

AN EFFICIENT ALGORITHM FOR DECIDING VANISHING OF SCHUBERT POLYNOMIAL COEFFICIENTS

ANSHUL ADVE, COLLEEN ROBICHAUX, AND ALEXANDER YONG

ABSTRACT. *Schubert polynomials* form a basis of all polynomials and appear in the study of cohomology rings of flag manifolds. The vanishing problem for Schubert polynomials asks if a coefficient of a Schubert polynomial is zero. We give a tableau criterion to solve this problem, from which we deduce the first polynomial time algorithm. These results are obtained from new characterizations of the *Schubertope*, a generalization of the permutahedron defined for any subset of the $n \times n$ grid. In contrast, we show that computing these coefficients explicitly is #P-complete.

1. INTRODUCTION

Schubert polynomials form a linear basis of all polynomials $\mathbb{Z}[x_1, x_2, x_3, \dots]$. They were introduced by A. Lascoux–M.-P. Schützenberger [7] to study the cohomology ring of the flag manifold. These polynomials represent the Schubert classes under the Borel isomorphism. A reference is the textbook [4].

If $w_0 = n \ n - 1 \ \dots \ 2 \ 1$ is the longest length permutation in S_n , then

$$\mathfrak{S}_{w_0}(x_1, \dots, x_n) := x_1^{n-1} x_2^{n-2} \dots x_{n-1}.$$

Otherwise, $w \neq w_0$ and there exists i such that $w(i) < w(i+1)$. Then one sets

$$\mathfrak{S}_w(x_1, \dots, x_n) = \partial_i \mathfrak{S}_{ws_i}(x_1, \dots, x_n),$$

where s_i is the transposition swapping i and $i+1$ and

$$\partial_i f := \frac{f(\dots, x_i, x_{i+1}, \dots) - f(\dots, x_{i+1}, x_i, \dots)}{x_i - x_{i+1}}.$$

Since ∂_i satisfies

$$\partial_i \partial_j = \partial_j \partial_i \text{ for } |i - j| > 1, \text{ and } \partial_i \partial_{i+1} \partial_i = \partial_{i+1} \partial_i \partial_{i+1},$$

the above description of \mathfrak{S}_w is well-defined. In addition, under the inclusion $\iota : S_n \hookrightarrow S_{n+1}$ defined by $w(1) \dots w(n) \mapsto w(1) \dots w(n) \ n+1$, $\mathfrak{S}_w = \mathfrak{S}_{\iota(w)}$. Thus one unambiguously refers to \mathfrak{S}_w for each $w \in S_\infty = \bigcup_{n \geq 1} S_n$.

The *graph* $G(w)$ of a permutation $w \in S_n$ is the $n \times n$ grid, with a \bullet placed in position $(i, w(i))$ (in matrix coordinates). The *Rothe diagram* of w is given by

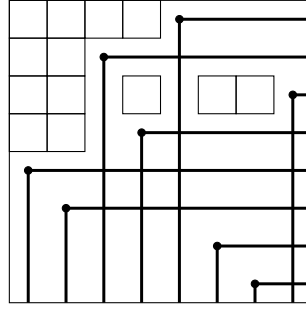
$$D(w) = \{(i, j) : 1 \leq i, j \leq n, j < w(i), i < w^{-1}(j)\}.$$

This is pictorially described with rays that strike out boxes south and east of each \bullet in $G(w)$. $D(w)$ are the remaining boxes.

The *code* of w , denoted $\text{code}(w)$ is the vector (c_1, c_2, \dots, c_L) where c_i is the number of boxes in the i -th row of $D(w)$ and L indexes the southmost row with a positive number of boxes. To each $w \in S_\infty$ there is a unique associated code; see [8, Proposition 2.1.2].

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Example 1.1. If $w = 53841267 \in S_8$ (in one line notation) then $D(w)$ is depicted by:



Here, $\text{code}(w) = (4, 2, 5, 2)$.

Consider the monomial expansion

$$\mathfrak{S}_w = \sum_{\alpha \in \mathbb{Z}_{\geq 0}^n} c_{\alpha,w} x^\alpha.$$

Now, $c_{\alpha,w} = 0$ unless $\alpha_i = 0$ for $i > L$, and moreover, $c_{\alpha,w} \in \mathbb{Z}_{\geq 0}$. Let Schubert be the problem of deciding $c_{\alpha,w} \neq 0$, as measured in the input size of α and w (under the assumption that arithmetic operations take constant time). The INPUT is $\text{code} = (c_1, \dots, c_L) \in \mathbb{Z}_{\geq 0}^L$ with $c_L > 0$ and $\alpha \in \mathbb{Z}_{\geq 0}^L$. Schubert returns YES if $c_{\alpha,w} > 0$ and NO otherwise.

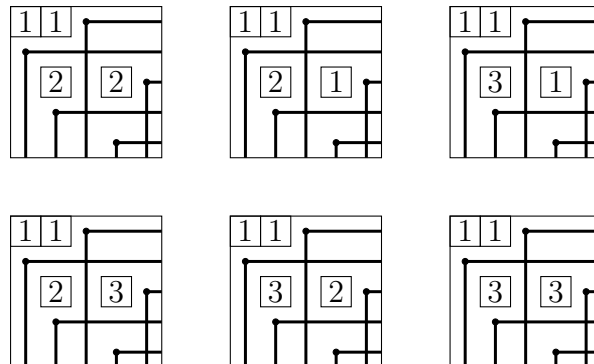
Theorem 1.2. Schubert \in P.

We prove Theorem 1.2 using another result. Fix $n \in \mathbb{Z}_{>0}$ and let $D \subseteq [n]^2$. We call D a *diagram* and visualize D as a subset of an $n \times n$ grid of boxes, oriented so that $(r, c) \in [n]^2$ represents the box in the r th row from the top and the c th column from the left. Let $\text{PerfectTab}(D, \alpha)$ be the fillings of D with α_k many k 's, where entries in each column are distinct, any entry in row i is $\leq i$, and each box contains exactly one entry. Let $\text{PerfectTab}_\downarrow(D, \alpha) \subseteq \text{PerfectTab}(D, \alpha)$ be fillings where entries in each column increase from top to bottom.

Theorem 1.3. $c_{\alpha,w} > 0 \iff \text{PerfectTab}(D(w), \alpha) \neq \emptyset \iff \text{PerfectTab}_\downarrow(D(w), \alpha) \neq \emptyset$

In general $\#\text{PerfectTab}(D(w), \alpha) \neq c_{\alpha,w}$ but rather $\#\text{PerfectTab}(D(w), \alpha) \geq c_{\alpha,w}$ (cf. [3]).

Example 1.4. Here are the tableaux in $\bigcup_\alpha \text{PerfectTab}_\downarrow(D(31524), \alpha)$:



Hence, for instance, $c_{(2,1,1),31524} > 0$ but $c_{(4),31524} = 0$.

To prove Theorems 1.2 and 1.3 we establish results about the *Schubitope* introduced in [9]. This polytope \mathcal{S}_D is defined with a halfspace description for any $D \subseteq [n]^2$. We prove (Theorem 2.13) that a lattice point α is in \mathcal{S}_D if and only if $\text{PerfectTab}(D, \alpha) \neq \emptyset$ where D is any diagram.

We then introduce the *indicator polytope* $\mathcal{P}(D, \alpha)$ whose lattice points $\mathcal{P}(D, \alpha)_{\mathbb{Z}}$ are in bijection with $\text{PerfectTab}(D, \alpha)$. We prove that $\mathcal{P}(D, \alpha) \neq \emptyset \iff \mathcal{P}(D, \alpha)_{\mathbb{Z}} \neq \emptyset$ (Theorem 2.27). Thus determining $\mathcal{P}(D, \alpha)_{\mathbb{Z}} \neq \emptyset$ (and equivalently $\alpha \in \mathcal{S}_D$) is in P using L. Khachiyan's ellipsoid method for linear programming, see [12]. We give two proofs of Theorem 2.27. The first shows $\mathcal{P}(D, \alpha)$ is totally unimodular. Hence $\mathcal{P}(D, \alpha) \neq \emptyset$ implies $\mathcal{P}(D, \alpha)$ has integral vertices. Our second proof obviates total unimodularity and is potentially adaptable to problems lacking that property. However, only the high-level structure of the second proof is easily generalizable — the rest is necessarily *ad hoc*.

For the special case of Rothe diagrams $D = D(w)$, using results of A. Fink-K. Mészáros-A. St. Dizier [2, Corollary 12 and Theorem 14] conjectured in [9, Conjectures 5.1 and 5.13],

$$(1) \quad \alpha \in \mathcal{S}_{D(w)} \iff c_{\alpha, w} > 0.$$

This, combined with our results on the Schubitope, proves Theorems 1.2 and 1.3.

The class #P in L. Valiant's complexity theory of counting problems are those that count the number of accepting paths of a nondeterministic Turing machine running in polynomial time. A problem $\mathcal{P} \in \#P$ is *complete* if for any problem $\mathcal{Q} \in \#P$ there exists a polynomial-time counting reduction from \mathcal{Q} to \mathcal{P} . These are the hardest of the problems in #P. There does not exist a polynomial time algorithm for such problems unless $P = NP$.

In contrast with Theorem 1.2, we prove:

Theorem 1.5. *Counting $c_{\alpha, w}$ is #P-complete.*

Given $\{c_{\alpha, w} \in \mathbb{Z}_{\geq 0}\}$ it is standard to ask for a counting rule for $c_{\alpha, w}$. A complexity motivation is an *appropriate* rule that establishes a counting problem is in #P with respect to given input (length). The rule of [1] establishes that counting $c_{\alpha, w}$ is in #P if the input is (w, α) but not if the input is $(\text{code}(w), \alpha)$. For the latter input assumption, we use the transition algorithm of [6] and its *graphical* reformulation from [5]. This allows us to give a polynomial time counting reduction to the #P-complete problem of counting Kostka coefficients [10], (see Section 5).

2. THE SCHUBITOPE

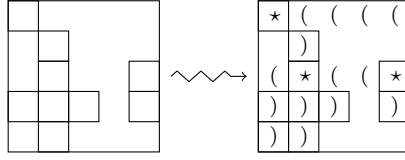
Consider a diagram $D \subseteq [n]^2$. Given $S \subseteq [n]$ and a column $c \in [n]$, construct a string denoted $\text{word}_{c, S}(D)$ by reading column c from top to bottom and recording

- (if $(r, c) \notin D$ and $r \in S$,
-) if $(r, c) \in D$ and $r \notin S$, and
- \star if $(r, c) \in D$ and $r \in S$.

Let $\theta_D^c(S) = \#\{\star\text{'s in } \text{word}_{c, S}(D)\} + \#\{\text{paired } ()\text{'s in } \text{word}_{c, S}(D)\}$ and

$$\theta_D(S) = \sum_{c=1}^n \theta_D^c(S).$$

Example 2.1. In the diagram D below, we labelled the corresponding strings for $\text{word}_{c, S}(D)$ for $S = \{1, 3\}$. For instance, we see $\text{word}_{5, \{1, 3\}}(D) = (\star)$.



The *Schubitope* \mathcal{S}_D , as defined in [9], is the polytope

$$(2) \quad \left\{ (\alpha_1, \dots, \alpha_n) \in \mathbb{R}_{\geq 0}^n : \alpha_1 + \dots + \alpha_n = \#D \text{ and } \sum_{i \in S} \alpha_i \leq \theta_D(S) \text{ for all } S \subseteq [n] \right\}.$$

2.1. **Characterizations via tableaux.** A *tableau* of shape D is a map

$$\tau : D \rightarrow [n] \cup \{\circ\},$$

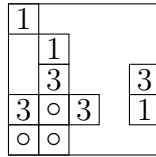
where $\tau(r, c) = \circ$ indicates that the box (r, c) is unlabelled. Let $\text{Tab}(D)$ denote the set of such tableaux.

It will be useful to reformulate the original definition of $\theta_D(S)$ into the language of tableaux. Given $S \subseteq [n]$, define $\pi_{D,S} \in \text{Tab}(D)$ by

$$(3) \quad \pi_{D,S}(r, c) = \begin{cases} r & \text{if } (r, c) \text{ contributes a ``*'' to } \text{word}_{c,S}(D), \\ s & \text{if } (r, c) \text{ contributes a ``)'' to } \text{word}_{c,S}(D) \text{ which is} \\ & \text{paired with an ``('' from } (s, c), \\ \circ & \text{otherwise.} \end{cases}$$

In (3) and throughout, we pair by the standard “inside-out” convention.

Example 2.2. Continuing Example 2.1, below is $\pi_{D,\{1,3\}}(D)$

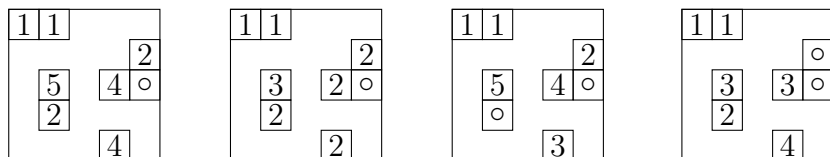


Proposition 2.3. For all $D \subseteq [n]^2$ and $S \subseteq [n]$, we have $\theta_D(S) = \#\pi_{D,S}^{-1}(S)$.

Proof. $\pi_{D,S}(r, c) \in S$ if and only if (r, c) falls into one of the first two cases in (3). □

Say $\tau \in \text{Tab}(D)$ is *flagged* if $\tau(r, c) \leq r$ whenever $\tau(r, c) \neq \circ$. It is *column-injective* if $\tau(r, c) \neq \tau(r', c)$ whenever $r \neq r'$ and $\tau(r, c) \neq \circ$. Let $\text{FCITab}(D) \subseteq \text{Tab}(D)$ be the set of tableaux of shape D which are flagged and column-injective.

Example 2.4. Of the tableaux of shape D below, only the second and fourth are flagged, and only the third and fourth are column-injective.



Proposition 2.5. $\pi_{D,S} \in \text{FCITab}(D)$ for all $D \subseteq [n]^2$ and $S \subseteq [n]$.

Proof. This is immediate from (3). □

A simple consequence of being flagged and column-injective is the following.

Proposition 2.6. *Let $\tau \in \text{FCITab}(D)$. Then for all $(r, c) \in [n]^2$ and $S \subseteq [n]$, we have*

$$(4) \quad \#\{(i, c) \in \tau^{-1}(S) : i < r\} \leq \#\{i \in S : i \leq r\},$$

with strict inequality whenever $(r, c) \in \tau^{-1}(S)$.

Proof. The map $(i, c) \mapsto \tau(i, c)$ from $\{(i, c) \in \tau^{-1}(S) : i \leq r\}$ to $\{i \in S : i \leq r\}$ is well-defined since τ is flagged. It is injective since τ is column-injective. Thus (4) holds, and

$$\#\{(i, c) \in \tau^{-1}(S) : i < r\} < \#\{(i, c) \in \tau^{-1}(S) : i \leq r\} \leq \#\{i \in S : i \leq r\}$$

whenever $(r, c) \in \tau^{-1}(S)$, establishing the strict inequality assertion. □

In fact, a stronger assertion holds when $\tau = \pi_{D,S}$.

