AN INTRODUCTION TO SCHUR FUNCTIONS

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1. **Symmetric Functions.** Consider polynomials in n variables $x_1, \ldots, x_n$. Given a multiindex $\alpha = (\alpha_1, \ldots, \alpha_n)$, let x^α denote the monomial $x_1^{\alpha_1} \cdots x_n^{\alpha_n}$. A *symmetric polynomial* is a polynomial of the form

$$f(x_1, \ldots, x_n) = \sum_\alpha c_\alpha x^\alpha,$$

where, for any permutation $w \in S_n$,

$$c_{(\alpha_1, \ldots, \alpha_n)} = c_{(\alpha_{w(1)}, \ldots, \alpha_{w(n)})}.$$

The integer partition λ obtained by sorting the coordinates of α is called the *shape* of α, denoted $\lambda(\alpha)$. The most obvious example of a symmetric polynomial in n variables is the *monomial symmetric function*, defined for each integer partition λ:

$$m_\lambda = \sum_{\lambda(\alpha) = \lambda} c_\alpha x^\alpha.$$

Note that m_λ is homogeneous of degree $|\lambda|$ (the sum of the parts of λ).

Exercise 1.1. *Take $n = 4$. Compute the monomial symmetric functions $m_{(3)}$, $m_{(2,1)}$, and $m_{(1^3)}$.*

Theorem 1.2. *The polynomials $m_\lambda(x_1, \ldots, x_n)$, as λ runs over all the integer partition of d, form a basis for the space of homogeneous symmetric polynomials of degree d in n variables.*

2. **Complete and Elementary Symmetric Polynomials.** Recall that the coefficients of a polynomial are symmetric polynomials in its roots:

(1) $(t - x_1)(t - x_2) \cdots (t - x_n)$
$$= t^n - e_1(x_1, \ldots, x_n)t^{n-1} + \cdots + (-1)^n e_n(x_1, \ldots, x_n),$$

where coefficient $e_i(x_1, \ldots, x_n)$ of t^{n-i} is given by:

(2) $$e_i(x_1, \ldots, x_n) = \sum_{1 \leq j_1 < \cdots < j_i \leq n} x_{j_1} x_{j_2} \cdots x_{j_i}.$$

The polynomial e_i is called the ith *elementary symmetric polynomial*. By convention, $e_i(x_1, \ldots, x_n) = 0$, for $i > n$.

The identity (1) can be written more elegantly as:

$$(1 + tx_1)\cdots(1 + tx_n) = \sum_{i=0}^{n} e_i(x_1,\cdots,x_n)t^i.$$

Dually[1], the *complete symmetric polynomials* are defined by the formal identity:

$$\frac{1}{(1 - x_1 t)\cdots(1 - x_n t)} = \sum_{i=0}^{\infty} h_i(x_1,\cdots,x_n)t^i.$$

Example 2.1. *In three variables:*

$$e_2(x_1, x_2, x_3) = x_1 x_2 + x_1 x_3 + x_2 x_3,$$
$$h_2(x_1, x_2, x_3) = x_1^2 + x_1 x_2 + x_1 x_3 + x_2^2 + x_2 x_3 + x_2^3.$$

Exercise 2.2. *Show that*

$$h_i(x_1,\cdots x_n) = \sum_{1 \leq j_1 \leq \cdots \leq j_i \leq n} x_{j_1}\cdots x_{j_i}.$$

More generally, for any integer partition $\lambda = (\lambda_1,\ldots,\lambda_l)$, define:

$$h_\lambda = h_{\lambda_1} h_{\lambda_2} \cdots h_{\lambda_l},$$
$$e_\lambda = e_{\lambda_1} e_{\lambda_2} \cdots e_{\lambda_l}.$$

Theorem 2.3. *Given partitions $\lambda = (\lambda_1,\ldots,\lambda_l)$ and $\mu = (\mu_1,\cdots,\mu_m)$ of d, let $M_{\lambda\mu}$ denote the number of matrices (a_{ij}) with non-negative integer entries whose ith row sums to λ_i for each i, and whose jth column sums to μ_j for each j. Then*

$$h_\lambda = \sum_\mu M_{\lambda\mu} m_\mu.$$

Dually, let $= N_{\lambda\mu}$ denote the number of integer matrices (a_{ij}) with entries 0 or 1, whose ith row sums to λ_i for each i, and whose jth column sums to μ_j for each j.

$$e_\lambda = \sum_\mu N_{\lambda\mu} m_\mu.$$

Proof. To prove the second identity involving elementary symmetric functions, note that a monomial in the expansion of

$$e_\lambda = \prod_{i=1}^{l} \sum_{j_1 < \cdots < j_{\lambda_j}} x_{j_1}\cdots x_{j_{\lambda_i}}$$

is a product of summands, one chosen from each of the l factors. Construct an $l \times m$ matrix (a_{ij}) corresponding to such a choice as follows: if the summand $x_{j_1}\cdots x_{j_{\lambda_i}}$ is chosen from the ith factor, then set the entries $a_{i,j_1},\ldots,a_{i,j_{\lambda_j}}$ to be 1 (the remaining entries of the ith row are 0). Clearly

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[1]We will refer to the replacing of $(1+u)$ by $(1-u)^{-1}$ in a formal identity as *dualization.*

the ith row of such a matrix sums to λ_i. The monomial corresponding to this choice is x^μ if, for each j, the the number of i for which x_j appears in the monomial corresponding to the $i$gth row is μ_j. This is just the sum of the jth column of the matrix (a_{ij}). It follows that the coefficient of x^μ, and hence the coefficient of m_μ in the expansion of e_λ in the basis of monomial symmetric functions of degree n, is $N_{\lambda\mu}$.

A similar proof can be given for the first identity involving complete symmetric functions. The only difference is that variables may be repeated in the monomials that appear in h_i. Counting the number of repetitions (instead of just recording 0 or 1) gives non-negative integer matrices. $\square$

3. **Alternating Polynomials.** An *alternating polynomial* in $x_1, \ldots, x_n$ is of the form:

$$(3) \qquad f(x_1, \ldots, x_n) = \sum_\alpha c_\alpha x_\alpha,$$

where, $c_{w(\alpha)} = \epsilon(w) c_\alpha$ for every multiindex α as in Section 1. Here $\epsilon : S_n \to \{\pm 1\}$ denotes the sign function. Equivalently, an alternating polynomial is one whose sign is reversed upon the interchange of any two variables.

Exercise 3.1. *If α is a multiindex where $\alpha_i = \alpha_j$ for some $i \neq j$, then $c_\alpha = 0$.*

In particular, every monomial in an alternating polynomial must be composed of distinct powers. Moreover, the polynomial is completely determined by the coefficients c_α of strictly decreasing multiindices, namely, multiindices of the form $\alpha = (\alpha_1, \ldots, \alpha_n)$ with $\alpha_1 > \cdots > \alpha_n$.

Exercise 3.2. *Let $\delta = (n-1, n-2, \ldots, 1, 0)$. Given an integer partition with at most n parts, we will pad it with 0's so that it can be regarded as a weakly decreasing multiindex of length n. Then $\lambda \mapsto \lambda + \delta$ is a bijection from the set of integer partitions with at most n onto the set of strictly decreasing multiindices.*

Example 3.3. *Let $\lambda = (\lambda_1, \ldots, \lambda_n)$ be a weakly decreasing multiindex. The polynomial:*

$$a_{\lambda+\delta} = \det(x_i^{\lambda_j + n - j})$$

is alternating, with unique strictly decreasing monomial $x^{\lambda+\delta}$.

Exercise 3.4. *The alternating polynomial of the form (3) is equal to*

$$\sum_\lambda c_\lambda a_{\lambda+\delta},$$

the sum being over all weakly decreasing multiindices.

4. **Interpretation of Alternants with Labeled Abaci.** A labeled abacus with n beads is a word $w = (w_k; k \geq 0)$ such that the subword of non-zero letters is a permutation of $1, 2, \ldots, n$. The sign $\epsilon(w)$ of the abacus is the sign of this permutation, the support is the set $\mathrm{supp}(w) = \{k \mid w_k > 0\}$, and the weight is defined as:

$$\mathrm{wt}(w) = \prod_k x_{w_k}^k.$$

The shape of the abacus, $\mathrm{shape}(w)$ is the unique partition λ such that the components of $\lambda + \delta$ form the support of w.

Example 4.1. *Consider the labeled abacus $w = 510032046000\cdots$. Its underlying permutation is 513246, which has sign -1, so $\epsilon(w) = -1$. Also, $\mathrm{supp}(w) = \{0, 1, 4, 5, 7, 8\}$, $\mathrm{shape}(w) = (3, 3, 2, 2)$ (indeed, $(3, 3, 2, 2, 0, 0) + (5, 4, 3, 2, 1, 0) = (8, 7, 5, 4, 1, 0)$) and $\mathrm{wt}(w) = x_5^0 x_1^1 x_3^4 x_2^5 x_4^7 x_6^8$. We visualize the abacus w as a configuration of beads on a single runner, with possible positions of beads numbered $1, 2, 3, \ldots$. If $w_k = i$ where $i > 0$, then a bead labeled i is placed in position k on the runner. If $w_k = 0$, then the position k is unoccupied. In the running example the visualization is:*

0	1	2	3	4	5	6	7	8	9
⑤	①	•	•	③	②	•	④	⑥	•

The first row shows the positions $k = 0, 1, \ldots$ on the runner and the second row shows the beads.

