \in \mathfrak{A}^* \) be a nonempty word. Prove that \( w \) is Lyndon if and only if for any two nonempty words \( u \in \mathfrak{A}^* \) and \( v \in \mathfrak{A}^* \) satisfying \( w = uv \), there exists at least one \( s \in u \sqcup v \) satisfying \( s > w \).

6.3. **Radford’s theorem on the shuffle algebra.** We recall that our goal in Chapter 6 is to exhibit an algebraically independent generating set of the k-algebra \( \text{QSym} \). Having the notion of Lyndon words – which will, to some extent, but not literally, parametrize this generating set – in place, we could start the construction of this generating set immediately. However, it might come off as rather unmotivated this way, and so we begin with some warmups. First, we shall prove Radford’s theorem on the shuffle algebra.

**Definition 6.60.** A polynomial algebra will mean a k-algebra which is isomorphic to the polynomial ring \( k[x_i \mid i \in I] \) as a k-algebra (for some indexing set \( I \)). Note that \( I \) needs not be finite.

\[\text{Proof.} \text{ Assume the contrary. Then, there exists at least one } w \in \mathcal{L} \text{ such that } \left( \frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u} \right) \neq 1. \text{ Consider this } w. \text{ Both } \text{mult}_w u \text{ and } \text{mult}_w v \text{ must be positive (since } \left( \frac{\text{mult}_w u + \text{mult}_w v}{\text{mult}_w u} \right) \neq 1). \text{ Since } \text{mult}_w u \text{ is positive, there must be at least one term in the CFL factorization of } u \text{ which is equal to } w. \text{ In other words, there is at least one } i \in \{1, 2, \ldots, p\} \text{ satisfying } a_i = w (\text{since } (a_1, a_2, \ldots, a_p) \text{ is the CFL factorization of } u). \text{ Similarly, there is at least one } j \in \{1, 2, \ldots, q\} \text{ satisfying } b_j = w. \text{ These } i \text{ and } j \text{ satisfy } a_i = w = b_j, \text{ which contradicts } a_i > b_j. \text{ This contradiction shows that our assumption was false, qed.}
\]

\[\text{Proof of (6.23): } \text{Let } w \in \mathcal{L} \text{ be such that } w \neq u. \text{ Then, the number of terms in the list } (u) \text{ which are equal to } w \text{ is } 0. \text{ Since } (a) \text{ is the CFL factorization of } u, \text{ this rewrites as follows: The number of terms in the CFL factorization of } u \text{ which are equal to } w \text{ is } 0. \text{ In other words, } \text{mult}_w u = 0. \text{ This proves (6.23).} \]
Equivalently, a polynomial algebra can be defined as a \( k \)-algebra which has an algebraically independent (over \( k \)) generating set. Yet equivalently, a polynomial algebra can be defined as a \( k \)-algebra which is isomorphic to the symmetric algebra of a free \( k \)-module.

Keep in mind that when we say that a certain bialgebra \( A \) is a polynomial algebra, we are making no statement about the coalgebra structure on \( A \). The isomorphism from \( A \) to the symmetric algebra of a free \( k \)-module needs not be a coalgebra isomorphism, and the algebraically independent generating set of \( A \) needs not consist of primitives. Thus, showing that a bialgebra \( A \) is a polynomial algebra does not trivialize the study of its bialgebraic structure.

**Remark 6.61.** Let \( V \) be a \( k \)-module, and let \( \mathfrak{A} \) be a totally ordered set. Let \( b_a \) be an element of \( V \) for every \( a \in \mathfrak{A} \). Consider the shuffle algebra \( Sh (V) \) (defined in Definition 1.73).

For every word \( w \in \mathfrak{A}^* \) over the alphabet \( \mathfrak{A} \), let us define an element \( b_w \) of \( Sh (V) \) by \( b_w = b_{w_1}b_{w_2}\cdots b_{w_{\ell}} \), where \( \ell \) is the length of \( w \). (The multiplication used here is that of \( T(V) \), not that of \( Sh(V) \); the latter is denoted by \( \shuffle \).)

Let \( u \in \mathfrak{A}^* \) and \( v \in \mathfrak{A}^* \) be two words over the alphabet \( \mathfrak{A} \). Let \( n = \ell(u) \) and \( m = \ell(v) \). Then,

\[
b_u \shuffle b_v = \sum_{\sigma \in Sh_{n,m}} b_{\sigma(u,v)}.
\]

**Exercise 6.62.** Prove Remark 6.61.

**Hint:** This follows from the definition of \( \shuffle \).

We can now state Radford’s theorem [131, Theorem 3.1.1(e)]:

**Theorem 6.63.** Assume that \( \mathbb{Q} \) is a subring of \( k \). Let \( V \) be a free \( k \)-module with a basis \( (b_a)_{a \in \mathfrak{A}} \), where \( \mathfrak{A} \) is a totally ordered set. Then, the shuffle algebra \( Sh (V) \) (defined in Definition 1.73) is a polynomial \( k \)-algebra. An algebraically independent generating set of \( Sh (V) \) can be constructed as follows:

For every word \( w \in \mathfrak{A}^* \) over the alphabet \( \mathfrak{A} \), let us define an element \( b_w \) of \( Sh (V) \) by \( b_w = b_{w_1}b_{w_2}\cdots b_{w_{\ell}} \), where \( \ell \) is the length of \( w \). (The multiplication used here is that of \( T(V) \), not that of \( Sh(V) \); the latter is denoted by \( \shuffle \).) Let \( L \) denote the set of all Lyndon words over the alphabet \( \mathfrak{A} \). Then, \((b_w)_{w \in L}\) is an algebraically independent generating set of the \( k \)-algebra \( Sh (V) \).

**Example 6.64.** For this example, let \( \mathfrak{A} \) be the alphabet \( \{1, 2, 3, \ldots\} \) with total order given by \( 1 < 2 < 3 < \cdots \), and assume that \( \mathbb{Q} \) is a subring of \( k \). Let \( V \) be the free \( k \)-module with basis \( (b_a)_{a \in \mathfrak{A}} \). We use the notations of Theorem 6.63. Then, Theorem 6.63 yields that \((b_w)_{w \in L}\) is an algebraically independent generating set of the \( k \)-algebra \( Sh (V) \). Here are some examples of elements of \( Sh (V) \) written as polynomials in this generating set:

\[
\begin{align*}
b_{12} &= b_{12} \quad \text{(the word 12 itself is Lyndon)}; \\
b_{21} &= b_{1}\shuffle b_{2} - b_{12}; \\
b_{11} &= \frac{1}{2} b_{1}\shuffle b_{1}; \\
b_{123} &= b_{123} \quad \text{(the word 123 itself is Lyndon)}; \\
b_{132} &= b_{132} \quad \text{(the word 132 itself is Lyndon)}; \\
b_{213} &= b_{213} - b_{123} - b_{132}; \\
b_{231} &= b_{231} - b_{2}\shuffle b_{13} + b_{132}; \\
b_{312} &= b_{312} - b_{2}\shuffle b_{12} - b_{123} - b_{132}; \\
b_{321} &= b_{321} - b_{2}\shuffle b_{2}\shuffle b_{3} - b_{23}\shuffle b_{1} + b_{3}\shuffle b_{12} + b_{123}; \\
b_{112} &= b_{112} \quad \text{(the word 112 itself is Lyndon)}; \\
b_{121} &= b_{121} - 2b_{112}; \\
b_{1212} &= \frac{1}{2} b_{12}\shuffle b_{12} - 2b_{1122}; \\
b_{2321} &= b_{2321} - b_{1}\shuffle b_{2}\shuffle b_{3}\shuffle b_{1} - b_{1}\shuffle b_{2}\shuffle b_{3}\shuffle b_{2}\shuffle b_{1} - b_{1}\shuffle b_{2}\shuffle b_{3}\shuffle b_{4} \\
&\quad + b_{1}\shuffle b_{2}\shuffle b_{3}\shuffle b_{4} + b_{12}\shuffle b_{3}\shuffle b_{4} + b_{123}\shuffle b_{4} - b_{1234}.
\end{align*}
\]
Note that Theorem 6.63 cannot survive without the condition that $\mathbb{Q}$ be a subring of $k$. For instance, for any $v \in V$, we have $v1w = 2v$ in $Sh(V)$, which vanishes if $2 = 0$ in $k$; this stands in contrast to the fact that polynomial $k$-algebras are integral domains when $k$ itself is one. We will see that QSym is less sensitive towards the base ring in this regard (although proving that QSym is a polynomial algebra is much easier when $\mathbb{Q}$ is a subring of $k$).

In our proof of Theorem 6.63 (but not only there), we will use part (a) of the following lemma\(^{227}\), which makes proving that certain families indexed by Lyndon words generate certain $k$-algebras more comfortable:

**Lemma 6.65.** Let $A$ be a commutative $k$-algebra. Let $\mathfrak{A}$ be a totally ordered set. Let $\mathcal{L}$ be the set of all Lyndon words over the alphabet $\mathfrak{A}$. Let $b_w$ be an element of $A$ for every $w \in \mathcal{L}$. For every word $u \in \mathfrak{A}^\ast$, define an element $b_u$ of $A$ by $b_u = b_{a_1}b_{a_2} \cdots b_{a_p}$, where $(a_1, a_2, \ldots, a_p)$ is the CFL factorization of $u$.

(a) The family $(b_w)_{w \in \mathcal{L}}$ is an algebraically independent generating set of the $k$-algebra $A$ if and only if the family $(b_u)_{u \in \mathfrak{A}^\ast}$ is a basis of the $k$-module $A$.

(b) The family $(b_w)_{w \in \mathcal{L}}$ generates the $k$-algebra $A$ if and only if the family $(b_u)_{u \in \mathfrak{A}^\ast}$ spans the $k$-module $A$.

(c) Assume that the $k$-algebra $A$ is graded. Let $\text{wt} : \mathfrak{A} \to \{1, 2, 3, \ldots\}$ be any map such that for every $N \in \{1, 2, 3, \ldots\}$, the set $\text{wt}^{-1}(N)$ is finite.

For every word $w \in \mathfrak{A}^\ast$, define an element $\text{Wt}(w) \in \mathbb{N}$ by $\text{Wt}(w) = \text{wt}(w_1) + \text{wt}(w_2) + \cdots + \text{wt}(w_k)$, where $k$ is the length of $w$.

Assume that for every $w \in \mathcal{L}$, the element $b_w$ of $A$ is homogeneous of degree $\text{Wt}(w)$.

Assume further that the $k$-module $A$ has a basis $(g_u)_{u \in \mathfrak{A}^\ast}$ having the property that for every $u \in \mathfrak{A}^\ast$, the element $g_u$ of $A$ is homogeneous of degree $\text{Wt}(u)$.

Assume also that the family $(b_w)_{w \in \mathcal{L}}$ generates the $k$-algebra $A$.

Then, this family $(b_w)_{w \in \mathcal{L}}$ is an algebraically independent generating set of the $k$-algebra $A$.

**Exercise 6.66.** Prove Lemma 6.65.

[Hint: For (a) and (b), notice that the $b_u$ are the “monomials” in the $b_w$. For (c), use Exercise 2.46(b) in every homogeneous component of $A$.]

The main workhorse of our proof of Theorem 6.63 will be the following consequence of Theorem 6.36(c):

**Proposition 6.67.** Let $V$ be a free $k$-module with a basis $(b_a)_{a \in \mathfrak{A}}$, where $\mathfrak{A}$ is a totally ordered set.

For every word $w \in \mathfrak{A}^\ast$ over the alphabet $\mathfrak{A}$, let us define an element $b_w$ of $Sh(V)$ by $b_w = b_{w_1}b_{w_2} \cdots b_{w_{\ell}}$, where $\ell$ is the length of $w$. (The multiplication used here is that of $T(V)$, not that of $Sh(V)$; the latter is denoted by $\odot$.)

For every word $u \in \mathfrak{A}^\ast$, define an element $b_u$ by $b_u = b_{u_1}b_{u_2} \cdots b_{u_p}$, where $(a_1, a_2, \ldots, a_p)$ is the CFL factorization of $u$.

If $\ell \in \mathbb{N}$ and if $x \in \mathfrak{A}^\ell$ is a word, then there is a family $(\eta_{x,y})_{y \in \mathfrak{A}^\ell} \in \mathbb{N}^{\mathfrak{A}^\ell}$ of elements of $\mathbb{N}$ satisfying

$$b_x = \sum_{y \in \mathfrak{A}^\ell, y \leq x} \eta_{x,y}b_y$$

and $\eta_{x,x} \neq 0$ (in $\mathbb{N}$).

Before we prove this, let us show a very simple lemma:

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\(^{226}\)A pattern emerges in the formulas for $b_{12}$, $b_{123}$ and $b_{1234}$: for every $n \in \mathbb{N}$, we have

$$b_{\alpha(n-1,\ldots,1)} = \sum_{\alpha \in \text{Comp}_n} (-1)^{\alpha_1}(\alpha) b_{d_1(\alpha)}b_{d_2(\alpha)}\cdots b_{d_{\ell(\alpha)}(\alpha)},$$

where $(d_1(\alpha)) \cdots (d_{\ell(\alpha)}(\alpha))$ is the factorization of the word $(1, 2, \ldots, n)$ into factors of length $\alpha_1, \alpha_2, \ldots, \alpha_{\ell}$ (where $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_{\ell})$). This can be proved by an application of the principle of inclusion and exclusion (as it is easy to see that for any composition $\alpha$ of $n$, the product $b_{d_1(\alpha)}b_{d_2(\alpha)}\cdots b_{d_{\ell(\alpha)}(\alpha)}$ equals the sum of $b_{\beta}$ for all words $\beta \in \mathfrak{C}$ satisfying $\text{Des}(\beta^{-1}) \subseteq D(\alpha)$).

\(^{227}\)And in a later proof, we will also use its part (c) (which is tailored for application to QSym).
Lemma 6.68. Let \( \mathfrak{A} \) be a totally ordered set. Let \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \). Let \( \sigma \in \text{Sh}_{n,m} \).

(a) If \( u, v \) and \( v' \) are three words satisfying \( \ell (u) = n, \ell (v) = m, \ell (v') = m \) and \( v' < v \), then \( u \cdot \sigma v' < u \cdot \sigma v \).

(b) If \( u, u' \) and \( v \) are three words satisfying \( \ell (u) = n, \ell (u') = n, \ell (v) = m \) and \( u' < u \), then \( u' \cdot \sigma v < u \cdot \sigma v \).

(c) If \( u, v \) and \( v' \) are three words satisfying \( \ell (u) = n, \ell (v) = m, \ell (v') = m \) and \( v' \leq v \), then \( u \cdot \sigma v' \leq u \cdot \sigma v \).

Exercise 6.69. Prove Lemma 6.68.


[Hint: Proceed by induction over \( \ell \). In the induction step, apply Theorem 6.36(c)\( ^{228} \) to \( u = a_1 \) and \( v = a_2 a_3 \cdots a_p \), where \((a_1, a_2, \ldots, a_p)\) is the CFL factorization of \( x \). Use Lemma 6.68 to get rid of smaller terms.]


[Hint: According to Lemma 6.65(a), it suffices to show that the family \((b_w)_{w \in \mathfrak{A}}\) defined in Proposition 6.67 is a basis of the \( k \)-module \( \text{Sh}(V) \). When \( \mathfrak{A} \) is finite, the latter can be proven by triangularity using Proposition 6.67. Reduce the general case to that of finite \( \mathfrak{A} \).]

6.4. Polynomial freeness of QSym: statement and easy parts.

Definition 6.72. For the rest of Chapter 6, we introduce the following notations: We let \( \mathfrak{A} \) be the totally ordered set \( \{1, 2, 3, \ldots\} \) with its natural order (that is, \( 1 < 2 < 3 < \cdots \)). Thus, the words over \( \mathfrak{A} \) are precisely the compositions. That is, \( \mathfrak{A}^* = \text{Comp} \). We let \( \mathfrak{L} \) denote the set of all Lyndon words over \( \mathfrak{A} \). These Lyndon words are also called Lyndon compositions.

A natural question is how many Lyndon compositions of a given size exist. While we will not use the answer, we nevertheless record it:

Exercise 6.73. Show that the number of Lyndon compositions of size \( n \) equals

\[
\frac{1}{n} \sum_{d \mid n} \mu (d) \left( 2^{n/d} - 1 \right) = \frac{1}{n} \sum_{d \mid n} \mu (d) 2^{n/d} - \delta_{n,1}
\]

for every positive integer \( n \) (where "\( \sum \)" means a sum over all positive divisors of \( n \), and where \( \mu \) is the number-theoretic Möbius function).

[Hint: One solution is similar to the solution of Exercise 6.29 using CFL factorization. Another proceeds by defining a bijection between Lyndon compositions and Lyndon words over a two-letter alphabet \( \{0, 1\} \) (with \( 0 < 1 \)) which are \( \neq 1 \).\( ^{229} \)]

Let us now state Hazewinkel’s result ([64, Theorem 8.1], [68, §6.7]) which is the main goal of Chapter 6:

Theorem 6.74. The \( k \)-algebra QSym is a polynomial algebra. It is isomorphic, as a graded \( k \)-algebra, to the \( k \)-algebra \( k[x_w \mid w \in \mathfrak{L}] \). Here, the grading on \( k[x_w \mid w \in \mathfrak{L}] \) is defined by setting \( \deg (x_w) = \sum_{i=1}^{\ell(w)} w_i \) for every \( w \in \mathfrak{L} \).

We shall prove Theorem 6.74 in the next section (Section 6.5). But the particular case of Theorem 6.74 when \( Q \) is a subring of \( k \) can be proven more easily; we state it as a proposition:

Proposition 6.75. Assume that \( Q \) is a subring of \( k \). Then, Theorem 6.74 holds.

We will give two proofs of Proposition 6.75 in this Section 6.4; a third proof of Proposition 6.75 will immediately result from the proof of Theorem 6.74 in Section 6.5. (There is virtue in giving three different proofs, as they all construct different isomorphisms \( k[x_w \mid w \in \mathfrak{L}] \rightarrow \text{QSym} \).

Our first proof – originating in Malvenuto’s [109, Corollaire 4.20] – can be given right away; it relies on Exercise 5.44:

\( ^{228} \)Or Theorem 6.36(e), if you prefer.

\( ^{229} \)This bijection is obtained by restricting the bijection

\[
\text{Comp} \rightarrow \{w \in \{0, 1\}^* \mid w \text{ does not start with } 1\},
\]

\[
(a_1, a_2, \ldots, a_p) \rightarrow 01^{a_1-1}01^{a_2-1} \cdots 01^{a_p-1}
\]

(where \( 01^k \) is to be read as \( 0 \langle 1 \rangle^k \), not as \( (01)^k \)) to the set of Lyndon compositions. The idea behind this bijection is well-known in the Grothendieck-Teichmüller community: see, e.g., [69, §3.1] (and see [57, Note 5.16] for a different appearance of this idea).
**First proof of Proposition 6.75.** Let $V$ be the free $\mathbb{k}$-module with basis $(b_n)_{n \in \{1,2,3,\ldots\}}$. Endow the $\mathbb{k}$-module $V$ with a grading by assigning to each basis vector $b_n$ the degree $n$. Exercise 5.44(k) shows that QSym is isomorphic to the shuffle algebra $Sh(V)$ (defined as in Proposition 1.73) as Hopf algebras. By being a bit more careful, we can obtain the slightly stronger result that QSym is isomorphic to the shuffle algebra $Sh(V)$ as graded Hopf algebras\textsuperscript{230}. In particular, QSym $\cong Sh(V)$ as graded $\mathbb{k}$-algebras.

Theorem 6.63 (applied to $b_n = b_n$) yields that the shuffle algebra $Sh(V)$ is a polynomial $\mathbb{k}$-algebra, and that an algebraically independent generating set of $Sh(V)$ can be constructed as follows:

For every word $w \in \mathfrak{A}^*$ over the alphabet $\mathfrak{A}$, let us define an element $b_w$ of $Sh(V)$ by $b_w = b_{w_1} b_{w_2} \cdots b_{w_\ell}$, where $\ell$ is the length of $w$. (The multiplication used here is that of $T(V)$, not that of $Sh(V)$; the latter is denoted by $\shuffle$.) Then, $(b_w)_{w \in \mathfrak{A}^*}$ is an algebraically independent generating set of the $\mathbb{k}$-algebra $Sh(V)$.

For every $w \in \mathfrak{A}^*$, we have $b_w = b_{w_1} b_{w_2} \cdots b_{w_{\ell(w)}}$ (by the definition of $b_w$). For every $w \in \mathfrak{A}^*$, the element $b_w$ of the polynomial algebra $k[b_{\ell(w)}]$ is homogeneous of degree $\sum_{i=1}^{\ell(w)} \deg (b_{w_i}) = \sum_{i=1}^{\ell(w)} w_i$.

Now, define a grading on the $\mathbb{k}$-algebra $k[x_w \mid w \in \mathfrak{A}^*]$ by setting $\deg (x_w) = \sum_{i=1}^{\ell(w)} w_i$ for every $w \in \mathfrak{L}$. By the universal property of the polynomial algebra $k[x_w \mid w \in \mathfrak{A}^*]$, we can define a $\mathbb{k}$-algebra homomorphism $\Phi : k[x_w \mid w \in \mathfrak{L}] \to Sh(V)$ by setting $\Phi (x_w) = b_w$ for every $w \in \mathfrak{L}$.

This homomorphism $\Phi$ is a $\mathbb{k}$-algebra isomorphism (since $(b_w)_{w \in \mathfrak{A}^*}$ is an algebraically independent generating set of the $\mathbb{k}$-algebra $Sh(V)$) and is graded (because for every $w \in \mathfrak{L}$, the element $b_w$ of $Sh(V)$ is homogeneous of degree $\sum_{i=1}^{\ell(w)} w_i = \deg (x_w)$). Thus, $\Phi$ is an isomorphism of graded $\mathbb{k}$-algebras. Hence, $Sh(V) \cong k[x_w \mid w \in \mathfrak{L}]$ as graded $\mathbb{k}$-algebras. Altogether, QSym $\cong Sh(V) \cong k[x_w \mid w \in \mathfrak{L}]$ as graded $\mathbb{k}$-algebras. Thus, QSym is a polynomial algebra. This proves Theorem 6.74 under the assumption that $\mathfrak{F}$ be a subring of $\mathbb{k}$. In other words, this proves Proposition 6.75. \hfill $\square$

Our second proof of Proposition 6.75 comes from Hazewinkel/Gubaren/i/Kirichenko \cite{Hazewinkel/Gubaren/i/Kirichenko} (where Proposition 6.75 appears as [68, Theorem 6.5.13]). This proof will construct an explicit algebraically independent family generating the $\mathbb{k}$-algebra QSym. \textsuperscript{231} The generating set will be very unsophisticated: it will be $(M_\alpha)_{\alpha \in \mathfrak{A}^*}$, where $\mathfrak{A}$ and $\mathfrak{L}$ are as in Theorem 6.74. Here, we are using the fact that words over the alphabet $\{1,2,3,\ldots\}$ are the same thing as compositions, so, in particular, a monomial quasisymmetric function $M_\alpha$ is defined for every such word $\alpha$.

It takes a bit of work to show that this family indeed fits the bill. We begin with a corollary of Proposition 5.3 that is essentially obtained by throwing away all non-bijeective maps $f$:

**Proposition 6.76.** Let $\alpha \in \mathfrak{A}^*$ and $\beta \in \mathfrak{A}^*$. Then,

$$M_\alpha M_\beta = \sum_{\gamma \in \alpha \shuffle \beta} M_\gamma + \text{(a sum of terms of the form } M_\delta \text{ with } \delta \in \mathfrak{A}^* \text{ satisfying } \ell(\delta) < \ell(\alpha) + \ell(\beta)). \text{ }$$

**Exercise 6.77.** Prove Proposition 6.76.

[Hint: Recall what was said about the $p = \ell + m$ case in Example 5.4.]

**Corollary 6.78.** Let $\alpha \in \mathfrak{A}^*$ and $\beta \in \mathfrak{A}^*$. Then, $M_\alpha M_\beta$ is a sum of terms of the form $M_\delta$ with $\delta \in \mathfrak{A}^*$ satisfying $\ell(\delta) \leq \ell(\alpha) + \ell(\beta)$.

