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## RESEARCH ARTICLE

# A Proof of the Extended Delta Conjecture 

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#### Abstract

We prove the Extended Delta Conjecture of Haglund, Remmel, and Wilson, a combinatorial formula for $\Delta_{h_{l}} \Delta_{e_{k}}^{\prime} e_{n}$, where $\Delta_{e_{k}}^{\prime}$ and $\Delta_{h_{l}}$ are Macdonald eigenoperators and $e_{n}$ is an elementary symmetric function. We actually prove a stronger identity of infinite series of $\mathrm{GL}_{m}$ characters expressed in terms of LLT series. This is achieved through new results in the theory of the Schiffmann algebra and its action on the algebra of symmetric functions.


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## 1. Introduction

We prove the Extended Delta Conjecture of Haglund, Remmel and Wilson [14] by adapting methods from our work in [2] on a generalized Shuffle Theorem and proving new results about the action of the elliptic Hall algebra on symmetric functions. As in [2], we reformulate the conjecture as the polynomial truncation of an identity of infinite series of $\mathrm{GL}_{m}$ characters, expressed in terms of LLT series. We then prove the stronger infinite series identity using a Cauchy identity for non-symmetric Hall-Littlewood polynomials.

The conjecture stemmed from studies of the diagonal coinvariant algebra $\mathrm{DR}_{n}$ in two sets of $n$ variables, whose character as a doubly graded $S_{n}$ module has remarkable links with both classical combinatorial enumeration and the theory of Macdonald polynomials. It was shown in [17] that this character is neatly given by the formula $\Delta_{e_{n-1}}^{\prime} e_{n}$, where $e_{n}$ is the $n$-th elementary symmetric function, and for any symmetric function $f, \Delta_{f}^{\prime}$ is a certain eigenoperator on Macdonald polynomials whose eigenvalues depend on $f$.

The Shuffle Theorem, conjectured in [13] and proven by Carlsson and Mellit in [4], gives a combinatorial expression for $\Delta_{e_{n-1}}^{\prime} e_{n}$ in terms of Dyck paths - that is, lattice paths from ( $0, n$ ) to $(n, 0)$ that lie weakly below the line segment connecting these two points.

An expanded investigation led Haglund, Remmel and Wilson [14 to the Delta Conjecture, a combinatorial prediction for $\Delta_{e_{k}}^{\prime} e_{n}$, for all $0 \leq k<n$. This led to a flurry of activity (e.g. [6, 11, 14, 15, 16, 21, 22, 23, 26, 28]), including a conjecture by Zabrocki 27] that $\Delta_{e_{k}}^{\prime} e_{n}$ captures the character of the super-diagonal coinvariant ring $\mathrm{SDR}_{n}$, a deformation of $\mathrm{DR}_{n}$ involving the addition of a set of anti-commuting variables.

The Delta Conjecture has been extended in two directions. One gives a Compositional generalization, recently proved by D'Adderio and Mellit [7]. The other involves a second eigenoperator $\Delta_{h_{l}}$, where $h_{l}$ is the $l$-th homogeneous symmetric function. The Extended Delta Conjecture [14, Conjecture 7.4] is, for $l \geq 0$ and $1 \leq k \leq n$,

$$
\begin{equation*}
\Delta_{h_{l}} \Delta_{e_{k-1}}^{\prime} e_{n}=\left\langle z^{n-k}\right\rangle \sum_{\substack{\lambda \in \mathbf{D}_{n+l} \\ P \in \mathbf{L}_{n+l, l}(\lambda)}} q^{\operatorname{dinv}(P)} t^{\operatorname{area}(\lambda)} x^{\mathrm{wt}+(P)} \prod_{i: r_{i}(\lambda)=r_{i-1}(\lambda)+1}\left(1+z t^{-r_{i}(\lambda)}\right), \tag{1}
\end{equation*}
$$

in which $\lambda$ is a Dyck path and $P$ is a certain type of labelling of $\lambda$ (see $\S 2$ for full definitions). D'Adderio, Iraci and Wyngaerd proved the Schröder case and the $t=0$ specialization of the conjecture [5, 6]; Qiu and Wilson [21] reformulated the conjecture and established the $q=0$ specialization as well.