Theorem 4.2. *For every partition λ the alternant in n variables,*

$$a_{\lambda+\delta} = (-1)^{\lfloor n/2 \rfloor} \sum_w \epsilon(w)\mathrm{wt}(w),$$

the sum being over all labeled abaci with n beads and shape λ.

Proof. The theorem follows from the expansion of the determinant. □

5. **Cauchy's Bialternant Form of a Schur Function.** The simplest polynomial of the form $a_{\lambda+\delta}$ arises when $\lambda = 0$; a_δ is the Vandermonde determinant:

$$a_\delta = \prod_{1 \leq i < j \leq n} (x_i - x_j).$$

Exercise 5.1. *Show that, for every weakly decreasing multiindex λ, $a_{\lambda+\delta}$ is divisible by a_δ in the ring of polynomials in $x_1, \ldots, x_n$.*

Exercise 5.2. *Show that $f \mapsto f a_\delta$ is an isomorphism of the space of symmetric polynomials in $x_1, \ldots, x_n$ of degree d onto the space of alternating polynomials of degree $d + \binom{n}{2}$.*

The above exercise gives the historically oldest definition of Schur functions—*Cauchy's bialternant formula*:

(4) $$s_\lambda(x_1, \ldots, x_n) = a_{\lambda+\delta}/a_\delta,$$

for any partition λ with at most n parts. If λ has more than n parts, set $s_\lambda(x_1, \ldots, x_n) = 0$. This is clearly a symmetric function of degree $|\lambda|$.

Theorem 5.3. *As λ runs over all integer partitions of d with at most n parts, the Schur functions $s_\lambda(x_1, \ldots, x_n)$ form a basis of the space of all homogeneous symmetric functions in $x_1, \ldots, x_n$ of degree d.*

Proof. This follows from Exercises 3.4 and 5.2. $\qquad\square$

Exercise 5.4 (Stability of Schur functions). *Show that substituting $x_n = 0$ in the Schur function $s_\lambda(x_1, \ldots, x_n)$ with n variables gives the corresponding Schur function $s_\lambda(x_1, \ldots, x_{n-1})$ with $n - 1$ variables.*

6. **Pieri's rule.** The set of integer partitions is endowed with the *containment order*. We say that a partition $\lambda = (\lambda_1, \ldots, \lambda_l)$ *contains* a partition $\mu = (\mu_1, \ldots, \mu_m)$ if $l \geq m$, and $\lambda_i \geq \mu_i$ for every $i = 1, \cdots, m$. We write $\lambda \supset \mu$ or $\mu \subset \lambda$. Recall that the Young diagram of the partition λ is the set of points

$$\{(i, j) \mid 1 \leq i \leq l, \ 1 \leq j \leq \lambda_i\}.$$

Visually, each node (i, j) of the Young diagram is replaced by a box, and the box corresponding to (i, j) is placed in the ith row and jth column (matrix notation). Thus, the Young diagram of $\lambda = (6, 5, 3, 3)$ is depicted by:
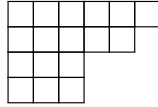
Note that containment of partitions is nothing but the containment relation on their Young diagrams. By abuse of notation, we will also use λ to denote the Young diagram of λ.

By a skew-shape, we mean a difference of Young diagrams $\lambda \setminus \mu$, where $\lambda \supset \mu$. We write λ/μ for this skew-shape. A skew-shape is called a *horizontal strip* (respectively, a *vertical strip*) if it has at most one element in each vertical column (respectively, horizontal row).

Theorem 6.1. *For every partition λ, and every positive integer k,*

$$s_\lambda h_k = \sum_\mu s_\mu,$$

where the sum runs over all partitions $\mu \supset \lambda$ such that μ/λ is a horizontal strip of size k. Dually,

$$s_\lambda e_k = \sum_\mu s_\mu,$$

where the sum runs over all partitions $\mu \supset \lambda$ such that μ/λ is a vertical strip of size k.

Proof. We reproduce the proof from Loehr [?]. Let $\mathrm{Abc}(\lambda)$ denote the set of all n-bead labeled abaci (see Section 4) of shape λ. Let $M(n, k)$ denote the set of all vectors $\alpha = (\alpha_1, \ldots, \alpha_n)$ with non-negative integer coordinates and

sum k. Set $\mathrm{wt}(\alpha) = x_1^{\alpha_1} \cdots x_n^{\alpha_n}$. The first identity is equivalent to showing that:

$$\sum_{w \in \mathrm{Abc}(\lambda)} \epsilon(w) \sum_{\alpha \in M(n,k)} \mathrm{wt}(\alpha) = \sum_{\mu} \sum_{w \in \mathrm{Abc}(\mu)} \epsilon(w)\mathrm{wt}(w),$$

the sum on the right being over all partitions $\mu \supset \lambda$ such that μ/λ is a horizontal strip. We will define an involution I on the $\mathrm{Abc}(\lambda) \times M(n,k)$ whose fixed points correspond to elements of $\coprod_{\mu/\lambda \text{ is a horiz. strip of size } k} \mathrm{Abc}(\mu) \times M(n,k)$ under a bijection that preserves weights and signs, and such that if $I(w, \alpha) = (w', \alpha')$ then $\mathrm{wt}(w)\mathrm{wt}(\alpha) = \mathrm{wt}(w')\mathrm{wt}(\alpha')$ and $\epsilon(w') = -\epsilon(w)$. Then all terms on the left hand side, except for those which do not correspond to fixed points, will cancel, and the surviving terms will give the right hand side.

To construct I, scan the abacus from left to right. Upon encountering a bead numbered j, move the bead α_j steps to the right, one step at a time. If this process completes without this bead colliding with another bead, (w, α) is a fixed point of I. The new abacus w^* has $\epsilon(w^*) = \epsilon(w)$ (the underlying permutation remains unchanged), and $\mathrm{shape}(w^*)/shape(w)$ is a horizontal strip of size k.

However, suppose a collision does occur, say the first collision is when bead j hits bead k that is located $p \leq \alpha_j$ position to the right of its initial position. Define $I(w, \alpha) = (w', \alpha')$, where w' is w with the beads i and j interchanged, $\alpha_j' = \alpha_j - p$, $\alpha_k' = \alpha_k + p$ and all other coordinates of α and α' are equal. Clearly w' has the opposite sign from w, and $\mathrm{wt}(w)\mathrm{wt}(\alpha) = \mathrm{wt}(w')\mathrm{wt}(\alpha')$. One checks that $I(w', \alpha') = (w, \alpha)$.

Example 6.2. *Let $n = 6$, $\lambda = (3, 3, 2, 2, 0, 0)$, $k = 3$, and*

$$(w, \alpha) = (51003204600 \cdots, (2, 1, 0, 0, 0, 0)).$$

Reading the abacus from left to right, the first bead encountered is numbered 1, which can be moved 2 places to the right without any collisions. After that the bead numbered 2 can be moved 1 place to the right, again without collisions. So (w, α) is a fixed point for I. The new abacus $50013024600 \cdots$ has shape $(3, 3, 3, 2, 2, 0)$ obtained by adding a horizontal 3-strip to $(3, 3, 2, 2, 0, 0)$.

On the other hand, if $\alpha = (1, 1, 1, 0, 0, 0)$, then the first collision is of the bead numbered 3 with the bead numbered 2 in the very first step. We have $I(w, \alpha) = (51002304600 \cdots, (1, 2, 0, 0, 0, 0))$.

Let $N(n, k)$ denote the set of vectors $\alpha = (\alpha_1, \ldots, \alpha_n)$ such that $\alpha_i \in \{0, 1\}$ for each i, and $\alpha_1 + \cdots + \alpha_n = k$. For elementary symmetric functions, we wish to prove:

$$\sum_{w \in \mathrm{Abc}(\lambda)} \epsilon(w) \sum_{\alpha \in N(n,k)} \mathrm{wt}(\alpha) = \sum_{\mu} \sum_{w \in \mathrm{Abc}(\mu)} \epsilon(w)\mathrm{wt}(w),$$

where μ runs over all partition such that μ/λ is a vertical strip of size k.