\textsuperscript{230}Proof. In the solution of Exercise 5.44(k), we have shown that QSym $\cong T(V)^o$ as graded Hopf algebras. But Remark 1.75(b) shows that the Hopf algebra $T(V)^o$ is naturally isomorphic to the shuffle algebra $Sh(V^o)$ as Hopf algebras; it is easy to see that the natural isomorphism $T(V)^o \to Sh(V^o)$ is graded (because it is the direct sum of the isomorphisms $(V^n)^o \to (V^n)^{\otimes n}$ over all $n \in \mathbb{N}$, and each of these isomorphisms is graded). Hence, $T(V)^o \cong Sh(V^o)$ as graded Hopf algebras. But $V^o \cong V$ as graded $\mathbb{k}$-modules (since $V$ is of finite type), and thus $Sh(V^o) \cong Sh(V)$ as graded Hopf algebras. Altogether, we obtain QSym $\cong T(V)^o \cong Sh(V^o) \cong Sh(V)$ as graded Hopf algebras, qed.

\textsuperscript{231}We could, of course, obtain such a family from our above proof as well (this is done by Malvenuto in \cite[Corollaire 4.20]{Malvenuto}), but it won’t be a very simple one.

\textsuperscript{232}The sum $\sum_{\gamma \in \alpha \shuffle \beta} M_\gamma$ ranges over the multiset $\alpha \shuffle \beta$; if an element appears several times in $\alpha \shuffle \beta$, then it has accordingly many addends corresponding to it.
**Exercise 6.79.** Prove Corollary 6.78.

We now define a partial order on the compositions of a given nonnegative integer:

**Definition 6.80.** Let \( n \in \mathbb{N} \). We define a binary relation \( \leq \) on the set \( \text{Comp}_n \) as follows: For two compositions \( \alpha \) and \( \beta \) in \( \text{Comp}_n \), we set \( \alpha \leq \beta \) if and only if \( \text{wll} \)

\[
\ell (\alpha) < \ell (\beta) \quad \text{or} \quad (\ell (\alpha) = \ell (\beta) \text{ and } \alpha \leq \beta \text{ in lexicographic order}).
\]

This binary relation \( \leq \) is the smaller-or-equal relation of a total order on \( \text{Comp}_n \); we refer to said total order as the \textit{wll-order} on \( \text{Comp}_n \), and we denote by \(<\) the smaller relation of this total order.

Notice that if \( \alpha \) and \( \beta \) are two compositions satisfying \( \ell (\alpha) = \ell (\beta) \), then \( \alpha \leq \beta \) in lexicographic order if and only if \( \alpha \leq \beta \) with respect to the relation \( \leq \) defined in Definition 6.1.

A remark about the name “wll-order” is in order. We have taken this notation from [64, Definition 6.7.14], where it is used for an extension of this order to the whole set \( \text{Comp} \). We will never use this extension, as we will only ever compare two compositions of the same integer.\(^{233}\)

We now state a fact which is similar (and plays a similar role) to Proposition 6.67:

**Proposition 6.81.** For every composition \( u \in \text{Comp} = \mathbb{A}^* \), define an element \( M_u \in \text{QSym} \) by \( M_u = M_{a_1} M_{a_2} \cdots M_{a_p} \), where \( (a_1, a_2, \ldots, a_p) \) is the CFL factorization of the word \( u \).

If \( n \in \mathbb{N} \) and if \( x \in \text{Comp}_n \), then there is a family \( \{\eta_{x,y}\}_{y \in \text{Comp}_n} \in \mathbb{N}^{\text{Comp}_n} \) of elements of \( \mathbb{N} \) satisfying

\[
M_x = \sum_{\substack{y \in \text{Comp}_n; \ y \leq x \text{ wll}}} \eta_{x,y} M_y
\]

and \( \eta_{x,x} \neq 0 \) (in \( \mathbb{N} \)).

Before we prove it, let us show the following lemma:

**Lemma 6.82.** Let \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \). Let \( u \in \text{Comp}_n \) and \( v \in \text{Comp}_m \). Let \( z \) be the lexicographically highest element of the multiset \( u \cup v \).

(a) We have \( z \in \text{Comp}_{n+m} \).

(b) There exists a positive integer \( h \) such that

\[
M_u M_v = h M_z + \left( \text{a sum of terms of the form } M_w \text{ with } w \in \text{Comp}_{n+m} \text{ satisfying } w < z \right).
\]

(c) Let \( v' \in \text{Comp}_m \) be such that \( v' \leq v \). Then,

\[
M_u M_{v'} = \left( \text{a sum of terms of the form } M_w \text{ with } w \in \text{Comp}_{n+m} \text{ satisfying } w < z \right).
\]

**Exercise 6.83.** Prove Lemma 6.82.

[Hint: For (b), set \( h \) to be the multiplicity with which the word \( z \) appears in the multiset \( u \cup v \), then use Proposition 6.76 and notice that \( M_u M_v \) is homogeneous of degree \( n + m \). For (c), use (b) for \( v' \) instead of \( v \) and notice that Lemma 6.68(a) shows that the lexicographically highest element of the multiset \( u \cup v' \) is \( z \).]

**Exercise 6.84.** Prove Proposition 6.81.

[Hint: Proceed by strong induction over \( n \). In the induction step, let \( (a_1, a_2, \ldots, a_p) \) be the CFL factorization of \( x \), and set \( u = a_1 \) and \( v = a_2 a_3 \cdots a_p \); then apply Proposition 6.81 to \( v \) instead of \( x \), and multiply the resulting equality \( M_v = \sum_{\substack{y \in \text{Comp}_m; \ y \leq v \text{ wll}}} \eta_{v,y} M_y \) with \( M_u \) to obtain an expression for \( M_u M_v = M_x \). Use

\(^{233}\) In [64, Definition 6.7.14], the name “wll-order” is introduced as an abbreviation for “weight first, then length, then lexicographic” (in the sense that two compositions are first compared by their weights, then, if the weights are equal, by their lengths, and finally, if the lengths are also equal, by the lexicographic order). For us, the alternative explanation “word length, then lexicographic” serves just as well.
Lemma 6.82 to show that this expression has the form $\sum_{y \in \text{Comp}_n} \eta_{x,y} M_y$ with $\eta_{x,x} \neq 0$; here it helps to remember that the lexicographically highest element of the multiset $u \cup v$ is $uv = x$ (by Theorem 6.36(c)).]

We are almost ready to give our second proof of Proposition 6.75; our last step is the following proposition:

**Proposition 6.85.** Assume that $Q$ is a subring of $k$. Then, $(M_w)_{w \in \mathcal{G}}$ is an algebraically independent generating set of the $k$-algebra $Q\text{Sym}$.

**Exercise 6.86.** Prove Proposition 6.85.

[Hint: Define $M_u$ for every $u \in \text{Comp}$ as in Proposition 6.81. Conclude from Proposition 6.81 that, for every $n \in \mathbb{N}$, the family $(M_u)_{u \in \text{Comp}_n}$ expands upper-triangularly (with respect to the total order $\leq$ on $\text{Comp}_n$) with respect to the basis $(M_u)_{u \in \text{Comp}_n}$ of $Q\text{Sym}$, and the diagonal entries of this expansion are $\neq 0$. Conclude that this family $(M_u)_{u \in \text{Comp}_n}$ is a basis of $Q\text{Sym}_n$ itself, and so the whole family $(M_u)_{u \in \text{Comp}}$ is a basis of $Q\text{Sym}$. Conclude using Lemma 6.65(a).]

**Second proof of Proposition 6.75.** Proposition 6.85 yields that $(M_w)_{w \in \mathcal{G}}$ is an algebraically independent generating set of the $k$-algebra $Q\text{Sym}$.

Define a grading on the $k$-algebra $k \left[ x_w \mid w \in \mathcal{G} \right]$ by setting $\deg (x_w) = \sum_{i=1}^{\ell(w)} w_i$ for every $w \in \mathcal{G}$. By the universal property of the polynomial algebra $k \left[ x_w \mid w \in \mathcal{G} \right]$, we can define a $k$-algebra homomorphism $\Phi : k \left[ x_w \mid w \in \mathcal{G} \right] \to Q\text{Sym}$ by setting

$$\Phi (x_w) = M_w \quad \text{for every } w \in \mathcal{G}.$$ 

This homomorphism $\Phi$ is a $k$-algebra isomorphism (since $(M_w)_{w \in \mathcal{G}}$ is an algebraically independent generating set of the $k$-algebra $Q\text{Sym}$) and is graded (because for every $w \in \mathcal{G}$, the element $M_w$ of $Q\text{Sym}$ is homogeneous of degree $|w| = \sum_{i=1}^{\ell(w)} w_i = \deg(x_w)$). Thus, $\Phi$ is an isomorphism of graded $k$-algebras. Hence, $Q\text{Sym} \cong k \left[ x_w \mid w \in \mathcal{G} \right]$ as graded $k$-algebras. In particular, this shows that $Q\text{Sym}$ is a polynomial algebra. This proves Theorem 6.74 under the assumption that $Q$ be a subring of $k$. Proposition 6.75 is thus proven again. 

\[\Box\]

### 6.5. Polynomial freeness of $Q\text{Sym}$: the general case.

We now will prepare for proving Theorem 6.74 without any assumptions on $k$. In our proof, we follow [64] and [68, §6.7], but without using the language of plethysm and Frobenius maps. We start with the following definition:

**Definition 6.87.** Let $\alpha$ be a composition. Write $\alpha$ in the form $\alpha = (\alpha_1, \alpha_2, \ldots, \alpha_{\ell})$ with $\ell = \ell(\alpha)$.

(a) Let $\text{SIS}(\ell)$ denote the set of all strictly increasing $\ell$-tuples $(i_1, i_2, \ldots, i_\ell)$ of positive integers.\(^{234}\) For every $\ell$-tuple $i = (i_1, i_2, \ldots, i_\ell) \in \text{SIS}(\ell)$, we denote the monomial $x_{i_1}^{\alpha_1} x_{i_2}^{\alpha_2} \cdots x_{i_\ell}^{\alpha_\ell}$ by $x_i^\alpha$. This $x_i^\alpha$ is a monomial of degree $\alpha_1 + \alpha_2 + \cdots + \alpha_{\ell} = |\alpha|$. Then,

$$M_{\alpha} = \sum_{i \in \text{SIS}(\ell)} x_i^\alpha.$$ 

\[234\]“Strictly increasing” means that $i_1 < i_2 < \cdots < i_\ell$ here. Of course, the elements of $\text{SIS}(\ell)$ are in 1-to-1 correspondence with $\ell$-element subsets of $\{1, 2, 3, \ldots\}$.

\[235\]Proof of (6.24): By the definition of $M_{\alpha}$, we have

$$M_{\alpha} = \sum_{\begin{subarray}{c}i_1 < i_2 < \cdots < i_\ell \text{ in } \{1, 2, 3, \ldots\} \end{subarray}} x_{i_1}^{\alpha_1} x_{i_2}^{\alpha_2} \cdots x_{i_\ell}^{\alpha_\ell} = \sum_{(i_1, i_2, \ldots, i_\ell) \in \text{SIS}(\ell)} x_{i_1}^{\alpha_1} x_{i_2}^{\alpha_2} \cdots x_{i_\ell}^{\alpha_\ell} = \sum_{(i_1, i_2, \ldots, i_\ell) \in \text{SIS}(\ell)} x_{(i_1, i_2, \ldots, i_\ell)}^{\alpha_1 \alpha_2 \cdots \alpha_\ell} = \sum_{i = (i_1, i_2, \ldots, i_\ell) \in \text{SIS}(\ell)} x_i^\alpha = \sum_{i \in \text{SIS}(\ell)} x_i^\alpha,$$

qed.
(b) Consider the ring $\mathbf{k}[[x]]$ endowed with the coefficientwise topology\textsuperscript{236}. The family $\left(x^i\right)_{i \in \text{SIS}(t)}$ of elements of $\mathbf{k}[[x]]$ is power-summable\textsuperscript{237}. Hence, for every $f \in \Lambda$, there is a well-defined power series $f \left(\left(x^i\right)_{i \in \text{SIS}(t)}\right) \in \mathbf{k}[[x]]$ obtained by “evaluating” $f$ at $\left(x^i\right)_{i \in \text{SIS}(t)}$ \textsuperscript{238}. In particular, for every $s \in \mathbb{N}$, we can evaluate the symmetric function $e_s \in \Lambda$ \textsuperscript{239} at $\left(x^i\right)_{i \in \text{SIS}(t)}$. The resulting power series $e_s \left(\left(x^i\right)_{i \in \text{SIS}(t)}\right) \in \mathbf{k}[[x]]$ will be denoted $M^{(s)}_\alpha$. Thus,

$$M^{(s)}_\alpha = e_s \left(\left(x^i\right)_{i \in \text{SIS}(t)}\right).$$

The power series $M^{(s)}_\alpha$ are the power series $e_s (\alpha)$ in \textcite{68}. We will shortly (in Corollary 6.94(a)) see that $M^{(s)}_\alpha \in Q\text{Sym}$ (although this is also easy to prove by inspection). Here are some examples of $M^{(s)}_\alpha$:

\textsuperscript{236}This topology is defined as follows:

We endow the ring $\mathbf{k}$ with the discrete topology. Then, we can regard the $\mathbf{k}$-module $\mathbf{k}[[x]]$ as a direct product of infinitely many copies of $\mathbf{k}$ (by identifying every power series in $\mathbf{k}[[x]]$ with the family of its coefficients). Hence, the product topology is a well-defined topology on $\mathbf{k}[[x]]$; this topology is denoted as the \textit{coefficientwise topology}. A sequence $(a_n)_{n \in \mathbb{N}}$ of power series converges to a power series $a$ with respect to this topology if and only if for every monomial $m$, all sufficiently high $n \in \mathbb{N}$ satisfy $(\text{the coefficient of } m \text{ in } a_n) = (\text{the coefficient of } m \text{ in } a)$.

Note that this is not the topology obtained by taking the completion of $\mathbf{k}[[x_1,x_2,x_3,\ldots]]$ with respect to the standard grading (in which all $x_i$ have degree 1). (The latter completion is actually a smaller ring than $\mathbf{k}[[x]]$.)

\textsuperscript{237}Let us define what “power-summable” means for us:

A family $(n_i)_{i \in I} \in \mathbb{N}^I$ (where $I$ is some set) is said to be \textit{finitely supported} if all but finitely many $i \in I$ satisfy $n_i = 0$. If $(n_i)_{i \in I} \in \mathbb{N}^I$ is a finitely supported family, then $\sum_{i \in I} n_i$ is a well-defined element of $\mathbb{N}$. If $N \in \mathbb{N}$, then a family $(n_i)_{i \in I} \in \mathbb{N}^I$ will be called $(\leq N)$-\textit{supported} if it is finitely supported and satisfies $\sum_{i \in I} n_i \leq N$.

We say that a family $(s_i)_{i \in I} \in \mathbb{N}^I$ of elements of a topological commutative $\mathbf{k}$-algebra $R$ is \textit{power-summable} if it satisfies the following property: For every $N \in \mathbb{N}$, the sum

$$\sum_{(n_i)_{i \in I} \in \mathbb{N}^I, (n_i)_{i \in I} \leq (\leq N)\text{-}\text{supported}} \sum_{(n_i)_{i \in I} \in \mathbb{N}^I, (n_i)_{i \in I} \leq (\leq N)\text{-}\text{supported}} \alpha(n_i)_{i \in I} \prod_{i \in I} s^j_i$$

converges in the topology on $R$ for every choice of scalars $\alpha(n_i)_{i \in I} \in \mathbf{k}$ corresponding to all $(\leq N)$-supported $(n_i)_{i \in I} \in \mathbb{N}^I$. In our specific case, we consider $\mathbf{k}[[x]]$ as a topological commutative $\mathbf{k}$-algebra, where the topology is the coefficientwise topology. The fact that the family $\left(x^i\right)_{i \in \text{SIS}(t)}$ is power-summable then can be proven as follows:

- If $\alpha \neq \emptyset$, then this fact follows from the (easily-verified) observation that every given monomial in the variables $x_1,x_2,x_3,\ldots$ can be written as a product of monomials of the form $x^i$ (with $i \in \text{SIS}(t)$) in only finitely many ways.
- If $\alpha = \emptyset$, then this fact follows by noticing that $\left(x^i\right)_{i \in \text{SIS}(t)}$ is a finite family (indeed, $\text{SIS}(t) = \text{SIS}(0) = \{(0)\}$, and every finite family is power-summable.

\textsuperscript{238}Here is how this power series $f \left(\left(x^i\right)_{i \in \text{SIS}(t)}\right)$ is formally defined:

Let $R$ be any topological commutative $\mathbf{k}$-algebra, and let $(s_i)_{i \in I} \in R^I$ be any power-summable family of elements of $R$. Assume that the indexing set $I$ is countably infinite, and fix a bijection $j : \{1, 2, 3, \ldots\} \rightarrow I$. Let $g \in R(\mathbf{k})$ be arbitrary. Then, we can substitute $s_{j(1)}, s_{j(2)}, s_{j(3)}, \ldots$ for the variables $x_1, x_2, x_3, \ldots$ in $g$, thus obtaining an infinite sum which converges in $R$ (in fact, its convergence follows from the fact that the family $(s_i)_{i \in I} \in R^I$ is power-summable). The value of this sum will be denoted by $g \left((s_i)_{i \in I}\right)$. In general, this value depends on the choice of the bijection $j$, so the notation $g \left((s_i)_{i \in I}\right)$ is unambiguous only if this bijection $j$ is chosen once and for all. However, when $g \in \Lambda$, one can easily see that the choice of $j$ has no effect on $g \left((s_i)_{i \in I}\right)$.

We can still define $g \left((s_i)_{i \in I}\right)$ when the set $I$ is finite instead of being countably infinite. In this case, we only need to modify our above definition as follows: Instead of fixing a bijection $j : \{1, 2, 3, \ldots\} \rightarrow I$, we now fix a bijection $j : \{1, 2, 3, \ldots, |I|\} \rightarrow I$, and instead of substituting $s_{j(1)}, s_{j(2)}, s_{j(3)}, \ldots$ for the variables $x_1, x_2, x_3, \ldots$ in $g$, we now substitute $s_{j(1)}, s_{j(2)}, \ldots, s_{j(|I|)}$, $0, 0, 0, \ldots$ for the variables $x_1, x_2, x_3, \ldots$ in $g$. Again, the same observations hold as before: $g \left((s_i)_{i \in I}\right)$ is independent on $j$ if $g \in \Lambda$.

Hence, $g \left((s_i)_{i \in I}\right)$ is well-defined for every $g \in R(\mathbf{x})$, every countable (i.e., finite or countably infinite) set $I$, every topological commutative $\mathbf{k}$-algebra $R$ and every power-summable family $(s_i)_{i \in I} \in R^I$ of elements of $R$, as long as a bijection $j$ is chosen. In particular, we can apply this to $g = f$, $\mathbf{I} = \text{SIS}(t)$, $R = \mathbf{k}[[x]]$ and $(s_i)_{i \in I} = \left(x^i\right)_{i \in \text{SIS}(t)}$ choosing $j$ to be the bijection which sends every positive integer $k$ to the $k$-th smallest element of $\text{SIS}(t)$ in the lexicographic order. (Of course, since $f \in \Lambda$, the choice of $j$ is irrelevant.)

\textsuperscript{239}Recall that $e_0 = 1$, and that $e_s = 0$ for $s < 0$. 
Example 6.88. If \( \alpha \) is a composition and \( \ell \) denotes its length \( \ell (\alpha) \), then
\[
M^{(0)}_{\alpha} = \frac{e_0}{\ell!} \left( \sum_{i \in SIS(\ell)} (x_i^\alpha) \right) = 1 \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right) = 1
\]
and
\[
M^{(1)}_{\alpha} = e_1 \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right) = \sum_{i \in SIS(\ell)} x_i^\alpha = M_{\alpha}
\]
(by (6.24))
and
\[
M^{(2)}_{\alpha} = e_2 \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right) = \sum_{i < j \in SIS(\ell)} x_i^\alpha x_j^\alpha
\]
(where the notation "\( i < j \)" should be interpreted with respect to an arbitrary but fixed total order on the set \( SIS(\ell) \) – for example, the lexicographic order). Applying the last of these three equalities to \( \alpha = (2, 1) \), we obtain
\[
M^{(2)}_{(2, 1)} = \sum_{(i, j) \in SIS(2): \ell (i, j) < \ell (j, j)} x_i^2 x_j^2 + \sum_{i < j} x_i^2 x_j^2 = M_{(2, 1, 1)} + M_{(2, 3, 1)} + 2M_{(2, 2, 1)} + M_{(2, 2, 2)} = M_{(4, 1, 1)}
\]
(here, we have WLOG assumed that the order on \( SIS(2) \) is lexicographic)
\[
= M_{(2, 1, 2, 1)} + M_{(2, 3, 1)} + 2M_{(2, 2, 1, 1)} + M_{(2, 2, 2)} + M_{(4, 1, 1)}.
\]
Of course, every negative integer \( s \) satisfies
\[
M^{(s)}_{\alpha} = e_s \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right) = 0.
\]

There is a determinantal formula for the \( s! M^{(s)}_{\alpha} \) (and thus also for \( M^{(s)}_{\alpha} \) when \( s! \) is invertible in \( k \), but in order to state it, we need to introduce one more notation:

Definition 6.89. Let \( \alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell) \) be a composition, and let \( k \) be a positive integer. Then, \( \alpha \{ k \} \) will denote the composition \( (k\alpha_1, k\alpha_2, \ldots, k\alpha_\ell) \). Clearly, \( \ell (\alpha \{ k \}) = \ell (\alpha) \) and \( |\alpha \{ k \}| = k |\alpha| \).

Exercise 6.90. Let \( \alpha \) be a composition. Write the composition \( \alpha \) in the form \( \alpha = (\alpha_1, \alpha_2, \ldots, \alpha_\ell) \) with \( \ell = \ell (\alpha) \).

(a) Show that the \( s \)-th power-sum symmetric function \( p_s \in \Lambda \)
\[
p_s \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right) = M^{(s)}_{\alpha}
\]
for every positive integer \( s \).

(b) Let us fix a total order on the set \( SIS(\ell) \) (for example, the lexicographic order). Show that the \( s \)-th elementary symmetric function \( e_s \in \Lambda \) satisfies
\[
M^{(s)}_{\alpha} = e_s \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right) = \sum_{(i_1, i_2, \ldots, i_s) \in SIS(\ell)^s} x_{i_1}^\alpha x_{i_2}^\alpha \cdots x_{i_s}^\alpha
\]
for every \( s \in \mathbb{N} \).

(c) Let \( s \in \mathbb{N} \), and let \( n \) be a positive integer. Let \( e^{(n)}_s \) be the symmetric function \( \sum_{i_1 < i_2 < \cdots < i_s} x_{i_1}^n x_{i_2}^n \cdots x_{i_s}^n \in \Lambda \). Then, show that
\[
M^{(s)}_{\alpha \{ n \}} = e^{(n)}_s \left( \sum_{i \in SIS(\ell)} x_i^\alpha \right).
\]
\footnote{This is not completely obvious, but easy to check (see Exercise 6.90(b)).}
(d) Let \( s \in \mathbb{N} \), and let \( n \) be a positive integer. Prove that there exists a polynomial \( P \in k[z_1, z_2, z_3, \ldots] \) such that \( M_{\alpha(n)}^{(s)} = P \left( M_{\alpha(1)}, M_{\alpha(2)}, M_{\alpha(3)}, \ldots \right) \).

[**Hint:** For (a), (b) and (c), apply the definition of \( f \left( (x_i^q)_{i \in \text{SIS}(\ell)} \right) \) with \( f \) a symmetric function\(^{241}\). For (d), recall that \( \Lambda \) is generated by \( e_1, e_2, e_3, \ldots \).]