Let us briefly outline the steps by which we prove (11).
Feigin-Tsymbaliuk [8] and Schiffmann-Vasserot [25] constructed an action of the elliptic Hall algebra $\mathcal{E}$ of Burban and Schiffmann [3] on the algebra of symmetric functions. The operators $\Delta_{f}$ and $\Delta_{f}^{\prime}$ are part of the $\mathcal{E}$ action. In Theorem4.4.1, we use this to reformulate the left hand side of (1) as the polynomial part of an explicit infinite series of virtual $\mathrm{GL}_{m}$ characters with coefficients in $\mathbb{Q}(q, t)$. The proof of Theorem 4.4.1 relies on a symmetry (Proposition 4.3.3) between distinguished elements of $\mathcal{E}$ introduced by Negut 19 and their transposes.

In Theorem 5.1.1, we also reformulate the right hand side of (1) as the polynomial part of an infinite series, in this case expressed in terms of the LLT series introduced by Grojnowski and Haiman in [12]. This given, we ultimately arrive at Theorem6.3.6 - an identity of infinite series of $\mathrm{GL}_{m}$ characters which implies the Extended Delta Conjecture by taking the polynomial part on each side.

Although the Extended Delta Conjecture and the Compositional Delta Conjecture both imply the Delta Conjecture, they generalize it in different directions and our methods are quite different from those of D'Adderio and Mellit. It would be interesting to know whether a common generalization is possible.

## 2. The Extended Delta Conjecture

The Extended Delta Conjecture equates a "symmetric function side," involving the action of a Macdonald operator on an elementary symmetric function, with a "combinatorial side." We begin by recalling the definitions of these two quantities.

### 2.1. Symmetric function side

Integer partitions are written $\lambda=\left(\lambda_{1} \geq \cdots \geq \lambda_{l}\right)$, sometimes with trailing zeroes allowed. We set $|\lambda|=\lambda_{1}+\cdots+\lambda_{l}$ and let $\ell(\lambda)$ be the number of non-zero parts. We identify a partition $\lambda$ with its French style Ferrers shape, the set of lattice squares (or boxes) with northeast corner in the set

$$
\begin{equation*}
\left\{(i, j) \mid 1 \leq j \leq \ell(\lambda), 1 \leq i \leq \lambda_{j}\right\} \tag{2}
\end{equation*}
$$

The shape generator of $\lambda$ is the polynomial

$$
\begin{equation*}
B_{\lambda}(q, t)=\sum_{(i, j) \in \lambda} q^{i-1} t^{j-1} \tag{3}
\end{equation*}
$$

Let $\Lambda=\Lambda_{\mathbf{k}}(X)$ be the algebra of symmetric functions in an infinite alphabet of variables $X=$ $x_{1}, x_{2}, \ldots$, with coefficients in the field $\mathbf{k}=\mathbb{Q}(q, t)$. We follow the notation of Macdonald 18 for the graded bases of $\Lambda$. Basis elements are indexed by a partition $\lambda$ and have homogeneous degree $|\lambda|$. Examples include the elementary symmetric functions $e_{\lambda}=e_{\lambda_{1}} \cdots e_{\lambda_{k}}$, complete homogeneous symmetric functions $h_{\lambda}=h_{\lambda_{1}} \cdots h_{\lambda_{k}}$, power-sums $p_{\lambda}=p_{\lambda_{1}} \cdots p_{\lambda_{k}}$, monomial symmetric functions $m_{\lambda}$, and Schur functions $s_{\lambda}$.