We construct an involution I on $\mathrm{Abc}(\lambda) \times N(n, k)$ as follows: scan the abacus from *right to left*. Upon encountering a bead numbered j, if $\alpha_j = 1$,

try to move the bead one step to the right. If this process completes without collisions, then (w, α) is a fixed point of I. Otherwise, if the first collision occurs with bead numbered j colliding with bead numbered k, then define w' to be w with beads j and k interchanged. Also, since the kth bead was adjacent to the jth bead, it could not have been moved in its turn. So $\alpha_k = 0$. Let α' be obtained from α by interchanging α_k and α_j. $\square$

Example 6.3. *The pair* $(51003204600\cdots, (1,1,1,0,0,0))$ *is a fixed point for* I, *and the shifted abacus is* $(50100324600\cdots)$ *of shape* $(3,3,3,3,1,0)$. *On the other hand*

$$I(51003204600\cdots, (0,0,1,0,1,1)) = (51002304600\cdots, (0,1,0,0,1,1)).$$

The following is a special case of Pieri's rule:

Corollary 6.4. *For every positive integer* k,

$$s_{(k)} = h_k, \ \ and \ s_{(1^k)} = e_k.$$

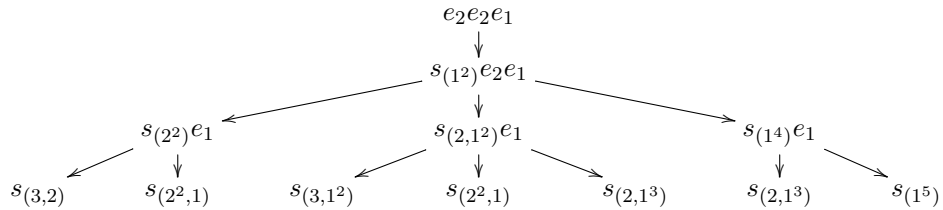**Exercise 6.5.** *Use Pieri's rule to show that:*

$$h_k e_l = s_{(k,1^l)} + s_{(k+1,1^{l-1})}.$$

Conclude that

$$s_{(j+1,1^k)} = \sum_{l=0}^{k}(-1)^l h_{j+l+1} e_{k-l}.$$

7. **Schur to Complete and Elementary via Tableaux.** Pieri's rule allows us to compute the complete and elementary symmetric functions h_λ and e_λ in terms of Schur functions.
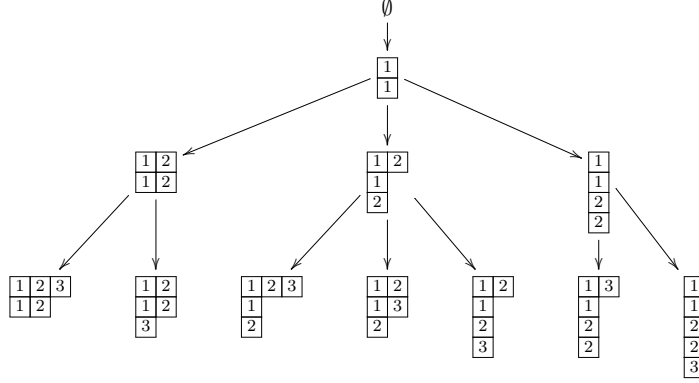
Example 7.1. *Repeated application of Pieri's rule gives an expansion of* $e_{(2,2,1)} = e_2 e_2 e_1$ *as:*
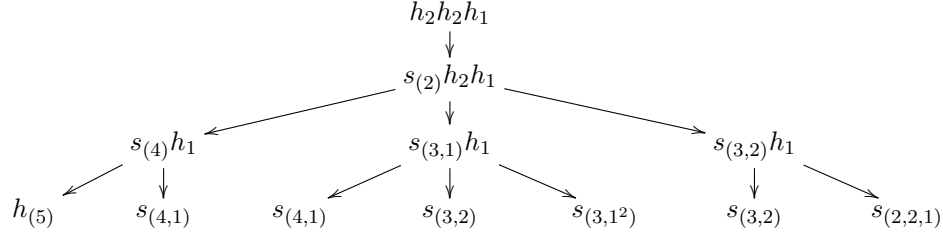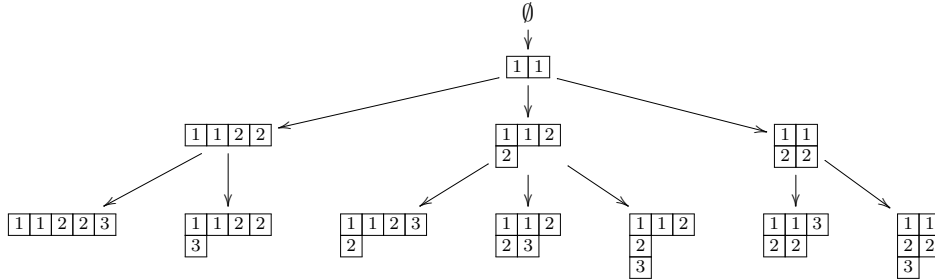

giving:

$$e_{(2^2,1)} = s_{(3,2)} + 2s_{(2^2,1)} + s_{(3,1^2)} + 2s_{(2,1^3)} + s_{(1^5)}.$$

The steps going from the first line of the above calculation to each term of the last line can be recorded by putting numbers into Young diagrams:

$\emptyset$

$\begin{array}{|c|}\hline 1\\\hline 1\\\hline\end{array}$

$\begin{array}{|c|c|}\hline 1&2\\\hline 1&2\\\hline\end{array}$
$\begin{array}{|c|c|}\hline 1&2\\\hline 1\\\cline{1-1} 2\\\cline{1-1}\end{array}$
$\begin{array}{|c|}\hline 1\\\hline 1\\\hline 2\\\hline 2\\\hline\end{array}$

$\begin{array}{|c|c|c|}\hline 1&2&3\\\hline 1&2\\\cline{1-2}\end{array}$
$\begin{array}{|c|c|}\hline 1&2\\\hline 1&2\\\hline 3\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|c|}\hline 1&2&3\\\hline 1\\\cline{1-1} 2\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|}\hline 1&2\\\hline 1&3\\\hline 2\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|}\hline 1&2\\\hline 1\\\cline{1-1} 2\\\cline{1-1} 3\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|}\hline 1&3\\\hline 1\\\cline{1-1} 2\\\cline{1-1} 2\\\cline{1-1}\end{array}$
$\begin{array}{|c|}\hline 1\\\hline 1\\\hline 2\\\hline 2\\\hline 3\\\hline\end{array}$

The boxes in the vertical strip added at the ith stage are filled with i.

Example 7.2. *Repeated application of Pieri's rule gives an expansion of $h_{(2,2,1)} = h_2 h_2 h_1$ as:*

$h_2 h_2 h_1$

$s_{(2)} h_2 h_1$

$s_{(4)} h_1 \qquad s_{(3,1)} h_1 \qquad s_{(3,2)} h_1$

$h_{(5)} \quad s_{(4,1)} \qquad s_{(4,1)} \quad s_{(3,2)} \quad s_{(3,1^2)} \qquad s_{(3,2)} \quad s_{(2,2,1)}$

giving:

$$h_{(2^2,1)} = s_{(5)} + 2s_{(4,1)} + 2s_{(3,2)} + s_{(3,1^2)} + s_{(2,2,1)}.$$

The steps going from the first line of the above calculation to each term of the last line can be recorded by putting numbers into Young diagrams:

$\emptyset$

$\begin{array}{|c|c|}\hline 1&1\\\hline\end{array}$

$\begin{array}{|c|c|c|c|}\hline 1&1&2&2\\\hline\end{array}$
$\begin{array}{|c|c|c|}\hline 1&1&2\\\hline 2\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|}\hline 1&1\\\hline 2&2\\\hline\end{array}$

$\begin{array}{|c|c|c|c|c|}\hline 1&1&2&2&3\\\hline\end{array}$
$\begin{array}{|c|c|c|c|}\hline 1&1&2&2\\\hline 3\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|c|c|}\hline 1&1&2&3\\\hline 2\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|c|}\hline 1&1&2\\\hline 2&3\\\hline\end{array}$
$\begin{array}{|c|c|c|}\hline 1&1&2\\\hline 2\\\cline{1-1} 3\\\cline{1-1}\end{array}$
$\begin{array}{|c|c|c|}\hline 1&1&3\\\hline 2&2\\\hline\end{array}$
$\begin{array}{|c|c|}\hline 1&1\\\hline 2&2\\\hline 3\\\cline{1-1}\end{array}$

The boxes in the horizontal strip added at the ith stage are filled with i.

Definition 7.3 (Semistandard tableau)**.** A semistandard tableau of shape $\lambda = (\lambda_1, \ldots, \lambda_l)$ and type $\mu = (\mu_1, \ldots, \mu_m)$ is the Young diagram of λ filled with numbers $1, \ldots, m$ such that the number i appears μ_i times, the numbers weakly increase along rows, and strictly increase along columns.

Exercise 7.4. *Semistandard tableaux of shape λ and type μ correspond to chains of integer partitions*

$$\emptyset = \lambda^{(0)} \subset \lambda^{(1)} \subset \lambda^{(2)} \subset \cdots \subset \lambda^{(m)} = \lambda$$

where $\lambda^{(i)}/\lambda^{(i-1)}$ is a horizontal strip of size μ_i.

Example 7.5. *The semistandard tableau of type $(3,2)$ and type $(2,2,1)$ are $\begin{array}{ccc}1&1&2\\2&3\end{array}$ and $\begin{array}{ccc}1&1&3\\2&2\end{array}$. They correspond to the chains:*

$$\square\square \subset \square\square\square \subset \square\square\square \quad and \quad \square\square \subset \square\square \subset \square\square\square,$$

respectively. As illustrated in Example 7.2, the coefficient of $s_{(3,2)}$ in the complete symmetric function $h_{(2,2,1)}$ is the number of semistandard tableau of shape $(3,2)$ and type $(2,2,1)$.

Definition 7.6 (Kostka number)**.** Given two partitions λ and μ, the Kostka number $K_{\lambda\mu}$ is the number of semistandard tableau of shape λ and type μ.