**Exercise 6.91.** Let \( s \in \mathbb{N} \). Show that the composition (1) satisfies \( M_{(1)}^{(s)} = e_s \).

**Proposition 6.92.** Let \( \alpha = (\alpha_1, \alpha_2, \ldots, \alpha_{\ell}) \) be a composition.

(a) Let \( n \in \mathbb{N} \). Define a matrix \( A_n^{(\alpha)} = \left( a_{i,j}^{(\alpha)} \right)_{i,j=1,2,\ldots,n} \) by

\[
a_{i,j}^{(\alpha)} = \begin{cases} M_{\alpha(i-j+1)}, & \text{if } i \geq j; \\ i, & \text{if } i = j - 1; \\ 0, & \text{if } i < j - 1 \end{cases}
\]

for all \( (i,j) \in \{1,2,\ldots,n\}^2 \).

This matrix \( A_n^{(\alpha)} \) looks as follows:

\[
A_n^{(\alpha)} = \begin{pmatrix}
M_{\alpha(1)} & 1 & 0 & \cdots & 0 & 0 \\
M_{\alpha(2)} & M_{\alpha(1)} & 2 & \cdots & 0 & 0 \\
M_{\alpha(3)} & M_{\alpha(2)} & M_{\alpha(1)} & \cdots & 0 & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
M_{\alpha(n-1)} & M_{\alpha(n-2)} & M_{\alpha(n-3)} & \cdots & M_{\alpha(1)} & n - 1 \\
M_{\alpha(n)} & M_{\alpha(n-1)} & M_{\alpha(n-2)} & \cdots & M_{\alpha(2)} & M_{\alpha(1)}
\end{pmatrix}.
\]

Then, \( \det \left( A_n^{(\alpha)} \right) = n! M_{\alpha(n)}^{(\alpha)} \).

(b) Let \( n \) be a positive integer. Define a matrix \( B_n^{(\alpha)} = \left( b_{i,j}^{(\alpha)} \right)_{i,j=1,2,\ldots,n} \) by

\[
b_{i,j}^{(\alpha)} = \begin{cases} i M_{\alpha(j)}, & \text{if } j = 1; \\
M_{\alpha(i-j+1)}, & \text{if } j > 1 \end{cases}
\]

for all \( (i,j) \in \{1,2,\ldots,n\}^2 \).

---

\(^{241}\)There are two subtleties that need to be addressed:

- the fact that the definition of \( f \left( (x_i^q)_{i \in \text{SIS}(\ell)} \right) \) distinguishes between two cases depending on whether or not \( \text{SIS}(\ell) \) is finite;
- the fact that the total order on the set \( \{1,2,3,\ldots\} \) (which appears in the summation subscript in the equality \( e_s = \sum_{i_1,i_2,\ldots,i_s \in \{1,2,3,\ldots\}^s} x_{i_1} x_{i_2} \cdots x_{i_s} \)) has nothing to do with the total order on the set \( \text{SIS}(\ell) \) (which appears in the summation subscript in \( \sum_{(i_1,i_2,\ldots,i_s) \in \text{SIS}(\ell)} x_{i_1}^\alpha x_{i_2}^\alpha \cdots x_{i_s}^\alpha \)). For instance, the former total order is well-founded, whereas the latter may and may not be. So there is (generally) no bijection between \( \{1,2,3,\ldots\} \) and \( \text{SIS}(\ell) \) preserving these orders (even if \( \text{SIS}(\ell) \) is infinite). Fortunately, this does not matter much, because the total order is only being used to ensure that every product of \( s \) pairwise distinct elements appears exactly once in the sum.
The matrix $B_{n}^{(\alpha)}$ looks as follows:

\[
B_{n}^{(\alpha)} = \begin{pmatrix}
M_{\alpha}^{(1)} & M_{\alpha}^{(0)} & M_{\alpha}^{(-1)} & \cdots & M_{\alpha}^{(1-n)} & M_{\alpha}^{(-n+2)} \\
2M_{\alpha}^{(2)} & M_{\alpha}^{(1)} & M_{\alpha}^{(0)} & \cdots & M_{\alpha}^{(-1-n)} & M_{\alpha}^{(-n+3)} \\
3M_{\alpha}^{(3)} & M_{\alpha}^{(2)} & M_{\alpha}^{(1)} & \cdots & M_{\alpha}^{(-2-n)} & M_{\alpha}^{(-n+4)} \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
(n-1)M_{\alpha}^{(n-1)} & M_{\alpha}^{(n-2)} & M_{\alpha}^{(n-3)} & \cdots & M_{\alpha}^{(1-n)} & M_{\alpha}^{(0)} \\
nM_{\alpha}^{(n)} & M_{\alpha}^{(n-1)} & M_{\alpha}^{(n-2)} & \cdots & M_{\alpha}^{(2-n)} & M_{\alpha}^{(1)} 
\end{pmatrix}
\]

Then, $\det(B_{n}^{(\alpha)}) = \mu_{\alpha} \langle n \rangle$.

**Exercise 6.93.** Prove Proposition 6.92.

[Hint: Substitute $(xi)_{i \in \text{SIS}(\ell)}$ for the variable set in Exercise 2.89, and recall Exercise 6.90(a).]

**Corollary 6.94.** Let $\alpha$ be a composition. Let $s \in \mathbb{Z}$.

(a) We have $M_{\alpha}^{(s)} \in \text{QSym}$.
(b) We have $M_{\alpha}^{(s)} \in \text{QSym}_{s[\alpha]}$.

**Exercise 6.95.** Prove Corollary 6.94.

We make one further definition:

**Definition 6.96.** Let $\alpha$ be a nonempty composition. Then, we denote by $\gcd \alpha$ the greatest common divisor of the parts of $\alpha$. (For instance, $\gcd(8,6,4) = 2$.) We also define $\text{red} \alpha$ to be the composition

\[
\left(\frac{\alpha_1}{\gcd \alpha}, \frac{\alpha_2}{\gcd \alpha}, \ldots, \frac{\alpha_\ell}{\gcd \alpha}\right),
\]

where $\alpha$ is written in the form $(\alpha_1, \alpha_2, \ldots, \alpha_\ell)$.

We say that a nonempty composition $\alpha$ is reduced if $\gcd \alpha = 1$.

We define $\mathcal{RL}$ to be the set of all reduced Lyndon compositions. In other words, $\mathcal{RL} = \{w \in \mathcal{L} \mid w \text{ is reduced}\}$ (since $\mathcal{L}$ is the set of all Lyndon compositions).

Hazewinkel, in [68, proof of Thm. 6.7.5], denotes $\mathcal{RL}$ by $\epsilon \text{LYN}$, calling reduced Lyndon compositions “elementary Lyndon words”.

**Remark 6.97.** Let $\alpha$ be a nonempty composition.

(a) We have $\alpha = (\text{red} \alpha) \{\gcd \alpha\}$.
(b) The composition $\alpha$ is Lyndon if and only if the composition $\text{red} \alpha$ is Lyndon.
(c) The composition $\text{red} \alpha$ is reduced.
(d) If $\alpha$ is reduced, then $\text{red} \alpha = \alpha$.
(e) If $s \in \{1, 2, 3, \ldots\}$, then the composition $\alpha \{s\}$ is nonempty and satisfies $\text{red} (\alpha \{s\}) = \text{red} \alpha$ and $\gcd (\alpha \{s\}) = s \gcd \alpha$.
(f) We have $(\gcd \alpha) |\text{red} \alpha| = |\alpha|$.

**Exercise 6.98.** Prove Remark 6.97.

Our goal in this section is now to prove the following result of Hazewinkel:

**Theorem 6.99.** The family \( M_{w}^{(s)} \) \((w,s) \in \mathcal{RL} \times \{1,2,3,\ldots\} \) is an algebraically independent generating set of the $k$-algebra $\text{QSym}$. 
This will (almost) immediately yield Theorem 6.74.

Our first step towards proving Theorem 6.99 is the following observation:

**Lemma 6.100.** The family \( \left( M^{(s)}_{\alpha} \right)_{(w,s)\in \mathbb{N} \times \{1,2,3,\ldots\}} \) is a reindexing of the family \( \left( M^{(\gcd \alpha)}_{\text{red} \alpha} \right)_{\alpha \in \mathcal{L}} \).

**Exercise 6.101.** Prove Lemma 6.100.

Next, we show a lemma:

**Lemma 6.102.** Let \( \alpha \) be a nonempty composition. Let \( s \in \mathbb{N} \). Then,

\[
\sum_{M_{\alpha} \in \text{Comp}_{s} \setminus \{M_{\alpha} \}} k M_{\beta}.
\]

(That is, \( s!M^{(s)}_{\alpha} - M^{s}_{\alpha} \) is a \( k \)-linear combination of terms of the form \( M_{\beta} \) with \( \beta \) ranging over the compositions of \( s | \alpha | \) satisfying \( \ell(\beta) \leq (s - 1) \ell(\alpha) \).

**Exercise 6.103.** Prove Lemma 6.102.

[Hint: There are two approaches: One is to apply Proposition 6.92(a) and expand the determinant; the other is to argue which monomials can appear in \( s!M^{(s)}_{\alpha} - M^{s}_{\alpha} \).]

We now return to studying products of monomial quasisymmetric functions:

**Lemma 6.104.** Let \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \). Let \( u \in \text{Comp}_{n} \) and \( v \in \text{Comp}_{m} \). Let \( z \) be the lexicographically highest element of the multiset \( u \cup v \). Let \( h \) be the multiplicity with which \( z \) appears in \( u \cup v \). Then,

\[
M_{u}M_{v} = h M_{z} + \left( \text{a sum of terms of the form } M_{w} \text{ with } w \in \text{Comp}_{n+m} \text{ satisfying } w \prec z \right).
\]

**Proof of Lemma 6.104.** Lemma 6.104 was shown during the proof of Lemma 6.82(b). \( \square \)

**Corollary 6.105.** Let \( n \in \mathbb{N} \) and \( m \in \mathbb{N} \). Let \( u \in \text{Comp}_{n} \) and \( v \in \text{Comp}_{m} \). Regard \( u \) and \( v \) as words in \( \mathbb{A}^{*} \). Assume that \( u \) is a Lyndon word. Let \( (b_{1}, b_{2}, \ldots, b_{q}) \) be the CFL factorization of the word \( v \).

Assume that \( u \geq b_{j} \) for every \( j \in \{1,2,\ldots,q\} \). Let

\[
h = 1 + |\{ j \in \{1,2,\ldots,q\} \mid b_{j} = u \}|.
\]

Then,

\[
M_{u}M_{v} = h M_{uv} + \left( \text{a sum of terms of the form } M_{w} \text{ with } w \in \text{Comp}_{n+m} \text{ satisfying } w \prec uv \right).
\]


[Hint: Apply Lemma 6.104, and notice that \( u v \) is the lexicographically highest element of the multiset \( u \cup v \) (by Theorem 6.36(e)), and that \( h \) is the multiplicity with which this word \( u v \) appears in the multiset \( u \cup v \) (this is a rewriting of Theorem 6.36(e)).]

**Corollary 6.107.** Let \( k \in \mathbb{N} \) and \( s \in \mathbb{N} \). Let \( x \in \text{Comp}_{k} \) be such that \( x \) is a Lyndon word. Then:

(a) The lexicographically highest element of the multiset \( x \cup x^{s} \) is \( x^{s+1} \).

(b) We have

\[
M_{x}M_{x^{s}} = (s+1) M_{x^{s+1}} + \left( \text{a sum of terms of the form } M_{w} \text{ with } w \in \text{Comp}_{(s+1)k} \text{ satisfying } w \prec x^{s+1} \right).
\]

(c) Let \( t \in \text{Comp}_{sk} \) be such that \( t \prec x^{s} \). Then,

\[
M_{x}M_{t} = \left( \text{a sum of terms of the form } M_{w} \text{ with } w \in \text{Comp}_{(s+1)k} \text{ satisfying } w \prec x^{s+1} \right).
\]

\[\text{The following equality makes sense because we have } z \in \text{Comp}_{n+m} \text{ (by Lemma 6.82(a)).}\]

[Hint: Notice that \((x, x, \ldots, x)\) is the CFL factorization of the word \(x^s\). Now, part (a) of Corollary 6.107 follows from Theorem 6.36(c), part (b) follows from Corollary 6.105, and part (c) from Lemma 6.82(c) (using part (a)).]

Corollary 6.109. Let \(n \in \mathbb{N}\) and \(m \in \mathbb{N}\). Let \(u \in \text{Comp}_n\) and \(v \in \text{Comp}_m\). Regard \(u\) and \(v\) as words in \(\mathbb{A}^*\). Let \((a_1, a_2, \ldots, a_p)\) be the CFL factorization of \(u\). Let \((b_1, b_2, \ldots, b_q)\) be the CFL factorization of the word \(v\). Assume that \(a_i > b_j\) for every \(i \in \{1, 2, \ldots, p\}\) and \(j \in \{1, 2, \ldots, q\}\). Then,

\[
M_u M_v = M_{uv} + \left( \text{a sum of terms of the form } M_w \text{ with } w \in \text{Comp}_{n+m} \text{ satisfying } w < uv \right).
\]


[Hint: Combine Lemma 6.104 with the parts (c) and (d) of Theorem 6.36.]

Corollary 6.111. Let \(n \in \mathbb{N}\). Let \(u \in \text{Comp}_n\) be a nonempty composition. Regard \(u\) as a word in \(\mathbb{A}^*\). Let \((a_1, a_2, \ldots, a_p)\) be the CFL factorization of \(u\). Let \(k \in \{1, 2, \ldots, p - 1\}\) be such that \(a_k > a_{k+1}\). Let \(x\) be the word \(a_1 a_2 \cdots a_k\), and let \(y\) be the word \(a_{k+1} a_{k+2} \cdots a_p\). Then,

\[
M_u = M_x M_y - \left( \text{a sum of terms of the form } M_w \text{ with } w \in \text{Comp}_n \text{ satisfying } w < u \right).
\]


[Hint: Apply Corollary 6.109 to \(x, y, |x|, |y|, k, p - k, (a_1, a_2, \ldots, a_k)\) and \((a_{k+1}, a_{k+2}, \ldots, a_p)\) instead of \(u, v, n, m, p, q, (a_1, a_2, \ldots, a_p)\) and \((b_1, b_2, \ldots, b_q)\); then, notice that \(xy = u\) and \(|x| + |y| = n\).]

Corollary 6.113. Let \(k \in \mathbb{N}\). Let \(x \in \text{Comp}_k\) be a composition. Assume that \(x\) is a Lyndon word. Let \(s \in \mathbb{N}\). Then,

\[
M_x^s - s! M_{x^s} \in \sum_{w \in \text{Comp}_k; w \ll x^s} k M_w.
\]

(Recall that \(x^s\) is defined to be the word \(xx \cdots x\).)


[Hint: Rewrite the claim of Corollary 6.113 in the form \(M_x^s \in s! M_{x^s} + \sum_{w \in \text{Comp}_k; w \ll x^s} k M_w\). This can be proven by induction over \(s\), where in the induction step we need the following two observations:

1. We have \(M_x M_{x^s} \in (s + 1) M_{x^s+1} + \sum_{w \in \text{Comp}_{w \ll x^{s+1}}} k M_w\).

2. For every \(t \in \text{Comp}_k\) satisfying \(t \ll x^s\), we have \(M_x M_{t} \in \sum_{w \in \text{Comp}_{w \ll x^{s+1}}} k M_w\).

These two observations follow from parts (b) and (c) of Corollary 6.107.]

Corollary 6.115. Let \(k \in \mathbb{N}\). Let \(x \in \text{Comp}_k\) be a composition. Assume that \(x\) is a Lyndon word. Let \(s \in \mathbb{N}\). Then,

\[
M_x^{(s)} - M_{x^s} \in \sum_{w \in \text{Comp}_k; w \ll x^s} k M_w.
\]

(Recall that \(x^s\) is defined to be the word \(xx \cdots x\).)

U = QSym. To this purpose, it is enough to prove that
\[ k \in \text{Comp}_k \]
\[ \ell(\beta) \leq (s-1)\ell(x) \]

For every reduced Lyndon composition \( \alpha \) and every \( j \in \{1, 2, 3, \ldots \} \), the quasisymmetric function \( M^{(j)}_w \) is an element of the family \( \left( M^{(\gcd w)}_{\text{red}_w} \right)_{w \in \mathbb{L}} \) and thus belongs to \( U \). Combine this with Exercise 6.90(d) to see that
\[ (6.26) \quad M_\beta \in U \quad \text{for every composition } \beta. \]

It remains to prove that \( (6.26) \) holds for \( \beta = \alpha \). In other words, it remains to prove that \( M_\alpha \in U \). Let \( (a_1, a_2, \ldots, a_p) \) be the CFL factorization of the word \( \alpha \). Assume WLOG that \( p \neq 0 \) (else, all is trivial). We are in one of the following two cases:

Case 1: All of the words \( a_1, a_2, \ldots, a_p \) are equal.

Case 2: Not all of the words \( a_1, a_2, \ldots, a_p \) are equal.
In Case 2, there exists a $k \in \{1, 2, \ldots, p-1\}$ satisfying $a_k > a_{k+1}$ (since $a_1 \geq a_2 \geq \cdots \geq a_p$), and thus Corollary 6.111 (applied to $u = \alpha$, $n = |\alpha|$, $x = a_1 a_2 \cdots a_k$ and $y = a_{k+1} a_{k+1} \cdots a_p$) shows that

$$M_{\alpha} = \sum_{\substack{w \in \text{Comp}_{|\alpha|} |\alpha|}} M_w \in U \quad \text{(by the induction hypothesis)}$$

$$- \left( \sum_{\substack{w \in \text{Comp}_{|\alpha|} |\alpha|}} M_w \in U \quad \text{(by the induction hypothesis)} \right)$$

$$\in UU - (\text{a sum of terms in } U) \subset U.$$ 

Hence, it only remains to deal with Case 1. In this case, set $x = a_1 = a_2 = \cdots = a_p$. Thus, $\alpha = a_1 a_2 \cdots a_p = x^p$, whence $|\alpha| = p |x|$. But Corollary 6.115 (applied to $s = p$ and $k = |x|$) yields

$$M_{x^p} - M_x \in \sum_{\substack{w \in \text{Comp}_{|\alpha|} |\alpha|; \ w < x^p \ \text{will}}} k M_w \in U \quad \text{(since } p |x| = |\alpha| \text{ and } x^p = \alpha)$$

$$\subset \sum_{\substack{w \in \text{Comp}_{\Lambda}; \ w \leq \alpha \ \text{will}}} kU \subset U,$$

so that $M_{x^p} \in \sum_{\substack{w \in \text{Comp}_{\Lambda}; \ w \leq \alpha \ \text{will}}} U \subset U - U \subset U$. This rewrites as $M_{\alpha} \in U$ (since $\alpha = x^p$). So $M_{\alpha} \in U$ is proven in both Cases 1 and 2, and thus the induction proof of (6.26) is finished.]

**Exercise 6.118.** Prove Theorem 6.74.

Of course, this proof of Theorem 6.74 yields a new (third) proof for Proposition 6.75.

We notice the following corollary of our approach to Theorem 6.74:

**Corollary 6.119.** The $\Lambda$-algebra $QSym$ is a polynomial algebra (over $\Lambda$).

**Exercise 6.120.** Prove Corollary 6.119.

[Hint: The algebraically independent generating set $\left( M_{w}^{(s)} \right)_{(w,s) \in \mathbb{N} \times \{1,2,3,\ldots\}}$ of $QSym$ contains the elements $M_{(1)}^{(s)} = e_s \in \Lambda$ for all $s \in \{1,2,3,\ldots\}$.]
7. Aguilar-Bergeron-Sottile character theory Part I: QSym as a terminal object

It turns out that the universal mapping property of NSym as a free associative algebra leads via duality to a universal property for its dual QSym, elegantly explaining several combinatorial invariants that take the form of quasisymmetric or symmetric functions:

- Ehrenborg’s quasisymmetric function of a ranked poset [47],
- Stanley’s chromatic symmetric function of a graph [153],
- the quasisymmetric function of a matroid considered in [19].

7.1. Characters and the universal property.

Definition 7.1. Given a Hopf algebra \( A \) over \( k \), a character is an algebra morphism \( A \xrightarrow{\zeta} k \), that is,

- \( \zeta(1_A) = 1_k \),
- \( \zeta \) is \( k \)-linear, and
- \( \zeta(ab) = \zeta(a)\zeta(b) \) for \( a, b \) in \( A \).

Example 7.2. A particularly important character for \( A = \text{QSym} \) is defined as follows:

\[
\text{QSym} \xrightarrow{\zeta_Q} k \quad f(x) \mapsto f(1, 0, 0, \ldots) = [f(x)]_{x_1=1, x_2=x_3=\cdots=0}.
\]

Hence,

\[
\zeta_Q(M_\alpha) = \zeta_Q(L_\alpha) = \begin{cases} 1 & \text{if } \alpha = (n) \text{ for some } n, \\ 0 & \text{otherwise}. \end{cases}
\]

In other words, the restriction \( \zeta_Q|_{\text{QSym}_n} \) coincides with the functional \( H_n \) in \( \text{NSym}_n = \text{Hom}_k(\text{QSym}_n, k) \): one has for \( f \) in \( \text{QSym}_n \) that

\[
\zeta_Q(f) = (H_n, f).
\]

It is worth remarking that there is nothing special about setting \( x_1 = 1 \) and \( x_2 = x_3 = \cdots = 0 \): for quasisymmetric \( f \), we could have defined the same character \( \zeta_Q \) by picking any variable, say \( x_n \), and sending

\[
f(x) \mapsto [f(x)]_{x_n=1, \text{ and } x_m=0 \text{ for } m \neq n}.
\]

This character \( \text{QSym} \xrightarrow{\zeta_Q} k \) has a certain universal property.

Theorem 7.3. A connected graded Hopf algebra \( A \) together with a character \( A \xrightarrow{\zeta} k \) induces a unique graded Hopf morphism \( A \xrightarrow{\Psi} \text{QSym} \) making this diagram commute:

\[
\begin{array}{ccc}
A & \xrightarrow{\Psi} & \text{QSym} \\
\downarrow{\zeta} & & \downarrow{\zeta_Q} \\
k & & 
\end{array}
\]

Furthermore, \( \Psi \) has this formula on elements of \( A_n \):

\[
\Psi(a) = \sum_{\alpha \in \text{Comp}_n} \zeta_\alpha(a)M_\alpha
\]

where for \( \alpha = (\alpha_1, \ldots, \alpha_\ell) \), the map \( \zeta_\alpha \) is the composite

\[
A_n \xrightarrow{\Delta^{(\ell-1)}} A^\otimes \ell \xrightarrow{\pi_\alpha} A_{\alpha_1} \otimes \cdots \otimes A_{\alpha_\ell} \xrightarrow{\zeta_\alpha^\otimes} k
\]

in which \( A^\otimes \ell \xrightarrow{\pi_\alpha} A_{\alpha_1} \otimes \cdots \otimes A_{\alpha_\ell} \) is the canonical projection.

\[\text{We are using the notation of Proposition 5.9 here, and we are still identifying QSym with QSym(\text{x}), where \text{x denotes the infinite chain } (x_1 < x_2 < \cdots).}\]
Proof. One argues that \( \Psi \) is unique, and has formula \((7.3)\), using only that \( \zeta \) is \( k \)-linear and sends 1 to 1 and that \( \Psi \) is a graded \( k \)-coalgebra map making \((7.2)\) commute. Equivalently, consider the adjoint \( k \)-algebra map\footnote{Here we are using the fact that there is a 1-to-1 correspondence between graded \( k \)-linear maps \( A \to \text{QSym} \) and graded \( k \)-linear maps \( \text{QSym}^o \to A^o \) given by \( f \mapsto f^* \), and this correspondence has the property that a given graded map \( f : A \to \text{QSym} \) is a \( k \)-coalgebra map if and only if \( f^* \) is a \( k \)-algebra map. This is a particular case of Exercise 1.67(f).} \[ \text{NSym} = \text{QSym}^o \xrightarrow{\Psi^*} A^o. \]

Commutativity of \((7.2)\) implies that for \( a \in A_n \),
\[ (\Psi^*(H_n), a) = (H_n, \Psi(a)) \overset{(7.1)}{=} \zeta_Q(\Psi(a)) = \zeta(a), \]

whereas gradedness of \( \Psi^* \) yields that \((\Psi^*(H_m), a) = 0 \) whenever \( a \in A_n \) and \( m \neq n \). In other words, \( \Psi^*(H_n) \)

is the element of \( A^o \) defined as the following functional on \( A \):
\[ \Psi^*(H_n)(a) = \begin{cases} \zeta(a) & \text{if } a \in A_n, \\ 0 & \text{if } a \in A_m \text{ for some } m \neq n. \end{cases} \tag{7.4} \]

By the universal property for \( \text{NSym} \cong k(H_1, H_2, \ldots) \) as free associative \( k \)-algebra, we see that any choice of a \( k \)-linear map \( A \xrightarrow{\zeta} k \) uniquely produces a \( k \)-algebra map \( \Psi^* : \text{QSym}^o \to A^o \), which satisfies \((7.4)\) for all \( n \geq 1 \). It is easy to see that this \( \Psi^* \) then automatically satisfies \((7.4)\) for \( n = 0 \) as well if \( \zeta \) sends 1 to 1 (it is here that we use \( \zeta(1) = 1 \) and the connectedness of \( A \)). Hence, any given \( k \)-linear map \( A \xrightarrow{\zeta} k \) sending 1 to 1 uniquely produces a \( k \)-algebra map \( \Psi^* : \text{QSym}^o \to A^o \) which satisfies \((7.4)\) for all \( n \geq 0 \). Formula \((7.3)\) follows as
\[ \Psi(a) = \sum_{\alpha \in \text{Comp}} (H_\alpha, \Psi(a)) M_\alpha \]
and for a composition \( \alpha = (\alpha_1, \ldots, \alpha_\ell) \), one has
\[ (H_\alpha, \Psi(a)) = (\Psi^*(H_\alpha), a) = (\Psi^*(H_{\alpha_1}) \cdots \Psi^*(H_{\alpha_\ell}), a) = (\Psi^*(H_{\alpha_1}) \otimes \cdots \otimes \Psi^*(H_{\alpha_\ell}), \Delta^{(\ell-1)}(a)) \overset{(7.4)}{=} (\zeta^\otimes \circ \pi_\alpha)(\Delta^{(\ell-1)}(a)) = \zeta_\alpha(a), \]

where the definition of \( \zeta_\alpha \) was used in the last equality.