As is conventional, $\omega: \Lambda \rightarrow \Lambda$ denotes the $\mathbf{k}$-algebra involution defined by $\omega s_{\lambda}=s_{\lambda^{*}}$, where $\lambda^{*}$ denotes the transpose of $\lambda$, and $\langle-,-\rangle$ denotes the symmetric bilinear inner product such that $\left\langle s_{\lambda}, s_{\mu}\right\rangle=\delta_{\lambda, \mu}$.

The basis of modified Macdonald polynomials, $\tilde{H}_{\mu}(X ; q, t)$, is defined 9$]$ from the integral form Macdonald polynomials $J_{\mu}(X ; q, t)$ of [18 using the device of plethystic evaluation. For an expression $A$ in terms of indeterminates, such as a polynomial, rational function, or formal series, $p_{k}[A]$ is defined to be the result of substituting $a^{k}$ for every indeterminate $a$ occurring in $A$. We define $f[A]$ for any $f \in \Lambda$ by substituting $p_{k}[A]$ for $p_{k}$ in the expression for $f$ as a polynomial in the power-sums $p_{k}$, so that $f \mapsto f[A]$ is a homomorphism. The variables $q, t$ from our ground field $\mathbf{k}$ count as indeterminates. The modified Macdonald polynomials are defined by

$$
\begin{equation*}
\tilde{H}_{\mu}(X ; q, t)=t^{n(\mu)} J_{\mu}\left[\frac{X}{1-t^{-1}} ; q, t^{-1}\right] \tag{4}
\end{equation*}
$$

where

$$
\begin{equation*}
n(\mu)=\sum_{i}(i-1) \mu_{i} \tag{5}
\end{equation*}
$$

For any symmetric function $f \in \Lambda$, let $f[B]$ denote the eigenoperator on the basis $\left\{\tilde{H}_{\mu}\right\}$ of $\Lambda$ such that

$$
\begin{equation*}
f[B] \tilde{H}_{\mu}=f\left[B_{\mu}(q, t)\right] \tilde{H}_{\mu} \tag{6}
\end{equation*}
$$

The left hand side of (1) is expressed in the notation of [14], where $\Delta_{f}=f[B]$ and $\Delta_{f}^{\prime}=$ $f[B-1]$. Hence, the symmetric function side of the Extended Delta Conjecture is

$$
\begin{equation*}
h_{l}[B] e_{k-1}[B-1] e_{n} \tag{7}
\end{equation*}
$$

### 2.2. The combinatorial side

The right hand side of the Extended Delta Conjecture (1) is a combinatorial generating function that counts labelled lattice paths.

Definition 2.2.1. A Dyck path is a south-east lattice path lying weakly below the line segment connecting the points $(0, N)$ and $(N, 0)$. The set of such paths is denoted $\mathbf{D}_{N}$. The staircase path $\delta$ is the Dyck path alternating between south and east steps.

Each $\lambda \in \mathbf{D}_{N}$ has area $(\lambda)=|\delta / \lambda|$ defined to be the number of lattice squares lying above $\lambda$ and below $\delta$. Let $r_{i}(\lambda)$ be the area contribution from squares in the $i$-th row, numbered from north to south; in other words, $r_{i}$ is the distance from the $i$-th south step of $\lambda$ to the $i$-th south step of $\delta$. Note that

$$
\begin{equation*}
r_{1}(\lambda)=0, \quad r_{i}(\lambda) \leq r_{i-1}(\lambda)+1 \quad \text { for } i>1, \quad \text { and } \quad \sum_{i=1}^{N} r_{i}(\lambda)=|\delta / \lambda| . \tag{8}
\end{equation*}
$$

Definition 2.2.2. A labelling $P=\left(P_{1}, \ldots, P_{N}\right) \in \mathbb{N}^{N}$ attaches a label in $\mathbb{N}=\{0,1, \ldots\}$ to each south step of $\lambda \in \mathbf{D}_{N}$ so that the labels increase from north to south along vertical runs of south steps, as shown in Figure 1 . The set of labellings is denoted by $\mathbf{L}_{N}(\lambda)$, or simply $\mathbf{L}(\lambda)$. Given $0 \leq l<N$, a partial labelling of $\lambda \in \mathbf{D}_{N}$ is a labelling where 0 occurs exactly $l$ times and never on the run at $x=0$. We denote the set of these partial labellings by $\mathbf{L}_{N, l}(\lambda)$.