Proposition 2.7. *If $(r, c) \in D \subseteq [n]^2$ and $S \subseteq [n]$, then*

$$(r, c) \in \pi_{D,S}^{-1}(S) \iff \#\{(i, c) \in \pi_{D,S}^{-1}(S) : i < r\} < \#\{i \in S : i \leq r\}.$$

Proof. (\Rightarrow) This direction follows from Propositions 2.5 and 2.6.

(\Leftarrow) If $r \in S$, then (r, c) contributes a “ \star ” to $\text{word}_{c,S}(D)$, so $\pi_{D,S}(r, c) = r \in S$, as desired. Thus we assume $r \notin S$. The hypothesis combined with this assumption says

$$\#\{(i, c) \in \pi_{D,S}^{-1}(S) : i < r\} < \#\{i \in S : i \leq r\} = \#\{i \in S : i < r\}.$$

Thus, there is a maximal $s \in S$ with $s < r$ such that $\pi_{D,S}(r', c) \neq s$ whenever $r' < r$. If $(s, c) \in D$, then (s, c) contributes a “ \star ” to $\text{word}_{c,S}(D)$, so $\pi_{D,S}(s, c) = s$, contradicting our choice of s . Therefore, (s, c) contributes an “(” to $\text{word}_{c,S}(D)$. If this “(” is paired by a “)” contributed by $(r', c) \in D$ with $r' < r$, then $\pi_{D,S}(r', c) = s$, again a contradiction. Thus, this “(” pairs the “)” from (r, c) , so $\pi_{D,S}(r, c) = s \in S$. Hence, $(r, c) \in \pi_{D,S}^{-1}(S)$ as desired. □

The previous two propositions combined assert that $\{(r, c) \in \pi_{D,S}^{-1}(S)\}$ is characterized by greedy selection as one moves down each column c . The next proposition shows that this greedy algorithm maximizes $\#\tau^{-1}(S)$ among all $\tau \in \text{FCITab}(D)$.

Proposition 2.8. *Let $D \subseteq [n]^2$ and $S \subseteq [n]$. Then $\#\pi_{D,S}^{-1}(S) \geq \#\tau^{-1}(S)$ for all $\tau \in \text{FCITab}(D)$.*

Proof. If not, then there exist $\tau \in \text{FCITab}(D)$ and $(r, c) \in [n]^2$ satisfying

$$\#\{(i, c) \in \pi_{D,S}^{-1}(S) : i \leq r\} < \#\{(i, c) \in \tau^{-1}(S) : i \leq r\}$$

and we can choose these such that r is minimized. Then because r is minimal,

$$\#\{(i, c) \in \pi_{D,S}^{-1}(S) : i < r\} = \#\{(i, c) \in \tau^{-1}(S) : i < r\}$$

and $(r, c) \in \tau^{-1}(S) \setminus \pi_{D,S}^{-1}(S)$, so in particular $(r, c) \in D$. Thus Proposition 2.6 implies

$$\#\{(i, c) \in \pi_{D,S}^{-1}(S) : i < r\} = \#\{(i, c) \in \tau^{-1}(S) : i < r\} < \#\{i \in S : i \leq r\}.$$

But then we must have $(r, c) \in \pi_{D,S}^{-1}(S)$ by Proposition 2.7, a contradiction. □

If τ has shape a subset of $[n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{R}_{\geq 0}^n$, say τ exhausts α over S if

$$\sum_{i \in S} \alpha_i \leq \#\tau^{-1}(S).$$

Example 2.9. Only the left tableau below exhausts $\alpha = (3, 2, 2, 4)$ over $S = \{1, 3\}$.

1				
	1	1		
		3		
		4	4	
4	o	4	3	

1				
	1	2		
		3		
		4	4	
4	o	4	2	

Theorem 2.10. Let $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$ with $\alpha_1 + \dots + \alpha_n = \#D$. Then $\alpha \in \mathcal{S}_D$ if and only if for each $S \subseteq [n]$, there exists $\tau_{D,S} \in \text{FCITab}(D)$ which exhausts α over S .

Proof of Theorem 2.10. (\Rightarrow) The inequalities in (2) combined with Proposition 2.3 imply

$$\sum_{i \in S} \alpha_i \leq \theta_D(S) = \#\pi_{D,S}^{-1}(S).$$

Thus, $\tau_{D,S} := \pi_{D,S}$ exhausts α over S .

(\Leftarrow) By Propositions 2.8 and 2.3,

$$\sum_{i \in S} \alpha_i \leq \#\tau_{D,S}^{-1}(S) \leq \#\pi_{D,S}^{-1}(S) = \theta_D(S),$$

so the inequalities in (2) hold. □

Remark 2.11. The proof of (\Rightarrow) shows that we can take $\tau_{D,S} = \pi_{D,S}$ in Theorem 2.10.

It would be nice if $\tau_{D,S}$ did not depend on S , i.e., if some τ_D exhausted α over all $S \subseteq [n]$, so we could take $\tau_{D,S} = \tau_D$ in Theorem 2.10. Indeed, this is shown in Theorem 2.13.

Say $\tau \in \text{Tab}(D)$ has content α if $\#\tau^{-1}(\{i\}) = \alpha_i$ for each $i \in [n]$. Let $\text{Tab}(D, \alpha)$ and $\text{FCITab}(D, \alpha)$ be the subsets of $\text{Tab}(D)$ and $\text{FCITab}(D)$, respectively, of those tableaux which have content α . In addition, call a tableau $\tau \in \text{Tab}(D)$ *perfect* if $\tau \in \text{FCITab}(D)$, and if no boxes are left unlabelled, i.e., $\tau^{-1}(\{o\}) = \emptyset$. Thus, the set of perfect tableaux of content α is precisely $\text{PerfectTab}(D, \alpha) \subseteq \text{FCITab}(D, \alpha)$ introduced in Section 1.

Proposition 2.12. Let $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$. Then $\text{PerfectTab}(D, \alpha) \neq \emptyset$ if and only if $\alpha_1 + \dots + \alpha_n = \#D$ and $\text{FCITab}(D, \alpha) \neq \emptyset$.

Proof. (\Rightarrow) Let $\tau \in \text{PerfectTab}(D, \alpha)$. Then $\tau \in \text{FCITab}(D, \alpha)$, and since τ has content α and satisfies $\tau^{-1}(\{o\}) = \emptyset$,

$$\alpha_1 + \dots + \alpha_n = \#\tau^{-1}(\{1\}) + \dots + \#\tau^{-1}(\{n\}) = \#D.$$

(\Leftarrow) Let $\tau \in \text{FCITab}(D, \alpha)$. Then since τ has content α ,

$$\#\tau^{-1}(\{o\}) = \#D - \#\tau^{-1}(\{1\}) - \dots - \#\tau^{-1}(\{n\}) = \#D - \alpha_1 - \dots - \alpha_n = 0.$$

Thus, $\tau \in \text{PerfectTab}(D, \alpha)$. □

Theorem 2.13. Let $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$. Then $\alpha \in \mathcal{S}_D$ if and only if $\text{PerfectTab}(D, \alpha) \neq \emptyset$.

The proof will require a lemma regarding tableaux of the form $\tau = \pi_{D,S}$.

Lemma 2.14. Let $D \subseteq [n]^2$, and $S, T \subseteq [n]$ be disjoint. Set

$$\tilde{D} = D \setminus \pi_{D,S}^{-1}(S) \text{ and } U = S \cup T.$$

Then

$$\pi_{D,U}^{-1}(U) = \pi_{D,S}^{-1}(S) \cup \pi_{\tilde{D},T}^{-1}(T).$$

Proof. Let $(r, c) \in D$, and assume by induction on r that

$$(5) \quad (i, c) \in \pi_{D,U}^{-1}(U) \iff (i, c) \in \pi_{D,S}^{-1}(S) \cup \pi_{\tilde{D},T}^{-1}(T)$$

whenever $i < r$. This clearly holds in the base case $r = 1$. By Proposition 2.7, $(r, c) \in \pi_{D,U}^{-1}(U)$ if and only if

$$(6) \quad \#\{(i, c) \in \pi_{D,U}^{-1}(U) : i < r\} < \#\{i \in U : i \leq r\}.$$

By (5) and the fact that

$$\pi_{D,S}^{-1}(S) \cap \tilde{D} = \emptyset = S \cap T,$$

(6) is equivalent to

$$\#\{(i, c) \in \pi_{D,S}^{-1}(S) : i < r\} + \#\{(i, c) \in \pi_{\tilde{D},T}^{-1}(T) : i < r\} < \#\{i \in S : i \leq r\} + \#\{i \in T : i \leq r\}.$$

By applying Proposition 2.6 twice, we see that this holds if and only if at least one of (i) and (ii) below hold.

- (i) $\#\{(i, c) \in \pi_{D,S}^{-1}(S) : i < r\} < \#\{i \in S : i \leq r\}$
- (ii) $\#\{(i, c) \in \pi_{\tilde{D},T}^{-1}(T) : i < r\} < \#\{i \in T : i \leq r\}$

By Proposition 2.7, (i) is equivalent to $(r, c) \in \pi_{D,S}^{-1}(S)$. If indeed $(r, c) \in \pi_{D,S}^{-1}(S)$ holds, then our induction step is complete. Otherwise, $(r, c) \notin \pi_{D,S}^{-1}(S)$, so by definition, $(r, c) \in \tilde{D}$. Thus, applying Proposition 2.7 to \tilde{D} , $T \subseteq [n]$ and $(r, c) \in \tilde{D}$, (ii) is equivalent to $(r, c) \in \pi_{\tilde{D},T}^{-1}(T)$. Hence, (5) holds for all $i \leq r$. \square

Corollary 2.15. Let $D \subseteq [n]^2$ and $S \subseteq U \subseteq [n]$. Then $\pi_{D,S}^{-1}(S) \subseteq \pi_{D,U}^{-1}(U)$.

Proof. Take $T = U \setminus S$ in Lemma 2.14. \square

Finally, we are ready to prove Theorem 2.13.

Proof of Theorem 2.13. (\Leftarrow) Let $\tau_D \in \text{PerfectTab}(D, \alpha)$. Then $\alpha_1 + \cdots + \alpha_n = \#D$ by Proposition 2.12. Also, for each $S \subseteq [n]$,

$$\sum_{i \in S} \alpha_i = \sum_{i \in S} \#\tau_D^{-1}(\{i\}) = \#\tau_D^{-1}(S),$$

so τ_D exhausts α over S . Thus, $\alpha \in \mathcal{S}_D$ by Theorem 2.10.

(\Rightarrow) We induct on the sum of the row indices of each box in D , i.e., $\sum_{(i,j) \in D} i$. The base case of an empty diagram is trivial, so we may assume $D \neq \emptyset$. Then since $\alpha \in \mathcal{S}_D$, (2) implies $\alpha_1 + \cdots + \alpha_n = \#D > 0$, so we can choose m maximal such that $\alpha_m > 0$.

Case 1: (D contains boxes below row m). Pick $(r, c) \in D$ below row m (so $r > m$).

Claim 2.16. There exists $r_1 < r$ such that $(r_1, c) \notin D$.

Proof of Claim 2.16. By Theorem 2.10, there exists $\tau_{D,[m]} \in \text{FCITab}(D)$ such that

$$(7) \quad \#\tau_{D,[m]}^{-1}([m]) \geq \alpha_1 + \cdots + \alpha_m = \alpha_1 + \cdots + \alpha_n = \#D.$$

Thus, $\tau_{D,[m]}(D) \subseteq [m]$. Consequently, by column-injectivity of $\tau_{D,[m]}$, there can be at most m boxes in each column of D . Since $(r, c) \in D$ with $r > m$, there are more than m boxes in column c if $(r_1, c) \in D$ for all $r_1 < r$. Hence there must be some $r_1 < r$ for which $(r_1, c) \notin D$, as asserted. \square

By Claim 2.16, we can choose $r_1 < r$ maximal such that $(r_1, c) \notin D$. Let

$$\tilde{D} = (D \setminus \{(r, c)\}) \cup \{(r_1, c)\}.$$

Claim 2.17. $\alpha \in \mathcal{S}_{\tilde{D}}$.

Proof of Claim 2.17. Since $\alpha \in \mathcal{S}_D$, $(r, c) \in D$, and $(r_1, c) \notin D$, we have

$$\alpha_1 + \cdots + \alpha_n = \#D = \#\tilde{D}.$$

Let $S \subseteq [n]$ and $T = S \cap [m]$. Then define $\tau_{\tilde{D},S} \in \text{Tab}(\tilde{D})$ by

$$\tau_{\tilde{D},S}(i, j) = \begin{cases} \pi_{D,T}(r, c) & \text{if } (i, j) = (r_1, c), \\ \pi_{D,T}(i, j) & \text{otherwise.} \end{cases}$$

If $\pi_{D,T}(r, c) = \circ$, then certainly $\tau_{\tilde{D},S} \in \text{FCITab}(\tilde{D})$. Otherwise, let $s = \pi_{D,T}(r, c)$. Since $(r, c) \in D$ but $r \notin T$, (r, c) contributes a “)” to $\text{word}_{c,S}(D)$. Thus, by (3), (s, c) contributes an “(”, so in particular $(s, c) \notin D$. From our choice of r_1 , we must therefore have $s \leq r_1$, so $\tau_{\tilde{D},S}$ is flagged. Hence, $\tau_{\tilde{D},S} \in \text{FCITab}(\tilde{D})$.

By construction,

$$\#\tau_{\tilde{D},S}^{-1}(\{i\}) = \#\pi_{D,T}^{-1}(\{i\})$$

for each $i \in [n]$, so $\tau_{\tilde{D},S}$ exhausts α over T by Theorem 2.10 and in particular Remark 2.11. Since $\alpha_i = 0$ for all $i > m$, we can write

$$\sum_{i \in S} \alpha_i = \sum_{i \in T} \alpha_i \leq \#\tau_{\tilde{D},S}^{-1}(T) \leq \#\tau_{\tilde{D},S}^{-1}(S).$$

Therefore, $\tau_{\tilde{D},S} \in \text{FCITab}(\tilde{D})$ exhausts α over S , so $\alpha \in \mathcal{S}_{\tilde{D}}$ by Theorem 2.10. \square

Since $r_1 < r$,

$$\sum_{(i,j) \in \tilde{D}} i < \sum_{(i,j) \in D} i.$$

Thus, Claim 2.17 and induction yields $\tau_{\tilde{D}} \in \text{PerfectTab}(\tilde{D}, \alpha)$. Define $\tau_D \in \text{Tab}(D)$ by

$$\tau_D(i, j) = \begin{cases} \tau_{\tilde{D}}(r_1, c) & \text{if } (i, j) = (r, c), \\ \tau_{\tilde{D}}(i, j) & \text{otherwise.} \end{cases}$$

Then it is easy to check that $\tau_D \in \text{PerfectTab}(D, \alpha)$, so Case 1 is complete.