We wish to show that if, in addition, \( A \) is a Hopf algebra and \( A \xrightarrow{\zeta} k \) is a character (algebra map), then \( A \xrightarrow{\Psi} \text{QSym} \) will be an algebra map, that is, the two maps \( A \otimes A \to \text{QSym} \) given by \( \Psi \circ m \) and \( m \circ (\Psi \otimes \Psi) \) coincide. To see this, consider these two diagrams having the two maps in question as the composites of their top rows:
\[ A \otimes A \xrightarrow{m} A \xrightarrow{\Psi} \text{QSym} \quad \text{and} \quad A \otimes A \xrightarrow{\Psi \otimes \Psi} \text{QSym} \otimes^2 \xrightarrow{m} \text{QSym} \]

The fact that \( \zeta, \zeta_Q \) are algebra maps makes the above diagrams commute, so that applying the uniqueness in the first part of the proof to the character \( A \otimes A \xrightarrow{\zeta \otimes \zeta} k \) proves the desired equality \( \Psi \circ m = m \circ (\Psi \otimes \Psi) \). \( \Box \)

Remark 7.4. When one assumes in addition that \( A \) is cocommutative, it follows that the image of \( \Psi \) will lie in the subalgebra \( \Lambda \subset \text{QSym} \), e.g. from the explicit formula \((7.3)\) and the fact that one will have \( \zeta_\alpha = \zeta_\beta \) whenever \( \beta \) is a rearrangement of \( \alpha \). In other words, the character \( A \xrightarrow{\zeta} k \) defined by restricting \( \zeta_Q \) to \( \Lambda \), or by
\[ \zeta_\Lambda(m_\lambda) = \begin{cases} 1 & \text{if } \lambda = (n) \text{ for some } n, \\ 0 & \text{otherwise}, \end{cases} \]

has a universal property as terminal object with respect to characters on cocommutative co- or Hopf algebras.
We close this section by discussing a well-known polynomiality and reciprocity phenomenon; see, e.g., Humper and Martin [77, Prop. 2.2], Stanley [153, §4].

**Definition 7.5.** The binomial Hopf algebra (over the commutative ring $k$) is the polynomial algebra $k \{m\}$ in a single variable $m$, with a Hopf algebra structure transported from the symmetric algebra $\text{Sym} (k^{1})$ (which is a Hopf algebra by virtue of Example 1.21, applied to $V = k^{1}$) along the isomorphism $\text{Sym} (k^{1}) \to k \{m\}$ which sends the standard basis element of $k^{1}$ to $m$. Thus the element $m$ is primitive; that is, $\Delta m = 1 \otimes m + m \otimes 1$ and $S(m) = -m$. As $S$ is an algebra anti-endomorphism by Proposition 1.32 and $k \{m\}$ is commutative, one has $S(g)(m) = g(-m)$ for all polynomials $g(m)$ in $k \{m\}$.

**Definition 7.6.** For an element $f(x)$ in $Q\text{Sym}$ and a nonnegative integer $m$, let $ps^{1}(f)(m)$ denote the element of $k$ obtained by principal specialization at $q = 1$

$$ps^{1}(f)(m) = \left[ f(x) \right]_{x_{1}=x_{2}=\cdots=x_{m}=1, \atop x_{m+1}=x_{m+2}=\cdots=0} = f(1,1,\ldots,1,0,0,\ldots).$$

**Proposition 7.7.** Assume that $\mathbb{Q}$ is a subring of $k$. The map $ps^{1}$ has the following properties.

(i) Let $f \in Q\text{Sym}$. There is a unique polynomial in $k \{m\}$ which agrees for each nonnegative integer $m$ with $ps^{1}(f)(m)$, and which, by abuse of notation, we will also denote $ps^{1}(f)(m)$. If $f$ lies in $Q\text{Sym}_{n}$, then $ps^{1}(f)(m)$ is a polynomial of degree at most $n$, taking these values on $M_{\alpha}, L_{\alpha}$ for $\alpha = (\alpha_{1}, \ldots, \alpha_{\ell})$ in $\text{Comp}_{n}$:

$$ps^{1}(M_{\alpha})(m) = \binom{m}{\ell},$$

$$ps^{1}(L_{\alpha})(m) = \binom{m-\ell+n}{n}.$$

(ii) The map $Q\text{Sym} \xrightarrow{ps^{1}} k \{m\}$ is a Hopf morphism into the binomial Hopf algebra.

(iii) For all $m$ in $\mathbb{Z}$ and $f$ in $Q\text{Sym}$ one has

$$\zeta_{Q}^{*m}(f) = ps^{1}(f)(m).$$

In particular, one also has

$$\zeta_{Q}^{*(-m)}(f) = ps^{1}(S(f))(m) = ps^{1}(f)(-m).$$

(iv) For a graded Hopf algebra $A$ with a character $A \xrightarrow{\zeta} k$, and any element $a$ in $A$, the polynomial $ps^{1}(\Psi(a))(m)$ in $k \{m\}$ has degree at most $n$, and when specialized to $m$ in $\mathbb{Z}$ satisfies

$$\zeta^{*m}(a) = ps^{1}(\Psi(a))(m).$$

**Proof.** To prove assertion (i), note that one has

$$ps^{1}(M_{\alpha})(m) = M_{\alpha}(1,1,\ldots,1,0,0,\ldots) = \sum_{1 \leq i_{1} \leq \cdots \leq i_{\ell} \leq m} [x_{i_{1}}^{\alpha_{1}} \cdots x_{i_{\ell}}^{\alpha_{\ell}}]_{x_{j}=1} = \binom{m}{\ell},$$

$$ps^{1}(L_{\alpha})(m) = L_{\alpha}(1,1,\ldots,1,0,0,\ldots) = \sum_{1 \leq i_{1} \leq \cdots \leq i_{\ell} \leq m; \atop i_{k} < \ell_{k+1} \text{ if } k \in D(\alpha)} [x_{i_{1}} \cdots x_{i_{n}}]_{x_{j}=1} = \binom{m-\ell+n}{n}.$$ 

As $\{M_{\alpha}\}_{\alpha \in \text{Comp}_{n}}$ form a basis for $Q\text{Sym}_{n}$, and $\binom{m}{\ell}$ is a polynomial function in $m$ of degree $\ell(\leq n)$, one concludes that for $f$ in $Q\text{Sym}_{n}$ one has that $ps^{1}(f)(m)$ is a polynomial function in $m$ of degree at most $n$. The polynomial giving rise to this function is unique, since infinitely many of its values are fixed.

To prove assertion (ii), note that $ps^{1}$ is an algebra morphism because it is an evaluation homomorphism. To check that it is a coalgebra morphism, it suffices to check $\Delta \circ ps^{1} = (ps^{1} \otimes ps^{1}) \circ \Delta$ on each $M_{\alpha}$ for
\( \alpha = (\alpha_1, \ldots, \alpha_\ell) \) in \( \text{Comp}_n \). Using the Vandemonde summation \( (A + B)_\ell = \sum_k \binom{A}{k} \binom{B}{\ell - k} \), one has
\[
(\Delta \circ \psi^1)(M_\alpha) = \Delta \left( \frac{m}{\ell} \right) = \left( \frac{m \otimes 1 + 1 \otimes m}{\ell} \right) = \sum_{k=0}^\ell \binom{m \otimes 1}{k} \left( \frac{1 \otimes m}{\ell - k} \right) = \sum_{k=0}^\ell \binom{m}{k} \otimes \left( \frac{m}{\ell - k} \right)
\]
while at the same time
\[
((\psi^1 \otimes \psi^1) \circ \Delta)(M_\alpha) = \sum_{k=0}^\ell \psi^1(M_{(\alpha_1, \ldots, \alpha_k)}) \otimes \psi^1(M_{(\alpha_{k+1}, \ldots, \alpha_\ell)}) = \sum_{k=0}^\ell \binom{m}{k} \otimes \left( \frac{m}{\ell - k} \right).
\]
Thus \( \psi^1 \) is a bialgebra map, and hence also a Hopf map, by Proposition 1.46(c).

For assertion (iii), first assume \( m \) lies in \( \{0, 1, 2, \ldots\} \). Since \( \zeta_Q(f) = f(1, 0, 0, \ldots) \), one has
\[
\zeta_Q^*(m)(f) = \zeta_Q^m \circ \Delta^{(m-1)} f(x) = \zeta_Q^m \left( f(x^{(1)}, x^{(2)}, \ldots, x^{(m)}) \right)
= \left[ f(x^{(1)}, x^{(2)}, \ldots, x^{(m)}) \right]_{x^{(1)} = x^{(2)} = \cdots = x^{(m)} = 1, \atop x_j^{(3)} = x_j^{(4)} = \cdots = 0 \text{ for all } j}
= f(1, 0, 0, \ldots, 1, 0, 0, \ldots, 1, 0, 0, \ldots) = f(\underbrace{1, 1, \ldots, 1, 0, 0, \ldots}_{m \text{ ones}}).
\]

But then Proposition 1.46(a) also implies
\[
\zeta_Q^{*(-m)}(f) = \left( \zeta_Q^{*(1)} \right)^m(f) = (\zeta_Q \circ S^m)(f) = \zeta_Q^m(S(f)) = \psi^1(S(f))(m) = S(\psi^1(f))(m) = \psi^1(f)(-m).
\]

For assertion (iv), note that
\[
\zeta^m(a) = (\zeta_Q \circ \Psi)^m(a) = (\zeta_Q^m)(\Psi(a)) = \psi^1(\Psi(a))(m).
\]
where the three equalities come from (7.2), Proposition 1.46(a), and assertion (iii) above, respectively. \( \square \)

Remark 7.8. Aguiar, Bergeron and Sottile give a very cute (third) proof of the QSym antipode formula Theorem 5.11, via Theorem 7.3, in [4, Example 4.8]. They apply Theorem 7.3 to the coopposite coalgebra QSym\(_{\text{cop}}\) and its character \( \zeta_Q^{*(1)} \). One can show that the map QSym\(_{\text{cop}}\) \( \psi \to \) QSym induced by \( \zeta_Q^{*(1)} \) is \( S = \psi \), the antipode of QSym, because \( S : \text{QSym} \to \text{QSym} \) is a coalgebra anti-endomorphism (by Exercise 1.47) satisfying \( \zeta_Q^{*(-1)} = \zeta_Q \circ S \). They then use the formula (7.3) for \( \Psi = S \) (together with the polynomiality Proposition 7.7) to derive Theorem 5.11.

Exercise 7.9. Show that \( \zeta_Q^m(f) = \psi^1(f)(m) \) for all \( f \in \text{QSym} \) and \( m \in \{0, 1, 2, \ldots\} \). (This was already proven in Proposition 7.7(iii); give an alternative proof using Proposition 5.7.)

7.2. Example: Ehrenborg’s quasisymmetric function of a ranked poset. Here we consider incidence algebras, coalgebras and Hopf algebras generally, and then particularize to the case of graded posets, to recover Ehrenborg’s interesting quasisymmetric function invariant via Theorem 7.3.

7.2.1. Incidence algebras, coalgebras, Hopf algebras.

Definition 7.10. Given a family \( \mathcal{P} \) of finite partially ordered sets \( P \), let \( \mathbf{k}[\mathcal{P}] \) denote the free \( \mathbf{k} \)-module whose basis consists of symbols \( [P] \) corresponding to isomorphism classes of posets \( P \) in \( \mathcal{P} \).

We will assume throughout that each \( P \) in \( \mathcal{P} \) is bounded, that is, it has a unique minimal element \( \hat{0} := \hat{0}_P \) and a unique maximal element \( \hat{1} := \hat{1}_P \). In particular, \( P \neq \emptyset \), although it is allowed that \( |P| = 1 \), so that \( \hat{0} = \hat{1} \); denote this isomorphism class of posets with one element by \( [\emptyset] \).

If \( \mathcal{P} \) is closed under taking intervals
\[
[x, y] := [x, y]_P := \{ z \in P : x \leq_P y \}
\]
247 See Exercise 7.9 for an alternative way to prove this, requiring less thought to verify its soundness.
then one can easily that the following coproduct and counit endow $k[P]$ with the structure of a coalgebra, called the \textit{(reduced) incidence coalgebra}:

$$\Delta[P] := \sum_{x \in P} [0, x] \otimes [x, \hat{1}],$$

$$\epsilon[P] := \begin{cases} 1 & \text{if } |P| = 1 \\ 0 & \text{otherwise.} \end{cases}$$

The dual algebra $k[P]^*$ is generally called the \textit{reduced incidence algebra (modulo isomorphism)} for the family $P$ (see, e.g., [142]). It contains the important element $k[P] \xrightarrow{\zeta} k$, called the \textit{$\zeta$-function} that takes the value $\zeta[P] = 1$ for all $P$.

If $P$ (is not empty and) satisfies the further property of being \textit{hereditary} in the sense that for every $P_1, P_2$ in $P$, the Cartesian product poset $P_1 \times P_2$ with componentwise partial order is also in $P$, then one can check that the following product and unit endow $k[P]$ with the structure of a (commutative) algebra:

$$[P_1] \cdot [P_2] := m([P_1] \otimes [P_2]) := [P_1 \times P_2],$$

$$1_{k[P]} := [o].$$

**Proposition 7.11.** For any hereditary family $P$ of finite posets, $k[P]$ is a bialgebra, and even a Hopf algebra with antipode $S$ given as in (1.16) (Takeuchi’s formula):

$$S[P] = \sum_{k \geq 0} (-1)^k \sum_{\hat{0} = x_0 < \cdots < x_k = \hat{1}} [x_0, x_1] \cdot \cdots \cdot [x_{k-1}, x_k].$$

**Proof.** Checking the commutativity of the pentagonal diagram in (1.8) amounts to the fact that, for any $(x_1, x_2) <_{P_1 \times P_2} (y_1, y_2)$, one has a poset isomorphism

$$[(x_1, x_2), (y_1, y_2)]_{P_1 \times P_2} \cong [x_1, y_1]_{P_1} \times [x_2, y_2]_{P_2}.$$

Commutativity of the remaining diagrams in (1.8) is straightforward, and so $k[P]$ is a bialgebra. But then Remark 1.45 implies that it is a Hopf algebra, with antipode $S$ as in (1.16), because the map $f := \text{id}_{k[P]} - u \epsilon$ (sending the class $[o]$ to $0$, and fixing all other $[P]$) is locally $*$-nilpotent:

$$f^k[P] = \sum_{\hat{0} = x_0 < \cdots < x_k = \hat{1}} [x_0, x_1] \cdot \cdots \cdot [x_{k-1}, x_k]$$

will vanish due to an empty sum whenever $k$ exceeds the maximum length of a chain in the finite poset $P$. \hfill $\square$

It is perhaps worth remarking how this generalizes the Möbius function formula of P. Hall. Note that the zeta function $k[P] \xrightarrow{\zeta} k$ is a character, that is, an algebra morphism. Proposition 1.46(a) then tells us that $\zeta$ should have a convolutional inverse $k[P] \xrightarrow{\mu := \zeta^{-1}} k$, traditionally called the \textit{Möbius function}, with the formula $\mu = \zeta^{-1} = \zeta \circ S$. Rewriting this via the antipode formula for $S$ given in Proposition 7.11 yields P. Hall’s formula.

**Corollary 7.12.** For a finite bounded poset $P$, one has

$$\mu[P] = \sum_{k \geq 0} (-1)^k |\text{chains } 0 = x_0 < \cdots < x_k = \hat{1} \text{ in } P|.$$

We can also notice that $S$ is an algebra anti-endomorphism (by Proposition 1.32), thus an algebra endomorphism (since $k[P]$ is commutative). Hence, $\mu = \zeta \circ S$ is a composition of two algebra homomorphisms, thus an algebra homomorphism itself. We therefore obtain the following classical fact:

**Corollary 7.13.** For two finite bounded posets $P$ and $Q$, we have $\mu[P \times Q] = \mu[P] \cdot \mu[Q]$. 
7.2.2. The incidence Hopf algebras for ranked posets and Ehrenborg’s function.

**Definition 7.14.** Take \( P \) to be the class of bounded ranked finite posets \( P \), that is, those for which all maximal chains from \( 
abla 0 \) to \( 
abla 1 \) have the same length \( r(P) \). This is a hereditary class, as it implies that any interval is \([x,y]_P\) is also ranked, and the product of two bounded ranked posets is also bounded and ranked. It also uniquely defines a rank function \( P \longrightarrow \mathbb{N} \) in which \( r(\nabla 0) = 0 \) and \( r(x) \) is the length of any maximal chain from \( 
abla 0 \) to \( x \).

**Example 7.15.** Consider a pyramid with apex vertex \( a \) over a square base with vertices \( b, c, d, e \):

![Diagram of a pyramid with apex vertex a over a square base with vertices b, c, d, e]

Ordering its faces by inclusion gives a bounded ranked poset \( P \), where the rank of an element is one more than the dimension of the face it represents:

**Definition 7.16.** Ehrenborg’s quasisymmetric function \( \Psi[P] \) for a bounded ranked poset \( P \) is the image of \([P]\) under the map \( k[P] \xrightarrow{\Psi} \text{QSym} \) induced by the zeta function \( k[P] \xrightarrow{\zeta} k \) as a character, via Theorem 7.3.

The quasisymmetric function \( \Psi[P] \) captures several interesting combinatorial invariants of \( P \); see Stanley [154, Chap. 3] for more background on these notions.

**Definition 7.17.** Let \( P \) be a bounded ranked poset \( P \) of rank \( r(P) := r(\nabla 1) \). Define its rank-generating function

\[
\text{RGF}(P,q) := \sum_{p \in P} q^{r(p)},
\]

its characteristic polynomial

\[
\chi(P,q) := \sum_{p \in P} \mu(\nabla 0,p)q^{r(p)}
\]

(where \( \mu(p,q) \) is shorthand for \( \mu([p,q]) \)), its zeta polynomial

\[
(7.6) \quad Z(P,m) = |\{\text{multichains } \nabla 0 \leq p_1 \leq \cdots \leq p_{m-1} \leq p \leq \nabla 1\}|
\]

\[
(7.7) \quad = \sum_{s=0}^{r(P)-1} \binom{m}{s+1} |\{\text{chains } \nabla 0 < p_1 < \cdots < p_s < \nabla 1\}|
\]
and for a subset $S \subset \{1, 2, \ldots, r(P) - 1\}$, its flag number $f_S$, as a component of its flag $f$-vector $(f_S)_{S \subset [r-1]}$
\[ f_S = |\{\text{chains } 0 < p_1 < p \cdots < p_s < p \text{ with } \{r(p_1), \ldots, r(p_s)\} = S\}|, \]
as well as the flag $h$-vector entry $h_T$ given by $f_S = \sum_{T \subset S} (-1)^{|S\setminus T|} f_T$.

**Example 7.18.** For the poset $P$ in Example 7.15, one has $\text{RGF}(P, q) = 1 + 5q + 8q^2 + 5q^3 + q^4$. Since $P$ is the poset of faces of a polytope, the Möbius function values for its intervals are easily predicted: $\mu(x, y) = (-1)^{r(x,y)}$, that is, $P$ is an Eulerian ranked poset; see Stanley [154, §3.16]. Hence its characteristic polynomial is trivially related to the rank generating function, sending $q \mapsto -q$, that is,
\[ \chi(P, q) = \text{RGF}(P, -q) = 1 - 5q + 8q^2 - 5q^3 + q^4. \]

Its flag $f$-vector and $h$-vector entries are given in the following table.

<table>
<thead>
<tr>
<th>$S$</th>
<th>$f_S$</th>
<th>$h_S$</th>
</tr>
</thead>
<tbody>
<tr>
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<td>1</td>
<td>1</td>
</tr>
<tr>
<td>${1}$</td>
<td>5</td>
<td>$5 - 1 = 4$</td>
</tr>
<tr>
<td>${2}$</td>
<td>8</td>
<td>$8 - 1 = 7$</td>
</tr>
<tr>
<td>${3}$</td>
<td>5</td>
<td>$5 - 1 = 4$</td>
</tr>
<tr>
<td>${1, 2}$</td>
<td>16</td>
<td>$16 - (5 + 8) + 1 = 4$</td>
</tr>
<tr>
<td>${1, 3}$</td>
<td>16</td>
<td>$16 - (5 + 5) + 1 = 7$</td>
</tr>
<tr>
<td>${2, 3}$</td>
<td>16</td>
<td>$16 - (5 + 8) + 1 = 4$</td>
</tr>
<tr>
<td>${1, 2, 3}$</td>
<td>32</td>
<td>$32 - (16 + 16 + 16) + (5 + 8 + 5) - 1 = 1$</td>
</tr>
</tbody>
</table>

and using (7.7), its zeta polynomial is
\[ Z(P, m) = 1 \binom{m}{1} + (5 + 8 + 5) \binom{m}{2} + (16 + 16 + 16) \binom{m}{3} + 32 \binom{m}{4} = \frac{m^2(2m - 1)(2m + 1)}{3}. \]

**Theorem 7.19.** Assume that $\mathbb{Q}$ is a subring of $\mathbb{k}$. Ehrenborg’s quasisymmetric function $\Psi[P]$ for a bounded ranked poset $P$ encodes

1. the flag $f$-vector entries $f_S$ and flag $h$-vector entries $h_S$ as its $M_\alpha$ and $L_\alpha$ expansion coefficients: \[ \Psi[P] = \sum_\alpha f_{D(\alpha)}(P) M_\alpha = \sum_\alpha h_{D(\alpha)}(P) L_\alpha, \]

2. the zeta polynomial as the specialization from Definition 7.6
\[ Z(P, m) = \psi_1^\alpha(\Psi[P]) | (m), \quad M_\alpha \neq 0, \quad \Psi[P] \text{ and } \psi_1^\alpha(\Psi[P]). \]

3. the characteristic polynomial as the specialization
\[ \chi(P, q) = \left(\psi_q \circ \chi \right) \circ \Psi[P] \]

where $Q\text{Sym} \xrightarrow{\psi_q} \mathbb{k}[q]$ maps $f(x) \mapsto f(q, 0, 0, \ldots)$.