Exercise 7.7. *For every partition λ, show that $K_{\lambda\lambda} = 1$.*

Exercise 7.8 (*f*-number)**.** *The f-number of a partition λ of n is defined to be the Kostka number $K_{\lambda,(1^n)}$, and is denoted f_λ.*

Exercise 7.9. *For a partition λ, let λ^- denote the set of all partitions whose Young diagram can be obtained by removing one box from the Young diagram of λ. Show that $f_\lambda = \sum_{\mu \in \lambda^-} f_\mu$.*

Exercise 7.10. *A hook is a partition of the form $h(a,b) = (a+1, 1^b)$. Show that $f_{h(a,b)} = \binom{a+b}{a}$.*

In order to understand the expansion of elementary symmetric functions we would need a variant of semistandard tableaux, one where the difference between successive shapes are vertical strips, rather than horizontal strips. However, it has become common practice to *conjugate* partitions instead:

Definition 7.11 (Conjugate of a partition)**.** The *conjugate* of a partition λ is the partition λ' whose Young diagram is given by:

$$\lambda' = \{(j,i) \mid (i,j) \in \lambda\}.$$

In other words, the Young diagram of λ' is the reflection of the Young diagram of λ about its principal diagonal.

Clearly $\lambda \mapsto \lambda'$ is an involution. For example, if $\lambda = (2,2,1)$, then $\lambda' = (3,2)$.

Exercise 7.12. *Semistandard tableaux of shape λ' and type μ correspond to chains of integer partitions*

$$\emptyset = \lambda^{(0)} \subset \lambda^{(1)} \subset \lambda^{(2)} \subset \cdots \subset \lambda^{(m)} = \lambda$$

where $\lambda^{(i)}/\lambda^{(i-1)}$ is a vertical strip of size μ_i.

Theorem 7.13. *The expansion of complete symmetric functions in terms of Schur functions is given by:*

$$h_\mu = \sum_\lambda K_{\lambda\mu} s_\lambda.$$

Dually, the extension of elementary symmetric functions in terms of Schur functions is given by:

$$e_\mu = \sum_\lambda K_{\lambda'\mu} s_\lambda.$$

8. **Triangularity of Kostka Numbers.** In this section we give a necessary and sufficient condition for the positivity of Kostka number. As usual, take $\lambda = (\lambda_1, \ldots, \lambda_l)$ and $\mu = (\mu_1, \ldots, \mu_m)$. Suppose $K_{\lambda\mu} > 0$. Then there exists a semistandard tableau t of shape λ and type μ. Since the columns of t are strictly increasing, all the 1's in t must occur in its first row, so $\lambda_1 \geq \mu_1$. Also, all the 2's must occur in the first two rows (along with all the 1's), so $\lambda_1 + \lambda_2 \geq \mu_1 + \mu_2$. More generally, all the numbers $1, \ldots, i$ for $i = 1, \ldots, m$ should occur in the first m rows of t. We have:

(5) $$\lambda_1 + \cdots + \lambda_i \geq \mu_1 + \cdots \mu_i \text{ for } i = 1, \ldots, m.$$

Definition 8.1. We say that an integer partition $\lambda = (\lambda_1, \ldots, \lambda_l)$ *dominates* the integer partition $\mu = (\mu_1, \ldots, \mu_m)$ $|\lambda| = |\mu|$ and if the condition (5) holds. When this happens we write $\lambda \rhd \mu$. This relation defines a partial order on the set of all integer partitions of n for any non-negative integer n.

Exercise 8.2. *Show that (n) is maximal and (1^n) is minimal among all the integer partitions of n. What is the smallest integer n for which the dominance order on partitions of n is not a linear order?*

Theorem 8.3 (Triangularity of Kostka Numbers)**.** *Given partition λ and μ of an integer n, $K_{\lambda\mu} > 0$ if and only if $\lambda \rhd \mu$.*

Proof. We have already seen that if $K_{\lambda\mu} > 0$, then $\lambda \rhd \mu$. While reading the proof of the converse, it is helpful to keep in mind Example 8.4 below. Suppose that $\lambda \rhd \mu$. Then $\lambda_1 \geq \mu_1 \geq \mu_m$. Therefore, the Young diagram of λ has at least μ_m cells in its first row, or in other words, it has at least μ_m columns. Choose the smallest integer i for which $\lambda_i \geq \mu_m$. Fill the bottommost box in the λ_{i+1} leftmost columns with m. Also, from the ith row, fill the rightmost $\mu_m - \lambda_{i+1}$ boxes with m. The remaining (unfilled) boxes in the Young diagram of λ now form the Young diagram of the partition

$$\eta = (\lambda_1, \ldots, \lambda_{i-1}, \lambda_i - \mu_m + \lambda_{i+1}, \lambda_{i+2}, \ldots, \lambda_l),$$

a partition with $l - 1$ parts. Writing $(\eta_1, \ldots, \eta_{l-1})$ for the parts of η, note that, since the first $i - 1$ parts of η are the same as those of λ, we have:

$$\eta_1 + \cdots + \eta_j \geq \mu_1 + \cdots + \mu_j$$

for $j \leq i - 1$. For $j \geq i$, we have

$$\eta_1 + \cdots + \eta_j = \lambda_1 + \cdots + \lambda_{j+1} - \mu_m$$
$$\geq \mu_1 + \cdots + \mu_j + \mu_{j+1} - \mu_m$$
$$\geq \mu_1 + \cdots + \mu_j.$$

It follows that $\eta \rhd (\mu_1, \ldots, \mu_{m-1})$. The result now follows by induction in m. $\qquad\square$

Example 8.4. *Consider the case where $\lambda = (7, 3, 2)$ and $\mu = (4, 4, 4)$. Then the smallest integer i such that $\lambda_i \geq 4$ is 1. Accordingly, we enter 3 into the bottom-most boxes in the three leftmost columns, and also into one rightmost box in the first row:*

We are left with the problem of finding a semistandard tableau of shape $(6, 2)$ and type $(4, 4)$. Recursively applying our process to this smaller problem gives:

and finally the desired tableau

Theorem 8.5. *The complete symmetric functions:*

$$\{h_\mu \mid \mu \text{ is a partition of } d \text{ with at most } n \text{ parts}\}$$

and the elementary symmetric functions:

$$\{e_\mu \mid \mu \text{ is a partition of } d \text{ with } \mu_1 \leq n\}$$

form bases of the space of homogeneous symmetric polynomials of degree d in variables $x_1, \ldots, x_n$.

Proof. In view of the triangularity of Kostka numbers (Theorem 8.3) and the fact that $K_{\lambda\lambda} = 1$ (Exercise 7.7) the theorem follows from Theorem 7.13. $\quad\square$

9. **The hook-length formula.** Let λ be a partition of n. The f-number $f_\lambda = K_{\lambda,(1^n)}$ is given by the hook-length formula:

$$(6) \qquad\qquad f_\lambda = \frac{n!}{\prod_{(i,j) \in \lambda} h_{ij}}.$$

Here h_{ij} denotes the number of cells in the Young diagram of λ that lie either to the right of (i, j) in the ith row, or below (i, j) in the jth column, including

the cell (i, j) itself. For example, the Young diagram of $\lambda = (6, 5, 3, 3)$ with the hook lengths are entered into the cells is:

9	8	7	4	3	1
7	6	5	2	1	
4	3	2			
3	2	1			

Theorem 9.1 (Frobenius dimension formula). *For every partition λ of n,*

$$f_\lambda = \frac{n!}{\prod_{i=1}^n (\lambda_1 + n - i)!} a_\delta(\lambda_1 + n - 1, \lambda_2 + n - 2, \ldots, \lambda_n).$$

Proof. Work with symmetric polynomials in n variables. The starting point of the proof of the Frobenius dimension formula is that f_λ is the coefficient of s_λ in the elementary symmetric function $e_{(1^n)} = (x_1 + \cdots + x_n)^n$ (by Theorem 7.13). In terms of alternating functions, we have that f_λ is the coefficient of $a_{\lambda+\delta}$ in $(x_1 + \cdots + x_n)^n a_\delta$. This is the same as the coefficients of the strictly decreasing monomial $x^{\lambda+\delta}$ in $(x_1 + \cdots + x_n)^n a_\delta$. Expanding out the latter sing the multinomial theorem give:

$$\sum_{\alpha_1 + \cdots + \alpha_n = n} \sum_{w \in S_n} \frac{n!}{\alpha_1! \cdots \alpha_n!} \epsilon(w) x_1^{\alpha_1 + \delta_{w(1)}} x_2^{\alpha_2 + \delta_{w(2)}} \cdots x_n^{\alpha_n + \delta_{w(n)}}.$$

The coefficient of $x^{\lambda+\delta}$ in the above sum is:

$$\sum_{\alpha, w} \frac{n!}{\alpha_1! \cdots \alpha_n!} \epsilon(w),$$

the sum being over all vectors $\alpha = (\alpha_1, \ldots, \alpha_n)$ with non-negative integer coordinates summing to n such that:

$$(7) \qquad\qquad\qquad \alpha_i + \delta_{w(i)} = \lambda_i + \delta_i.$$

Mutliplying and dividing by the desired denominator of the Frobenius dimension formula gives:

$$f_\lambda = \frac{n!}{\prod_{i=1}^n (\lambda_1 + n - i)!} \sum_{\alpha, w} \epsilon(w) \prod_{i=1}^n \frac{(\lambda_i + \delta_i)!}{\alpha_i!}$$

$$= \frac{n!}{\prod_{i=1}^n (\lambda_1 + n - i)!} \sum_{w \in S_n} \epsilon(w) \prod_{i=1}^n (\lambda_i + \delta_i)(\lambda_i + \delta - 1) \cdots (\lambda_i + \delta_i - \delta_{w(i)} + 1),$$

where the second step uses the identity (7) relation λ and α. The condition that $\alpha \geq 0$ for each i is now automatically taken care of, since one of the factors in the product

$$(\lambda_i + \delta_i)(\lambda_i + \delta - 1) \cdots (\lambda_i + \delta_i - \delta_{w(i)} + 1)$$

will be zero if $\alpha_i = \lambda_i + \delta_i - \delta_{w(i)} < 0$.