Case 2: (D does not contain boxes below row m). We say an inequality $\sum_{i \in S} \alpha_i \leq \theta_D(S)$ from (2) is *nontrivial* if

$$(8) \quad \sum_{i \in S} \alpha_i > 0 \quad \text{and} \quad \theta_D(S) < \#D.$$

Case 2a: (All nontrivial inequalities from (2) are strict). Thus if (8) holds, then

$$(9) \quad \sum_{i \in S} \alpha_i < \theta_D(S).$$

Claim 2.18. *There exists $c \in [n]$ such that $(m, c) \in D$.*

Proof of Claim 2.18. By Theorem 2.10, there exists some $\tau_{D, \{m\}} \in \text{FCITab}(D)$ which exhausts α over $\{m\}$. Then

$$\#\tau_{D, \{m\}}^{-1}(\{m\}) \geq \alpha_m > 0,$$

so $\tau_{D, \{m\}}(r, c) = m$ for some $(r, c) \in D$. Since $\tau_{D, \{m\}}$ is flagged, we must have $r \geq m$. But by the assumption of Case 2, there are no boxes below row m , so $r = m$. \square

Pick $c \in [n]$ as in Claim 2.18. Then let $\tilde{D} = D \setminus \{(m, c)\}$ and $\tilde{\alpha} = (\tilde{\alpha}_1, \dots, \tilde{\alpha}_n) := (\alpha_1, \dots, \alpha_{m-1}, \alpha_m - 1, 0, \dots, 0)$.

Claim 2.19. $\tilde{\alpha} \in \mathcal{S}_{\tilde{D}}$.

Proof of Claim 2.19. Since $\alpha_i = 0$ for all $i > m$, and $(m, c) \in D$, we have

$$(10) \quad \tilde{\alpha}_1 + \dots + \tilde{\alpha}_n = \alpha_1 + \dots + \alpha_n - 1 = \#D - 1 = \#\tilde{D}.$$

For each $S \subseteq [n]$, let

$$\tau_{\tilde{D}, S} = \pi_{D, S}|_{\tilde{D}} \in \text{FCITab}(\tilde{D})$$

be the restriction of $\pi_{D, S}$ to \tilde{D} . Then by Proposition 2.3,

$$(11) \quad \#\tau_{\tilde{D}, S}^{-1}(S) \geq \#\pi_{D, S}^{-1}(S) - 1 = \theta_D(S) - 1.$$

If $\sum_{i \in S} \alpha_i = 0$, then

$$\sum_{i \in S} \tilde{\alpha}_i = 0 \leq \#\tau_{\tilde{D}, S}^{-1}(S).$$

If $\theta_D(S) = \#D$, then by (10) and (11),

$$\sum_{i \in S} \tilde{\alpha}_i \leq \tilde{\alpha}_1 + \dots + \tilde{\alpha}_n = \#D - 1 = \theta_D(S) - 1 \leq \#\tau_{\tilde{D}, S}^{-1}(S).$$

Finally, if $\sum_{i \in S} \alpha_i > 0$ and $\theta_D(S) < \#D$, then (9) must hold, so by (9) and (11),

$$\sum_{i \in S} \tilde{\alpha}_i \leq \sum_{i \in S} \alpha_i \leq \theta_D(S) - 1 \leq \#\tau_{\tilde{D}, S}^{-1}(S).$$

In all three cases, $\tau_{\tilde{D}, S}$ exhausts $\tilde{\alpha}$ over S , so $\tilde{\alpha} \in \mathcal{S}_{\tilde{D}}$ by Theorem 2.10. \square

By construction,

$$\sum_{(i, j) \in \tilde{D}} i < \sum_{(i, j) \in D} i.$$

Thus, Claim 2.19 and induction yield $\tau_{\tilde{D}} \in \text{PerfectTab}(\tilde{D}, \tilde{\alpha})$. Define $\tau_D \in \text{Tab}(D)$ by

$$\tau_D(i, j) = \begin{cases} m & \text{if } (i, j) = (m, c), \\ \tilde{\tau}(i, j) & \text{otherwise.} \end{cases}$$

Clearly, τ_D is flagged, has content α , and satisfies $\tau_D^{-1}(\{\circ\}) = \emptyset$. The only potential obstruction to column-injectivity is that there could be some $r \neq m$ for which $\tau_D(r, c) = m$.

This is impossible, since τ_D is flagged, so such an r must be greater than m , but by the assumption of Case 2 there are no boxes below row m . Thus, $\tau_D \in \text{PerfectTab}(D, \alpha)$, so Case 2a is complete.

Case 2b: (There exists a tight, nontrivial inequality in (2)). Thus, there exists $A \subseteq [n]$ satisfying

$$(12) \quad 0 < \sum_{i \in A} \alpha_i = \theta_D(A) < \#D.$$

Let $D^{(1)} = \pi_{D,A}^{-1}(A)$ and $D^{(2)} = D \setminus D^{(1)}$. Then for each $i \in [n]$, set

$$\alpha_i^{(1)} = \begin{cases} \alpha_i & \text{if } i \in A, \\ 0 & \text{if } i \notin A \end{cases} \quad \text{and} \quad \alpha_i^{(2)} = \begin{cases} \alpha_i & \text{if } i \notin A, \\ 0 & \text{if } i \in A. \end{cases}$$

Claim 2.20. $\alpha^{(1)} := (\alpha_1^{(1)}, \dots, \alpha_n^{(1)}) \in \mathcal{S}_{D^{(1)}}$.

Proof of Claim 2.20. By (12) and Proposition 2.3, we have

$$\alpha_1^{(1)} + \dots + \alpha_n^{(1)} = \sum_{i \in A} \alpha_i = \theta_D(A) = \#\pi_{D,A}^{-1}(A) = \#D^{(1)}.$$

Let $S \subseteq [n]$ and $T = S \cap A$. Then set

$$\tau_{D^{(1)},S} = \pi_{D,T}|_{D^{(1)}} \in \text{FCITab}(D^{(1)}).$$

By Corollary 2.15, $\pi_{D,T}^{-1}(T) \subseteq D^{(1)}$, so $\tau_{D^{(1)},S}^{-1}(T) = \pi_{D,T}^{-1}(T)$. Thus, by Remark 2.11, $\tau_{D^{(1)},S}$ exhausts α over T . Hence,

$$\sum_{i \in S} \alpha_i^{(1)} = \sum_{i \in T} \alpha_i \leq \#\tau_{D^{(1)},S}^{-1}(T) \leq \#\tau_{D^{(1)},S}^{-1}(S),$$

so $\tau_{D^{(1)},S}$ exhausts $\alpha^{(1)}$ over S , and consequently $\alpha^{(1)} \in \mathcal{S}_{D^{(1)}}$ by Theorem 2.10. \square

Claim 2.21. $\alpha^{(2)} := (\alpha_1^{(2)}, \dots, \alpha_n^{(2)}) \in \mathcal{S}_{D^{(2)}}$.

Proof of Claim 2.21. By (12) and Proposition 2.3,

$$\alpha_1^{(2)} + \dots + \alpha_n^{(2)} = \alpha_1 + \dots + \alpha_n - \sum_{i \in A} \alpha_i = \#D - \theta_D(A) = \#D - \#\pi_{D,A}^{-1}(A) = \#D^{(2)}.$$

Let $S \subseteq [n]$, $T = S \setminus A$, and $U = A \cup T$. Then by Theorem 2.10, Remark 2.11, (12), Proposition 2.3, and Lemma 2.14, we can write

$$\begin{aligned} \sum_{i \in S} \alpha_i^{(2)} &= \sum_{i \in U} \alpha_i - \sum_{i \in A} \alpha_i \leq \#\pi_{D,U}^{-1}(U) - \theta_D(A) \\ &= \#\pi_{D,U}^{-1}(U) - \#\pi_{D,A}^{-1}(A) = \#\pi_{D^{(2)},T}^{-1}(T) \leq \#\pi_{D^{(2)},T}^{-1}(S). \end{aligned}$$

Thus, $\tau_{D^{(2)},S} := \pi_{D^{(2)},T}$ exhausts $\alpha^{(2)}$ over S , so $\alpha^{(2)} \in \mathcal{S}_{D^{(2)}}$ by Theorem 2.10. \square

By (12) and Proposition 2.3, we have

$$0 < \#\pi_{D,A}^{-1}(A) < \#D,$$

so $D^{(1)}, D^{(2)} \subsetneq D$. Thus, by Claims 2.20 and 2.21 and induction, there exist

$$\tau_{D^{(1)}} \in \text{PerfectTab}(D^{(1)}, \alpha^{(1)}) \quad \text{and} \quad \tau_{D^{(2)}} \in \text{PerfectTab}(D^{(2)}, \alpha^{(2)}).$$

Define $\tau_D = \tau_{D^{(1)}} \cup \tau_{D^{(2)}} \in \text{Tab}(D)$ by

$$\tau_D(i, j) = \begin{cases} \tau_{D^{(1)}}(i, j) & \text{if } (i, j) \in D^{(1)}, \\ \tau_{D^{(2)}}(i, j) & \text{if } (i, j) \in D^{(2)}. \end{cases}$$

Clearly τ_D is flagged and satisfies $\tau_D^{-1}(\{\circ\}) = \emptyset$. It has content α because $\alpha = \alpha^{(1)} + \alpha^{(2)}$, and it is column-injective because the images of $\tau_{D^{(1)}}$ and $\tau_{D^{(2)}}$ are disjoint. Therefore, $\tau_D \in \text{PerfectTab}(D, \alpha)$ and Case 2b is complete.

This completes the proof of Theorem 2.13. \square

2.2. Polytopal descriptions of perfect tableaux. Given $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$, define the indicator polytope

$$\mathcal{P}(D, \alpha) \subseteq \mathbb{R}^{n^2}$$

to be the polytope with points of the form $(\alpha_{ij})_{i,j \in [n]} = (\alpha_{11}, \dots, \alpha_{n1}, \dots, \alpha_{1n}, \dots, \alpha_{nn})$ governed by the inequalities (A)-(C) below.

(A) Column-Injectivity Conditions: For all $i, j \in [n]$,

$$0 \leq \alpha_{ij} \leq 1.$$

(B) Content Conditions: For all $i \in [n]$,

$$\sum_{j=1}^n \alpha_{ij} = \alpha_i.$$

(C) Flag Conditions: For all $s, j \in [n]$,

$$\sum_{i=1}^s \alpha_{ij} \geq \#\{(i, j) \in D : i \leq s\}.$$

Proposition 2.22. *Let $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$ with $\alpha_1 + \dots + \alpha_n = \#D$. If $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$, then for each $j \in [n]$, we have*

$$\sum_{i=1}^n \alpha_{ij} = \#\{(i, j) \in D : i \in [n]\}.$$

Proof. From the flag conditions (C) where $s = n$, we have that

$$\sum_{i=1}^n \alpha_{ij} \geq \#\{(i, j) \in D : i \in [n]\}.$$

If this inequality is strict for any j , then using the content conditions (B), we can write

$$\#D = \alpha_1 + \dots + \alpha_n = \sum_{i=1}^n \sum_{j=1}^n \alpha_{ij} = \sum_{j=1}^n \sum_{i=1}^n \alpha_{ij} > \sum_{j=1}^n \#\{(i, j) \in D : i \in [n]\} = \#D,$$

a contradiction. \square

Theorem 2.23. *Let $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$. Then $\text{PerfectTab}(D, \alpha) \neq \emptyset$ if and only if $\alpha_1 + \dots + \alpha_n = \#D$ and $\mathcal{P}(D, \alpha) \cap \mathbb{Z}^{n^2} \neq \emptyset$.*

Proof. (\Rightarrow) By Proposition 2.12, we have $\alpha_1 + \cdots + \alpha_n = \#D$. Let $\tau \in \text{PerfectTab}(D, \alpha)$. Then for each $i, j \in [n]$, set

$$\alpha_{ij} = \#\{r \in [n] : \tau(r, j) = i\} = \begin{cases} 1 & \text{if } \tau(r, j) = i \text{ for some } r \in [n], \\ 0 & \text{otherwise,} \end{cases}$$

where the second equality follows from the fact that τ is column-injective.

Claim 2.24. $(\alpha_{ij}) \in \mathcal{P}(D, \alpha) \cap \mathbb{Z}^{n^2}$.

Proof of Claim 2.24. Clearly $(\alpha_{ij}) \in \mathbb{Z}^{n^2}$ and the column-injectivity conditions (A) hold. Since τ has content α ,

$$\sum_{j=1}^n \alpha_{ij} = \sum_{j=1}^n \#\{r \in [n] : \tau(r, j) = i\} = \#\tau^{-1}(\{i\}) = \alpha_i$$

for each $i \in [n]$, so the content conditions (B) hold. Finally, for each $s, j \in [n]$, we have

$$\sum_{i=1}^s \alpha_{ij} = \#\{r \in [n] : \tau(r, j) \leq s\} \geq \#\{(r, j) \in D : r \leq s\}$$

since τ is flagged. Thus, the flag conditions (C) also hold. \square

(\Leftarrow) Let $(\alpha_{ij}) \in \mathcal{P}(D, \alpha) \cap \mathbb{Z}^{n^2}$. By the column-injectivity conditions (A), $\alpha_{ij} \in \{0, 1\}$. Thus, by Proposition 2.22, there exists for each $j \in [n]$ a bijection

$$\varphi_j : \{i \in [n] : (i, j) \in D\} \rightarrow \{i \in [n] : \alpha_{ij} = 1\}$$

that is order-preserving, i.e., φ_j satisfies $\varphi_j(i) < \varphi_j(i')$ whenever $i < i'$. Define $\tau \in \text{Tab}(D)$ by $\tau(i, j) = \varphi_j(i)$.

Claim 2.25. $\tau \in \text{PerfectTab}(D, \alpha)$.

Proof of Claim 2.25. By construction, $\tau^{-1}(\{\circ\}) = \emptyset$. Since φ_j is injective and order-preserving, τ is strictly increasing along columns, hence column-injective. For each $i \in [n]$, the content conditions (B) imply

$$\tau^{-1}(\{i\}) = \sum_{j=1}^n \#\varphi_j^{-1}(\{i\}) = \sum_{j=1}^n \alpha_{ij} = \alpha_i,$$

so τ has content α . Finally, the flag conditions (C) show that for each $s, j \in [n]$,

$$\#\{i \leq s : (i, j) \in D\} \leq \sum_{i=1}^s \alpha_{ij} = \#\{i \leq s : \alpha_{ij} = 1\},$$

so $\varphi_j(i) \leq i$ for each $(i, j) \in D$ since φ_j is order-preserving. Thus, $\tau(i, j) = \varphi_j(i) \leq i$ and τ is flagged. Hence, $\tau \in \text{PerfectTab}(D, \alpha)$. \square

This shows that $\text{PerfectTab}(D, \alpha) \neq \emptyset$ and completes the proof of the theorem. \square

Remark 2.26. The proof of Claim 2.25 shows that if $\text{PerfectTab}(D, \alpha) \neq \emptyset$, then we can find $\tau \in \text{PerfectTab}(D, \alpha)$ which is not only column-injective, but also strictly increasing along columns, so $\tau(i, j) < \tau(i', j)$ whenever $i < i'$. Thus $\text{PerfectTab}(D, \alpha) \neq \emptyset$ if and only if $\text{PerfectTab}(D, \alpha)_\downarrow \neq \emptyset$.