**Proof.** In assertion (i), the expansion $\Psi[P] = \sum_\alpha f_{D(\alpha)}(P) M_\alpha$ is (7.3), since $\zeta_0^\alpha[P] = f_{D(\alpha)}(P)$. The $L_\alpha$ expansion follows by inclusion-exclusion, as $L_\alpha = \sum_{\beta \supset D(\alpha)} M_\beta$ and $f_S(P) = \sum_{T \subset S} h_T$.

Assertion (ii) is immediate from Proposition 7.7(iv), since $Z(P, m) = \zeta^m[P]$.

Assertion (iii) can be deduced from assertion (i), but it is perhaps more fun and in the spirit of things to proceed as follows. Note that $\psi_q(M_\alpha) = q^n$ for $\alpha = (n)$, and $\psi_q(M_\alpha)$ vanishes for all other $\alpha \neq (n)$ in $\text{Comp}_n$. Hence for a bounded ranked poset $P$ one has
\[ (\psi_q \circ \Psi)[P] = q^{r(P)}. \]

---

\[ ^{248} \text{Actually, (7.7) is false if } |P| = 1 \text{ (but only then). We use (7.6) to define } Z(P, m) \text{ in this case.} \]

\[ ^{249} \text{In fact, Ehrenborg defined } \Psi[P] \text{ in [37, Dehn. 4.1] via this } M_\alpha \text{ expansion, and then showed that it gave a Hopf morphism.} \]
Consequently, using (1.17) one can compute

\[ RGF(P, q) = \sum_{p \in P} q^{r(p)} \cdot 1 = \sum_{p \in P} q^{r([0,p])} \cdot \zeta[p, \hat{1}] \]

\[ = \sum_{p \in P} (\psi_q \circ \Psi)[0, p] \cdot (\zeta Q \circ \Psi)[p, \hat{1}] \]

\[ = \sum_{p \in P} (\psi_q \circ \zeta Q)(\Psi[P]) = (\psi_q \otimes \zeta Q)(\Delta \Psi[P]) \]

Similarly, for assertion (iv) first note that Proposition 7.11 and Corollary 7.12 let one calculate that

\[ (\psi_q \circ S)[P] = \sum_k (-1)^k \sum_{\hat{0}=x_0 < \cdots < x_k=1} (\psi_q \circ \Psi)([x_0, x_1]) \cdots (\psi_q \circ \Psi)([x_{k-1}, x_k]) \]

\[ = \sum_k (-1)^k \sum_{\hat{0}=x_0 < \cdots < x_k=1} q^{r(p)} = \mu(\hat{0}, \hat{1}) q^{r(P)}. \]

This is used in the penultimate equality here:

\[ ((\psi_q \circ S) \ast \zeta Q) \circ \Psi[P] \overset{(1.17)}{=} (\psi_q \circ S \circ \Psi) \ast (\zeta Q \circ \Psi)[P] = (\psi_q \circ S \circ \Psi) \ast \zeta[P] \]

\[ = \sum_{p \in P} (\psi_q \circ S \circ \Psi)[0, p] \cdot \zeta[p, \hat{1}] = \sum_{p \in P} \mu([0, p]) q^{r(p)} = \chi(P, q). \]

7.3. Example: Stanley’s chromatic symmetric function of a graph. We introduce the chromatic Hopf algebra of graphs and an associated character \( \zeta \) so that the map \( \Psi \) from Theorem 7.3 sends a graph \( G \) to Stanley’s chromatic symmetric function of \( G \). Then principal specialization \( ps^1 \) sends this to the chromatic polynomial of the graph.

7.3.1. The chromatic Hopf algebra of graphs.

Definition 7.20. The chromatic Hopf algebra (see Schmitt [144, §3.2]) \( G \) is a free \( k \)-module whose \( k \)-basis elements \([G]\) are indexed by isomorphism classes of (finite) simple graphs \( G = (V, E) \). Define for \( G_1 = (V_1, E_1), G_2 = (V_2, E_2) \) the multiplication

\[ [G_1] \cdot [G_2] := [G_1 \sqcup G_2] \]

where \([G_1 \sqcup G_2] \) denote the isomorphism class of the disjoint union, on vertex set \( V = V_1 \sqcup V_2 \) which is a disjoint union of copies of their vertex sets \( V_1, V_2 \), with edge set \( E = E_1 \sqcup E_2 \). For example,

\[ \begin{bmatrix} \bullet & \bullet & \bullet \\ \bullet & \bullet & \bullet \end{bmatrix} \cdot \begin{bmatrix} \bullet \\ \bullet \end{bmatrix} = \begin{bmatrix} \bullet & \bullet & \bullet \\ \bullet & \bullet & \bullet \end{bmatrix} \]

Thus the class \([\emptyset] \) of the empty graph \( \emptyset \) having \( V = \emptyset, E = \emptyset \) is a unit element.

Given a subset \( V' \subset V \), the subgraph induced on vertex set \( V' \) is defined as the graph \( G|_{V'} := (V', E') \) with edge set \( E' = \{ e \in E : e \in \{v_1, v_2\} \subset V' \} \). This lets one define a comultiplication

\[ \Delta[G] := \sum_{(V_1, V_2) : V_1 \sqcup V_2 = V} [G|_{V_1}] \otimes [G|_{V_2}] \]

Define a counit

\[ \epsilon[G] := \begin{cases} 1 & \text{if } G = \emptyset \\ 0 & \text{otherwise.} \end{cases} \]

Proposition 7.21. The above maps endow \( G \) with the structure of a graded connected finite type Hopf algebra over \( k \), which is both commutative and cocommutative.
Example 7.22. Here are some examples of these structure maps:

\[
\Delta \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] = 1 \otimes \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] + 2 \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] \otimes \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] + 2 \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] \otimes \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] + \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] \otimes \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] + \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] \otimes \left[ \begin{array}{c}
\bullet \\
\bullet
\end{array} \right] \otimes 1
\]

Proof of Proposition 7.21. The associativity of the multiplication and comultiplication should be clear as

\[
m^{(2)}([G_1] \otimes [G_2] \otimes [G_3]) = [G_1 \sqcup G_2 \sqcup G_3]
\]

\[
\Delta^{(2)}[G] = \sum_{(V_1,V_2,V_3); \ V = V_1 \sqcup V_2 \sqcup V_3} [G|_{V_1}] \otimes [G|_{V_2}] \otimes [G|_{V_3}].
\]

Checking the unit and counit conditions are straightforward. Commutativity of the pentagonal bialgebra diagram in (1.8) comes down to check that, given graphs \(G_1, G_2\) on disjoint vertex sets \(V_1, V_2\), when one applies to \([G_1] \otimes [G_2]\) either the composite \(\Delta \circ m\) or the composite \((m \otimes m) \circ (id \otimes T \otimes id) \circ (\Delta \otimes \Delta)\), the result is the same:

\[
\sum_{(V_1,V_2,V_3); \ V = V_1 \sqcup V_2 \sqcup V_3} [G_1|_{V_1} \sqcup G_2|_{V_2}] \otimes [G_1|_{V_2} \sqcup G_2|_{V_3}].
\]

Letting \(G_n\) be the \(k\)-span of \([G]\) having \(n\) vertices makes \(G\) a bialgebra which is graded and connected, and hence also a Hopf algebra by Proposition 1.36. Cocommutativity should be clear, and commutativity follows from the graph isomorphism \(G_1 \sqcup G_2 \cong G_2 \sqcup G_1\). Finally, \(G\) is of finite type since there are only finitely many isomorphism classes of simple graphs on \(n\) vertices for every given \(n\). \(\square\)

Remark 7.23. Humpert and Martin [77, Theorem 3.1] gave the following expansion for the antipode in the chromatic Hopf algebra, containing fewer terms than Takeuchi’s general formula (1.16): given a graph \(G = (V,E)\), one has

\[
S[G] = \sum_F (-1)^{|V| - \text{rank}(F)} \text{acyc}(G/F)[G_{V,F}].
\]

Here \(F\) runs over all subsets of edges that form flats in the graphic matroid for \(G\), meaning that if \(e = \{v,v'\}\) is an edge in \(E\) for which one has a path of edges in \(F\) connecting \(v\) to \(v'\), then \(e\) also lies in \(F\). Here \(G/F\) denotes the quotient graph in which all of the edges of \(F\) have been contracted, while \(\text{acyc}(G/F)\) denotes its number of acyclic orientations, and \(G_{V,F} := (V,F)\) as a simple graph.\(^{250}\)

Remark 7.24. In [13], Benedetti, Hallam and Machacek define a Hopf algebra of simplicial complexes, which contains \(G\) as a Hopf subalgebra (and also has \(G\) as a quotient Hopf algebra). They compute a formula for its antipode similar to (and generalizing) (7.9).

Remark 7.25. The chromatic Hopf algebra \(G\) is used in [90] and [31, §14.4] to study Vassiliev invariants of knots. In fact, a certain quotient of \(G\) (named \(F\) in [90] and \(L\) in [31, §14.4]) is shown to naturally host invariants of chord diagrams and therefore Vassiliev invariants of knots.

Remark 7.26. The \(k\)-algebra \(G\) is isomorphic to a polynomial algebra (in infinitely many indeterminates) over \(k\). Indeed, every finite graph can be uniquely written as a disjoint union of finitely many connected finite graphs (up to order). Therefore, the basis elements \([G]\) of \(G\) corresponding to connected finite graphs \(G\) are

\(^{250}\) The notation \(\text{rank}(F)\) denotes the rank of \(F\) in the graphic matroid of \(G\). We can define it without reference to matroid theory as the maximum cardinality of a subset \(F'\) of \(F\) such that the graph \(G_{V,F'}\) is acyclic. Equivalently, \(\text{rank}(F)\) is \(|V| - c(F)|\), where \(c(F)\) denotes the number of connected components of the graph \(G_{V,F}\). Thus, the equality (7.9) can be rewritten as \(S[G] = \sum_F (-1)^{c(F)} \text{acyc}(G/F)[G_{V,F}]\).

In this form, this equality is also proven in [14, Thm. 3.1].
algebraically independent in \( G \) and generate the whole \( k \)-algebra \( G \) (indeed, the disjoint unions of connected finite graphs are precisely the monomials in these elements). Thus, \( G \) is isomorphic to a polynomial \( k \)-algebra with countably many generators (one for each isomorphism class of connected finite graphs). As a consequence, for example, we see that \( G \) is an integral domain if \( k \) is an integral domain.

7.3.2. A “ribbon basis” for \( G \) and selfduality. In this subsection, we shall explore a second basis of \( G \) and a bilinear form on \( G \). This material will not be used in the rest of these notes (except in Exercise 7.44), but it is of some interest and provides an example of how a commutative cocommutative Hopf algebra can be studied.

First, let us define a second basis of \( G \), which is obtained by Möbius inversion (in an appropriate sense) from the standard basis \(( [G] )\) (\( [G] \) is an isomorphism class of finite graphs):

**Definition 7.27.** For every finite graph \( G = (V,E) \), set

\[
[G]^\# = \sum_{\substack{H = (V,E'); \\
E' \supset E^c}} (-1)^{|E' \setminus E^c|} [H] \in G,
\]

where \( E^c \) denotes the complement of the subset \( E \) in the set of all two-element subsets of \( V \). Clearly, \( [G]^\# \) depends only on the isomorphism class \( [G] \) of \( G \), not on \( G \) itself.

**Proposition 7.28.**

(a) Every finite graph \( G = (V,E) \) satisfies

\[
[G] = \sum_{\substack{H = (V,E'); \\
E \cap E = \emptyset}} [H]^\#.
\]

(b) The elements \([G]^\#\), where \([G] \) ranges over all isomorphism classes of finite graphs, form a basis of the \( k \)-module \( G \).

(c) For any graph \( H = (V,E) \), we have

\[
\Delta [H]^\# = \sum_{\substack{V = V_1 \sqcup V_2; \\
H = H_{V_1} \sqcup H_{V_2}}} [H_{V_1}]^\# \otimes [H_{V_2}]^\#.
\]

(d) For any two graphs \( H_1 = (V_1,E_1) \) and \( H_2 = (V_2,E_2) \), we have

\[
[H_1]^\# [H_2]^\# = \sum_{\substack{H = (V_1 \sqcup V_2,E); \\
H_{V_1} = H_1; \\
H_{V_2} = H_2}} [H]^\#.
\]

For example,

\[
\begin{bmatrix}
\bullet \\
\bullet \\
\bullet
\end{bmatrix}^\# = \begin{bmatrix}
\bullet \\
\bullet \\
-1
\end{bmatrix} - \begin{bmatrix}
\bullet \\
\bullet \\
-2
\end{bmatrix} + \begin{bmatrix}
\bullet \\
\bullet \\
\bullet
\end{bmatrix}.
\]

Proving Proposition 7.28 is part of Exercise 7.33 further below.

The equalities that express the elements \([G]^\#\) in terms of the elements \([H]\) (as in Definition 7.27), and vice versa (Proposition 7.28(a)), are reminiscent of the relations (5.21) and (5.20) between the bases \((R_\alpha)\) and \((H_\alpha)\) of \( \text{NSym} \). In this sense, we can call the basis of \( G \) formed by the \([G]^\#\) a “ribbon basis” of \( G \).

We now define a \( k \)-bilinear form on \( G \):
Definition 7.29. For any two graphs $G$ and $H$, let $\text{Iso}(G, H)$ denote the set of all isomorphisms from $G$ to $H$. Let us now define a $k$-bilinear form $(\cdot, \cdot) : G \times G \to k$ on $G$ by setting

$$([G]^\sharp, [H]) = |\text{Iso}(G, H)|.$$  

Proposition 7.30. The form $(\cdot, \cdot) : G \times G \to k$ is symmetric.

Again, we refer to Exercise 7.33 for a proof of Proposition 7.30.

The basis of $G$ constructed in Proposition 7.28(b) and the bilinear form $(\cdot, \cdot)$ defined in Definition 7.29 can be used to construct a Hopf algebra homomorphism from $G$ to its graded dual $G^\circ$.

Definition 7.31. For any finite graph $G$, let $\text{aut}(G)$ denote the number $|\text{Iso}(G, G)|$. Notice that this is a positive integer, since the set $\text{Iso}(G, G)$ is nonempty (it contains $\text{id}_G$).

Now, recall that the Hopf algebra $G$ is a connected graded Hopf algebra of finite type. The $n$-th homogeneous component is spanned by the $[G]$ where $G$ ranges over the graphs with $n$ vertices. Since $G$ is of finite type, its graded dual $G^\circ$ is defined. Let $([G]^\ast)[G]$ is an isomorphism class of finite graphs be the basis of $G^\circ$ dual to the basis $([G])_{[G]}$ is an isomorphism class of finite graphs of $G$. Define a $k$-linear map $\psi : G \to G^\circ$ by

$$\psi ([G]^\sharp) = \text{aut}(G) \cdot [G]^\ast$$

for every finite graph $G$.

Proposition 7.32. Consider the map $\psi : G \to G^\circ$ defined in Definition 7.31.

(a) This map $\psi$ satisfies $(\psi(a))(b) = (a, b)$ for all $a \in G$ and $b \in G$.

(b) The map $\psi : G \to G^\circ$ is a Hopf algebra homomorphism.

(c) If $Q$ is a subring of $k$, then the map $\psi$ is a Hopf algebra isomorphism $G \to G^\circ$.

Exercise 7.33. Prove Proposition 7.28, Proposition 7.30 and Proposition 7.32.

Remark 7.34. Proposition 7.32(c) shows that the Hopf algebra $G$ is self-dual when $Q$ is a subring of $k$. On the other hand, if $k$ is a field of positive characteristic, then $G$ is never self-dual. Here is a quick way to see this: The elements $[G]^\ast$ of $G^\circ$ defined in Definition 7.31 have the property that

$$([\circ]^\ast)^n = n! \sum_{[G]} [G]^\ast$$

for every $n \in \mathbb{N}$, where $\circ$ denotes the graph with one vertex. Thus, if $p$ is a prime and $k$ is a field of characteristic $p$, then $([\circ]^\ast)^p = 0$. Hence, the $k$-algebra $G^\circ$ has nilpotents in this situation. However, the $k$-algebra $G$ does not (indeed, Remark 7.26 shows that it is an integral domain whenever $k$ is an integral domain). Thus, when $k$ is a field of characteristic $p$, then $G$ and $G^\circ$ are not isomorphic as $k$-algebras (let alone as Hopf algebras).

---

251 We recall that if $G = (V, E)$ and $H = (W, F)$ are two graphs, then an isomorphism from $G$ to $H$ means a bijection $\varphi : V \to W$ such that $\varphi^*(E) = F$. Here, $\varphi^*$ denotes the map from the powerset of $V$ to the powerset of $W$ which sends every $T \subseteq V$ to $\varphi(T) \subseteq W$.

252 This is well-defined, because:

- the number $|\text{Iso}(G, H)|$ depends only on the isomorphism classes $[G]$ and $[H]$ of $G$ and $H$, but not on $G$ and $H$ themselves;
- the elements $[G]^\sharp$, where $[G]$ ranges over all isomorphism classes of finite graphs, form a basis of the $k$-module $G$ (because of Proposition 7.28(b));

253 This is well-defined, since $([G]^\sharp)[G]$ is an isomorphism class of finite graphs is a basis of the $k$-module $G$ (because of Proposition 7.28(b)).

254 To see this, observe that the tensor $[\circ] \otimes^n$ appears in the iterated coproduct $\Delta^{(n-1)}([G])$ exactly $n!$ times whenever $G$ is a graph on $n$ vertices.
### 7.3.3. Stanley’s chromatic symmetric function of a graph.

**Definition 7.35.** Stanley’s chromatic symmetric function $\Psi[G]$ for a simple graph $G = (V, E)$ is the image of $[G]$ under the map $G \xrightarrow{\psi} \text{QSym}$ induced via Theorem 7.3 from the edge-free character $G \xrightarrow{\zeta} k$ defined by

$$(7.12) \quad \zeta[G] = \begin{cases} 1 & \text{if } G \text{ has no edges, that is, } G \text{ is an independent/stable set of vertices,} \\ 0 & \text{otherwise.} \end{cases}$$

Note that, because $G$ is cocommutative, $\Psi[G]$ is symmetric and not just quasisymmetric; see Remark 7.4.

Recall that for a graph $G = (V, E)$, a (vertex-)coloring $f : V \to \{1, 2, \ldots\}$ is called proper if no edge $e = \{v, v'\}$ in $E$ has $f(v) = f(v')$.

**Proposition 7.36.** For a graph $G = (V, E)$, the symmetric function $\Psi[G]$ has the expansion \(^{255}\)

$$\Psi[G] = \sum_{\text{proper colorings } f : V \to \{1, 2, \ldots\}} x_f$$

where $x_f := \prod_{v \in V} x_{f(v)}$. In particular, its specialization from Proposition 7.6 gives the chromatic polynomial of $G$:

$$\text{ps}^1\Psi[G](m) = \chi_G(m) = |\{\text{proper colorings } f : V \to \{1, 2, \ldots, m\}\}|$$

**Proof.** The iterated coproduct $G \xrightarrow{\Delta^{(e-1)}} G^{\otimes e}$ sends

$$[G] \mapsto \sum_{(V_1, \ldots, V_e) : V = V_1 \sqcup \cdots \sqcup V_e} [G|V_1] \otimes \cdots \otimes [G|V_e]$$

and the map $\zeta^{\otimes e}$ sends each addend on the right to 1 or 0, depending upon whether each $V_i \subset V$ is a stable set or not, that is, whether the assignment of color $i$ to the vertices in $V_i$ gives a proper coloring of $G$. Thus formula (7.3) shows that the coefficient $\zeta_\alpha$ of $x_1^{\alpha_1} \cdots x_e^{\alpha_e}$ in $\Psi[G]$ counts the proper colorings $f$ in which $|f^{-1}(i)| = \alpha_i$ for each $i$. \qed

**Example 7.37.** For the complete graph $K_n$ on $n$ vertices, one has

$$\Psi[K_n] = n!e_n$$

$$\text{ps}^1(\Psi[K_n])(m) = n!e_n(1, 1, \ldots, 1) = n! \binom{m}{n} \left( \underbrace{m \text{ ones}}_{\text{m ones}} \right)$$

$$= m(m-1) \cdots (m-(n-1)) = \chi(K_n, m).$$

In particular, the single vertex graph $K_1$ has $\Psi[K_1] = e_1$, and since the Hopf morphism $\Psi$ is in particular an algebra morphism, a graph $K_1^{\otimes n}$ having $n$ isolated vertices and no edges will have $\Psi[K_1^{\otimes n}] = e_1^n$.

As a slightly more interesting example, the graph $P_3$ which is a path having three vertices and two edges will have

$$\Psi[P_3] = m(2, 1) + 6m(1, 1, 1) = e_2e_1 + 3e_3$$

One might wonder, based on the previous examples, when $\Psi[G]$ is $e$-positive, that is, when does its unique expansion in the $\{e_\lambda\}$ basis for $\Lambda$ have nonnegative coefficients? This is an even stronger assertion than $s$-positivity, that is, having nonnegative coefficients for the expansion in terms of Schur functions $\{s_\lambda\}$, since each $e_\lambda$ is $s$-positive. This weaker property fails, starting with the claw graph $K_{3,1}$, which has

$$\Psi[K_{3,1}] = s(3,1) - s(2,2) + 5s(2,1,1) + 8s(1,1,1,1).$$

On the other hand, a result of Gasharov [55, Theorem 2] shows that one at least has $s$-positivity for $\Psi[\text{inc}(P)]$ where $\text{inc}(P)$ is the incomparability graph of a poset which is $(3 + 1)$-free; we refer the reader to Stanley [153, §5] for a discussion of the following conjecture, which remains open\(^{256}\):

\(^{255}\)In fact, Stanley defined $\Psi[G]$ in [153, Defn. 2.1] via this expansion.

\(^{256}\)A recent refinement for incomparability graphs of posets which are both $(3 + 1)$- and $(2 + 2)$-free, also known as unit interval orders is discussed by Shareshian and Wachs [147].
Conjecture 7.38. For any $3 + 1$-free poset $P$, the incomparability graph $\text{inc}(P)$ has $\Psi[\text{inc}(P)]$ an $e$-positive symmetric function.

Here is another question about $\Psi[G]$: how well does it distinguish nonisomorphic graphs? Stanley gave this example of two graphs $G_1, G_2$ having $\Psi[G_1] = \Psi[G_2]$:

$$G_1 = \begin{array}{c}
\bullet \\
\bullet \\
\end{array} \quad G_2 = \begin{array}{c}
\bullet \\
\bullet \\
\end{array}$$

At least $\Psi[G]$ appears to do better at distinguishing trees, much better than its specialization, the chromatic polynomial $\chi(G, m)$, which takes the same value $m(m - 1)^{n-1}$ on all trees with $n$ vertices.

Question 7.39. Does the chromatic symmetric function (for $k = \mathbb{Z}$) distinguish trees?

It has been checked that the answer is affirmative for trees on 23 vertices or less. There are also interesting partial results on this question by Martin, Morin and Wagner [119].

We close this section with a few other properties of $\Psi[G]$ proven by Stanley which follow easily from the theory we have developed. For example, his work makes no explicit mention of the chromatic Hopf algebra $G$, and the fact that $\Psi$ is a Hopf morphism (although he certainly notes the trivial algebra morphism property $\Psi[G_1 \sqcup G_2] = \Psi[G_1] \Psi[G_2]$). One property he proves is implicitly related to $\Psi$ as a coalgebra morphism: he considers (in the case when $Q$ is a subring of $k$) the effect on $\Psi$ of the operator $\frac{\partial}{\partial p_1} : \Lambda_Q \rightarrow \Lambda_Q$ which acts by first expressing a symmetric function $f = f(p_1, p_2, \ldots)$ as a polynomial in the power sums $\{p_n\}$, and then applies $\frac{\partial}{\partial p_1}$. It is not hard to see that $\frac{\partial}{\partial p_1}$ is the same as the skewing operator $s_{(1)}^1 = p_1^1$: both act as derivations on $\Lambda_Q = \mathbb{Q}[p_1, p_2, \ldots]$, and agree in their effect on each $p_n$, in that both send $p_1 \mapsto 1$, and both annihilate $p_2, p_3, \ldots$.