Each labelling $P \in \mathbf{L}(\lambda)$ is assigned a statistic $\operatorname{dinv}(P)$, defined to be the number of pairs $(i<j)$ such that either

$$
\left\{\begin{array}{l}
r_{i}(\lambda)=r_{j}(\lambda) \text { and } P_{i}<P_{j} \text { or }  \tag{9}\\
r_{i}(\lambda)=r_{j}(\lambda)+1 \text { and } P_{i}>P_{j}
\end{array}\right.
$$

The weight of a labelling $P$ is defined so zero labels do not contribute, by

$$
\begin{equation*}
x^{\mathrm{wt}+(P)}=\prod_{i \in[N]: P_{i} \neq 0} x_{P_{i}} \tag{10}
\end{equation*}
$$

This is equivalent to letting $x_{0}=1$ in $x^{\mathrm{wt}(P)}:=\prod_{i \in[N]} x_{P_{i}}$.
The above defines the right hand side of 11, with $\left\langle z^{n-k}\right\rangle$ denoting the coefficient of $z^{n-k}$. Remark 2.2.3. In [14, a Dyck path is a north-east lattice path lying weakly above the line segment connecting $(0,0)$ and ( $N, N$ ), and labellings increase from south to north along vertical runs. After reflecting the picture about a horizontal line, our conventions on paths, labellings, and the definition of $\operatorname{dinv}(P)$ match those in [14. Separately, [13] uses the same conventions


Figure 1. A path $\lambda$ and partial labelling $P \in \mathbf{L}_{11,2}(\lambda)$, with $\operatorname{area}(\lambda)=10, \operatorname{dinv}(P)=15$, $x^{\mathrm{wt}+(P)}=x_{1}^{2} x_{2} x_{3}^{2} x_{4}^{2} x_{5} x_{6}$, and $x^{\mathrm{wt}(P)}=x_{0}^{2} x_{1}^{2} x_{2} x_{3}^{2} x_{4}^{2} x_{5} x_{6}$.
that we do for Dyck paths, but defines labellings to increase from south to north, and defines $\operatorname{dinv}(P)$ with the inequalities in (9) reversed. However, since the sum

$$
\begin{equation*}
\sum_{P \in \mathbf{L}(\lambda)} q^{\operatorname{dinv}(P)} x^{\mathrm{wt}(P)} \tag{11}
\end{equation*}
$$

is a symmetric function [13], it is unchanged if we reverse the ordering on labels, after which the conventions in [13] agree with those used here.

We prefer another slight modification based on the following lemma which was mentioned in [14] without details.

Lemma 2.2.4. For any Dyck path $\lambda \in \mathbf{D}_{N}$, we have

$$
\begin{equation*}
\prod_{\substack{1<i \leq N \\ r_{i}(\lambda)=r_{i-1}(\lambda)+1}}\left(1+z t^{-r_{i}(\lambda)}\right)=\prod_{\substack{1<i \leq N \\ c_{i}(\lambda)=c_{i-1}(\lambda)+1}}\left(1+z t^{-c_{i}(\lambda)}\right), \tag{12}
\end{equation*}
$$

where $c_{i}(\lambda)=r_{i}\left(\lambda^{*}\right)$ is the contribution to $|\delta / \lambda|$ from boxes in the $i$-th column, numbered from right to left.

Proof. The condition $r_{i}(\lambda)=r_{i-1}(\lambda)+1$ means that $\lambda$ has consecutive south steps in rows $i-1$ and $i$ with no intervening east step. Similarly, $c_{i}(\lambda)=c_{i-1}(\lambda)+1$ if and only if $\lambda$ has consecutive east steps in columns $i-1$ and $i$ (numbered right to left). Consider the word formed by listing the steps in $\lambda$ in the southeast direction from $(0, N)$ to $(N, 0)$, as shown here for the example in Figure 1.