Theorem 2.23 formulates the problem of determining if $\text{PerfectTab}(D, \alpha) \neq \emptyset$ in terms of feasibility of an integer linear programming problem. In general, integral feasibility is NP-complete. We now show that in our case, feasibility of the problem is equivalent to feasibility of its LP-relaxation:

Theorem 2.27. *Let $D \subseteq [n]^2$ and $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}^n$ with $\alpha_1 + \dots + \alpha_n = \#D$. Then $\mathcal{P}(D, \alpha) \cap \mathbb{Z}^{n^2} \neq \emptyset$ if and only if $\mathcal{P}(D, \alpha) \neq \emptyset$.*

For reasons given in the Introduction, we provide two proofs of this fact.

Proof 1 of Theorem 2.27. We write the constraints (A)-(C) in the form $M\vec{x} \leq \vec{b}$ where M is a $(3n^2 + n) \times n^2$ block matrix and \vec{b} is a vector of length $3n^2 + n$ of the form

$$M = \begin{pmatrix} M_{A_1} \\ M_{A_2} \\ M_B \\ M_C \end{pmatrix} \text{ and } \vec{b} = (b_i)_{i=1}^{3n^2+n}.$$

Let \vec{b}_I denote the subvector of \vec{b} containing those b_i with $i \in I \subseteq [3n^2 + n]$. Also, we use the following coordinatization:

$$\vec{x} = (\alpha_{11}, \dots, \alpha_{n1}, \alpha_{12}, \dots, \alpha_{n2}, \dots, \alpha_{nn})^T.$$

- M_{A_1} is the $n^2 \times n^2$ block corresponding to the condition $0 \leq \alpha_{ij}$ from (A). Thus, $M_{A_1} = -I_{n^2}$ and $b_r = 0$ for $r \in [1, n^2]$.
- M_{A_2} is the $n^2 \times n^2$ block corresponding to $\alpha_{ij} \leq 1$ from (A). Hence, $M_{A_2} = I_{n^2}$ and $b_r = 1$ for $r \in [n^2 + 1, 2n^2]$.
- M_C is the $n^2 \times n^2$ matrix for (C). Thus,

$$M_C = \begin{pmatrix} M_{C_T} & 0 & \dots & 0 \\ 0 & M_{C_T} & \dots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \dots & M_{C_T} \end{pmatrix}$$

where $M_{C_T} = (c_{ij})_{1 \leq i, j \leq n}$ is lower triangular such that $c_{ij} = -1$ for $i \geq j$. Also,

$$b_{(2n^2+n)+n(j-1)+s} = -\#\{(i, j) \in D : i \leq s\}, \text{ for } s, j \in [n].$$

- M_B is the $n \times n^2$ block encoding (B). Take $M_B = (I_n \ I_n \ \dots \ I_n)$ and $\vec{b}_{[2n^2+1, 2n^2+n]} = (\alpha_i)_{i \in [n]}$. Clearly $M_B \vec{x} \leq (\alpha_i)_{i \in [n]}$ encodes the inequalities $\sum_{j=1}^n \alpha_{ij} \leq \alpha_i$. Now, (B) requires $\sum_{j=1}^n \alpha_{ij} = \alpha_i$. However, $\alpha_1 + \dots + \alpha_n = \#D$ ensures that

$$\begin{pmatrix} M_B \\ M_C \end{pmatrix} \vec{x} \leq \vec{b}_{[2n^2+1, 3n^2+n]} \text{ only if } M_B \vec{x} = (\alpha_i)_{i \in [n]}.$$

Summarizing, $M\vec{x} \leq \vec{b}$ indeed encodes (A)-(C).

Example 2.28. For $n = 2$ consider $\vec{x} = (\alpha_{11}, \alpha_{21}, \alpha_{12}, \alpha_{22})^T$ with $D = \{(1, 1), (1, 2), (2, 2)\} \subset [2] \times [2]$ and $\alpha = (2, 1)$.

We have

$$M_{A_1} \vec{x} = \begin{pmatrix} -1 & 0 & 0 & 0 \\ 0 & -1 & 0 & 0 \\ 0 & 0 & -1 & 0 \\ 0 & 0 & 0 & -1 \end{pmatrix} \begin{pmatrix} \alpha_{11} \\ \alpha_{21} \\ \alpha_{12} \\ \alpha_{22} \end{pmatrix} \leq \begin{pmatrix} 0 \\ 0 \\ 0 \\ 0 \end{pmatrix}$$

$$\begin{aligned}
M_{A_2}\vec{x} &= \begin{pmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} \alpha_{11} \\ \alpha_{21} \\ \alpha_{12} \\ \alpha_{22} \end{pmatrix} \leq \begin{pmatrix} 1 \\ 1 \\ 1 \\ 1 \end{pmatrix} \\
M_B\vec{x} &= \begin{pmatrix} 1 & 0 & 1 & 0 \\ 0 & 1 & 0 & 1 \end{pmatrix} \begin{pmatrix} \alpha_{11} \\ \alpha_{21} \\ \alpha_{12} \\ \alpha_{22} \end{pmatrix} \leq \begin{pmatrix} \alpha_1 \\ \alpha_2 \end{pmatrix} = \begin{pmatrix} 2 \\ 1 \end{pmatrix} \\
M_C\vec{x} &= \begin{pmatrix} -1 & 0 & 0 & 0 \\ -1 & -1 & 0 & 0 \\ 0 & 0 & -1 & 0 \\ 0 & 0 & -1 & -1 \end{pmatrix} \begin{pmatrix} \alpha_{11} \\ \alpha_{21} \\ \alpha_{12} \\ \alpha_{22} \end{pmatrix} \leq \begin{pmatrix} -\#\{(i, 1) \in D : i \leq 1\} \\ -\#\{(i, 1) \in D : i \leq 2\} \\ -\#\{(i, 2) \in D : i \leq 1\} \\ -\#\{(i, 2) \in D : i \leq 2\} \end{pmatrix} = \begin{pmatrix} -1 \\ -1 \\ -1 \\ -2 \end{pmatrix}
\end{aligned}$$

Theorem 2.29. M is a totally unimodular matrix; that is, every minor of M equals 0, 1, or -1 .

Proof. Suppose M' is a square submatrix of M with k rows from M_{A_1} or M_{A_2} . We show by induction on k that $\det(M') \in \{0, \pm 1\}$.

For the base case $k = 0$, consider M' an $\ell \times \ell$ submatrix of M with only rows from M_B and M_C . Let M'_B, M'_C be the corresponding blocks of M' , i.e. $M' = \begin{pmatrix} M'_B \\ M'_C \end{pmatrix}$ where M'_B , or M'_C , is the submatrix of M_B , or M_C respectively, using the rows and columns of M' . Since M_B has one 1 per column, M'_B has at most one 1 per column. By the form of M_C , it is straightforward to row reduce M'_C to obtain a $(0, -1)$ -matrix M''_C with at most one -1 in each column. Let $M'' = \begin{pmatrix} M'_B \\ M''_C \end{pmatrix}$, an $\ell \times \ell$ matrix. It is textbook (see [11, Theorem 13.3]) that if a $(0, \pm 1)$ -matrix N has at most one 1 and at most one -1 in each column, N is totally unimodular; hence $\det(M') = \pm \det(M'') \in \{0, -1, 1\}$ as desired. Thus the base case holds.

Now suppose M' is a square submatrix of M that contains $k \geq 1$ rows from M_{A_1} or M_{A_2} . Let R be such a row from M_{A_1} or M_{A_2} . If R contains all 0's, $\det(M') = 0$, and we are done. Otherwise R contains a single ± 1 . Hence the cofactor expansion for $\det(M')$ along R gives $\det(M') = \pm \det(M'')$ where M'' is a submatrix of M with $k - 1$ rows from M_{A_1} or M_{A_2} . So by induction, $\det(M') \in \{0, \pm 1\}$, as required. \square

Since M is totally unimodular then any vertices of $M\vec{x} \leq \vec{b}$ are integral [11, Theorem 13.2]. Thus, if $\mathcal{P}(D, \alpha) \neq \emptyset$ then its vertices are integral, i.e., $\mathcal{P}(D, \alpha) \cap \mathbb{Z}^{n^2} \neq \emptyset$. \square

Proof 2 of Theorem 2.27. Given a point $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$, we say a pair of sequences

$$(r_1, \dots, r_{m+1}; c_1, \dots, c_m) \in [n]^{m+1} \times [n]^m,$$

for some $m \in \mathbb{Z}_{>0}$, is *stable* at (α_{ij}) if the properties (i)-(iv) below hold. The purpose of each property will become clear later.

- (i) $r_{m+1} = r_1$.
- (ii) For all $k \in [m]$, $\alpha_{r_k c_k}, \alpha_{r_{k+1} c_k} \notin \mathbb{Z}$.
- (iii) For all $k \in [m]$, if $i > r_{k+1}$ and $\alpha_{i c_k} \notin \mathbb{Z}$, then $i = r_k$.

(iv) There exists $(r, c) \in [n]^2$ such that

$$\#\{k \in [m] : (r, c) = (r_k, c_k)\} \neq \#\{k \in [m] : (r, c) = (r_{k+1}, c_k)\}.$$

Claim 2.30. For any $(\alpha_{ij}) \in \mathcal{P}(D, \alpha) \setminus \mathbb{Z}^{n^2}$, there exists $(r_1, \dots, r_{m+1}; c_1, \dots, c_m)$ stable at (α_{ij}) .

Proof of Claim 2.30. Choose r_1, c_1 such that $\alpha_{r_1 c_1} \notin \mathbb{Z}$, and assume that we have fixed r_k, c_k such that $\alpha_{r_k c_k} \notin \mathbb{Z}$. By Proposition 2.22, we have

$$\sum_{i=1}^n \alpha_{i c_k} = \#\{(i, c_k) \in D : i \in [n]\} \in \mathbb{Z}.$$

Thus, as $\alpha_{r_k c_k} \notin \mathbb{Z}$, it makes sense to set

$$(13) \quad r_{k+1} = \max\{i \neq r_k : \alpha_{i c_k} \notin \mathbb{Z}\}.$$

If $r_{k+1} = r_\ell$ for some $\ell \in [k]$, then end the construction of these sequences. Otherwise, the content conditions (B) say that

$$\sum_{j=1}^n \alpha_{r_{k+1} j} = \alpha_{r_{k+1}} \in \mathbb{Z},$$

and since $\alpha_{r_{k+1} c_k} \notin \mathbb{Z}$, we can choose $c_{k+1} \neq c_k$ such that $\alpha_{r_{k+1} c_{k+1}} \notin \mathbb{Z}$, completing the recursive definition. By the pigeonhole principle, this process must halt, yielding sequences $r_1, \dots, r_\ell, \dots, r_{m+1}$ and $c_1, \dots, c_\ell, \dots, c_m$ with $r_{m+1} = r_\ell$.

By disregarding the first $\ell - 1$ terms of each sequence, we may assume $\ell = 1$ without loss of generality. Then we assert that $(r_1, \dots, r_{m+1}; c_1, \dots, c_m)$ is stable at (α_{ij}) . Indeed, (i) and (ii) are immediate from the construction, (iii) follows from (13), and (iv) holds because $(r, c) := (r_2, c_2)$ exists and satisfies

$$\#\{k \in [m] : (r, c) = (r_k, c_k)\} = 1 \quad \text{and} \quad \#\{k \in [m] : (r, c) = (r_{k+1}, c_k)\} = 0,$$

since $c_2 \neq c_1$ and $r_2 \neq r_k$ for all $k \neq 2$. □

We now fix a pair of sequences $(r_1, \dots, r_{m+1}; c_1, \dots, c_m)$. Given (α_{ij}) and $\delta > 0$, set

$$(14) \quad \alpha_{ij}^\delta = \alpha_{ij} + \delta[\#\{k \in [m] : (i, j) = (r_k, c_k)\} - \#\{k \in [m] : (i, j) = (r_{k+1}, c_k)\}].$$

Claim 2.31. If $(r_1, \dots, r_{m+1}; c_1, \dots, c_m)$ is stable at $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$, then $(\alpha_{ij}^\delta) \in \mathcal{P}(D, \alpha)$ for some $\delta > 0$.

Proof of Claim 2.31. First, note that the content conditions (B) are preserved regardless of our choice of δ . Indeed, for each $i \in [n]$,

$$\begin{aligned} \sum_{j=1}^n \alpha_{ij}^\delta &= \sum_{j=1}^n [\alpha_{ij} + \delta[\#\{k \in [m] : (i, j) = (r_k, c_k)\} - \#\{k \in [m] : (i, j) = (r_{k+1}, c_k)\}]] \\ &= \alpha_i + \delta[\#\{k \in [m] : i = r_k\} - \#\{k \in [m] : i = r_{k+1}\}], \end{aligned}$$

and the term in brackets vanishes by (i).

We next check the flag conditions (C). For each $s, j \in [n]$, we can write

$$\begin{aligned}
\sum_{i=1}^s \alpha_{ij}^\delta &= \sum_{i=1}^s [\alpha_{ij} + \delta[\#\{k \in [m] : (i, j) = (r_k, c_k)\} - \#\{k \in [m] : (i, j) = (r_{k+1}, c_k)\}]] \\
&= \sum_{i=1}^s \alpha_{ij} + \delta[\#\{k \in [m] : s \geq r_k \text{ and } j = c_k\} - \#\{k \in [m] : s \geq r_{k+1} \text{ and } j = c_k\}] \\
(15) \quad &\geq \sum_{i=1}^s \alpha_{ij} - \delta[\#\{k \in [m] : r_{k+1} \leq s < r_k \text{ and } j = c_k\}].
\end{aligned}$$

Thus, if $\#\{k \in [m] : r_{k+1} \leq s < r_k \text{ and } j = c_k\} = 0$, then the flag condition (C) for these s, j is preserved.