Proposition 7.40. (Stanley [153, Cor. 2.12(a)]) For any graph $G = (V, E)$, one has

$$\frac{\partial}{\partial p_1} \Psi[G] = \sum_{v \in V} \Psi[G|_{V \setminus v}].$$

Proof. One first computes

$$\Delta \Psi[G] = (\Psi \otimes \Psi) \Delta[G] = \sum_{(V_1, V_2) : V = V_1 \sqcup V_2} \Psi[G|_{V_1}] \otimes \Psi[G|_{V_2}].$$

Since degree considerations force $(s_{(1)}, \Psi[G|_{V_1}]) = 0$ unless $|V_1| = 1$, in which case $\Psi[G|_{V_1}] = s_{(1)}$, one has

$$\frac{\partial}{\partial p_1} \Psi[G] = s_{(1)}^1 \Psi[G] = \sum_{(V_1, V_2) : V = V_1 \sqcup V_2} (s_{(1)}, \Psi[G|_{V_1}]) \cdot \Psi[G|_{V_2}] = \sum_{v \in V} \Psi[G|_{V \setminus v}].$$

$\square$

Definition 7.41. Given a graph $G = (V, E)$, an acyclic orientation $\Omega$ of the edges $E$ (that is, an orientation of each edge such that the resulting directed graph has no cycles), and a vertex-coloring $f : V \rightarrow \{1, 2, \ldots\}$, say that the pair $(\Omega, f)$ are weakly compatible if whenever $\Omega$ orients an edge $\{v, v'\}$ in $E$ as $v \rightarrow v'$, one has $f(v) \leq f(v')$. Note that a proper vertex-coloring $f$ of a graph $G = (V, E)$ is weakly compatible with a unique acyclic orientation $\Omega$.

Proposition 7.42. (Stanley [153, Prop. 4.1, Thm. 4.2]) The involution $\omega$ of $\Lambda$ sends $\Psi[G]$ to $\omega(\Psi[G]) = \sum_{(\Omega, f)} x_f$ in which the sum runs over weakly compatible pairs $(\Omega, f)$ of an acyclic orientation $\Omega$ and vertex-coloring $f$.

Furthermore, the chromatic polynomial $\chi_G(m)$ has the property that $(-1)^{|V|} \chi(G, -m)$ counts all such weakly compatible pairs $(\Omega, f)$ in which $f : V \rightarrow \{1, 2, \ldots, m\}$ is a vertex-$m$-coloring.
Proof. As observed above, a proper coloring \( f \) is weakly compatible with a unique acyclic orientation \( \Omega \) of \( G \). Denote by \( P_\Omega \) the poset on \( V \) which is the transitive closure of \( \Omega \), endowed with a strict labelling by integers, that is, \( i <_\Omega j \) implies \( i < Z j \). Then proper colorings \( f \) that induce \( \Omega \) are the same as \( P_\Omega \)-partitions, so that

\[
\Psi[G] = \sum_{\Omega} F_{P_\Omega}(x).
\]

Applying the antipode \( S \) and using Corollary 5.27 gives

\[
\omega(\Psi[G]) = (-1)^{|V|} S(\Psi[G]) = \sum_{\Omega} F_{P_\Omega^\text{opp}}(x) = \sum_{(\Omega,f)} x_f
\]

where in the last line one sums over weakly compatible pairs as in the proposition. The last equality comes from the fact that since each \( P_\Omega \) has been given a strict labelling, \( P_\Omega^\text{opp} \) acquires a weak (or natural) labelling, that is \( i <_{P_\Omega^\text{opp}} j \) implies \( i < Z j \).

The last assertion follows from Proposition 7.7(iii). \( \square \)

Remark 7.43. The interpretation of \( \chi(G,-m) \) in Proposition 7.42 is a much older result of Stanley \[152\]. The special case interpreting \( \chi(G,-1) \) as \( (-1)^{|V|} \) times the number of acyclic orientations of \( G \) has sometimes been called Stanley’s \(-1\)-color theorem. It also follows (via Proposition 7.7) from Humpert and Martin’s antipode formula for \( \mathcal{G} \) discussed in Remark 7.23: taking \( \zeta \) to be the character of \( \mathcal{G} \) given in (7.12),

\[
\chi(G,-1) = \zeta^*(-1)|G| = \zeta(S(G)) = \sum_F (-1)^{|V|-\text{rank}(F)} \text{acyc}(G/F) \zeta(G_{V,F}) = (-1)^{|V|}\text{acyc}(G)
\]

where the last equality uses the vanishing of \( \zeta \) on graphs that have edges, so only the \( F = \emptyset \) term survives.

Exercise 7.44. If \( V \) and \( X \) are two sets, and if \( f : V \to X \) is any map, then eqs \( f \) will denote the set

\[
\{ \{ u, u' \} \mid u \in V, u' \in V, u \neq u' \text{ and } f(u) = f(u') \}.
\]

This is a subset of the set of all two-element subsets of \( V \).

If \( G = (V,E) \) is a finite graph, then show that the map \( \Psi \) introduced in Definition 7.35 satisfies

\[
\Psi([G]^f) = \sum_{f : V \to \{1,2,3,...\}; eqs f = E} x_f,
\]

where \( x_f := \prod_{v \in V} x_{f(v)} \). Here, \([G]^f\) is defined as in Definition 7.27.

7.4. Example: The quasisymmetric function of a matroid. We introduce the matroid-minor Hopf algebra of Schmitt \[141\], and studied extensively by Crapo and Schmitt \[33, 34, 35\]. A very simple character \( \zeta \) on this Hopf algebra will then give rise, via the map \( \Psi \) from Theorem 7.3, to the quasisymmetric function invariant of matroids from the work of Billera, Jia and the second author \[19\].

7.4.1. The matroid-minor Hopf algebra. We begin by reviewing some notions from matroid theory; see Oxley \[121\] for background, undefined terms and unproven facts.

Definition 7.45. A matroid \( M \) of rank \( r \) on a (finite) ground set \( E \) is specified by a nonempty collection \( \mathcal{B}(M) \) of \( r \)-element subsets of \( E \) with the following exchange property:

For any \( B, B' \) in \( \mathcal{B}(M) \) and \( b \) in \( B \), there exists \( b' \) in \( B' \) with \( (B \setminus \{b\}) \cup \{b'\} \) in \( \mathcal{B}(M) \).

Example 7.46. A matroid \( M \) with ground set \( E \) is represented by a family of vectors \( S = (v_e)_{e \in E} \) in a vector space if \( \mathcal{B}(M) \) is the collection of subsets \( B \subseteq E \) having the property that the subfamily \( (v_e)_{e \in B} \) is a basis for the span of all of the vectors in \( S \).

For example, if \( M \) is the matroid with \( \mathcal{B}(M) = \{\{a,b\}, \{a,c\}, \{a,d\}, \{b,c\}, \{b,d\}\} \) on the ground set \( E = \{a,b,c,d\} \), then \( M \) is represented by the family \( S = (v_a, v_b, v_c, v_d) \) of the four vectors \( v_a = (1,0), v_b = (1,1), v_c = (0,1) = v_d \) in \( \mathbb{R}^2 \) depicted here:

\[
\begin{array}{c}
\text{vector space} \\
\text{vectors:} \\
\text{basis:}
\end{array}
\]

\[
\begin{array}{c}
v_c, v_d \\
v_b \\
v_a
\end{array}
\]
Conversely, whenever \( E \) is a finite set and \( S = \{v_e\}_{e \in E} \) is a family of vectors in a vector space, then the set
\[
\{ B \subseteq E : \text{the subfamily } \{v_e\}_{e \in B} \text{ is a basis for the span of all of the vectors in } S \}
\]
is a matroid on the ground set \( E \).

A matroid is said to be linear if there exists a family of vectors in a vector space representing it. Not all matroids are linear, but many important ones are.

**Example 7.47.** A special case of matroids \( M \) represented by vectors are graphic matroids, coming from a graph \( G = (V, E) \), with parallel edges and self-loops allowed. One represents these by vectors in \( \mathbb{R}^V \) with standard basis \( \{e_v\}_{v \in V} \) by associating the vector \( e_v - e_v' \) to any edge connecting a vertex \( v \) with a vertex \( v' \). One can check (or see [121, §1.2]) that the bases \( B \) in \( \mathcal{B}(M) \) correspond to the edge sets of spanning forests for \( G \), that is, edge sets which contain one spanning tree for each connected component of \( G \). For example, the matroid \( \mathcal{B}(M) \) corresponding to the graph \( G = (V, E) \) shown below:

```
\begin{tikzpicture}
  \node (a) at (0,1) {a};
  \node (b) at (1,0) {b};
  \node (c) at (1,1) {c};
  \node (d) at (2,0) {d};
  \draw (a) -- (b);
  \draw (b) -- (c);
  \draw (c) -- (d);
  \draw (d) -- (a);
\end{tikzpicture}
```

is exactly the matroid represented by the vectors in Example 7.46; indeed, the spanning trees of this graph \( G \) are the edge sets \( \{a, b\}, \{a, c\}, \{a, d\}, \{b, c\}, \{b, d\} \).

To define the matroid-minor Hopf algebra one needs the basic matroid operations of deletion and contraction. These model the operations of deleting or contracting an edge in a graph. For configurations of vectors they model the deletion of a vector, or the passage to images in the quotient space modulo the span of a vector.

**Definition 7.48.** Given a matroid \( M \) of rank \( r \) and an element \( e \) of its ground set \( E \), say that \( e \) is loop (resp. coloop) of \( M \) if \( e \) lies in no basis (resp. every basis) \( B \) in \( \mathcal{B}(M) \). If \( e \) is not a coloop, the deletion \( M \setminus e \) is a matroid of rank \( r \) on ground set \( E \setminus \{e\} \) having bases
\[
\mathcal{B}(M \setminus e) := \{ B \in \mathcal{B}(M) : e \notin B \}.
\]
If \( e \) is not a loop, the contraction \( M/e \) is a matroid of rank \( r \) on ground set \( E \setminus \{e\} \) having bases
\[
\mathcal{B}(M/e) := \{ B \setminus \{e\} : e \in B \in \mathcal{B}(M) \}.
\]
When \( e \) is a loop of \( M \), then \( M/e \) has rank \( r \) instead of \( r - 1 \) and one defines its bases as in (7.14) rather than (7.15); similarly, if \( e \) is a coloop of \( M \) then \( M \setminus e \) has rank \( r - 1 \) instead of \( r \) and one defines its bases as in (7.15) rather than (7.14).

**Example 7.49.** Starting with the graph \( G \) and its graphic matroid \( M \) from Example 7.47, the deletion \( M \setminus a \) and contraction \( M/c \) correspond to the graphs \( G \setminus a \) and \( G/c \) shown here:

```
\begin{tikzpicture}
  \node (a) at (0,1) {a};
  \node (b) at (1,0) {b};
  \node (c) at (1,1) {c};
  \node (d) at (2,0) {d};
  \draw (a) -- (b);
  \draw (b) -- (c);
  \draw (c) -- (d);
  \draw (d) -- (a);
\end{tikzpicture}
```

One has

- \( \mathcal{B}(M \setminus a) = \{\{b, c\}, \{b, d\}\} \), so that \( b \) has become a coloop in \( M \setminus a \), and
- \( \mathcal{B}(M/c) = \{\{a\}, \{b\}\} \), so that \( d \) has become a loop in \( M/c \).

**Definition 7.50.** Deletions and contractions commute with each other, leading to well-defined operations for subsets \( A \subseteq E \) of the

- restriction \( M|_A \) on ground set \( A \), obtained by deleting all \( e \in E \setminus A \) in any order, and
- quotient/contraction \( M/A \) on ground set \( E \setminus A \), obtained by contracting all \( e \in A \) in any order.
We will also need the direct sum $M_1 \oplus M_2$, whose ground set $E = E_1 \sqcup E_2$ is the disjoint union of a copy of the ground sets $E_1, E_2$ for $M_1, M_2$, and having bases

$$ \mathcal{B}(M_1 \oplus M_2) := \{ B_1 \sqcup B_2 : B_i \in \mathcal{B}(M_i) \text{ for } i = 1, 2 \}.$$

Lastly, say that two matroids $M_1, M_2$ are isomorphic if there is a bijection of their ground sets $E_1 \xrightarrow{\varphi} E_2$ having the property that $\varphi \mathcal{B}(M_1) = \mathcal{B}(M_2)$.

Now one can define the matroid-minor Hopf algebra, originally introduced by Schmitt [141, §15], and studied further by Crapo and Schmitt [33, 34, 35].

**Definition 7.51.** Let $\mathcal{M}$ have $k$-basis elements $[M]$ indexed by isomorphism classes of matroids. Define the multiplication via

$$ [M_1] \cdot [M_2] := [M_1 \oplus M_2]$$

so that the class $[\emptyset]$ of the empty matroid $\emptyset$ having empty ground set gives a unit. Define the comultiplication for $M$ a matroid on ground set $E$ via

$$ \Delta[M] := \sum_{A \subseteq E} [M|_A] \otimes [M/A],$$

and a counit

$$ \epsilon[M] := \begin{cases} 1 & \text{if } M = \emptyset \\ 0 & \text{otherwise.} \end{cases}$$

**Proposition 7.52.** The above maps endow $\mathcal{M}$ with the structure of a graded connected finite type Hopf algebra over $k$, which is commutative.

**Proof.** Checking the unit and counit conditions are straightforward. Associativity and commutativity of the multiplication follow because the direct sum operation $\oplus$ for matroids is associative and commutative up to isomorphism. Coassociativity follows because for a matroid $M$ on ground set $E$, one has this equality between the two candidates for $\Delta^{(2)}[M]$

$$ \sum_{\emptyset \subseteq A_1 \subseteq A_2 \subseteq E} [M|_{A_1}] \otimes [(M|_{A_2})/A_1] \otimes [M/A_2] = \sum_{\emptyset \subseteq A_1 \subseteq A_2 \subseteq E} [M|_{A_1}] \otimes [(M/A_1)|_{A_2 \setminus A_1}] \otimes [M/A_2]$$

due to the matroid isomorphism $(M|_{A_2})/A_1 \cong (M/A_1)|_{A_2 \setminus A_1}$. Commutativity of the bialgebra diagram in (1.8) amounts to the fact that for a pair of matroids $M_1, M_2$ and subsets $A_1, A_2$ of their (disjoint) ground sets $E_1, E_2$, one has isomorphisms

$$ M_1|_{A_1} \oplus M_2|_{A_2} \cong (M_1 \oplus M_2)|_{A_1 \sqcup A_2}, $$

$$ M_1/A_1 \oplus M_2/A_2 \cong (M_1 \oplus M_2)/(A_1 \sqcup A_2). $$

Letting $\mathcal{M}_n$ be the $k$-span of $[M]$ for matroids whose ground set $E$ has cardinality $|E| = n$, one can then easily check that $\mathcal{M}$ becomes a bialgebra which is graded, connected, and of finite type, hence also a Hopf algebra by Proposition 1.36. \qed

### 7.4.2. A quasisymmetric function for matroids.

**Definition 7.53.** Define a character $\mathcal{M} \xrightarrow{\zeta} k$ by

$$ \zeta[M] = \begin{cases} 1 & \text{if } M \text{ has only one basis,} \\ 0 & \text{otherwise.} \end{cases} $$

It is easily checked that this is a character, that is, an algebra map $\mathcal{M} \xrightarrow{\zeta} k$. Note that if $M$ has only one basis, say $\mathcal{B}(M) = \{ B \}$, then $B := \text{coloops}(M)$ is the set of coloops of $M$, and $E \setminus B = \text{loops}(M)$ is the set of loops of $M$. Equivalently, $M = \bigoplus_{e \in E} M_{|\{e\}}$ is the direct sum of matroids each having one element, each a coloop or loop.

Define $\Psi[M]$ for a matroid $M$ to be the image of $[M]$ under the map $\mathcal{M} \xrightarrow{\Psi} \text{QSym}$ induced via Theorem 7.3 from the above character $\zeta$. 

It turns out that $\Psi[M]$ is intimately related with greedy algorithms and finding minimum cost bases. A fundamental property of matroids (and one that characterizes them, in fact; see [121, §1.8]) is that no matter how one assigns costs $f : E \to \mathbb{R}$ to the elements of $E$, the following greedy algorithm (generalizing Kruskal’s algorithm for finding minimum cost spanning trees) always succeeds in finding one basis $B$ in $\mathcal{B}(M)$ achieving the minimum total cost $f(B) := \sum_{b \in B} f(b)$:

**Algorithm 7.54.** Start with the empty subset $I_0 = \emptyset$ of $E$. For $j = 1, 2, \ldots, r$, having already defined the set $I_{j-1}$, let $e$ be the element of $E \setminus I_{j-1}$ having the lowest cost $f(e)$ among all those for which $I_{j-1} \cup \{e\}$ is independent, that is, still a subset of at least one basis $B$ in $\mathcal{B}(M)$. Then define $I_j := I_{j-1} \cup \{e\}$. Repeat this until $j = r$, and $B = I_r$ will be among the bases that achieve the minimum cost.

**Definition 7.55.** Say that a cost function $f : E \to \{1, 2, \ldots\}$ is $M$-generic if there is a unique basis $B$ in $\mathcal{B}(M)$ achieving the minimum cost $f(B)$.

**Example 7.56.** For the graphic matroid $M$ of Example 7.47, this cost function $f_1 : E \to \{1, 2, \ldots\}$

$$f_1(a) = 1 \quad f_1(b) = 3 \quad f_1(c) = 2$$

is $M$-generic, as it minimizes uniquely on the basis $\{a, d\}$, whereas this cost function $f_2 : E \to \{1, 2, \ldots\}$

$$f_2(a) = 1 \quad f_2(b) = 3 \quad f_2(c) = 2$$

is not $M$-generic, as it achieves its minimum value on the two bases $\{a, c\}, \{a, d\}$.

**Proposition 7.57.** For a matroid $M$ on ground set $E$, one has this expansion

$$\Psi[M] = \sum_{\text{$M$-generic } f : E \to \{1, 2, \ldots\}} x_f$$

where $x_f := \prod_{e \in E} x_{f(e)}$. In particular, for $m \geq 0$, its specialization $p^1$ from Definition 7.6 has this interpretation:

$$p^1\Psi[M](m) = \left|\left\{ \text{$M$-generic } f : E \to \{1, 2, \ldots, m\}\right\}\right|.$$ 

**Proof.** The iterated coproduct $\Delta^{(t-1)} : \mathcal{M} \to \mathcal{M} \otimes^\ell$ sends

$$\Psi[M] \mapsto \sum [M|A_1] \otimes [(M|A_2)/A_1] \otimes \cdots \otimes [(M|A_t)/A_{t-1}]$$

where the sum is over flags of nested subsets

$$(7.16) \quad \emptyset = A_0 \subseteq A_1 \subseteq \cdots \subseteq A_{t-1} \subseteq A_t = E.$$ 

The map $\zeta_{\otimes^\ell}$ sends each summand to 1 or 0, depending upon whether each $(M|A_j)/A_{j-1}$ has a unique basis or not. Thus formula (7.3) shows that the coefficient $\zeta_{\alpha}$ of $x_{i_1}^{\alpha_{i_1}} \cdots x_{i_\ell}^{\alpha_{i_\ell}}$ in $\Psi[M]$ counts the flags of subsets in (7.16) for which $|A_j \setminus A_{j-1}| = \alpha_j$ and $(M|A_j)/A_{j-1}$ has a unique basis, for each $j$.

Given a flag as in (7.16), associate the cost function $f : E \to \{1, 2, \ldots\}$ whose value on each element of $A_j \setminus A_{j-1}$ is $i_j$; conversely, given any cost function $f$, say whose distinct values are $i_1 < \ldots < i_\ell$, one associates the flag having $A_j \setminus A_{j-1} = f^{-1}(i_j)$ for each $j$.

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257 In fact, this expansion was the original definition of $\Psi[M]$ in [19, Defn. 1.1].
We will prove below, using induction on \( s = 0, 1, 2, \ldots, \ell \) the following claim: After having completed \( \alpha_1 + \alpha_2 + \cdots + \alpha_s \) steps in the greedy algorithm (Algorithm 7.54), there is a unique choice for the independent set produced thus far, namely

\[
I_{\alpha_1 + \alpha_2 + \cdots + \alpha_s} = \bigcup_{j=1}^{s} \text{coloops}(\langle M|_{A_j} \rangle/\langle A_{j-1} \rangle),
\]

if and only if each of the matroids \( \langle M|_{A_j} \rangle/\langle A_{j-1} \rangle \) for \( j = 1, 2, \ldots, s \) has a unique basis.

The case \( s = \ell \) in this claim would show what we want, namely that \( f \) is \( M \)-generic, minimizing uniquely on the basis shown in (7.17) with \( s = \ell \), if and only if each \( \langle M|_{A_j} \rangle/\langle A_{j-1} \rangle \) has a unique basis.

The assertion of the claim is trivially true for \( s = 0 \). In the inductive step, one may assume that

- the independent set \( I_{\alpha_1 + \alpha_2 + \cdots + \alpha_{s-1}} \) takes the form in (7.17), replacing \( s \) by \( s - 1 \),
- it is the unique \( f \)-minimizing basis for \( M|_{A_{s-1}} \), and
- \( \langle M|_{A_s} \rangle/\langle A_{s-1} \rangle \) has a unique basis for \( j = 1, 2, \ldots, s - 1 \).

Since \( A_{s-1} \) exactly consists of all of the elements \( e \) of \( E \) whose costs \( f(e) \) lie in the range \( \{i_1, i_2, \ldots, i_{s-1}\} \), in the next \( \alpha_s \) steps the algorithm will work in the quotient matroid \( M/A_{s-1} \) and attempt to augment \( I_{\alpha_1 + \alpha_2 + \cdots + \alpha_{s-1}} \) using the next-cheapest elements, namely the elements of \( A_s \setminus A_{s-1} \), which all have cost \( f \) equal to \( i_s \). Thus the algorithm will have no choices about how to do this augmentation if and only if \( \langle M|_{A_s} \rangle/\langle A_{s-1} \rangle \) has a unique basis, namely its set of coloops, in which case the algorithm will choose to add all of these coloops, giving \( I_{\alpha_1 + \alpha_2 + \cdots + \alpha_s} \) as described in (7.17). This completes the induction.

The last assertion follows from Proposition 7.7. \( \square \)

Example 7.58. If \( M \) has one basis then every function \( f : E \to \{1, 2, \ldots\} \) is \( M \)-generic, and

\[
\Psi[M] = \sum_{f : E \to \{1, 2, \ldots\}} x_f = (x_1 + x_2 + \cdots)|E| = M^{(1)}.
\]

Example 7.59. Let \( U_{r,n} \) denote the uniform matroid of rank \( r \) on \( n \) elements \( E \), having \( B(U_{r,n}) \) equal to all of the \( r \)-element subsets of \( E \).

As \( U_{1,2} \) has \( E = \{1, 2\} \) and \( B = \{\{1\}, \{2\}\} \), genericity means \( f(1) \neq f(2) \), so

\[
\Psi[U_{1,2}] = \sum_{(f(1), f(2)): f(1) \neq f(2)} x_{f(1)}x_{f(2)} = x_1x_2 + x_2x_1 = 2M^{(1,1)}.
\]

Similarly \( U_{1,3} \) has \( E = \{1, 2, 3\} \) with \( B = \{\{1\}, \{2\}, \{3\}\} \), and genericity means either that \( f(1), f(2), f(3) \) are all distinct, or that two of them are the same and the third is smaller. This shows

\[
\Psi[U_{1,3}] = 3\sum_{i<j} x_ix_j^2 + 2\sum_{i<j<k} x_ix_jx_k = 3M^{(1,2)} + 6M^{(1,1,1)}
\]

\[
\psi^1\Psi[U_{1,3}](m) = 3\left(\binom{m}{2}\right) + 6\left(\binom{m}{3}\right) = \frac{m(m-1)(2m-1)}{2}
\]

One can similarly analyze \( U_{2,3} \) and check that

\[
\Psi[U_{2,3}] = 3M^{(2,1)} + 6M^{(1,1,1)}
\]

\[
\psi^1\Psi[U_{2,3}](m) = 3\left(\binom{m}{2}\right) + 6\left(\binom{m}{3}\right) = \frac{m(m-1)(2m-1)}{2}
\]

These last examples illustrate the behavior of \( \Psi \) under the duality operation on matroids.