Setting $f_j(x) = x(x-1) \prod(x - j + 1)$, the above expression for f_λ becomes the determinant:

□

10. **Schensted's insertion algorithm.** Let t be a semistandard tableau, and x be a positive integer. Schensted's insertion algorithm is a method of inserting a box with the number x into t, resulting in a new tableau $\mathrm{INSERT}(t, x)$. Applied repeatedly, it gives a way to convert any word into a tableau. This tableau succinctly expresses some combinatorial properties of the original word.

First consider the case where t has a single row, with entries $a_1 \leq \cdots \leq a_k$. Use $\emptyset$ to denote the empty word. The algorithm ι takes as input the single row t and a letter x, and returns a pair (b, t'), where b' is either the empty word, or a single letter, and t' is a row:

$$\iota(a_1 a_2 \cdots a_k, x) = \begin{cases} (\emptyset, a_1 \cdots a_k x) & \text{if } x \geq a_k, \\ (a_j, a_1 \cdots a_{j-1} x a_j \cdots a_k) & \text{if } j = \min\{r \mid a_r > x\}. \end{cases}$$

In the second case, one says that x *has been inserted into* $t = a_1 \cdots a_k$, *obtaining* $t' = a_1 \cdots a_{j-1} x a_j \cdots a_k$, *and* ***bumping out*** a_j. Also, it is notationally convenient to write $\iota(t, \emptyset) = (\emptyset, t)$ (when nothing is inserted, t remains unchanged, and nothing is bumped out).

Now suppose t is a tableau, with rows $r_1, r_2, \ldots, r_l$. The reading word of t is $r_l r_{l-1} \cdots r_1$. Suppose that $\iota(r_1, x) = (y, r_1')$. Recursively define:

$$\mathrm{INSERT}(t, x) = \mathrm{INSERT}(r_l \cdots r_2, y) r_1'$$

Example 10.1. *Consider the insertion of* 3 *into the tableau:*

$$t = \begin{array}{|c|c|c|c|c|} \hline 1 & 3 & 3 & 5 & 8 \\ \hline 2 & 4 & 6 & 6 \\ \cline{1-4} 3 & 5 & 8 \\ \cline{1-3} 4 \\ \cline{1-1} \end{array}.$$

We have $\iota(13358, 3) = (5, 13338)$; $\iota(2466, 5) = (6, 2456)$; $\iota(358, 6) = (8, 356)$; $\iota(4, 8) = (\emptyset, 48)$. *Thus,* $\mathrm{INSERT}(t, 3)$ *is the tableau:*

$$\begin{array}{|c|c|c|c|c|} \hline 1 & 3 & 3 & 3 & 8 \\ \hline 2 & 4 & 5 & 6 \\ \cline{1-4} 3 & 5 & 6 \\ \cline{1-3} 4 & 8 \\ \cline{1-2} \end{array}.$$

In general, it is not possible to recover t *and* x *from* $\mathrm{INSERT}(t, x)$, *even if we know* x. *For example, the above tableau can be obtained by inserting* 3 *into a different tableau:*

$$\mathrm{INSERT}\left(\begin{array}{|c|c|c|c|c|} \hline 1 & 3 & 3 & 6 & 8 \\ \hline 2 & 4 & 5 \\ \cline{1-3} 3 & 5 & 6 \\ \cline{1-3} 4 & 8 \\ \cline{1-2} \end{array}, 3 \right) = \begin{array}{|c|c|c|c|c|} \hline 1 & 3 & 3 & 3 & 8 \\ \hline 2 & 4 & 5 & 6 \\ \cline{1-4} 3 & 5 & 6 \\ \cline{1-3} 4 & 8 \\ \cline{1-2} \end{array}.$$

Clearly, the shape of $\mathrm{INSERT}(t, x)$ can be obtained by adding one box to the shape of t. If we know the row into which the new box was added, and the value of x, then t can be recovered from $\mathrm{INSERT}(t, x)$. This recovery is based on the fact that ι can be inverted: define

$$\delta(a, a_1 a_2 \cdots a_k) = (a_1 \cdots a_{j-1} a a_{j+1} \cdots a_k, a_j),$$

where $j = k$ if $a_r \leq a$ for all $r = 1, \ldots, k$ and $j = \min\{r \mid a_{r+1} > a\}$. To recover t and x from $s = \text{INSERT}(t, x)$ and r, delete the last entry of the rth row of s, say x_r. Let u_{r-1} denote the $(r-1)$st row of s. Suppose $\delta(x_r, u_{r-1}) = (v_{r-1}, x_{r-1})$, replace the $(r-1)$st row of s with v_{r-1}. Continue this process until $\delta(x_2, u_1) = (v_1, x_1)$ is obtained and the first row of s is replaced with v_1. The tableau obtained at the end of this process is t, and $x = x_1$. Write $\text{DELETE}(t, r) = (s, x)$. The preceding discussion shows:

Theorem 10.2. *If* $\text{DELETE}(t, r) = (s, x)$*, then* $\text{INSERT}(s, x) = t$*, and* $\text{shape}(t)$ *is obtained from* $\text{shape}(s)$ *by adding a cell to its* r*th row.*

Exercise 10.3. *Verify Theorem 10.2 for the insertions in Example 10.1.*

11. **Tableaux and Words.** Let A_n^* denote the concatenation monoid of all words in the alphabet $\{1, \ldots, n\}$. For any $w = a_1 \cdots a_k \in A_n^*$, Schensted's insertion algorithm allows us to associate a unique semistandard tableau $P(w)$ as follows:

- If $w = a$ has only one letter, then $P(a)$ is the single-cell tableau with entry a.
- If $w = ua$, where $u \in A_n^*$ and $a \in \{1, \ldots, n\}$, then $P(w) = \text{INSERT}(P(u), a)$.

Example 11.1. *If* $w = 1374433254$*, then* $P(w)$ *is the tableau:*

$$\begin{array}{|c|c|c|c|c|}
\hline
1 & 2 & 3 & 3 & 4 \\
\hline
3 & 4 & 5 \\
\cline{1-3}
4 \\
\cline{1-1}
7 \\
\cline{1-1}
\end{array}$$

Given a semistandard tableau t, its reading word w is defined to be the sequence of numbers obtained from reading its rows from left to right, starting with the bottom row, and moving up sequentially to the top row. Since the first entry of each row is strictly smaller than the last entry of the row below it, the tableau t can be recovered from w by chopping it up into segments with a cut after each a_i with $a_{i+1} < a_i$ (we say that w has a descent at i). The resulting segments, taken from right to left, form the rows of t.

Example 11.2. *The reading word of the tableau* t *formed at the end of Example 8.4 is:*

$$w = 332231111223.$$

The tableau t *is recovered by marking off the descents* $w = 333|223|1111223$*, and then rearranging the segments into a tableau.*

Exercise 11.3. *Let* w *denote the reading word of a tableau* t*. Show that* $P(w) = t$*.*

Not every word comes from a tableau; for example the word 132, when broken up at descents gives rise to $\begin{array}{|c|c|}\hline 2 & \\\hline 1 & 3 \\\hline\end{array}$. We shall say that a word is a tableau if it is the reading word of a semistandard tableau.

Call the word $w = a_1 \cdots a_k$ a *row* if $a_1 \leq \cdots \leq a_k$. Call it a *column* if $a_1 > \cdots > a_k$. Write x^w for the monomial $x_{a_1} x_{a_2} \cdots x_{a_k}$.

Exercise 11.4. *Show that, for every positive integer* i,

$$h_k(x_1, \ldots, x_n) = \sum_{w \,\in\, A_n^* \ is \ a \ row \ of \ length \ k} x^w,$$

and

$$e_k(x_1, \ldots, x_n) = \sum_{w \,\in\, A_n^* \ is \ a \ column \ of \ length \ k} x^w.$$

If w_1 and w_2 are words, and $w_1 w_2$ is their concatenation, then

$$x^{w_1} x^{w_2} = x^{w_1 w_2}.$$

This gives rise to an algebra homomorphism ev $: \mathbf{Z}[A_n^*] \to \mathbf{Z}[x_1, \ldots, x_n]$ called the *evaluation map*.

$$\mathbf{H}_k := \sum_{w \,\in\, A_n^* \ is \ a \ row \ of \ length \ k} w$$

or the elements

$$\mathbf{E}_k = \sum_{w \,\in\, A_{n}* \ is \ a \ column \ of \ length \ k} w$$

Then Exercise 11.4 can be restated as the identities:

$$e_k = \mathrm{ev}(\mathbf{E}_k) \ \text{and} \ h_k = \mathrm{ev}(\mathbf{H}_k).$$

Obviously, ev has a large kernel; is domain is the free algebra, and it maps onto the polynomial algebra. The algebra on the right contains our primary object of interest—the algebra $\mathbf{Z}[x_1, \ldots, x_n]^{S_n}$ of symmetric polynomials. In the next few sections, we shall learn about an equivalence relation "≡" on A_n^* such that the resulting quotient monoid $\mathrm{Pl}_n := A_n^* / \equiv$ (called the *plactic monoid*) has the property that the subalgebra of $\mathbf{Z}[\mathrm{Pl}_n]$ generated by the elements is isomorphic to $\mathbf{Z}[x_1, \ldots, x_n]^{S_n}$ under the map $w \mapsto x^w$.