Otherwise, $r_{k+1} \leq s < r_k$ and $j = c_k$ for some $k \in [m]$, so (ii) and (iii) tell us that there is exactly one $i > s$ for which $\alpha_{ij} \notin \mathbb{Z}$, namely $i = r_k$. This, combined with Proposition 2.22, shows that

$$(16) \quad \sum_{i=1}^s \alpha_{ij} = \sum_{i=1}^n \alpha_{ij} - \sum_{i=s+1}^n \alpha_{ij} = \#\{(i, j) \in D : i \in [n]\} - \sum_{i=s+1}^n \alpha_{ij} \notin \mathbb{Z}.$$

By the nonintegrality from (16), the flag inequalities (C) for $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$ are strict:

$$(17) \quad \sum_{i=1}^s \alpha_{ij} > \#\{(i, j) \in D : i \leq s\}.$$

Hence, by taking δ sufficiently small and applying (15) and (17), we can ensure

$$\sum_{i=1}^s \alpha_{ij}^\delta \geq \sum_{i=1}^s \alpha_{ij} - \delta[\#\{k \in [m] : r_{k+1} \leq s < r_k \text{ and } j = c_k\}] \geq \#\{(i, j) \in D : i \leq s\}$$

for all $s, j \in [n]$, so the flag conditions (C) will be preserved. If $\alpha_{ij} \neq \alpha_{ij}^\delta$ then by (14) we must have $(i, j) = (r_k, c_k)$ or $(i, j) = (r_{k+1}, c_k)$ for some k , which by (ii) implies $0 < \alpha_{ij} < 1$. So we can require in addition that δ be small enough that $0 \leq \alpha_{ij}^\delta \leq 1$ for all $i, j \in [n]$. For such δ , the conditions (A)-(C) all hold, so $(\alpha_{ij}^\delta) \in \mathcal{P}(D, \alpha)$. \square

Finally, choose a point $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$ with the maximum number of integer coordinates. If $(\alpha_{ij}) \in \mathbb{Z}^{n^2}$, then we are done. Otherwise, there exists $(r_1, \dots, r_{m+1}; c_1, \dots, c_m)$ that is stable at (α_{ij}) by Claim 2.30. By (iv), there exists $(r, c) \in [n]^2$ such that $|\alpha_{rc}^\delta| \rightarrow \infty$ as $\delta \rightarrow \infty$, so α_{rc}^δ violates the column-injectivity conditions (A) for large δ . This, combined with Claim 2.31, shows that the set $S = \{\delta > 0 : (\alpha_{ij}^\delta) \in \mathcal{P}(D, \alpha)\}$ is nonempty and bounded above. Thus, we can define $\eta = \sup S$ and set $(\tilde{\alpha}_{ij}) = (\alpha_{ij}^\eta)$. Since $\mathcal{P}(D, \alpha)$ is closed and the map $\delta \mapsto (\alpha_{ij}^\delta)$ from S to $\mathcal{P}(D, \alpha)$ is continuous, this supremum is in fact a maximum, and $(\tilde{\alpha}_{ij}) \in \mathcal{P}(D, \alpha)$. By our choice of (α_{ij}) , we cannot have $\tilde{\alpha}_{r_k c_k} \in \mathbb{Z}$ or $\tilde{\alpha}_{r_{k+1} c_k} \in \mathbb{Z}$ for any $k \in [m]$, since then $(\tilde{\alpha}_{ij})$ has more integer coordinates than (α_{ij}) . Thus, $(r_1, \dots, r_{m+1}; c_1, \dots, c_m)$ is stable at $(\tilde{\alpha}_{ij})$, so by Claim 2.31, there exists $\delta > 0$ for which $(\tilde{\alpha}_{ij}^\delta) \in \mathcal{P}(D, \alpha)$. But then $(\alpha_{ij}^{\eta+\delta}) = (\tilde{\alpha}_{ij}^\delta) \in \mathcal{P}(D, \alpha)$, contradicting the maximality of η . \square

In summary, applying the results of this section to $D = D(w)$,

$$(18) \quad c_{\alpha, w} > 0 \stackrel{[2]}{\iff} \alpha \in \mathcal{S}_D \stackrel{\text{Thm 2.13}}{\iff} \text{PerfectTab}(D, \alpha) \neq \emptyset \stackrel{\text{Thm 2.23}}{\iff} \mathcal{P}(D, \alpha) \cap \mathbb{Z}^{n^2} \neq \emptyset \stackrel{\text{Thm 2.27}}{\iff} \mathcal{P}(D, \alpha) \neq \emptyset.$$

If $D \subseteq [n]^2$ has many identical columns, then many of the flag conditions (C) will look essentially the same. Thus, for efficiency of computation, we construct a “compressed” version of $\mathcal{P}(D, \alpha)$ that removes some of the repetitive inequalities.

A tuple $\mathcal{C} = (m, \{P_k\}_{k=1}^\ell, \{p_k\}_{k=1}^\ell, \{\lambda_k\}_{k=1}^\ell)$ is a *compression* of $D \subseteq [n]^2$ if:

- $m \leq n$ is a nonnegative integer such that $(r, p) \notin D$ whenever $r > m$ and $p \in [n]$,
- $P = P_1 \dot{\cup} \dots \dot{\cup} P_\ell \subseteq [n]$ such that if $p, p' \in P_k$ then

$$\{r \in [n] : (r, p) \in D\} = \{r \in [n] : (r, p') \in D\},$$

and moreover if D is nonempty in column p then $p \in P_k$ for some $k \in [\ell]$.

- $p_k \in P_k$ a representative for each $k \in [\ell]$, and
- $\lambda_k = \#P_k$ for each $k \in \ell$.

For $D \subseteq [n]^2$, a compression \mathcal{C} of D , and $\tilde{\alpha} = (\tilde{\alpha}_1, \dots, \tilde{\alpha}_m) \in \mathbb{Z}_{\geq 0}^m$ define

$$(19) \quad \mathcal{Q}(D, \mathcal{C}, \tilde{\alpha}) \subseteq \mathbb{R}^{m\ell}$$

to be the polytope with points of the form $(\tilde{\alpha}_{ik})_{i \in [m], k \in [\ell]}$ satisfying (A')-(C') below.

(A') Column-Injectivity Conditions: For all $i \in [m], k \in [\ell]$,

$$0 \leq \tilde{\alpha}_{ik} \leq 1.$$

(B') Content Conditions: For all $i \in [m]$,

$$\sum_{k=1}^{\ell} \lambda_k \tilde{\alpha}_{ik} = \alpha_i.$$

(C') Flag Conditions: For all $s \in [m], k \in [\ell]$,

$$\sum_{i=1}^s \tilde{\alpha}_{ik} \geq \#\{(i, p_k) \in D : i \leq s\}.$$

Remark 2.32. We can always take $m = \ell = n$ and $P_k = \{k\}$ for each $k \in [\ell]$, in which case $\mathcal{Q}(D, \mathcal{C}, \tilde{\alpha}) = \mathcal{P}(D, \alpha) \subseteq \mathbb{R}^{n^2}$.

Theorem 2.33. Let $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$ and $\tilde{\alpha} = (\tilde{\alpha}_1, \dots, \tilde{\alpha}_m) := (\alpha_1, \dots, \alpha_m)$. Then $\alpha_1 + \dots + \alpha_n = \#D$ and $\mathcal{P}(D, \alpha) \neq \emptyset$ if and only if $\alpha_1 + \dots + \alpha_m = \#D$, $\alpha_{m+1} = \dots = \alpha_n = 0$, and $\mathcal{Q}(D, \mathcal{C}, \tilde{\alpha}) \neq \emptyset$.

Proof. (\Rightarrow) Let $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$. Then by the content and flag conditions (B) and (C),

$$\begin{aligned} \#D &= \alpha_1 + \dots + \alpha_n \geq \alpha_1 + \dots + \alpha_m = \sum_{i=1}^m \sum_{j=1}^n \alpha_{ij} \\ &= \sum_{j=1}^n \sum_{i=1}^m \alpha_{ij} \geq \sum_{j=1}^n \#\{(i, j) \in D : i \leq m\} = \#D. \end{aligned}$$

Thus, $\alpha_1 + \dots + \alpha_m = \#D$ and $\alpha_{m+1} = \dots = \alpha_n = 0$. Now, for each $i \in [m]$ and $k \in [\ell]$, set

$$\tilde{\alpha}_{ik} = \frac{1}{\lambda_k} \sum_{j \in P_k} \alpha_{ij}.$$

We claim that $(\tilde{\alpha}_{ik}) \in \mathcal{Q}(D, \mathcal{C}, \alpha)$. First, for each $i \in [m]$ and $k \in [\ell]$, we have

$$0 \leq \tilde{\alpha}_{ik} = \frac{1}{\lambda_k} \sum_{j \in P_k} \alpha_{ij} \leq \frac{1}{\lambda_k} \sum_{j \in P_k} 1 = 1,$$

so the column-injectivity conditions (A') are satisfied. Next, for each $i \in [m]$, (B) implies

$$\sum_{k=1}^{\ell} \lambda_k \tilde{\alpha}_{ik} = \sum_{k=1}^{\ell} \sum_{j \in P_k} \alpha_{ij} = \sum_{j=1}^n \alpha_{ij} = \alpha_i,$$

so the content conditions (B') are satisfied. Finally, for each $s \in [m]$ and $k \in [\ell]$, (C) implies

$$\sum_{i=1}^s \tilde{\alpha}_{ik} = \frac{1}{\lambda_k} \sum_{j \in P_k} \sum_{i=1}^s \alpha_{ij} \geq \frac{1}{\lambda_k} \sum_{j \in P_k} \#\{(i, j) \in D : i \leq s\} = \#\{(i, p_k) \in D : i \leq s\},$$

so the flag conditions (C') are satisfied.

(\Leftarrow) Clearly $\alpha_1 + \dots + \alpha_n = \#D$. Let $(\tilde{\alpha}_{ik}) \in \mathcal{Q}(D, \mathcal{C}, \tilde{\alpha})$. For each $i, j \in [n]$, set

$$\alpha_{ij} = \begin{cases} 0 & \text{if } i > m, \\ \tilde{\alpha}_{ik} & \text{if } i \leq m \text{ and } j \in P_k. \end{cases}$$

We claim that $(\alpha_{ij}) \in \mathcal{P}(D, \alpha)$. The column-injectivity conditions (A) are clear. If $i > m$,

$$\sum_{j=1}^n \alpha_{ij} = 0 = \alpha_i.$$

Otherwise $i \leq m$, and (B') implies

$$\sum_{j=1}^n \alpha_{ij} = \sum_{k=1}^{\ell} \sum_{j \in P_k} \tilde{\alpha}_{ik} = \sum_{k=1}^{\ell} \lambda_k \tilde{\alpha}_{ik} = \alpha_i.$$

Thus, the content conditions (B) hold. Finally, if $s \in [n]$ and $j \in P_k$, then (C') implies

$$\sum_{i=1}^s \alpha_{ij} = \sum_{i=1}^{\min\{s, m\}} \tilde{\alpha}_{ik} \geq \#\{(i, p_k) \in D : i \leq \min\{s, m\}\} = \#\{(i, j) \in D : i \leq s\}.$$

Hence, the flag conditions (C) hold as well. \square

2.3. Deciding membership in the Schubitope. We use the above results of this section to give a polynomial time algorithm to check if a lattice point is in the Schubitope.

Let $D \subseteq [n]^2$, and fix a compression $\mathcal{C} = (m, \{P_k\}_{k=1}^{\ell}, \{p_k\}_{k=1}^{\ell}, \{\lambda_k\}_{k=1}^{\ell})$ of D (as in Section 2.2).

Theorem 2.34. *Let $\alpha = (\alpha_1, \dots, \alpha_n) \in \mathbb{Z}_{\geq 0}^n$. Then $\alpha \in \mathcal{S}_D$ if and only if $\alpha_1 + \dots + \alpha_m = \#D$, $\alpha_{m+1} = \dots = \alpha_n = 0$, and $\mathcal{Q}(D, \mathcal{C}, \tilde{\alpha}) \neq \emptyset$, where $\tilde{\alpha} = (\tilde{\alpha}_1, \dots, \tilde{\alpha}_m) := (\alpha_1, \dots, \alpha_m)$.*

Proof. This follows from Theorems 2.13, 2.23, 2.27, and 2.33. \square

For each $k \in [\ell]$, let $R_k(\mathcal{C}) = \{r \in [n] : (r, p_k) \in D\} \subseteq [m]$.

Theorem 2.35. *Given as input $\{R_k(\mathcal{C})\}_{k=1}^{\ell}, \{\lambda_k\}_{k=1}^{\ell}$, and $\tilde{\alpha} = (\tilde{\alpha}_1, \dots, \tilde{\alpha}_m) \in \mathbb{Z}_{\geq 0}^m$ satisfying $\tilde{\alpha}_1 + \dots + \tilde{\alpha}_m = \#D$, one can decide if $\alpha := (\tilde{\alpha}_1, \dots, \tilde{\alpha}_m, 0, \dots, 0) \in \mathbb{Z}_{\geq 0}^n$ lies in \mathcal{S}_D in polynomial time in m and ℓ .*

Remark 2.36. In view of Theorem 2.34, this input is most natural, because the conditions $\alpha_1 + \dots + \alpha_m = \#D$ and $\alpha_{m+1} = \dots = \alpha_n = 0$ are clearly necessary, and it contains the minimum amount of information we need to compute $\mathcal{Q}(D, \mathcal{C}, \tilde{\alpha})$.

Remark 2.37. As in Remark 2.32, we can take $m = \ell = n$ and $P_k = \{k\}$ for each $k \in [\ell]$, so we can check if α is in \mathcal{S}_D in polynomial time in n regardless of the structure of D .

Proof of Theorem 2.35. Since $R_k(\mathcal{C})$ takes m bits to encode for each $k \in [\ell]$, and $\mathcal{Q}(D, \mathcal{C}, \tilde{\alpha}) \subseteq \mathbb{R}^{m\ell}$ is governed by $O(m\ell)$ constraints, $\mathcal{Q}(D, \mathcal{C}, \tilde{\alpha})$ can be constructed in polynomial time in m and ℓ . By Theorem 2.34, we are done using L. Khachiyan's ellipsoid method [12]. \square

3. COMPUTING ROTHE DIAGRAMS

We will repeatedly use the following to establish the complexity of computing preliminary data of $D(w)$ given $\text{code}(w)$.

Proposition 3.1. *There exists an $O(L^2)$ -time algorithm to compute $(w(1), \dots, w(L))$ from the input $\text{code}(w) = (c_1, \dots, c_L)$.*

Proof. Clearly $w(1) = c_1 + 1$. After determining $w(1), \dots, w(i-1)$, we determine (in $O(L)$ -time) $\pi := \pi^{(i)} \in S_{i-1}$ such that $(w(\pi(1)) < w(\pi(2)) < \dots < w(\pi(i-1)))$. Next, set

$$B := (w(\pi(1)), w(\pi(2)) - w(\pi(1)), w(\pi(3)) - w(\pi(2)), \dots, w(\pi(i-1)) - w(\pi(i-2))).$$

Let

$$V_t := \sum_{j=1}^t (B_j - 1), \quad \text{for } 0 \leq t \leq i-1.$$

Set $w(i) := c_i + T + 1$ where $T := \max_{t \in [0, i-1]} \{t : c_i \geq V_t\}$. By construction, $w(1), \dots, w(i)$ is a partial permutation with code $(c_1, \dots, c_{i-1}, c_i)$. Each stage $1 \leq i \leq L$ takes $O(i)$ -time. \square

The *essential set* of w consists of the maximally southeast boxes of each connected component of $D(w)$, i.e.,

$$(20) \quad \text{Ess}(w) = \{(i, j) \in D(w) : (i+1, j), (i, j+1) \notin D(w)\}.$$

If it exists, we call the connected component of $D(w)$ involving $(1, 1)$ the *dominant component* and denote it by $\text{Dom}(w)$. For instance, in Example 1.1, $\text{Dom}(w)$ has shape $(4, 2, 2, 2)$. Further, if it exists, the *accessible box* \mathbf{z}_w is the southmost then eastmost box in $\text{Ess}(w) \setminus \text{Dom}(w)$. In Example 1.1,

$$\text{Ess}(w) = \{(1, 4), (3, 4), (3, 7), (4, 2)\} \text{ and } \mathbf{z}_w = (3, 7).$$

(Although $(4, 2)$ is the southmost box of $\text{Ess}(w)$, it is in $\text{Dom}(w)$, and hence not the accessible.)