SSSESEESSEEESESESSESEE


Treating south and east steps as left and right parentheses, each south step pairs with an east step to its right, and we have $r_{i}(\lambda)=c_{j}(\lambda)$ if the $i$-th south step (numbered left to right) pairs with the $j$-th east step (numbered right to left). Furthermore, the leftmost member of each
double south step pairs with the rightmost member of a double east step, as indicated in the word displayed above.

Since each index $i-1$ such that $r_{i}(\lambda)=r_{i-1}(\lambda)+1$ pairs with an index $j-1$ such that $c_{j}(\lambda)=c_{j-1}(\lambda)+1$, we have

$$
\begin{equation*}
\prod_{\substack{1<i \leq N \\ r_{i}(\lambda)=r_{i-1}(\lambda)+1}}\left(1+z t^{-r_{i-1}(\lambda)-1}\right)=\prod_{\substack{1<j \leq N \\ c_{j}(\lambda)=c_{j-1}(\lambda)+1}}\left(1+z t^{-c_{j-1}(\lambda)-1}\right) \tag{13}
\end{equation*}
$$

Now (12) follows.
Setting $N=n+l$ and $m=k+l$, the right hand side of (1), or the combinatorial side of the Extended Delta Conjecture, is equal to

$$
\begin{equation*}
\left\langle z^{N-m}\right\rangle \sum_{\substack{\lambda \in \mathbf{D}_{N} \\ P \in \mathbf{L}_{N, l}(\lambda)}} t^{|\delta / \lambda|} q^{\operatorname{dinv}(P)} x^{\mathrm{wt}+(P)} \prod_{\substack{1<i \leq N \\ c_{i}(\lambda)=c_{i-1}(\lambda)+1}}\left(1+z t^{-c_{i}(\lambda)}\right) \tag{14}
\end{equation*}
$$

## 3. Background on the Schiffmann algebra $\mathcal{E}$

From work of Feigin and Tsymbaliuk [8] and Schiffmann and Vasserot [25], we know that the operators $f[B]$ in (7) form part of an action of the elliptic Hall algebra $\mathcal{E}$ of Burban and Schiffmann [3, 24], or Schiffmann algebra for short, on the algebra of symmetric functions. In [2], we used this action to express the symmetric function side of a generalized Shuffle Theorem as the polynomial part of an explicit infinite series of $\mathrm{GL}_{l}$ characters. Here we derive a similar expression (Theorem4.4.1) for the symmetric function side (7) of the Extended Delta Conjecture.

For this purpose, we need a deeper study of the Schiffmann algebra than we did in [2], where a fragment of the theory was enough. We start with a largely self-contained description of $\mathcal{E}$ and its action on $\Lambda$, although we occasionally refer to [2] for the restatements of results from [3, 24, 25] in our notation, and for some proofs. A precise translation between our notation and that of [3, 24, 25] can be found in [2, eq. (25)]. In the presentation of $\mathcal{E}$ and its action on $\Lambda$, we freely use plethystic substitution, defined in 2.1 . Indeed, the ability to do so is a principal reason why we prefer the notation used here to that in the foundational papers on the Schiffmann algebra.

### 3.1. Description of $\mathcal{E}$

Let $\mathbf{k}=\mathbb{Q}(q, t)$, as in 42 . The Schiffmann algebra $\mathcal{E}$ is generated by a central Laurent polynomial subalgebra $F=\mathbf{k}\left[c_{1}^{ \pm}, c_{2}^{ \pm 1}\right]$ and a family of subalgebras $\Lambda_{F}\left(X^{m, n}\right)$ isomorphic to the algebra of symmetric functions $\Lambda_{F}(X)$ over $F$, one for each pair of coprime integers $(m, n)$. These are subject to defining relations spelled out below.