Definition 7.60. Given a matroid \( M \) of rank \( r \) on ground set \( E \), its dual or orthogonal matroid \( M^\perp \) is a matroid of rank \( |E| - r \) on the same ground set \( E \), having

\[ B(M^\perp) := \{ E \setminus B \}_{B \in B(M)}. \]

Here are a few examples of dual matroids.

Example 7.61. The dual of a uniform matroid is another uniform matroid:

\[ U_{r,n}^\perp = U_{n-r,n}. \]
Example 7.62. If $M$ is matroid of rank $r$ represented by family of vectors $\{e_1, \ldots, e_n\}$ in a vector space over some field $k$, one can find a family of vectors $\{e_1^+, \ldots, e_n^+\}$ that represent $M^\perp$ in the following way. Pick a basis for the span of the vectors $\{e_i\}_{i=1}^n$, and create a matrix $A$ in $k^{r \times n}$ whose columns express the $e_i$ in terms of this basis. Then pick any matrix $A^\perp$ whose row space is the null space of $A$, and one finds that the columns $\{e_i^+\}_{i=1}^n$ of $A^\perp$ represent $M^\perp$. See Oxley [121, §2.2].

Example 7.63. Let $G = (V,E)$ be a graph embedded in the plane with edge set $E$, giving rise to a graphic matroid $M$ on ground set $E$. Let $G^\perp$ be a planar dual of $G$, so that, in particular, for each edge $e$ in $E$, the graph $G^\perp$ has one edge $e^\perp$, crossing $e$ transversely. Then the graphic matroid of $G^\perp$ is $M^\perp$. See Oxley [121, §2.3].

Proposition 7.64. If $\Psi[M] = \sum_\alpha c_\alpha M_\alpha$ then $\Psi[M^\perp] = \sum_\alpha c_\alpha M_\alpha^{\text{rev}(\alpha)}$.

Proof. This amounts to showing that for any composition $\alpha = (\alpha_1, \ldots, \alpha_\ell)$, the cardinality of the set of $M$-generic $f$ having $x_f = x^\alpha$ is the same as the cardinality of the set of $M^\perp$-generic $f^\perp$ having $x_{f^\perp} = x^{\text{rev}(\alpha)}$. We claim that the map $f \mapsto f^\perp$ in which $f^\perp(e) = \ell + 1 - f(e)$ gives a bijection between these sets. To see this, note that any basis $B$ of $M$ satisfies

\begin{align}
(7.18) & \quad f(B) + f(E \setminus B) = \sum_{e \in E} f(e) \\
(7.19) & \quad f(E \setminus B) + f^\perp(E \setminus B) = (\ell + 1)(|E| - r),
\end{align}

where $r$ denotes the rank of $M$. Thus $B$ is $f$-minimizing if and only if $E \setminus B$ is $f$-maximizing (by (7.18)) if and only if $E \setminus B$ is $f^\perp$-minimizing (by (7.19)). Consequently $f$ is $M$-generic if and only if $f^\perp$ is $M^\perp$-generic.

The last assertion follows, for example, from the calculation in Proposition 7.7(i) that $\text{ps}^1(M_\alpha)(m) = \binom{m}{\ell(\alpha)}$ together with the fact that $\ell(\text{rev}(\alpha)) = \ell(\alpha)$.

Just as (7.13) showed that Stanley’s chromatic symmetric function of a graph has an expansion as a sum of $P$-partition enumerators for certain strictly labelled posets $P$, the same holds for $\Psi[M]$.

Definition 7.65. Given a matroid $M$ on ground set $E$, and a basis $B$ in $B(M)$, define the base-cobase poset $P_B$ to have $b < b'$ whenever $b$ lies in $B$ and $b'$ lies in $E \setminus B$ and $(B \setminus \{b\}) \cup \{b'\}$ is in $B(M)$.

Proposition 7.66. For any matroid $M$, one has $\Psi[M] = \sum_{B \in B(M)} F_{(P_B, \text{strict})}(x)$ where $F_{(P, \text{strict})}(x)$ for a poset $P$ means the $P$-partition enumerator for any strict labelling of $P$, i.e. a labelling such that the $P$-partitions satisfy $f(i) < f(j)$ whenever $i <_P j$.

In particular, $\Psi[M]$ expands nonnegatively in the $\{L_\alpha\}$ basis.

Proof. A basic result about matroids, due to Edmonds [46], describes the edges in the matroid base polytope which is the convex hull of all vectors $\sum_{b \in B} t_b e_b$ for $B \in B(M)$ inside $\mathbb{R}^E$ with standard basis $\{e_b\}_{b \in E}$. He shows that all such edges connect two bases $B, B'$ that differ by a single basis exchange, that is, $B' = (B \setminus \{b\}) \cup \{b'\}$ for some $b$ in $B$ and $b'$ in $E \setminus B$.

Polyhedral theory then says that a cost function $f$ on $E$ will minimize uniquely at $B$ if and only if one has a strict increase $f(B) < f(B')$ along each such edge $B \to B'$ emanating from $B$, that is, if and only if $f(b) < f(b')$ whenever $b <_B b'$ in the base-cobase poset $P_B$, that is, $f$ lies in $A(P_B, \text{strict})$. 

Example 7.67. The graphic matroid from Example 7.47 has this matroid base polytope, with the bases $B$ in $B(M)$ labelling the vertices:
The base-cobase posets $P_B$ for its five vertices $B$ are as follows:

\[
\begin{array}{cccc}
  a & b \\
  \uparrow & \downarrow \\
  c & d
\end{array} \quad \begin{array}{cccc}
  b & d \\
  \uparrow & \downarrow \\
  a & c
\end{array} \quad \begin{array}{cccc}
  a & d \\
  \uparrow & \downarrow \\
  b & c
\end{array} \quad \begin{array}{cccc}
  a & c \\
  \uparrow & \downarrow \\
  b & d
\end{array} \quad \begin{array}{cccc}
  b & c \\
  \uparrow & \downarrow \\
  a & d
\end{array}
\]

One can label the first of these five strictly as

\[
1, 2, 3, 4
\]

and compute its strict $P$-partition enumerator from the linear extensions $\{3412, 3421, 4312, 4321\}$ as

\[
L(2, 2) + L(2, 1, 1) + L(1, 1, 1, 1) + L(1, 1, 1, 1)
\]

while any of the last four can be labelled strictly as

\[
1, 2, 3, 4
\]

and they each have an extra linear extension $3142$ giving their strict $P$-partition enumerators as

\[
L(2, 2) + L(2, 1, 1) + L(1, 1, 1, 1) + L(1, 1, 1, 1) + L(1, 2, 1)
\]

Hence one has

\[
\Psi[M] = 5L(2, 2) + 5L(1, 1, 1, 1) + 4L(1, 1, 2, 1) + 5L(2, 1, 1, 1, 1) + 5L(1, 1, 1, 1, 1)
\]

As $M$ is a graphic matroid for a self-dual planar graph, one has a matroid isomorphism $M \cong M^\perp$ (see Example 7.63), reflected in the fact that $\Psi[M]$ is invariant under the symmetry swapping $M_\alpha \leftrightarrow M_{\text{rev}(\alpha)}$ (and simultaneously swapping $L_\alpha \leftrightarrow L_{\text{rev}(\alpha)}$).

This $P$-partition expansion for $\Psi[M]$ also allows us to identify its image under the antipode of QSym.

**Proposition 7.68.** For a matroid $M$ on ground set $E$, one has

\[
S(\Psi[M]) = (-1)^{|E|} \sum_{f : E \to \{1, 2, \ldots\}} |\{f \text{-maximizing bases } B\}| \cdot x_f
\]

and

\[
\text{ps}^1\Psi[M](m) = (-1)^{|E|} \sum_{f : E \to \{1, 2, \ldots, m\}} |\{f \text{-maximizing bases } B\}|.
\]

In particular, the expected number of $f$-maximizing bases among all cost functions $f : E \to \{1, 2, \ldots, m\}$ is

\[
(-m)^{-|E|} \text{ps}^1\Psi[M](m).
\]

**Proof.** Corollary 5.27 implies

\[
S(\Psi[M]) = \sum_{B \in \mathcal{B}(M)} S(F_{(P_B, \text{strict})}(x)) = (-1)^{|E|} \sum_{B \in \mathcal{B}(M)} F_{(P_B^\text{opp, natural})}(x)
\]

where $F_{(P, \text{natural})}(x)$ is the enumerator for $P$-partitions in which $P$ has been naturally labelled, so that they satisfy $f(i) \leq f(j)$ whenever $i <_P j$. When $P = P_B^\text{opp}$, this is exactly the condition for $f$ to achieve its maximum value at $f(B)$ (possibly not uniquely), that is, for $f$ to lie in the closed normal cone to the vertex indexed by $B$ in the matroid base polytope; compare this with the discussion in the proof of Proposition 7.66.

Thus one has

\[
S(\Psi[M]) = (-1)^{|E|} \sum_{f \text{ maximizing at } B} x_f
\]

which agrees with the statement of the proposition, after reversing the order of the summation.

The rest follows from Proposition 7.7. \(\square\)
Example 7.69. We saw in Example 7.67 that the matroid $M$ from Example 7.47 has
\[
\Psi[M] = 5L_{(2,2)} + 5L_{(1,1,2)} + 4L_{(1,2,1)} + 5L_{(2,1,1)} + 5L_{(1,1,1,1)},
\]
and therefore will have
\[
\psi_1^1 \Psi[M](m) = 5\binom{m - 2 + 4}{4} + (5 + 4 + 5)\binom{m - 3 + 4}{4} + 5\binom{m - 4 + 4}{4} = \frac{m(m - 1)(2m^2 - 2m + 1)}{2}.
\]
using $\psi_1^1(L_\alpha)(m) = \binom{m - |\alpha|}{|\alpha|}$ from Proposition 7.7 (i). Let us first do a reality-check on a few of its values
with $m \geq 0$ using Proposition 7.57, and for negative $m$ using Proposition 7.68:

<table>
<thead>
<tr>
<th>$m$</th>
<th>$\psi_1^1 \Psi<a href="m">M</a>$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$-1$</td>
<td>5</td>
</tr>
<tr>
<td>0</td>
<td>0</td>
</tr>
<tr>
<td>1</td>
<td>1</td>
</tr>
<tr>
<td>2</td>
<td>2</td>
</tr>
</tbody>
</table>

When $m = 0$, interpreting the set of cost functions $f : E \to \{1, 2, \ldots, m\}$ as being empty explains why
the value shown is 0. When $m = 1$, there is only one function $f : E \to \{1\}$, and it is not $M$-generic; any
of the 5 bases in $\mathcal{B}(M)$ will minimize $f(B)$, explaining both why the value for $m = 1$ is 0, but also explaining
the value of 5 for $m = -1$. The value of 5 for $m = 2$ counts these $M$-generic cost functions $f : E \to \{1, 2\}$:

Lastly, Proposition 7.68 predicts the expected number of $f$-minimizing bases for $f : E \to \{1, 2, \ldots, m\}$ as
\[
(-m)^{-|E|} \psi_1^1 \Psi[M](-m) = (-m)^{-|E|} \frac{m(m + 1)(2m^2 + 2m + 1)}{2} = \frac{(m + 1)(2m^2 + 2m + 1)}{2m^3},
\]
whose limit as $m \to \infty$ is 1, consistent with the notion that “most” cost functions should be generic with
respect to the bases of $M$, and maximize/minimize on a unique basis.

Remark 7.70. It is not coincidental that there is a similarity of results for Stanley’s chromatic symmetric
function of a graph $\Psi[G]$ and for the matroid quasisymmetric function $\Psi[M]$, such as the $P$-partition expansions (7.13) versus Proposition 7.66, and the reciprocity results Proposition 7.42 versus Proposition 7.68. It
was noted in [19, §9] that one can associate a similar quasisymmetric function invariant to any generalized
permutohedra in the sense of Postnikov [129]. Furthermore, recent work of Ardila and Aguiar [3] has shown
that there is a Hopf algebra of such generalized permutohedra, arising from a Hopf monoid in the sense
of Aguiar and Mahajan [6]. This Hopf algebra generalizes the chromatic Hopf algebra of graphs and the
matroid-minor Hopf algebra, and its quasisymmetric function invariant derives as usual from Theorem 7.3.
Their work [3] also provides a generalization of the chromatic Hopf algebra antipode formula of Humpert
and Martin [77] discussed in Remark 7.23 above.
8. The Malvenuto-Reutenauer Hopf Algebra of Permutations

Like so many Hopf algebras we have seen, the *Malvenuto-Reutenauer Hopf algebra* \( \text{FQSym} \) can be thought of fruitfully in more than one way. One is that it gives a natural noncommutative lift of the quasisymmetric \( P \)-partition enumerators and the fundamental basis \( \{ L_\alpha \} \) of \( \text{QSym} \), rendering their product and coproduct formulas even more natural.

### 8.1. Definition and Hopf Structure

**Definition 8.1.** Define \( \text{FQSym} = \bigoplus_{n \geq 0} \text{FQSym}_n \) to be a graded \( k \)-module in which \( \text{FQSym}_n \) has \( k \)-basis \( \{ F_w \} \) indexed by the permutations \( w = (w_1, \ldots, w_n) \) in \( S_n \).

We first attempt to lift the product and coproduct formulas (5.8), (5.7) in the \( \{ L_\alpha \} \) basis of \( \text{QSym} \). We attempt to define a product for \( u \in S_k, v \in S_\ell \) as follows:

\[
F_u F_v := \sum_{w \in u \shuffle v[k]} F_w
\]

where we regard permutations as words (namely, every \( \pi \in S_n \) is identified with the word \( (\pi_1, \pi_2, \ldots, \pi_n) \)), and where for \( v = (v_1, \ldots, v_\ell) \) one sets \( v[k] := (k + v_1, \ldots, k + v_\ell) \). Note that the multiset \( u \shuffle v[k] \) is an actual set in this situation (i.e., has each element appear only once) and is a subset of \( S_{k+\ell} \).

The coproduct will be defined using the notation of standardization \( \text{std}(i) \) of a word \( i = (i_1, \ldots, i_n) \) in some linearly ordered alphabet. This standardization is the permutation in \( S_n \) obtained by replacing all the occurrences of the smallest letter in \( i \) by the numbers \( 1, 2, \ldots, m_1 \) from left to right (where \( m_1 \) is the number of times this smallest letter occurs in \( i \)), then replacing all occurrences of the next smallest letter by the numbers \( m_1 + 1, m_1 + 2, \ldots, m_1 + m_2 \), from left to right (where \( m_2 \) is how often this letter occurs in \( i \)), etc.

**Example 8.2.** Considering words in the Roman alphabet \( a < b < c < \cdots \)

\[
\text{std}(b \ a \ c \ c \ b \ a \ a \ b \ a \ c \ b) = (5 \ 1 \ 9 \ 10 \ 6 \ 2 \ 3 \ 7 \ 4 \ 11 \ 8).
\]

Using this, define for \( w = (w_1, \ldots, w_n) \) in \( S_n \)

\[
\Delta F_w := \sum_{k=0}^n F_{\text{std}(w_1, w_2, \ldots, w_k)} \otimes F_{\text{std}(w_{k+1}, w_{k+2}, \ldots, w_n)}.
\]

It is possible to check directly that the maps defined in (8.1) and (8.2) endow \( \text{FQSym} \) with the structure of a graded connected finite type Hopf algebra; see Hazewinkel, Gubareni, Kirichenko [68, Thm. 7.1.8]. However in justifying this here, we will follow the approach of Duchamp, Hivert and Thibon [43, §3], which exhibits \( \text{FQSym} \) as a subalgebra of a larger ring of (noncommutative) power series of bounded degree in a totally ordered alphabet.

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258 Alternatively, one can define \( \text{std}(i) \) as the unique permutation \( \sigma \in S_n \) such that for every \( 1 \leq a < b \leq n \), we have \( \sigma(a) < \sigma(b) \) if and only if \( i_a \leq i_b \).
Definition 8.3. Given a totally ordered set $I$, create a totally ordered variable set $\{x_i\}_{i \in I}$, and the ring $R(\{X_i\}_{i \in I})$ of noncommutative power series of bounded degree in this alphabet. Many times, we will use a variable set $X := (x_1 < x_2 < \cdots)$, and call the ring $R(X)$.

We first identify the algebra structure for FQSym as the subalgebra of finite type within $R(\{X_i\}_{i \in I})$ spanned by the elements

$$F_w = F_w(\{X_i\}_{i \in I}) := \sum_{i = \{i_1, \ldots, i_n\}; \text{std}(i) = w^{-1}} x_i$$

where $X_i := x_{i_1} \cdots x_{i_n}$, as $w$ ranges over $\bigcup_{n \geq 0} \mathfrak{S}_n$.

Example 8.4. For the alphabet $X = (x_1 < x_2 < \cdots)$, in $R(X)$ one has

$$F_1 = \sum_{1 \leq i} x_i = x_1 + x_2 + \cdots$$

$$F_{12} = \sum_{1 \leq i < j} x_i x_j = x_1^2 + x_2^2 + \cdots + x_1 x_2 + x_2 x_3 + x_2 x_3 + x_1 x_4 + \cdots$$

$$F_{21} = \sum_{1 \leq i < j} x_j x_i = x_2 x_1 + x_3 x_1 + x_4 x_2 + x_4 x_1 + \cdots$$

$$F_{312} = \sum_{\text{std}(i) = 231} x_i = \sum_{1 \leq i < j \leq 3} x_j x_k x_i$$

$$= x_2^2 x_1 + x_3^2 x_1 + x_3^2 x_2 + \cdots + x_2 x_3 x_1 + x_2 x_4 x_1 + \cdots$$

$$= \sum_{\text{std}(i) = 132} x_i + \sum_{\text{std}(i) = 231} x_i + \sum_{\text{std}(i) = 321} x_i = F_{132} + F_{312} + F_{321} = \sum_{w \in 1 \cup 32} F_w$$

Proposition 8.5. For any totally ordered infinite set $I$, the elements $\{F_w\}$ as $w$ ranges over $\bigcup_{n \geq 0} \mathfrak{S}_n$ form the $k$-basis for a subalgebra FQSym$(\{X_i\}_{i \in I})$ of $R(X)$, which is graded connected and of finite type, having multiplication defined $k$-linearly by (8.1).

Consequently all such algebras are isomorphic to a single algebra FQSym, having basis $\{F_w\}$ and multiplication given by the rule (8.1), with the isomorphism mapping $F_w \mapsto F_w(\{X_i\}_{i \in I})$.

For example,

$$F_1 F_{21} = (x_1 + x_2 + x_3 + \cdots)(x_2 x_1 + x_3 x_1 + x_3 x_2 + x_4 x_1 + \cdots)$$

$$= x_1 \cdot x_3 x_2 + x_1 \cdot x_4 x_2 + \cdots + x_1 \cdot x_3 x_2 + x_2 \cdot x_3 x_2 + x_2 \cdot x_4 x_2 + \cdots$$

$$+ x_2 \cdot x_3 x_1 + x_2 \cdot x_4 x_1 + \cdots + x_2 \cdot x_3 x_1 + x_3 \cdot x_3 x_1 + x_3 \cdot x_3 x_2 + \cdots$$

$$+ x_3 \cdot x_2 x_1 + x_4 \cdot x_2 x_1 + \cdots$$

$$= \sum_{\text{std}(i) = 132} x_i + \sum_{\text{std}(i) = 231} x_i + \sum_{\text{std}(i) = 321} x_i = F_{132} + F_{312} + F_{321} = \sum_{w \in 1 \cup 32} F_w$$

Proof of Proposition 8.5. The elements $\{F_w(\{X_i\}_{i \in I})\}$ are linearly independent as they are supported on disjoint monomials, and so form a $k$-basis for their span. The fact that they multiply via rule (8.1) is the equivalence of conditions (i) and (iii) in the following Lemma 8.6, from which all the remaining assertions follow. \[\text{Let } N \text{ denote the free monoid on the alphabet } \{X_i\}_{i \in I}; \text{ it consists of words } X_{i_1} X_{i_2} \cdots X_{i_k}. \text{ We define a topological } k\text{-module } k(\langle\{X_i\}_{i \in I}\rangle) \text{ to be the Cartesian product } k^N \text{ equipped with the product topology, but we identify its element } (a_{w,u})_{u \in N} \text{ with the word } w \text{ for every } w \in N. \text{ Thus, every element } (\lambda_w)_{w \in N} \in k^N = k(\langle\{X_i\}_{i \in I}\rangle) \text{ can be rewritten as the convergent sum } \sum_{w \in N} \lambda_w w. \text{ We call } \lambda_w \text{ the coefficient of } w \text{ in this element (or the coefficient of this element before } w). \text{ The elements of } k(\langle\{X_i\}_{i \in I}\rangle) \text{ will be referred to as noncommutative power series. We define a multiplication on } k(\langle\{X_i\}_{i \in I}\rangle) \text{ by the formula}\]

\[
\left( \sum_{w \in N} \lambda_w w \right) \left( \sum_{w \in N} \mu_w w \right) = \sum_{w \in N} \left( \sum_{(u,v) \in N^2, w = u v} \lambda_u \mu_v \right) w.
\]

(This is well-defined thanks to the fact that, for each $w \in N$, there are only finitely many $(u,v) \in N^2$ satisfying $w = uv$.) Thus, $k(\langle\{X_i\}_{i \in I}\rangle)$ becomes a $k$-algebra with unity $1$ (the empty word). (It is similar to the monoid algebra $kN$ of $N$ over $k$, with the only difference that infinite sums are allowed.)

Now, we define $R(\{X_i\}_{i \in I})$ to be the $k$-subalgebra of $k(\langle\{X_i\}_{i \in I}\rangle)$ consisting of all noncommutative power series $\sum_{w \in N} \lambda_w w \in k(\langle\{X_i\}_{i \in I}\rangle)$ of bounded degree (i.e., such that all words $w \in N$ of sufficiently high length satisfy $\lambda_w = 0$).
Lemma 8.6. For a triple of permutations

\[
\begin{align*}
  u &= (u_1, \ldots, u_k) \text{ in } S_k, \\
  v &= (v_1, \ldots, v_{n-k}) \text{ in } S_{n-k}, \\
  w &= (w_1, \ldots, w_n) \text{ in } S_n,
\end{align*}
\]

the following conditions are equivalent:

(i) \( w^{-1} \) lies in the set \( u^{-1} \uplus v^{-1}[k] \),

(ii) \( u = \text{std}(w_1, \ldots, w_k) \) and \( v = \text{std}(w_{k+1}, \ldots, w_n) \),

(iii) for some word \( \mathbf{i} = (i_1, \ldots, i_n) \) with \( \text{std}(\mathbf{i}) = w \) one has \( u = \text{std}(i_1, \ldots, i_k) \) and \( v = \text{std}(i_{k+1}, \ldots, i_n) \).