12. **The Plactic Monoid.** The plactic monoid $\mathrm{Pl}(A_n)$ is the quotient of A_n^* by the equivalence relation generated by the Knuth relations:

$(K1)$ \hspace{3em} $xzy \equiv zxy$ if $x \le y < z$,

$(K2)$ \hspace{3em} $yxz \equiv yzx$ if $x < y \le z$.

Two words are said to be in the same plactic class if each can be obtained from the other by a sequence of moves of the form $(K1)$ and $(K2)$. Since both sides of the Knuth relations have the same evaluation, it follows that the evaluation map ev $: \mathbf{Z}[A_n^*] \to \mathbf{Z}[x_1, \ldots, x_n]$ factors through the plactic monoid algebra $\mathbf{Z}[\mathrm{Pl}(A_n)]$. Let E_k denote the image of $\mathbf{E}_k$ and H_k denote the image of $\mathbf{H}_k$ in $\mathbf{Z}[\mathrm{Pl}(A_n)]$.

Exercise 12.1. *Take* $n = 2$. *Show that* E_1 *and* E_2 *commute in* $\mathbf{Z}[\mathrm{Pl}(A_2)]$. *Show that they commute in* $\mathbf{Z}[\mathrm{Pl}(A_3)]$.

Exercise 12.2. *Define Schützenberger's forgotten relations by:*

$(F1)$ $\qquad\qquad\qquad\qquad xzy \cong yxz \ \ if \ x < y < z,$

$(F2)$ $\qquad\qquad\qquad\qquad zxy \cong yzx \ \ if \ x \le y \le z.$

Let $F(A_n)$ denote the monoid $A_n^*/\cong$. Show that the images of $\mathbf{E}_1$ and $\mathbf{E}_2$ commute in $F(A_3)$.

Exercise 12.3. *Show that any evaluation-preserving equivalence on A_3^* under which the images of $\mathbf{E}_1$ and $\mathbf{E}_2$ commute must include either the Knuth equivalences or Schützenberger's forgotten equivalences.*

13. **The plactic Pieri rules.** Observe that, if for any tableau t and $x \in A_n$, if $t' = \mathrm{INSERT}(t, x)$, then $\mathrm{shape}(t')$ is obtained by adding one box to $\mathrm{shape}(t)$.

Lemma 13.1. *Let t be the (the reading word of) a semistandard tableau in A_n^* and x, y be letters in A_n. Let $t' = \mathrm{INSERT}(t, x)$ and $t'' = \mathrm{INSERT}(t', y)$. Let a be the box added to $\mathrm{shape}(t')$ to obtain $\mathrm{shape}(t'')$, and b be the box added to $\mathrm{shape}(t)$ to obtain $\mathrm{shape}(t')$. If $x \le y$, then b lies in a column strictly to the right of the column of a. If $x > y$, then b lies in a row strictly below the row of a.*

Proof. Let $a_1 \cdots a_k$ be the first row of t. Consider first the case where $a \le y$. If $a_k \le x \le y$, then the result is obvious. If $x < a_k$ and $y \ge a_k$, then b lies in the $(k+1)$st column, whereas x bumps some letter $x' \le y$ to a lower row. This letter cannot come to rest in the $(k+1)$st column because that would violate the fact that columns increase strictly in a semistandard tableau. If $y < a_k$ as well, then x bumps out x' and y bumps out y' with $x' \le y'$. The problem is now reduced to the tableau obtained by removing the top row of t, allowing for the application of induction. In the base case (where the original tableau t is a row), a and b are the first and second boxes in the second row of t''.

Now suppose $x > y$. If $x \ge a_k$, then x first comes to rest at the end of the first row in t', but then y bumps some element of the first row of t' up to a lower row in t''. So a lies in the first row and b in a lower row. If $x \le a_k$, then the elements x' and y' bumped out from the first row by x and y respectively again satisfy $x' > y'$, allowing for an inductive argument. $\square$

Theorem 13.2 (Plactic Pieri rules). *Let $\mathrm{Tab}_n(\lambda)$ denote the set of all semistandard tableaux of shape λ and entries in A_n. Let $R_k(A_n)$ denote the set of rows of length k in A_n^*. Then the map $(t, r) \mapsto P(tr)$ defines a bijection:*

$$\mathrm{Tab}_n(\lambda) \times R_k(A_n) \to \coprod_{\mu/\lambda \ is \ a \ horizontal \ strip \ of \ size \ k} \mathrm{Tab}_n(\mu).$$

Let $C_k(A_n)$ denote the set of columns of length k in A_n^*. Then the map $(t, c) \mapsto P(tc)$ defines a bijection:

$$\mathrm{Tab}_n(\lambda) \times R_k(A_n) \to \coprod_{\mu/\lambda \ is \ a \ vertical \ strip \ of \ size \ k} \mathrm{Tab}_n(\mu).$$

Proof. Lemma 13.1 implies that shape($P(wr)$) is obtained from shape($P(w)$) by adding a horizontal strip, and that shape($P(wc)$) is obtained from shape($P(w)$) by adding a vertical strip. The bijectivity can be shown by repeated used of the DELETE algorithm (because we know which box from μ has to be removed at each step). $\qquad\square$

14. **The Lindström-Gessel-Viennot Lemma.** Let R be a commutative ring. Let S be any set of points, and $v : S \times S \to R$ be any function (we think of w as a *weight function*. Given $s, t \in S$, a path in S from s to t is is a sequence $\omega = (s = s_0, s_1, \ldots, s_k = t)$ of distinct points in S. We denote this by $\omega : s \to t$. The weight of the path ω is defined to be:

$$v(\omega) = v(s_0, s_1) v(s_1, s_2) \cdots v(s_{k-1}, s_k).$$

Definition 14.1 (Crossing paths). Two paths $\omega = (s_0, \ldots, s_k)$ and $\eta = (t_0, \ldots, t_l)$ are said to cross if $s_i = t_j$ for some $0 \leq i \leq k$ and $0 \leq j \leq l$.

Definition 14.2 (Crossing condition). Given a set S of points, a weight function $v : S \times S \to R$, and a points $A_1, \ldots, A_n$, and $B_1, \ldots, B_n$, we say that the *crossing condition* is satisfied if, whenever $1 \leq i < j \leq n$ and $1 \leq i' < j' \leq n$, and $\omega : i \to j'$ and $\eta : j \to i'$ are paths such that $v(\omega) \neq 0$ and $v(\eta) \neq 0$, then the paths ω and η cross.

Fix points $A_1, \ldots, A_n$ and $B_1, \ldots, B_n$ in S, and define an $n \times n$ matrix (a_{ij}) by:

$$a_{ij} = \sum_{\omega : A_i \to B_j} v(\omega).$$

Theorem 14.3 (Lindström-Gessel-Viennot Lemma). *Assume that the crossing condition (Definition 14.2) holds. Then the determinant of the matrix (a_{ij}) defined above is given by:*

$$\text{(8)} \qquad \det(a_{ij}) = \sum_{\omega_i : A_i \to B_i} v(\omega_1) \cdots v(\omega_n),$$

where the sum is over all n-tuples $(\omega_1, \ldots, \omega_n)$ of pairwise non-crossing paths $\omega_i : A_i \to B_i$.