We will need the following in Section 5:

Proposition 3.2. *Given $\text{code}(w)$, there exists an $O(L^2)$ -time algorithm to compute $\mathbf{z}_w = (r, c)$ or determine it does not exist.*

Proof. Use Proposition 3.1 to find $(w(1), \dots, w(L))$ in $O(L^2)$ -time. Next, compute

$$w_{NW}(i) := \{w(j) : w(j) \leq w(i), j \leq i\}.$$

Then take

$$Y(i) := \{q - 1 : q \in w_{NW}(i)\} \setminus w_{NW}(i), \quad \text{for } i \in [L].$$

Compute $k_i := \max Y(i)$ for $i \in [L]$ in $O(L^2)$ -time (if $k_i \geq 1$, then k_i is the column index of the eastmost box of $D(w)$ in row i). In $O(L^2)$ -time, calculate

$$I := \{i \in [2, \dots, L] : k_i > \min_{j < i} w(j)\}.$$

Let $Y := \{(i, k_i) : i \in I\}$. Hence, $Y \cap \text{Dom}(w) = \emptyset$. Thus, if $Y = \emptyset$, \mathbf{z}_w does not exist. Otherwise, $\mathbf{z}_w \in Y$. Thus, in $O(L)$ -time, determine $r := \max\{i : (i, k_i) \in Y\}$. Output $\mathbf{z}_w = (r, k_r)$. \square

The *pivots* of \mathbf{z}_w denoted $\text{Piv}(\mathbf{z}_w)$ are the \bullet 's of $D(w)$ that are maximally southeast, among those northwest of \mathbf{z}_w . In Example 1.1, $\text{Piv}((3, 7)) = \{(2, 3), (1, 5)\}$.

4. PROOFS OF THEOREMS 1.2 AND 1.3

4.1. Proof of Theorem 1.2. By (1) combined with Theorem 2.35, it remains to establish the complexity of computing a compression of $D(w)$. For this, we need the following lemmas and propositions. Fix $w \in S_\infty$ with $\text{code}(w) = (c_1, \dots, c_L)$. Let $\sigma \in S_L$ be such that $\{w(\sigma(1)) < w(\sigma(2)) < \dots < w(\sigma(L))\}$. For convenience, set $w(\sigma(0)) := 0$.

Lemma 4.1. *For $1 \leq h \leq L$, and for all*

$$j_1, j_2 \in \{w(\sigma(h-1)) + 1, w(\sigma(h-1)) + 2, \dots, w(\sigma(h)) - 1\},$$

we have $(i, j_1) \in D(w)$ if and only if $(i, j_2) \in D(w)$.

Proof. For each k , let $u_1^{(k)} < \dots < u_k^{(k)}$ be $w(1), w(2), \dots, w(k)$ sorted in increasing order. Set $u_0^{(k)} := 0$. The lemma follows from the inductive claim that in the first k rows of $D(w)$, the columns $u_{h-1}^{(k)} + 1, u_{h-1}^{(k)} + 2, \dots, u_h^{(k)} - 1$ are the same. The base case $k = 1$ is clear. The inductive step is straightforward by considering how, in row $k + 1$ of $D(w)$, the \bullet and its ray emanating east affects the columns. \square

Define a collection of intervals in $[n]$ by

$$P'_{2k-1} := [w(\sigma(k-1)) + 1, w(\sigma(k)) - 1] \text{ and } P'_{2k} := \{w(\sigma(k))\}, \text{ for } 1 \leq k \leq L.$$

Let $1 \leq h_1 < h_2 < \dots < h_\ell \leq 2L$ be indices of the intervals P'_h that are nonempty. Set $P_i := P'_{h_i}$.

Lemma 4.2. *If $j_1, j_2 \in P_k$ for some k , then $(i, j_1) \in D(w) \iff (i, j_2) \in D(w)$.*

Proof. This follows by the definition of $\{P_k\}_{k=1}^\ell$ together with Lemma 4.1. \square

Let $p_k := \min\{p \in P_k\}$ for each $k \in [\ell]$.

Proposition 4.3. *There exists an $O(L^2)$ -time algorithm to compute $\{P_k\}_{k=1}^\ell$, $\{p_k\}_{k=1}^\ell$, and $\{\#P_k\}_{k=1}^\ell$ from the input $\text{code}(w) = (c_1, \dots, c_L)$.*

Proof. Proposition 3.1 computes $(w(1), \dots, w(L))$ in $O(L^2)$ -time. It takes $O(L \log(L))$ -time to sort $(w(1), \dots, w(L))$, i.e., to compute $\sigma \in S_L$. Computing the endpoints, and thus cardinalities, of the P'_k takes $O(L)$ -time as there are at most $2L$ of them. Then we reindex $\{\#P'_k\}_{k=1}^{2L}$ to obtain $\{\#P_k\}_{k=1}^\ell$ in $O(L)$ -time. \square

For each $k \in [\ell]$, let

$$R_k := \{r \in [L] : (r, p_k) \in D(w)\}.$$

Proposition 4.4. *Computing $\{R_k\}_{k=1}^\ell$ from $\text{code}(w)$ takes $O(L^2)$ -time.*

Proof. By $D(w)$'s definition, $r \in R_k$ if and only if $w(r) > p_k$ and $p_k \notin \{w(i) : i < r\}$. Propositions 4.3 and 3.1 give $\{P_k\}_{k=1}^\ell$, $\{p_k\}_{k=1}^\ell$ and $\{w(1), \dots, w(L)\}$ in $O(L^2)$ -time. \square

Conclusion of proof of Theorem 1.2: Proposition 4.3 computes $\{P_k\}_{k=1}^\ell$, $\{p_k\}_{k=1}^\ell$, and $\{\#P_k\}_{k=1}^\ell$ in $O(L^2)$ -time. Proposition 4.4 finds $\{R_k\}_{k=1}^\ell$ in $O(L^2)$ -time. One checks, using Lemma 4.2, that $\mathcal{C} = (L, \{P_k\}_{k=1}^\ell, \{p_k\}_{k=1}^\ell, \{\#P_k\}_{k=1}^\ell)$ is a compression of $D(w)$. Hence we may apply Theorem 2.35 by taking $D := D(w)$, $R_k(\mathcal{C}) := R_k$, $\lambda_k := \#P_k$ for $k \in [\ell]$ and $m := L$. Thus the result follows by (1). \square

4.2. Proof of Theorem 1.3; an application. Remark 2.26 combined with (18) proves the theorem. \square

Let $n_{132}(w)$ be the number of 132-patterns in $w \in S_n$, that is,

$$n_{132}(w) = \#\{(i, j, k) : 1 \leq i < j < k \leq n, w(i) < w(k) < w(j)\}.$$

Corollary 4.5. *There are at least $n_{132}(w) + 1$ distinct vectors α such that $c_{\alpha, w} > 0$.*

Proof. Suppose $i < j < k$ index a 132 pattern in w . There is a box b of $D(w)$ in row j and column $w(k)$. There are $N := n_{132}(w)$ many such boxes, b_1, \dots, b_N (all distinct), listed in English language reading order. Let M_i be boxes in the same column and connected component as b_i that are weakly north of b_i and strictly south of any b_j , where $j < i$. Iteratively define fillings $F_0, F_1, F_2, \dots, F_N$ of $D(w)$:

(F_0) Fill each box c of $D(w)$ with the row number of c .

(F_i) For $1 \leq i \leq N$, F_i is the same as F_{i-1} except that $F_i(c) := F_{i-1}(c) - 1$ if $c \in M_i$.

Clearly, $F_0 \in \text{PerfectTab}_\downarrow(D(w)) := \bigcup_\alpha \text{PerfectTab}_\downarrow(D(w), \alpha)$. Inductively assume $F_{i-1} \in \text{PerfectTab}_\downarrow(D(w))$. Since labels only decrease, F_i satisfies the row bound condition. Next we check that each column is strictly increasing. Let m_i be the northmost box of M_i . If m_i is adjacent and directly below some b_j (for a $j < i$) then

$$F_i(b_j) = F_0(b_j) - 1 < F_0(m_i) - 1 = F_i(m_i),$$

as needed. Otherwise suppose m_i is adjacent and south of a non-diagram position. Let d_i (if it exists) be the first diagram box directly north of m_i . Then $F_0(d_i) < F_0(m_i) - 1$. Hence

$$F_i(d_i) \leq F_0(d_i) < F_0(m_i) - 1 = F_i(m_i),$$

verifying column increasingness here as well. That F_i is column increasing elsewhere is clear since F_{i-1} is column increasing (by induction) and only labels of M_i are changed.

It remains to check that every label of F_i is in $\mathbb{Z}_{>0}$. Since each box of $D(w)$ is decremented at most once, the only concern is there is a box x in the first row that appears in some M_i , since then $F_0(x) = 1$ and $F_i(x) = 0$. However, in this case b_i must be in $\text{Dom}(w)$, which implies that the "1" in the 132-pattern associated to b_i could not exist, a contradiction. Thus $F_i \in \text{PerfectTab}_\downarrow(D(w))$, completing the induction.

Finally, under Theorem 1.3, each F_i corresponds to a distinct exponent vector since the sum of the labels is strictly decreasing at each step $F_{i-1} \mapsto F_i$. \square

From Corollary 4.5, this result of A. Weigandt [15] is immediate:

Corollary 4.6 (A. Weigandt's 132-bound). $\mathfrak{S}_w(1, 1, 1, \dots, 1) \geq n_{132}(w) + 1$.

As shown in [15], Corollary 4.6 in turn implies $\mathfrak{S}_w(1, 1, \dots, 1) \geq 3$ if $n_{132}(w) \geq 2$, a recent conjecture of R. P. Stanley [13].

5. COUNTING $c_{\alpha, w}$ IS IN #P

5.1. Vexillary permutations. A permutation $w \in S_n$ is *vexillary* if there does not exist a 2143 *pattern*, i.e., indices $i < j < k < l$ such that w has the pattern $w(j) < w(i) < w(l) < w(k)$. For example, $w = \underline{53841267}$ is not vexillary; we underlined the positions of a 2143 pattern. *Fulton's criterion* states that w is vexillary if and only if there do not exist $(a, b), (c, d) \in \text{Ess}(w)$ such that $a < c$ and $b < d$. In Example 1.1, w is not vexillary due to $(1, 4)$ and $(3, 7)$. Our main reference for this subsection is [8, Chapter 2].

We will also use this characterization of vexillary permutations:

Theorem 5.1. [6] *Given $\text{code}(w) = (c_1, \dots, c_L) \in \mathbb{Z}_{\geq 0}^n$, w vexillary if and only if*

- (i) *if i is such that $c_i > c_{i+1}$, then $c_i > c_j$ for any $j > i$, and*
- (ii) *if i, h are such that $c_i \geq c_h$, then $\#\{j : i < j < h, c_j < c_h\} \leq c_i - c_h$.*

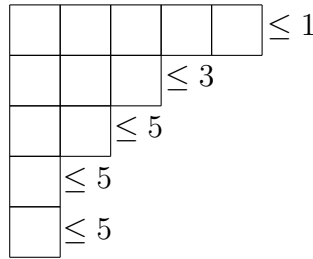
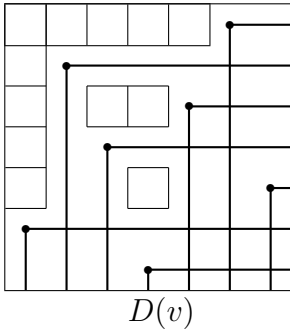
The *shape* of a vexillary permutation v is the partition $\lambda(v)$ formed by sorting $\text{code}(v) = (c_1, c_2, \dots)$ into decreasing order. Now, if $c_i \neq 0$, let e_i be the greatest integer $j \geq i$ such that $c_j \geq c_i$. The *flag*

$$\phi(v) = (\phi_1 \leq \phi_2 \leq \dots \leq \phi_m)$$

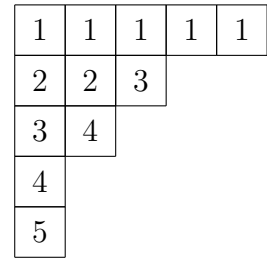
for v is the sequence of e_i 's sorted into increasing order; see, e.g., [8, Definition 2.2.9].

Example 5.2. Consider $\text{code}(v) = (5, 1, 3, 1, 2)$ for the vexillary $v = 6253714$. Here

$$e = (1, 5, 3, 5, 5), \phi(v) = (1, 3, 5, 5, 5) \text{ and } \lambda(v) = (5, 3, 2, 1, 1).$$



$\lambda(v)$ flagged by $\phi(v)$



$T \in \text{SSYT}(\lambda(v), \phi(v))$

For a partition $\lambda = (\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_m \geq 0)$ and a flag $\phi = (\phi_1 \leq \phi_2 \leq \dots \leq \phi_m)$ of positive integers, define the *flagged Schur function*

$$S_\lambda(\phi) = \det |h_{\lambda_i - i + j}(\phi_i)|_{i,j=1, \dots, m},$$

where

$$h_k(n) = \sum_{1 \leq i_1 \leq \dots \leq i_k \leq n} x_{i_1} \cdots x_{i_k}$$

is the complete homogeneous symmetric polynomial of degree k . Furthermore,

$$(21) \quad \mathfrak{S}_v = S_{\lambda(v)}(\phi(v)), \text{ for } v \text{ vexillary.}$$

A semistandard Young tableau of shape λ is *flagged* by ϕ if its entries in row i are $\leq \phi_i$; see Example 5.2. Denote the set of such tableaux by $\text{SSYT}(\lambda, \phi)$. Then

$$(22) \quad S_\lambda(\phi) = \sum_{T \in \text{SSYT}(\lambda, \phi)} x^{\text{content}(T)}.$$

where $\text{content}(T) = (\mu_1, \dots, \mu_\ell(\lambda))$ such that μ_i is the number of i 's in T .

5.2. Graphical transition. The transition recurrence for \mathfrak{S}_w was found by A. Lascoux and M.-P. Schützenberger [6]. This is transition for the case discussed in [5]:

Theorem 5.3 ([6], cf. [5]). *Let $\mathbf{z}_w = (r, c)$ and $w' = w \cdot (r \ k)$ where $k = w^{-1}(c)$. Then*

$$(23) \quad \mathfrak{S}_w = x_r \mathfrak{S}_{w'} + \sum_{w'' = w' \cdot (i \ k)} \mathfrak{S}_{w''},$$

where the summation is over $\{i : (i, w(i)) \in \text{Piv}(\mathbf{z}_w)\}$.