Proof of Lemma 8.6. The implication (ii) \( \Rightarrow \) (iii) is clear since \( \text{std}(w) = w \). The reverse implication (iii) \( \Rightarrow \) (ii) is best illustrated by example, e.g., considering Example 8.2 as concatenated, with \( n = 11, k = 6, n-k = 5 \):

\[
\begin{array}{c|c}
  w &= \text{std} \begin{pmatrix} b & a & c & c & b & a \end{pmatrix} = (1 3 2 5 4) \\
  v &= \text{std} \begin{pmatrix} a & b & a & c & b \end{pmatrix} = (3 7 4 11 8) \\
  u &= \text{std} \begin{pmatrix} 5 & 1 & 9 & 10 & 6 & 2 \end{pmatrix} = (3 1 5 6 4 2) \\
  v &= \text{std} \begin{pmatrix} a & b & a & c & b \end{pmatrix} = (1 3 2 5 4)
\end{array}
\]

The equivalence of (i) and (ii) is a fairly standard consequence of unique parabolic factorization for cosets \( xW_j \) (that is, the permutations \( x \in S_n \) satisfying \( x_1 < \cdots < x_k \) and \( x_{k+1} < \cdots < x_n \)). One can uniquely express any \( w \) in \( W \) as \( w = xy \) with \( x \) in \( W_j \) and \( y \) in \( W_j \), which here means that \( y = u \cdot v[k] = v[k] \cdot u \) for some \( u \in S_k \) and \( v \in S_{n-k} \). Therefore \( w = xw[k] \), if and only if \( w^{-1} = u^{-1}v^{-1}[k]x^{-1} \), which means that \( w^{-1} \) is the shuffle of the sequences \( u^{-1} \) in positions \( \{x_1, \ldots, x_k\} \) and \( v^{-1}[k] \) in positions \( \{x_{k+1}, \ldots, x_n\} \). \( \square \)

Example 8.7. To illustrate the equivalence of (i) and (ii) and the parabolic factorization in the preceding proof, let \( n = 9 \) and \( k = 5 \) with

\[
\begin{align*}
  w &= \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 \\
    4 & 9 & 6 & 1 & 5 & 8 & 2 & 3 & 7 \end{pmatrix} \\
  &= x \cdot u \cdot v[k] \\
  w^{-1} &= \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 \\
    4 & 9 & 6 & 1 & 5 & 8 & 2 & 3 & 7 \end{pmatrix} \\
  &= u^{-1} \cdot v^{-1}[k] \cdot x^{-1}
\end{align*}
\]

Proposition 8.5 yields that \( \text{FQSym} \) is isomorphic to the \( k \)-subalgebra \( \text{FQSym}(\mathbf{X}) \) of the \( k \)-algebra \( R(\mathbf{X}) \) when \( \mathbf{X} \) is the variable set \( \{X_1 < X_2 < \cdots\} \). We identify \( \text{FQSym} \) with \( \text{FQSym}(\mathbf{X}) \) along this isomorphism. For any infinite alphabet \( \{X_i\}_{i \in I} \) and any \( f \in \text{FQSym} \), we denote by \( f(\{x_i\}_{i \in I}) \) the image of \( f \) under the algebra isomorphism \( \text{FQSym} \rightarrow \text{FQSym}(\{X_i\}_{i \in I}) \) defined in Proposition 8.5.

One can now use this to define a coalgebra structure on \( \text{FQSym} \). Roughly speaking, one wants to first evaluate an element \( f \) in \( \text{FQSym} \cong \text{FQSym}(\mathbf{X}) \cong \text{FQSym}(\mathbf{X}, \mathbf{Y}) \) as \( f(\mathbf{X}, \mathbf{Y}) \), using the linearly ordered variable set \( \{\mathbf{X}, \mathbf{Y}\} := (X_1 < X_2 < \cdots < Y_1 < Y_2 < \cdots) \). Then one should take the image of \( f(\mathbf{X}, \mathbf{Y}) \) after imposing the partial commutativity relations

\[
(8.3) \quad X_iY_j = Y_jX_i \text{ for every pair } (X_i, Y_j) \in \mathbf{X} \times \mathbf{Y},
\]

and hope that this image lies in a subalgebra isomorphic to

\[
\text{FQSym}(\mathbf{X}) \otimes \text{FQSym}(\mathbf{Y}) \cong \text{FQSym} \otimes \text{FQSym}.
\]
We argue this somewhat carefully. Start by considering the canonical monoid epimorphism

\[ F(\mathbf{X}, \mathbf{Y}) \xrightarrow{\rho} M, \]

where \( F(\mathbf{X}, \mathbf{Y}) \) denotes the free monoid on the alphabet \( \mathbf{X} \) and \( \mathbf{Y} \) and \( M \) denotes its quotient monoid imposing the partial commutativity relations (8.3). Let \( k^M \) denote the \( k \)-module of all functions \( f : M \to k \), with pointwise addition and scalar multiplication; similarly define \( k^{F(\mathbf{X}, \mathbf{Y})} \). As both monoids \( F(\mathbf{X}, \mathbf{Y}) \) and \( M \) enjoy the property that an element \( m \) has only finitely many factorizations as \( m = m_1m_2 \), one can define a convolution algebra structure on both \( k^{F(\mathbf{X}, \mathbf{Y})} \) and \( k^M \) via

\[ (f_1 \ast f_2)(m) = \sum_{(m_1, m_2) \in N \times N} f_1(m_1)f_2(m_2), \]

where \( N \) is respectively \( F(\mathbf{X}, \mathbf{Y}) \) or \( M \). As fibers of the map \( \rho \) in (8.4) are finite, it induces a map of convolution algebras, which we also call \( \rho \):

\[ k^{F(\mathbf{X}, \mathbf{Y})} \xrightarrow{\rho} k^M. \]

Now recall that \( R(\mathbf{X}) \) denotes the algebra of noncommutative formal power series in the variable set \( \mathbf{X} \), of bounded degree, with coefficients in \( k \). One similarly has the ring \( R(\mathbf{X}, \mathbf{Y}) \), which can be identified with the subalgebra of \( k^{F(\mathbf{X}, \mathbf{Y})} \) consisting of the functions \( f : F(\mathbf{X}, \mathbf{Y}) \to k \) having a bound on the length of the words in their support (the value of \( f \) on a word in \( \mathbf{X}, \mathbf{Y} \) gives its power series coefficient corresponding to said word). We let \( R(M) \) denote the analogous subalgebra of \( k^M \); this can be thought of as the algebra of bounded degree “partially commutative power series” in the variable sets \( \mathbf{X} \) and \( \mathbf{Y} \). Note that \( \rho \) restricts to a map

\[ R(\mathbf{X}, \mathbf{Y}) \xrightarrow{\rho} R(M). \]

Finally, we claim (and see Proposition 8.9 below for a proof) that this further restricts to a map

\[ FQSym(\mathbf{X}, \mathbf{Y}) \xrightarrow{\rho} FQSym(\mathbf{X}) \otimes FQSym(\mathbf{Y}) \]

in which the target is identified with its image under the (injective\(^{260}\)) multiplication map

\[ FQSym(\mathbf{X}) \otimes FQSym(\mathbf{Y}) \to R(M), \quad f(\mathbf{X}) \otimes g(\mathbf{Y}) \mapsto f(\mathbf{X})g(\mathbf{Y}). \]

Using the identification of \( FQSym \) with all three of \( FQSym(\mathbf{X}), FQSym(\mathbf{Y}), FQSym(\mathbf{X}, \mathbf{Y}) \), the map \( \rho \) in (8.7) will then define a coproduct structure on \( FQSym \). Abusing notation, for \( f \) in \( FQSym \), we will simply write \( \Delta(f) = f(\mathbf{X}, \mathbf{Y}) \) instead of \( \rho(f(\mathbf{X}, \mathbf{Y})) \).

**Example 8.8.** Recall from Example 8.4 that one has

\[ F_{312} = \sum_{i \text{ std}(i) = 231} X_i = \sum_{1 \leq i < j \leq k} X_j X_k X_i \]

and therefore its coproduct is

\[ \Delta F_{312} = F_{312}(X_1, X_2, \ldots, Y_1, Y_2, \ldots) \quad (\text{by our abuse of notation}) \]

\[ = \sum_{i<j \leq k} X_j X_k X_i + \sum_{i \leq j, k} X_j Y_i \cdot X_k + \sum_{i \leq j \leq k} Y_j Y_i \cdot X_k + \sum_{i \leq j \leq k} 1 \cdot Y_j Y_i \cdot X_k \]

\[ = F_{312}(\mathbf{X}) \cdot 1 + F_{21}(\mathbf{X}) \cdot F_1(\mathbf{Y}) + F_1(\mathbf{X}) \cdot F_{12}(\mathbf{Y}) + 1 \cdot F_{312}(\mathbf{Y}) \]

\[ = F_{312} \otimes 1 + F_{21} \otimes F_1 + F_1 \otimes F_{12} + 1 \otimes F_{312}. \]

**Proposition 8.9.** The map \( \rho \) in (8.6) does restrict as claimed to a map as in (8.7), and hence defines a coproduct on \( FQSym \), acting on the \( \{F_w\} \) basis by the rule (8.2). This endows \( FQSym \) with the structure of a graded connected finite type Hopf algebra.

\(^{260}\) as images of the basis \( F_u(\mathbf{X}) \otimes F_v(\mathbf{Y}) \) of \( FQSym(\mathbf{X}) \otimes FQSym(\mathbf{Y}) \) are supported on disjoint monomials in \( R(M) \), so linearly independent.
Proof. Let $I$ be the totally ordered set $\{1 < 2 < 3 < \cdots \}$. Let $J$ be the totally ordered set $\{1 < 2 < 3 < \cdots < \bar{1} < \bar{2} < \bar{3} < \cdots \}$. We set $X_i = Y_i$ for every positive integer $i$. Then, the alphabet $(X, Y)$ can be written as $\{X_i\}_{i \in I}$.

If $i$ is a word over the alphabet $I = \{1 < 2 < 3 < \cdots \}$, then we denote by $\bar{i}$ the word over $J$ obtained from $i$ by replacing every letter $i$ by $\bar{i}$.

For the first assertion of Proposition 8.9, it suffices to check that $F_w$ indeed has the image under $\Delta$ claimed in (8.2). Let $n \in \mathbb{N}$ and $w \in \mathfrak{S}_n$. Then,

$$
\Delta F_w = F_w(X, Y) \quad \text{(by our abuse of notation)}
$$

$$
= \sum_{i \in J^n; \text{std}(i) = w^{-1}} (X, Y)_i = \sum_{t \in J^n; \text{std}(t) = w^{-1}} (X, Y)_t
$$

(8.8)

(since for every $t \in J^n$, there exists exactly one choice of $k \in \{0, 1, \ldots, n\}$ and $(i, j) \in I^k \times I^{n-k}$ satisfying $t \in i \uplus \bar{j}$; namely, $i$ is the restriction of $t$ to the subalphabet $I$ of $J$, whereas $j$ is the restriction of $t$ to $J \setminus I$, and $k$ is the length of $i$).

We now fix $k$ and $(i, j)$, and try to simplify the inner sum $\sum_{t \in J^n: \text{std}(t) = w^{-1}; \text{std}(\bar{t}) = w^{-1}} (X, Y)_t$ on the right hand side of (8.8). First we notice that this sum is nonempty if and only if there exists some $t \in i \uplus \bar{j}$ satisfying std$(t)$ = $w^{-1}$. This existence is easily seen to be equivalent to $w^{-1} \in \text{std}(i) \uplus \text{std}(j)[k]$ (since the standardization of any shuffle in $i \uplus \bar{j}$ is the corresponding shuffle in std$(i) \uplus \text{std}(j)[k]$). This, in turn, is equivalent to std$(i) = (\text{std}(w_1, \ldots, w_k))^{-1}$ and std$(j) = (\text{std}(w_{k+1}, \ldots, w_n))^{-1}$ (according to the equivalence (i) $\iff$ (ii) in Lemma 8.6). Hence, the inner sum on the right hand side of (8.8) is nonempty if and only if std$(i) = (\text{std}(w_1, \ldots, w_k))^{-1}$ and std$(j) = (\text{std}(w_{k+1}, \ldots, w_n))^{-1}$. When it is nonempty, it has only one addend, and this addend is $(X, Y)_t = X_iY_j$ (since $t \in i \uplus \bar{j}$). Summarizing, we see that the inner sum on the right hand side of (8.8) equals $X_iY_j$ when std$(i) = (\text{std}(w_1, \ldots, w_k))^{-1}$ and std$(j) = (\text{std}(w_{k+1}, \ldots, w_n))^{-1}$, and is empty otherwise. Thus, (8.8) simplifies to

$$
\Delta F_w = \sum_{k=0}^n \sum_{(i, j) \in I^k \times I^{n-k}: \text{std}(i) = (\text{std}(w_1, \ldots, w_k))^{-1}, \text{std}(j) = (\text{std}(w_{k+1}, \ldots, w_n))^{-1}} X_iY_j
$$

$$
= \sum_{k=0}^n \text{std}(w_1, \ldots, w_k)(X) \text{std}(w_{k+1}, \ldots, w_n)(Y)
$$

$$
= \sum_{k=0}^n \text{std}(w_1, \ldots, w_k) \otimes \text{std}(w_{k+1}, \ldots, w_n) \in \text{FQSym} \otimes \text{FQSym}.
$$

This proves (8.2), and thus the first assertion of Proposition 8.9.

From this, it is easy to derive that $\Delta$ satisfies coassociativity (i.e., the diagram (1.3) holds for $C = \text{FQSym}$). (Alternatively, one can obtain this from the associativity of multiplication using Corollary 8.11.) We have already verified the rule (8.2). The connected graded structure on $\text{FQSym}$ gives a counit and an antipode for free. \hfill \square

Exercise 8.10. We say that a permutation $w \in \mathfrak{S}_n$ is connected if $n$ is a positive integer and if there exists no $i \in \{1, 2, \ldots, n-1\}$ satisfying $f(\{1, 2, \ldots, i\}) = \{1, 2, \ldots, i\}$. Let $\mathfrak{C}$ denote the set of all connected permutations of all $n \in \mathbb{N}$. Show that $\text{FQSym}$ is a free (noncommutative) $k$-algebra with generators $(F_w)_{w \in \mathfrak{C}}$. (This statement means that $(F_{w_1}F_{w_2} \cdots F_{w_k})_{k \in \mathbb{N}; \,(w_1, w_2, \ldots, w_k) \in \mathfrak{C}^k}$ is a basis of the $k$-module $\text{FQSym}$.)

\footnote{In fact, the elements std$(t)$ for $t \in i \uplus \bar{j}$ are pairwise distinct, and thus only one of them can equal $w^{-1}$.}
Corollary 8.11. The Hopf algebra $\text{FQSym}$ is self-dual: Let $\{G_w\}$ be the dual $k$-basis to the $k$-basis $\{F_w\}$ for $\text{FQSym}$. Then, the $k$-linear map sending $G_w \mapsto F_{w^{-1}}$ is a Hopf algebra isomorphism $\text{FQSym}^\circ \rightarrow \text{FQSym}$.

Proof. For any $0 \leq k \leq n$, any $u \in \mathcal{S}_k$ and any $v \in \mathcal{S}_{n-k}$, one has

$$F_{u^{-1}}F_{v^{-1}} = \sum_{w^{-1} \in u^{-1} \sqcup v^{-1}} F_{w^{-1}}$$

via the equivalence of (i) and (ii) in Lemma 8.6. On the other hand, in $\text{FQSym}$ for any $0 \leq k \leq n$, any $u \in \mathcal{S}_k$ and any $v \in \mathcal{S}_{n-k}$, one has

$$F_{u^{-1}}F_{v^{-1}} = \sum_{w^{-1} \in u^{-1} \sqcup v^{-1}} F_{w^{-1}}$$

coming from the coproduct formula (8.2) for $\text{FQSym}$ in the $\{F_w\}$-basis. Comparing these equalities, we see that the $k$-linear map $\tau$ sending $G_w \mapsto F_{w^{-1}}$ is an isomorphism $\text{FQSym}^\circ \rightarrow \text{FQSym}$ of $k$-algebras. Hence, the adjoint $\tau^* : \text{FQSym}^\circ \rightarrow (\text{FQSym}^\circ)^\circ$ of this map is an isomorphism of $k$-coalgebras. But identifying $(\text{FQSym}^\circ)^\circ$ with $\text{FQSym}$ in the natural way (since $\text{FQSym}$ is of finite type), we easily see that $\tau^* = \tau$, whence $\tau$ itself is an isomorphism of both $k$-algebras and $k$-coalgebras, hence of $k$-bialgebras, hence of Hopf algebras.

We can now be a bit more precise about the relations between the various algebras $\Lambda, \text{QSym}, \text{NSym}, \text{FQSym}, R(X), R(x)$.

Not only does $\text{FQSym}$ allow one to lift the Hopf structure of $\text{QSym}$, it dually allows one to extend the Hopf structure of $\text{NSym}$. To set up this duality, note that Corollary 8.11 motivates the choice of an inner product on $\text{FQSym}$ in which

$$(F_u, F_v) := \delta_{u^{-1}, v^{-1}}.$$

We wish to identify the images of the ribbon basis $\{R_\alpha\}$ of $\text{NSym}$ when included in $\text{FQSym}$.

Definition 8.12. For any composition $\alpha$, define an element $R_\alpha$ of $\text{FQSym}$ by

$$R_\alpha := \sum_{w : \text{Des}(w) = \text{D}(\alpha)} F_{w^{-1}}$$

where the $w$ in the sums are supposed to belong to $\mathcal{S}_{|\alpha|}$, and where the descent set of a sequence $i = (i_1, \ldots, i_n)$ is defined by

$$\text{Des}(i) := \{j \in \{1, 2, \ldots, n-1\} : i_j > i_{j+1}\} = \text{Des}(\text{std}(i)).$$

Alternatively, $R_\alpha = \sum_T X_T$ in which the sum is over column-strict tableaux of the ribbon skew shape $\alpha$, and $X_T = X_i$ in which $i$ is the sequence of entries of $T$ read in order from the southwest toward the northeast.

Example 8.13. Taking $\alpha = (1, 3, 2)$, with ribbon shape and column-strict fillings $T$ as shown

$$\begin{array}{cccc}
\square & \square & \square & \square \\
\square & \square & \square & i_5 \leq i_6 \\
\square & \square & \square & i_2 \leq i_3 \leq i_4 \\
\square & \square & i_1 \\
\end{array}$$

one has that

$$R_{(1,3,2)} = \sum_{i = (i_1, i_2, i_3, i_4, i_5) : \text{Des}(i) = \text{D}(\alpha) = \{1, 4\}} X_i = \sum_{i_1 > i_2 \leq i_3 \leq i_4 \leq i_5 \leq i_6} X_{i_1}X_{i_2}X_{i_3}X_{i_4}X_{i_5}X_{i_6} = \sum_T X_T$$
Corollary 8.14. For every $n \in \mathbb{N}$ and $w \in \mathfrak{S}_n$, we let $\gamma(w)$ denote the unique composition $\alpha$ of $n$ satisfying $D(\alpha) = \text{Des}(w)$. The following $\mathbb{Z}$-linear surjection and injection

\[
\begin{align*}
\pi & : \text{FQSym} \twoheadrightarrow \text{QSym} \\
& F_w \mapsto L_{\gamma(w)} \\
\iota & : \text{NSym} \hookrightarrow \text{FQSym} \\
& R_\alpha \mapsto R_\alpha
\end{align*}
\]

are Hopf morphisms, and are also adjoint maps with respect to the above choice of inner product on FQSym and the usual dual pairing between NSym and QSym. Furthermore, the second map $\iota$ lets one factor the surjection $\text{NSym} \twoheadrightarrow \Lambda$ as follows

\[
\begin{align*}
\text{NSym} & \to \text{FQSym} \hookrightarrow \text{R}(\mathbf{X}) \to \text{R}(x) \\
R_\alpha & \mapsto R_\alpha \mapsto s_\alpha(x)
\end{align*}
\]

through the abelianization $\text{R}(\mathbf{X}) \to \text{R}(x)$ sending the noncommutative variable $X_i$ to the commutative $x_i$.

Proof. The fact that $\text{FQSym} \xrightarrow{\pi} \text{QSym}$ is a Hopf map comes from checking that it respects the product (compare (5.8) and (8.1)) and the coproducts (compare (5.7) and (8.2)) then applying Proposition 1.46(c). It follows that $\text{NSym} \xleftarrow{\iota} \text{FQSym}$ is a Hopf map once we check the adjointness assertion, via the following calculation:

\[
(\iota(R_\alpha), F_w) = (R_\alpha, F_w) = \sum_{u: \text{Des}(u) = D(\alpha)} (F_{w^{-1}}, F_w) = \begin{cases} 
1 & \text{if } \text{Des}(w) = D(\alpha) \\
0 & \text{otherwise}
\end{cases}
\]

\[= (R_\alpha, L_{\gamma(w)}) = (R_\alpha, \pi(F_w)).\]

The last assertion is clear: the abelianization map sends the noncommutative tableau monomial $X_T$ to the commutative tableau monomial $x_T$. \qed

We summarize some of this picture as follows:

\[
\begin{array}{c}
\text{FQSym} \xrightarrow{\text{dual}} \text{FQSym} \\
\downarrow \iota \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \downarrow \pi \\
\text{NSym} \xrightarrow{\text{dual}} \text{QSym} \\
\downarrow \pi \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \quad \downarrow \text{dual} \\
\Lambda \xrightarrow{\text{dual}} \Lambda
\end{array}
\]

9. Further topics

We summarize some of this picture as follows:

- **Review of representation theory of finite-dimensional algebras.**
  Review the notions of indecomposables, simples, projectives, along with the theorems of Krull-Remak-Schmidt, of Jordan-Hölder, and the two kinds of Grothendieck groups dual to each other.

- **$0$-Hecke algebra representation theory.**
  Describe the simples and projectives, following Denton, Hivert, Schilling, Thiery on $J$-trivial monoids.

- **NSym and Qsym as Grothendieck groups.**
  Give Krob and Thibon’s interpretation of
  - QSym and the Grothendieck group of composition series, and
  - NSym and the Grothendieck group of projectives.

Remark 9.1. Mention P. McNamara’s interpretation, in the case of supersolvable lattices, of the Ehrenborg quasisymmetric function as the composition series enumerator for an $H_n(0)$-action on the maximal chains.
9.0.2. Aguiar-Bergeron-Sottile character theory Part II: Odd and even characters, subalgebras.

9.0.3. Face enumeration, Eulerian posets, and cd-indices. Borrowing from Billera’s ICM notes.
- f-vectors, h-vectors
- flag f-vectors, flag h-vectors
- ab-indices and cd-indices

9.0.4. Other topics.
- Loday-Ronco Hopf algebra of planar binary trees
- Poirier-Reutenauer Hopf algebra of tableaux
- Reading Hopf algebra of Baxter permutations
- Hopf monoids, e.g. of Hopf algebra of generalized permutohedra, of matroids, of graphs, Stanley chromatic symmetric functions and Tutte polynomials
- Lam-Pylyavskyy Hopf algebra of set-valued tableaux
- Connes-Kreimer Hopf algebra and renormalization
- Noncommutative symmetric functions and \( \Omega \Sigma P^\infty \)
- Maschke’s theorem and “integrals” for Hopf algebras
- Nichols-Zoeller structure theorem and group-like elements
- Cartier-Milnor-Moore structure theorem and primitive elements
- Quasi-triangular Hopf algebras and quantum groups
- The Steenrod algebra, its dual, and tree Hopf algebras
- Ringel-Hall algebras of quivers
- Ellis-Khovanov odd symmetric function Hopf algebras (see also Lauda-Russell)

Student talks given in class were:
1. Al Garver, on Maschke’s theorem for finite-dimensional Hopf algebras
2. Jonathan Hahn, on the paper by Humpert and Martin.
3. Emily Gunawan, on the paper by Lam, Lauve and Sottile.
4. Jonas Karlsson, on the paper by Connes and Kreimer
5. Thomas McConville, on Butcher’s group and generalized Runge-Kutta methods.
7. Theodosios Douvropoulos, on the Cartier-Milnor-Moore theorem.
8. Alex Csar, on the Loday-Ronco Hopf algebra of binary trees
9. Kevin Dilks, on Reading’s Hopf algebra of (twisted) Baxter permutations
10. Becky Patrias, on the paper by Lam and Pylyavskyy
11. Meng Wu, on multiple zeta values and Hoffman’s homomorphism from QSym

10. Some open problems and conjectures
- Is there a proof of the Assaf-McNamara skew Pieri rule that gives a resolution of Specht or Schur/Weyl modules whose character corresponds to \( s_{\lambda/\mu} h_n \), whose terms model their alternating sum?
- Explicit antipodes in the Lam-Pylyavskyy Hopf algebras?
- P. McNamara’s question: are \( P \)-partition enumerators irreducible for connected posets \( P \)?
- Stanley’s question: are the only \( P \)-partition enumerators which are symmetric (not just quasisymmetric) those for which \( P \) is a skew-shape with a column-strict labelling?
- Does Stanley’s chromatic symmetric function distinguish trees?
- Hoffman’s stuffle conjecture
- Billera-Brenti’s nonnegativity conjecture for the total cd-index of Bruhat intervals

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