Proof. Let P be the set of all n-tuples of paths of the form:

$$\text{(9)} \qquad \bar{\omega} = (\omega_i : A_i \to B_{w(i)}, i = 1, \ldots, n),$$

where w is a permutation of $\{1, \ldots, n\}$. Define the weight of $\bar{\omega} \in P$ by:

$$v(\bar{\omega}) = \prod_{i=1}^{n} v(\omega_i)$$

and its sign by $\epsilon(\bar{\omega}) = \epsilon(w)$. Then the determinant on the left hand side of (8) expands to the sum:

$$\text{(10)} \qquad \sum_{\bar{\omega} \in P} \epsilon(\bar{\omega}) v(\bar{\omega}).$$

Define an involution $I : P \to P$ as follows: given $\bar\omega$ as in (9), if the paths are pairwise non-crossing, then $\bar\omega$ is a fixed point for I. In this case the crossing condition implies that w is the identity permutation. Otherwise, take the least i such that the path $\omega_i = (s_0, \ldots, s_k)$ crosses another path $\omega_j = (t_0, \ldots, t_l)$. Let j be the smallest number such that a point s_j of ω_i lies in the path ω_j, say $s_j = t_r$. Let $I(\bar\omega)$ be the family of paths obtained from $\bar\omega$ by modifying ω_i and ω_j to ω_i' and ω_j' as follows:

$$\omega_i' = (s_0, \ldots, s_j, t_{r+1}, \ldots, t_l),$$
$$\omega_j' = (t_0, \ldots, t_r, s_{j+1}, \ldots, s_k).$$

Clearly, $v(I(\bar\omega)) = v(\bar\omega)$ and $\epsilon(I(\bar\omega)) = -\epsilon(\bar\omega)$. This involution cancels out all the terms in (10) except those that occur on the right hand side of (8). □

15. **The Jacobi-Trudi Identities.** We have seen that the Kostka numbers can be used to express complete and elementary symmetric functions in terms of Schur functions. The reverse operation—that of expressing Schur functions in terms of complete or elementary symmetric functions—is done by the Jacobi-Trudi identities:

Theorem 15.1 (Jacobi-Trudi identities). *For every integer partition $\lambda = (\lambda_1, \ldots, \lambda_l)$ form the $l \times l$ matrices with (i,j)th entry $h_{\lambda_i - i + j}$ and $e_{\lambda_i - i + j}$ respectively. Then*

$$s_\lambda = \det(h_{\lambda_i - i + j}), \quad s_{\lambda'} = \det(e_{\lambda_i - i + j}).$$

Proof. The Jacobi-Trudi identities can be proved using the Lindström-Gessel-Viennot lemma (Theorem 14.3). For the first identity take S to be the positive cone in the the integer lattice:

$$S = \{(i,j) \mid i \geq 0, \ j > 0 \text{ are integers}\}.$$

Set the weight $v((i,j), (i+1,j))$ of each rightward horizontal edge to be x_j for $j = 1, \ldots, n$, the weight of each upward vertical edge $v((i,j), (i,j+1))$ to be 1 for all $j = 1, \ldots, n-1$. The remaining weights are all zero.

Lemma 15.2. *For all integers $i > 0$ and $j \geq 0$, we have:*

$$\sum_{\omega : (i,1) \to (i+j,n)} v(\omega) = h_j(x_1, \ldots, x_n).$$

Proof. Only rightward or upward steps have non-zero weights. So every path with non-zero weight is composed of unit upward and rightward steps. A path with non-zero weight from $(i,1)$ to $(i+j,n)$ must have exactly j rightward steps, say in rows $1 \leq i_1 \geq i_2 \cdots \leq i_j \leq n$. The weight of such a path is $x_{i_1} \cdots x_{i_j}$, and hence, the sum of the weights of all such paths is $h_j(x_1, \ldots, x_n)$. For an example, see Fig. 1. □
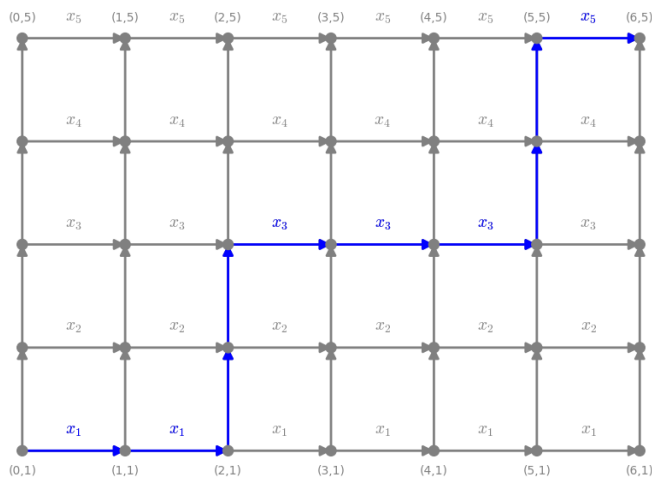
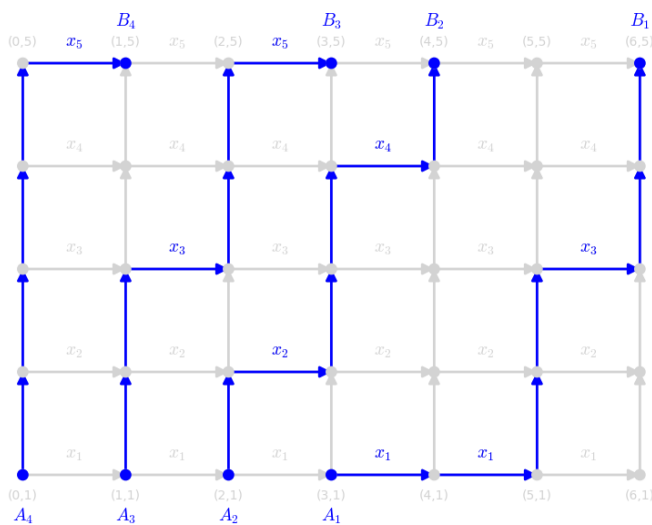FIGURE 1. A path from $(0,5)$ to $(6,1)$ whose weight is the monomial $x_1^2 x_3^3 x_5$ in $h_6(x_1, \ldots, x_5)$.

FIGURE 2. Non-crossing paths corresponding to the tableau

1	1	3
2	4	
3	5	
5		

Given $\lambda = (\lambda_1, \ldots, \lambda_l)$, and working with n variables $x_1, \ldots, x_n$, let $A_i = (l - i, 1)$ and $B_i = (\lambda_i + l - i, n)$ for $i = 1, \ldots, l$. Then by Lemma 15.2,

$$\sum_{\omega : A_i \to B_j} v(\omega) = h_{\lambda_j + i - j}.$$

So the left-hand-side of the first Jacobi-Trudi identity is the left-hand-side of the Lindström-Gessel-Viennot lemma. The right hand side of the

Linström-Gessel-Viennot lemma consists of a sequence of non-crossing paths $(\omega_1, \ldots, \omega_n)$, where $\omega_i : A_i \to B_i$. Reading the row numbers of the horizontal steps in ω_i gives a weakly increasing sequence of integers $1 \leq k_1 \leq \cdots \leq k_{\lambda_i} \leq n$. Enter these numbers into the ith row of the Young diagram of λ for $i = 1, \ldots, n$. Since the paths are non-crossing, the jth rightward step of ω_i must be strictly higher than the jth rightward step of ω_{i+1}. This means that the columns of the resulting numbering are strictly increasing, resulting in a semistandard tableau of shape λ.

Figure 2 shows the non-crossing path configuration corresponding to $n = 5$, $\lambda = (3, 2, 2, 1)$ which corresponds to the semistandard tableau

$$\begin{array}{|c|c|c|}\hline 1 & 1 & 3 \\\hline 2 & 4 \\\cline{1-2} 3 & 5 \\\cline{1-2} 5 \\\cline{1-1}\end{array}$$

. Thus, it follows from the Lindström-Gessel-Viennot lemma that

$$\det(h_{\lambda_j + i - j}) = \sum_{t \in \mathrm{Tab}(\lambda)} x^t.$$

For the second Jacobi-Trudi identity take

$$S = \{(i, j) \mid i \geq 0, \ j \geq 0\}.$$

Define $v((i, j), (i + 1, j)) = 1$ (as before) and $v((i - 1, j), (i, j + 1)) = x_i$; all other weights are zero. For the new weights, the analog of Lemma 15.2 is:

Lemma 15.3. *For all integers $i > 0$ and $j > 0$, we have:*

$$\sum_{\omega : (i, 0) \to (i + j, n)} v(\omega) = e_j(x_1, \ldots, x_n).$$

Proof. Every path with non-zero weights consists of unit upward or upper-rightward diagonal steps. A path with non-zero weight from $(i, 0)$ to $(i + j, n)$ must have n such steps, of which j must be diagonal. If the steps numbered $i_1, \ldots, i_j$ are the diagonal steps, then the path has weight $x_{i_1} \cdots x_{i_j}$. For an example of such a path, see Fig. 3. Summing over all possible paths gives $e_j(x_1, \ldots, x_n)$. $\square$

Suppose that the conjugate partition of λ is $\lambda' = (\lambda'_1, \ldots, \lambda'_k)$. In order to apply the Lindström-Gessel-Viennot lemma to obtain the second Jacobi-Trudi identity, take $A_i = (k - i, 0)$ and $B_i = (\lambda'_i + k - i, n)$ for $i = 1, \ldots, k$. Then by Lemma 15.3,

$$\sum_{\omega_i : A_i \to B_j} v(\omega) = e_{\lambda_j + i - j}.$$

So the left-hand-side of the second Jacobi-Trudi identity is the left-hand-side of the Lindström-Gessel-Viennot lemma.