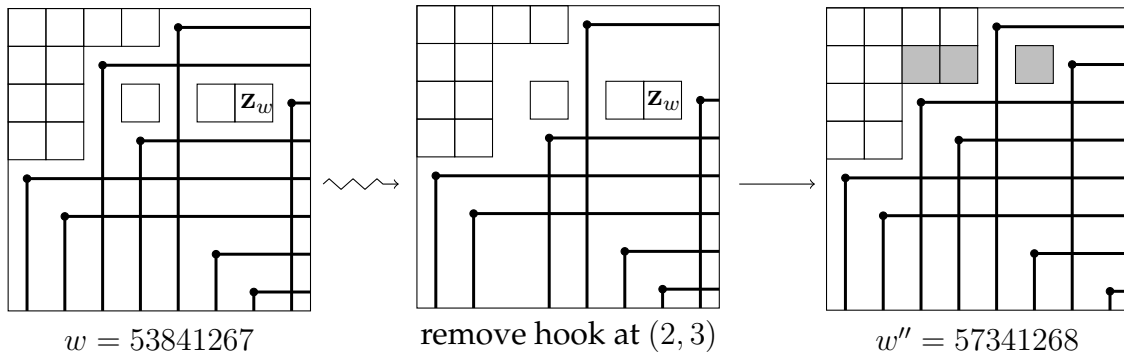
We will use the *graphical* transition tree $\mathcal{T}(w)$ of [5]. This reformulates (23) in terms of Rothe diagrams and certain moves on these diagrams. By definition, $D(w)$ (equivalently w) will label the root of $\mathcal{T}(w)$. If w is vexillary, stop. Otherwise, there exists an accessible box $\mathbf{z}_w = (r, c) \in D(w)$ (if not, $D(w) = \text{Dom}(w)$, contradicting w is not vexillary).

The children of $D(w)$ are Rothe diagrams resulting from two types of moves:

- (T.1) *Deletion moves:* remove \mathbf{z}_w from $D(w)$. The resulting diagram is $D(w')$. Add an edge $D(w) \xrightarrow{x_r} D(w')$.
- (T.2) *March moves:* There is a move for each $\mathbf{x}^{(i)} = (i, w(i)) \in \text{Piv}(\mathbf{z}_w)$. Let \mathcal{R} be the rectangle with corners \mathbf{z}_w and $\mathbf{x}^{(i)}$. Remove $\mathbf{x}^{(i)}$ and its rays from $G(w)$ to form $G^{(i)}(w)$. Order the boxes $\{b_i\}_{i=1}^r$ in \mathcal{R} in English reading order. Move b_1 strictly north and strictly west to the closest position not occupied by other boxes of $D(w)$ or rays from $G^{(i)}(w)$. Repeat with b_2, b_3, \dots where b_j may move into a square left unoccupied by earlier moves. The resulting diagram will be $D(w'')$ where $w'' = w' \cdot (i \ k)$. Add an edge $D(w) \xrightarrow{i} D(w'')$.

Repeat for each child $D(u)$. Stop when u vexillary; these permutations are the leaves $\mathcal{L}(w)$ of $\mathcal{T}(w)$. (Multiple leaves may be labelled by the same permutation.)

Example 5.4. Let $w = 53841267$. We compute the march move 2 for the pivot $(2, 3)$:



The moved boxes during $D(w) \mapsto D(w'')$ are shaded gray.

Example 5.5. Let $w = 53861247$. Using $\mathcal{T}(w)$ from Figure 1, we compute

$$\begin{aligned} \mathfrak{S}_w = & x_4 \cdot \mathfrak{S}_{73541268} + x_4 \cdot \mathfrak{S}_{57341268} + x_3^2 x_4 \cdot \mathfrak{S}_{53641278} + x_3 x_4 \cdot \mathfrak{S}_{63541278} + x_3 x_4 \cdot \mathfrak{S}_{56341278} \\ & + \mathfrak{S}_{74531268} + \mathfrak{S}_{57431268} + x_3^2 \cdot \mathfrak{S}_{54631278} + x_3 \cdot \mathfrak{S}_{64531278} + x_3 \cdot \mathfrak{S}_{56431278}. \end{aligned}$$

For instance, $c_{(4,2,5,3),w} := [x_1^4 x_2^2 x_3^5 x_4^3] \mathfrak{S}_w = 1$ is witnessed by

- the path $w \xrightarrow{x_4} \bullet \xrightarrow{x_3} \bullet \xrightarrow{x_3} u = 53641278$, and
- the semistandard tableau

$$T = \begin{array}{|c|c|c|c|} \hline 1 & 1 & 1 & 1 \\ \hline 2 & 2 & 3 & \\ \hline 3 & 3 & & \\ \hline 4 & 4 & & \\ \hline \end{array} \quad \text{of shape } \lambda(u), \text{ flagged by } \phi(u) = (1, 3, 4, 4).$$

Proposition 5.9 below formalizes a rule for $c_{\alpha,w}$ in terms of such pairs.

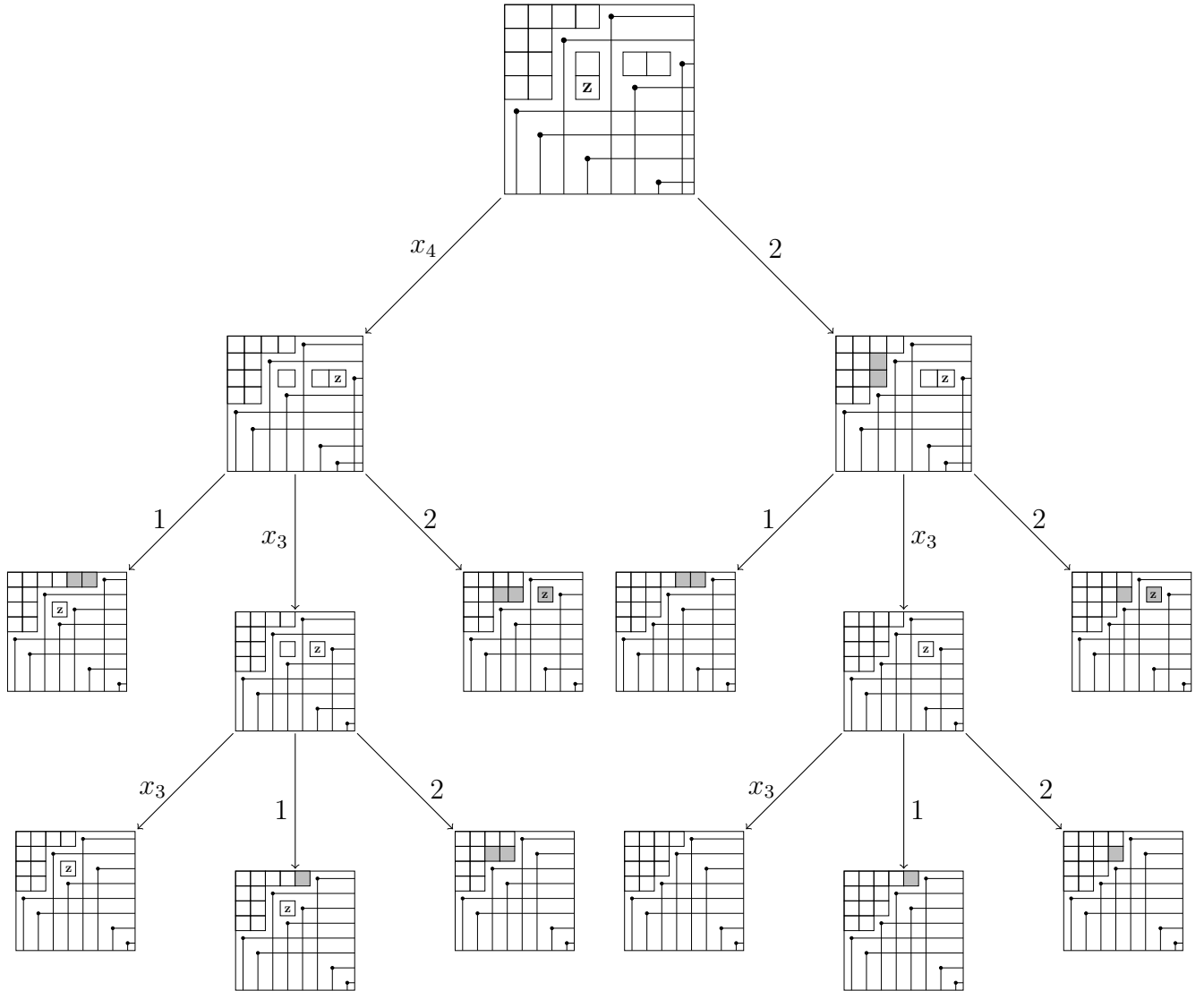


FIGURE 1. $\mathcal{T}(w)$ for $w = 53861247$ where the accessible boxes are marked with z and those boxes of the parent which moved are shaded gray.

5.3. **Proof of #P-ness.** The technical core of our proof of Theorem 1.5 is to show:

Theorem 5.6. *The problem of computing $c_{\alpha,w}$, given input α and $\text{code}(w)$, is in #P.*

Define X to be the set consisting of pairs (S, R) where:

- (X.1) $S = (s_1, \dots, s_h)$, $s_t \in [L] \cup \{(x_k, m_t) : k \in [L], m_t \in \mathbb{Z}_{>0}\}$ such that if $s_t = (x_k, m_t)$ then $s_{t+1} \neq (x_k, m_{t+1})$ for $t < h$, and
- (X.2) $R = (r_{ij})_{1 \leq i, j \leq L}$, where $r_{ij} \in \mathbb{Z}_{\geq 0}$.

Fix $w \in S_\infty$ and a vexillary permutation $v \in S_\infty$. A (w, v) -transition string is a sequence $S = (s_1, \dots, s_h)$ satisfying (X.1) such that if we interpret i as $\bullet \xrightarrow{i} \bullet$ and (x_k, m_t) as $\bullet \xrightarrow{x_k} \bullet \dots \bullet \xrightarrow{x_k} \bullet$ (m_t -times) then S describes a path from w to (a leaf labelled by) v in $\mathcal{T}(w)$. Let $\text{Trans}(w, v)$ be the set of such sequences.

The *deletion weight* of $S \in \text{Trans}(w, v)$ is

$$\text{delwt}(S) = \sum m_t \cdot \vec{e}_r,$$

where the summation is over $1 \leq t \leq h$ such that $s_t = (x_r, m_t) \in S$ for some $r \in [L]$ (depending on t). Here $\vec{e}_r \in \mathbb{Z}_{\geq 0}^L$ is the r -th standard basis vector and L is the length of $\text{code}(w) = (c_1, c_2, \dots, c_L)$.

Example 5.7. In Figure 1 we read the $(w = 53861247, v = 54631278)$ -transition string $S = (2, (x_3, 2))$ as the path $w \xrightarrow{2} \bullet \xrightarrow{x_3} \bullet \xrightarrow{x_3} v$. Here, $\text{delwt}(S) = (0, 0, 2, 0)$.

Suppose T is a tableau of shape $\lambda = (\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_L \geq 0)$, with entries in $[L]$ and weakly increasing along rows. Define

$$R(T) = (r_{ij})_{1 \leq i, j \leq L}$$

to be the $L \times L$ matrix where r_{ij} is the number of j 's in row i of T . $R(T)$ encodes T . As pointed out in (a preprint version of) [10], T might have exponentially many (in L) boxes, whereas $R(T)$ is a $O(L^2)$ description of T .

Example 5.8. If $\lambda = (4, 3, 1, 0, 0)$ and

$$T = \begin{array}{|c|c|c|c|} \hline 1 & 1 & 2 & 3 \\ \hline 2 & 4 & 5 & \\ \hline 4 & & & \\ \hline \end{array} \longleftrightarrow R(T) = \begin{pmatrix} 2 & 1 & 1 & 0 & 0 \\ 0 & 1 & 0 & 1 & 1 \\ 0 & 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 \end{pmatrix}$$

Let $X_{\alpha,w} = \{(S, R(T))\} \subseteq X$ such that the following hold:

- (X.1') $S \in \text{Trans}(w, v)$,
- (X.2') $T \in \text{SSYT}(\lambda(v), \phi(v))$, and
- (X.3') $\text{delwt}(S) + \text{content}(T) = \alpha$.

Proposition 5.9. $c_{\alpha,w} = \#X_{\alpha,w}$.

Proof. Iterating (23),

$$\mathfrak{S}_w = \sum_{\text{vexillary } v \in S_\infty} \sum_{S \in \text{Trans}(w, v)} x^{\text{delwt}(S)} \mathfrak{S}_v.$$

Hence

$$(24) \quad c_{\alpha, w} = \sum_{\text{vexillary } v \in S_\infty} \sum_{S \in \text{Trans}(w, v)} [x^\alpha] x^{\text{delwt}(S)} \mathfrak{S}_v.$$

The result then follows from by (21), (22), and (24) combined. \square

Proposition 5.10 (cf. [6]). *Let $\text{code}(w) = (c_1, \dots, c_L)$. Suppose $D(w')$ is obtained from $D(w)$ using move (T.1) and $D(w'')$ is obtained from $D(w)$ with move (T.2) for a pivot in row i . There is an $O(L^2)$ -time algorithm to compute*

- (I) $\text{code}(w') = (c_1, \dots, c_{r-1}, c_r - 1, c_{r+1}, \dots, c_L)$ and
- (II) $\text{code}(w'') = (c_1, \dots, c_{i-1}, c_i + b, c_{i+1}, \dots, c_{r-1}, c_r - b, c_{r+1}, \dots, c_L)$, for some $b \in \mathbb{Z}_{>0}$.

Proof. By Proposition 3.2, determine $\mathbf{z}_w := (r, c)$ in $O(L^2)$ -time.

For (I), $D(w')$ is obtained from $D(w)$ by deleting \mathbf{z}_w ; so the expression in (I) is clear.

For (II), using Proposition 3.1, we can find $\mathbf{x} = (i, w(i))$, in $O(L^2)$ -time; this is our (T.2) pivot. Notice that row r of $D(w) \cap \mathcal{R}$ is nonempty (it contains $\mathbf{z}_w = (r, c)$); let b be the number of boxes in this row. It is straightforward from the graphical description of \mathcal{R} in terms of Rothe diagrams that each row of $D(w) \cap \mathcal{R}$ either has zero boxes or $b > 0$ boxes. Moreover, the d -th box (say, from the left) of each row are in the same column.

Suppose $j_1, \dots, j_m \in [i + 1, r]$ index the rows where $D(w) \cap \mathcal{R} \neq \emptyset$ (and thus has b boxes). (T.2) moves the b boxes of j_1 to row i and moves the b boxes of j_q to row j_{q-1} for $q = 2, \dots, m$. As explained above $j_m = r$, so (T.2) moves no boxes into row r . Thus row r of $D(w'') \cap \mathcal{R}$ has zero boxes.