The right hand side of the Linström-Gessel-Viennot lemma consists of a sequence of non-crossing paths $(\omega_1, \ldots, \omega_n)$, where $\omega_i : A_i \to B_i$. Reading the row numbers where the upper-rightward steps in ω_i originate gives a strictly increasing sequence of integers $1 \leq k_1 \leq \cdots \leq k_{\lambda'_i} \leq n$. Enter these numbers into the ith column of the Young diagram of λ. Since the paths are non-crossing, the jth upper-rightward step of ω_i must be no lower than the jth upper-rightward step of ω_{i+1}. This means that the rows of the resulting numbering are weakly increasing, resulting in a semi-standard tableau of shape λ. Figure 2 shows the non-crossing path configuration
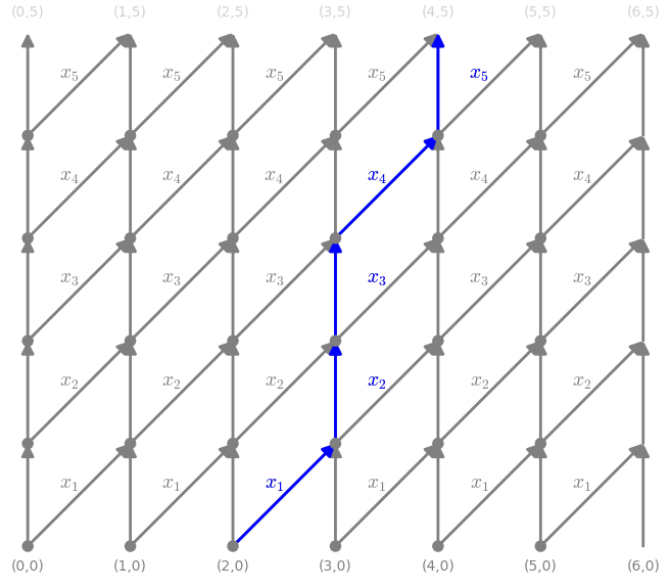
FIGURE 3. A path from $(2,0)$ to $(4,5)$ whose weight is the monomial $x_1 x_4$ in $e_2(x_1, \ldots, x_5)$.
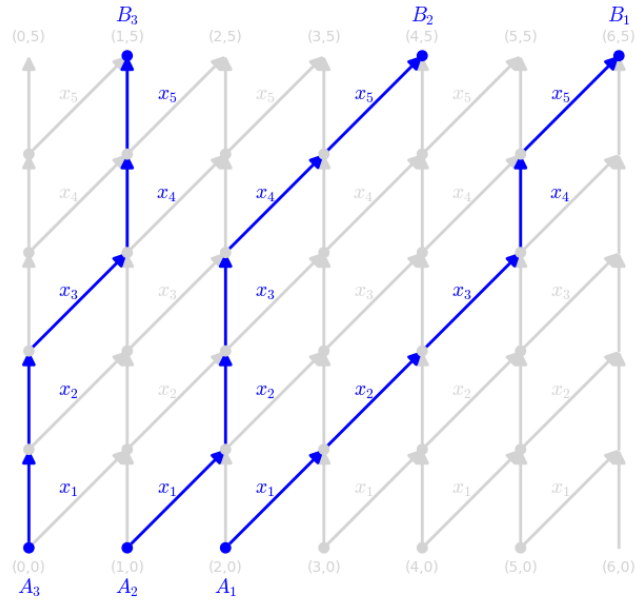

FIGURE 4. Non-crossing paths corresponding to the tableau

1	1	3
2	4	
3	5	
5		

.

corresponding to $n = 5$, $\lambda = (3, 2, 2, 1)$ which corresponds to the semistandard tableau $\begin{array}{|c|c|c|}\hline 1 & 1 & 3 \\\hline 2 & 4 \\\cline{1-2} 3 & 5 \\\cline{1-2} 5 \\\cline{1-1}\end{array}$. Thus, it follows from the Lindström-Gessel-Viennot lemma that

$$\det(e_{\lambda_j + i - j}) = \sum_{t \in \mathrm{Tab}(\lambda')} x^t,$$

proving the second Jacobi-Trudi identity. $\square$

16. Giambelli's identity.

Definition 16.1 (Frobenius coordinates)**.** Let $\lambda = (\lambda_1, \ldots, \lambda_l)$ be a partition. Its *Drufee rank* d is defined to be the largest integer i such that (i, i) lies in the Young diagram of λ. Let α_i denote the number of cells in the ith row that lie strictly to the right of (i, i) in the Young diagram of λ. Similarly let β_i denote the number of cells in the ith column that lie strictly below (i, i). Clearly $\alpha_1 > \cdots > \alpha_d$, $\beta_1 > \cdots > \beta_d)$, and the Young diagram of λ can be recovered from the data $(\alpha | \beta) = (\alpha_1, \ldots, \alpha_d | \beta_1, \ldots, \beta_d)$, which are called the *Frobenius coordinates* of λ[2].

Example 16.2. *The hook partition $(a + 1, 1^b)$ has Frobenius coordinates $(a|b)$. Hook partitions are precisely those partitions which have Durfee rank 1. The partition with frobenius coordinates $(5, 2, 1 | 4, 3, 0)$ is $(6, 4, 4, 2, 2)$. If λ has Frobenius coordinates $(\alpha | \beta)$, then its conjugate λ' has Frobenius coordinates $(\beta | \alpha)$. The size of a partition with Durfee rank d and Frobenius cooedinates $(\alpha | \beta)$ is $d + |\alpha| + |\beta|$.*

Schur functions of hook partitions can be calculated using Exercise 6.5, which, when written is terms of Frobenius coordinates, becomes:

$$(11) \qquad s_{(a|b)} = \sum_{l=0}^{b} (-1)^l h_{a+l+1} e_{b-l}$$

Theorem 16.3 (Giambelli's formula)**.** *For a partition $(\alpha_1, \ldots, \alpha_d | \beta_1, \ldots, \beta_d)$ in Frobenius coordinates,*

$$(12) \qquad s_{(\alpha | \beta)} = \det(s_{(\alpha_i | \beta_j)})_{d \times d}.$$

Note that the determinant on the right consists of hook-partition Schur functions, which are given by (11)*.*

Example 16.4. *The Schur function for $\lambda = (6, 4, 4, 2, 2)$ can be computed as:*

$$s_{(5,2,1|4,3,0)} = \det \begin{pmatrix} s_{(5|4)} & s_{(5|3)} & s_{(5|0)} \\ s_{(2|4)} & s_{(2|3)} & s_{(2|0)} \\ s_{(1|4)} & s_{(1|3)} & s_{(1|0)} \end{pmatrix}$$

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[2]While constructing the character tables of symmetric groups, Frobenius used these coordinates to index the irreducible representation, while he used the ordinary coordinates to index the congugacy classes.

Proof. The determinant on the right hand side of (12) can be written as:

$$(13) \qquad \sum_{w \in S_d} \epsilon(w) \prod_{t^i \in \mathrm{Tab}(\alpha_i, \beta_{w(i)})} x^{t_i}.$$

Let P be the set of all d-tuples of semistandard hook-shaped tableaux $t^1, \ldots, t^d$, where $t^i \in \mathrm{Tab}(\alpha_i, \beta_{w(i)})$. Then we can assemble $(t^1, \ldots, t^d)$ into a tableau t of shape $(\alpha|\beta)$ as follows: The entries in the arm of t^i go into the ith arm of t. The corner entry of t^i goes into the $(w(i), w(i))$th node of t. The entries in the leg of t_i go into the $w(i)$th leg of t. $\square$

17. **The Robinson-Schensted-Knuth Correspondences.** For an $m \times n$ matrix $A = (a_{ij})$, the column word u_A, row word v_A and their duals $\bar{u}_A$ and $\bar{v}_A$ are defined as follows:

$$u_A = 1^{a_{11}} 2^{a_{12}} \cdots n^{a_{1n}} \ 1^{a_{21}} 2^{a_{22}} \cdots n^{a_{2n}} \ \cdots \ 1^{a_{m1}} 2^{a_{m2}} \cdots n^{a_{mn}}$$

$$v_A = 1^{a_{11}} 2^{a_{21}} \cdots m^{a_{m1}} \ 1^{a_{12}} 2^{a_{22}} \cdots m^{a_{m2}} \ \cdots \ 1^{a_{1n}} 2^{a_{2n}} \cdots m^{a_{mn}}$$

$$\bar{u}_A = n^{a_{1n}} \cdots 2^{a_{12}} 1^{a_{11}} \ n^{a_{2n}} \cdots 2^{a_{22}} 1^{a_{21}} \ n^{a_{mn}} \cdots 2^{a_{m2}} 1^{a_{11}}$$

$$\bar{v}_A = m^{a_{m1}} \cdots 2^{a_{21}} 1^{a_{11}} \ m^{a_{m2}} \cdots 2^{a_{22}} 1^{a_{12}} \ \cdots \ m^{a_{mn}} \cdots 2^{a_{2n}} \cdots 1^{a_{1n}}$$

Definition 17.1 (Robinson-Schensted-Knuth Correspondences)**.** Define functions from integer matrices onto pairs of semistandard tableaux by:

$$\mathrm{RSK}(A) = (P(u_A), P(v_A)),$$

$$\mathrm{RSK}^*(A) = (P^*(u_A), P(\bar{v}_A)).$$

Lemma 17.2. *For every integer matrix A with non-negative entries, the tableaux $P(u_A)$ and $P(v_A)$ have the same shape. For every zero-one matrix A, the tableaux $P^*(u_A)$ and $P(\bar{v}_A)$ have the same shape.*

Proof. This follows from Greene's theorem. $\square$

Theorem 17.3 (Knuth's theorem)**.** *Let $\mathbf{M}_{\mu\nu}$ denote the set of integer matrices with non-negative entries, row sums $(\mu_1, \ldots, \mu_m)$, column sums $(\nu_1, \ldots, \nu_n)$. Then RSK gives rise to a bijection:*

$$_{\mu\nu} \xrightarrow{\sim} \coprod_{\lambda} \mathrm{Tab}(\lambda, \nu) \times \mathrm{Tab}(\lambda, \mu).$$