It remains to compute b in $O(L^2)$ -time. Using Proposition 3.1 compute, in $O(L^2)$ -time,

$$m := \#\{h < r : w(h) < w(i)\}.$$

Clearly $b = c_r - [(w(i) - 1) - m]$. \square

Let $s_t = (x_r, m_t)$, as in (X.1), be a valid (multi)-deletion move on $u \in \mathcal{T}(w)$. Let $u^{(m_t)} \in \mathcal{T}(w)$ be defined by $u \xrightarrow{x_k} \bullet \dots \bullet \xrightarrow{x_k} u^{(m_t)}$ (m_t -times).

Proposition 5.11. *Suppose $u \in \mathcal{T}(w)$ where $\text{code}(u) = (\tilde{c}_1, \dots, \tilde{c}_{L'})$. Let $s_t = (x_k, m_t)$ or $s_t = i$ be as in (X.1). Given input $\text{code}(u)$ and s_t , there is an $O(L^2)$ algorithm to respectively determine if $u \xrightarrow{x_k} \bullet \dots \bullet \xrightarrow{x_k} u^{(m_t)}$ (m_t -times) or $u \xrightarrow{i} u''$ occurs in $\mathcal{T}(w)$ and (if yes) to compute*

- $\text{code}(u^{(m_t)})$ in the case $s_t = (x_k, m_t)$ (a multi-deletion move (T.1)), or
- $\text{code}(u'')$ in the case $s_t = i$ (a march move (T.2)).

Proof. By Proposition 5.10, $L' \leq L$. Thus in our run-time analysis, we replace L' by L .

Proposition 3.2 finds $\mathbf{z}_u := (r, c)$ (or determines it does not exist) in $O(L^2)$ -time. If \mathbf{z}_u does not exist then u is dominant and thus vexillary; output s_t is invalid. Thus we assume henceforth that \mathbf{z}_u exists.

Case 1: ($s_t = (x_k, m_t)$.) Proposition 3.1 finds $u(1), \dots, u(L')$ in $O(L^2)$ -time. Determine (taking $O(L^2)$ time) if

$$(25) \quad c_r - \left(\left(\min_{i \in [r]} u(i) \right) - 1 \right) \geq m_t,$$

holds. We claim that s_t is valid if and only if (25) holds and $k = r$. Indeed, observe

$$(26) \quad \#\{\text{boxes in row } r \text{ of } \text{Dom}(u)\} = \left(\min_{i \in [r]} u(i) \right) - 1.$$

Thus, (25) is equivalent to the existence of m_t boxes in row r of $D(u) \setminus \text{Dom}(u)$. By (T.1), if $k = r$ this is equivalent to being able to apply $\bullet \xrightarrow{x_r} \bullet$ successively m_t -times.

Finally, if s_t is valid, by m_t applications of Proposition 5.10 (I),

$$(27) \quad \text{code}(u^{(m_t)}) = (\tilde{c}_1, \dots, \tilde{c}_{r-1}, \tilde{c}_r - m_t, \tilde{c}_{r+1}, \dots, \tilde{c}_{L'}).$$

Hence we can output (27) in $O(L^2)$ -time.

Case 2: ($s_t = i$.) By Proposition 3.1, determine $u(1), \dots, u(L')$ from $\text{code}(u)$ in $O(L^2)$ -time. In particular this computes $\mathbf{x} := (i, u(i))$ in $O(L^2)$ -time. To decide if s_t is valid we must determine if $\mathbf{x} \in \text{Piv}(\mathbf{z}_u)$. To do this, first calculate (in $O(L)$ -time)

$$u_{NW}(\mathbf{z}_u) := \{(j, u(j)) : j < r, u(j) < c\}.$$

By definition,

$$\text{Piv}(\mathbf{z}_u) = \{(j, u(j)) \in u_{NW}(\mathbf{z}_u) : \nexists (h, u(h)) \in u_{NW}(\mathbf{z}_u) \text{ with } h > j, u(h) > u(j)\}.$$

$\text{Piv}(\mathbf{z}_u)$ takes $O(L)$ -time to compute since $\#u_{NW}(\mathbf{z}_u) \leq r - 1 \leq L - 1$. Hence we check if $\mathbf{x} \in \text{Piv}(\mathbf{z}_u)$ in $O(L)$ -time. If this is false, we output a rejection. Otherwise, Proposition 5.10 outputs $\text{code}(u'')$ in $O(L^2)$ -time. \square

Proposition 5.12. *If $S = (s_1, \dots, s_h) \in \text{Trans}(w, v)$ then $h \leq L^2$.*

Proof. Let $w := w_0 \xrightarrow{s_1} w_1 \xrightarrow{s_2} \dots \xrightarrow{s_h} w_h = v$ be the path in $\mathcal{T}(w)$ associated to S . By (T.1) and (T.2), $\mathbf{z}_{w_{t+1}}$ is weakly northwest of \mathbf{z}_{w_t} . Hence, for any fixed r , those $t \in [0, h - 1]$ with \mathbf{z}_{w_t} in row r form an interval $I^{(r)} \subseteq [0, h - 1]$. Since $1 \leq r \leq L$, it suffices to prove

$$(28) \quad \#I^{(r)} \leq 2(r - 1).$$

By (X.1) the transition moves acting on row r alternate between multi-(T.1) moves (x_r, m_t) and (T.2) moves. Thus to show (28), it is enough to prove

$$(29) \quad \#\{t \in I^{(r)} : w_{t-1} \rightarrow w_t \text{ is a (T.2) move}\} \leq r - 1.$$

Consider a march move i with $\mathbf{z}_{w_{t-1}} = (r, c)$ and $\mathbf{x} = (i, w_{t-1}(i)) \in \text{Piv}(\mathbf{z}_{w_{t-1}})$. By (T.2), if $(r, c') \in D(w_{t-1})$ is in the same connected component as $\mathbf{z}_{w_{t-1}}$, the move i takes (r, c') strictly north of row r . Thus, each march move strictly reduces the number of components in row r . Let $t_0 = \min\{t \in I^{(r)}\}$. Since there are at most r \bullet 's weakly above row r , $D(w_{t_0})$ has at most $r - 1$ (non-dominant) components in row r . Hence (29) holds, as desired. \square

Proposition 5.13. *Let v be vexillary with $\text{code}(v) = (c_1, \dots, c_{L'})$ and $L' \leq L$. There exists an $O(L^2)$ -time algorithm to check if $R = (r_{ij})_{1 \leq i, j \leq L'}$ is $R = R(T)$ for some $T \in \text{SSYT}(\lambda(v), \phi(v))$.*

Proof. Since $L' \leq L$, it is $O(L^2)$ -time to calculate $\phi(v), \lambda(v)$. Let

$$\lambda_i := \sum_{j=1}^{L'} r_{ij}, \text{ for } 1 \leq i \leq L'.$$

First verify (in $O(L)$ -time) that $\lambda_i \geq \lambda_{i+1}$ for $1 \leq i \leq L' - 1$. Then $R = R(T)$ where T is the (unique) row weakly increasing tableau of shape λ with r_{ij} many j 's in row i .

To verify $T \in \text{SSYT}(\lambda(v), \phi(v))$ we must check that it is (i) is flagged by $\phi(v)$, (ii) has shape $\lambda(v)$, and (iii) is semistandard. For (i), we need

$$(30) \quad r_{ij} = 0 \text{ if } j > \phi(v)_i, \text{ for all } i, j \in [L'].$$

For (ii), we need

$$(31) \quad \lambda_i = \lambda(v)_i \text{ for each } i \in [L'].$$

For (iii), it remains to ensure that T is column strict, i.e.,

$$(32) \quad \sum_{j' \leq j} r_{i+1, j'} \leq \sum_{j' < j} r_{i, j'} \text{ for each } i \in [L' - 1], j \in [L'].$$

We found the inequalities (31) and (32) from a (preprint) version of [10]. The inequalities (30), (31), and (32) can be checked in $O(L^2)$ -time since $i, j \in [L'] \subseteq [L]$. \square

The following completes our proof that we can check that $(S, R) \in X_{\alpha, w}$ in $L^{O(1)}$ -time.

Proposition 5.14. *Given $(S, R) \in X$ and $(\text{code}(w), \alpha)$, one can determine if $(S, R) \in X_{\alpha, w}$ in $L^{O(1)}$ -time.*

Proof. By Propositions 5.11 and 5.12 combined, one determines in $O(L^4)$ -time if S encodes a path $w := w_0 \xrightarrow{s_1} w_1 \xrightarrow{s_2} \dots \xrightarrow{s_h} w_h = v$ in $\mathcal{T}(w)$. If so, the length of $\text{code}(v)$ is at most L . Thus, using Theorem 5.1, one checks v is vexillary in $O(L^3)$ -time. This decides if S satisfies (X.1'). Proposition 5.13 checks R satisfies (X.2') in $O(L^2)$ -time. Finally since $h \leq L^2$, computing $\text{delwt}(S)$ takes $O(L^2)$ -time. Hence (X.3') is checkable in $O(L^2)$ time. \square

Proof of Theorem 5.6: By Proposition 5.9, $\#X_{\alpha, w} = c_{\alpha, w}$. By Proposition 5.12, $(S, R) \in \#X_{\alpha, w}$ only if the list S has at most L^2 elements. Assuming this, we check (S, R) satisfies (X.1) and (X.2) in $O(L^2)$ -time. Using Proposition 5.14, we can verify $(S, R) \in X_{\alpha, w}$ in $L^{O(1)}$ -time. Thus, given input α and $\text{code}(w)$, computing $c_{\alpha, w}$ is in $\#\text{P}$. \square

5.4. Hardness, and the conclusion of the proof of Theorem 1.5. *Schur polynomials* are an important basis of the vector space of symmetric polynomials. The Schur polynomial $s_\lambda = a_{\lambda+\delta}/a_\delta$ where $\lambda = (\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_n \geq 0)$, $a_\gamma := \det(x_i^{\gamma_j})_{i,j=1}^n$, and $\delta = (n-1, n-2, \dots, 2, 1, 0)$. The flagged Schur function of Section 5.1 is a generalization of the Schur polynomial.

A permutation w is *grassmannian* if it has at most one *descent* i , i.e., where $w(i) > w(i+1)$. Given a partition $(\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_L \geq 0)$ define a grassmannian permutation w_λ by setting

$$w_\lambda(i) = i + \lambda_{L-i+1} \text{ for } 1 \leq i \leq L.$$

For w_λ grassmannian, it is well-known (see, e.g., [8]) that

$$(33) \quad \text{code}(w_\lambda) = (\lambda_L, \lambda_{L-1}, \dots, \lambda_1).$$

Moreover,

$$(34) \quad \mathfrak{S}_{w_\lambda} = s_\lambda(x_1, \dots, x_L) = \sum_{\alpha \in \mathbb{Z}_{\geq 0}^L} K_{\lambda, \alpha} x^\alpha,$$

where $K_{\lambda, \alpha}$ is the *Kostka coefficient*. This number counts semistandard tableaux of shape λ with content α .

By (34),

$$(35) \quad c_{\alpha, w_\lambda} = K_{\lambda, \alpha}.$$

By Theorem 5.6, counting $c_{\alpha, w}$ is in #P. Suppose there is an oracle to compute $c_{\alpha, w}$ in polynomial time in the input length of $(\text{code}(w), \alpha)$. This input length is the same as for the input λ, α for $K_{\lambda, \alpha}$. Hence (33) and (35) combined imply a polynomial-time counting reduction from $\{c_{\alpha, w}\}$ to Kostka coefficients. Now H. Narayanan [10] proved that counting $K_{\lambda, \alpha}$ is a #P-complete problem. Thus counting $c_{\alpha, w}$ is a #P-complete problem. \square

Remark 5.15. Suppose the input for counting $c_{\alpha, w}$ is (α, w) where $w \in S_n$ (in one-line notation). Then the above counting reduction is not polynomial time in the input length of the Kostka problem. For example, suppose $\lambda = \alpha = (2^L, 2^L, \dots, 2^L)$ (L -many). Then the input length of this instance of the Kostka problem is $2L^2 \in O(L^2)$. On the other hand, $w_\lambda \in S_{L+2^L}$. Therefore, a polynomial time algorithm for the Schubert coefficient problem in n would have $\Omega(2^L)$ run time for the Kostka problem.

It seems unlikely that there is a polynomial-time reduction under this input assumption. This is our justification to encode w via $\text{code}(w)$ rather than one line notation. \square

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REFERENCES

- [1] S. Billey, W. Jockusch and R. P. Stanley, *Some combinatorial properties of Schubert polynomials*, J. Algebraic Combin. **2**(1993), no. 4, 345–374.
- [2] A. Fink, K. Mészáros, and A. St. Dizier, *Schubert polynomials as integer point transforms of generalized permutahedra*, Adv. Math. **332** (2018), 465–475.
- [3] S. Fomin, C. Greene, V. Reiner, and M. Shimozono, *Balanced labellings and Schubert polynomials*, European J. Combin. **18** (1997), no. 4, 373–389.
- [4] W. Fulton, *Young tableaux. With applications to representation theory and geometry*. London Mathematical Society Student Texts, 35. Cambridge University Press, Cambridge, 1997.
- [5] A. Knutson and A. Yong, *A formula for K -theory truncation Schubert calculus*, Int. Math. Res. Not. 2004, no. 70, 3741–3756.
- [6] A. Lascoux and M. -P. Schützenberger, *Schubert polynomials and the Littlewood-Richardson rule*, Letters in Math. Physics **10** (1985), 111–124.
- [7] ———, *Polynômes de Schubert*, C. R. Acad. Sci. Paris Sér. I Math. **294** (1982), 447–450.
- [8] L. Manivel, *Symmetric functions, Schubert polynomials and degeneracy loci*. Translated from the 1998 French original by John R. Swallow. SMF/AMS Texts and Monographs, American Mathematical Society, Providence, 2001.
- [9] C. Monical, N. Tokcan and A. Yong, *Newton polytopes in algebraic combinatorics*, Sel. Math. **25**(5) (2019), no. 66, 37 pp.
- [10] H. Narayanan, *On the complexity of computing Kostka numbers and Littlewood-Richardson coefficients*, J. Alg. Comb., Vol. 24, N. 3, 2006, 347–354.
- [11] C. H. Papadimitriou and K. Steiglitz, Kenneth, *Combinatorial optimization: algorithms and complexity*. Corrected reprint of the 1982 original. Dover Publications, Inc., Mineola, NY, 1998. xvi+496 pp.

- [12] A. Schrijver, *Theory of Linear and Integer Programming*, John Wiley & sons, 1998.
- [13] R. P. Stanley, *Some Schubert shenanigans*, preprint, 2017. arXiv:1704.00851
- [14] L. G. Valiant, *The complexity of computing the permanent*, Theoret. Comput. Sci., 8(2):189–201, 1979.
- [15] A. Weigandt, *Schubert polynomials, 132-patterns, and Stanley's conjecture*, Algebr. Comb. 1 (2018), no. 4, 415–423.

UCLA, LOS ANGELES, CA 90095

Email address: aadve@g.ucla.edu

DEPT. OF MATHEMATICS, U. ILLINOIS AT URBANA-CHAMPAIGN, URBANA, IL 61801, USA

Email address: cer2@illinois.edu, ayong@illinois.edu