For any algebra $A$ containing a copy of $\Lambda$, there is an adjoint action of $\Lambda$ on $A$ arising from the Hopf algebra structure of $\Lambda$. Using two formal alphabets $X$ and $Y$ to distinguish between the tensor factors in $\Lambda \otimes \Lambda \cong \Lambda(X) \Lambda(Y)$, the coproduct and antipode for the Hopf algebra structure are given by the plethystic substitutions

$$
\begin{equation*}
\Delta f=f[X+Y], \quad S(f)=f[-X] \tag{15}
\end{equation*}
$$

The adjoint action of $f \in \Lambda$ on $\zeta \in A$ is then given by

$$
\begin{equation*}
(\operatorname{Ad} f) \zeta=\sum_{i} f_{i} \zeta g_{i}, \quad \text { where } \quad f[X-Y]=\sum_{i} f_{i}(X) g_{i}(Y) \tag{16}
\end{equation*}
$$

since the formula on the right is another way to write $(1 \otimes S) \Delta f=\sum_{i} f_{i} \otimes g_{i}$. More explicitly, we have

$$
\begin{equation*}
\left(\operatorname{Ad} p_{n}\right) \zeta=\left[p_{n}, \zeta\right] \quad \text { and } \quad\left(\operatorname{Ad} h_{n}\right) \zeta=\sum_{j+k=n}(-1)^{k} h_{j} \zeta e_{k} \tag{17}
\end{equation*}
$$

The last formula can be expressed for all $n$ at once as a generating function identity

$$
\begin{equation*}
(\operatorname{Ad} \Omega[z X]) \zeta=\Omega[z X] \zeta \Omega[-z X] \tag{18}
\end{equation*}
$$

where

$$
\begin{equation*}
\Omega(X)=\sum_{n=0}^{\infty} h_{n}(X) . \tag{19}
\end{equation*}
$$

We fix notation for the quantities

$$
\begin{equation*}
M=(1-q)(1-t), \quad \widehat{M}=\left(1-(q t)^{-1}\right) M \tag{20}
\end{equation*}
$$

which play a role in the presentation of $\mathcal{E}$ and will be referred to again later.

### 3.1.1. Basic structure and symmetries

The algebra $\mathcal{E}$ is $\mathbb{Z}^{2}$ graded with the central subalgebra $F$ in degree $(0,0)$ and $f\left(X^{m, n}\right)$ in degree $(d m, d n)$ for $f(X)$ of degree $d$ in $\Lambda(X)$.

The universal central extension $\widehat{\mathrm{SL}_{2}(\mathbb{Z})} \rightarrow \mathrm{SL}_{2}(\mathbb{Z})$ acts on the set of tuples

$$
\begin{equation*}
\left\{(m, n, \theta) \in\left(\mathbb{Z}^{2} \backslash \mathbf{0}\right) \times \mathbb{R} \mid \theta \text { is a value of } \arg (m+i n)\right\} \tag{21}
\end{equation*}
$$

lifting the $\mathrm{SL}_{2}(\mathbb{Z})$ action on pairs $(m, n)$, with the central subgroup $\mathbb{Z}$ generated by the 'rotation by $2 \pi^{\prime} \operatorname{map}(m, n, \theta) \mapsto(m, n, \theta+2 \pi)$. The group $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ acts on $\mathcal{E}$ by $\mathbf{k}$-algebra automorphisms, compatibly with the action of $\mathrm{SL}_{2}(\mathbb{Z})$ on the grading group $\mathbb{Z}^{2}$. Before giving the defining relations of $\mathcal{E}$, we specify how $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ acts on the generators.

For each pair of coprime integers $(m, n)$, we introduce a family of alphabets $X_{\theta}^{m, n}$, one for each value $\theta$ of $\arg (m+i n)$, related by

$$
\begin{equation*}
X_{\theta+2 \pi}^{m, n}=c_{1}^{m} c_{2}^{n} X_{\theta}^{m, n} \tag{22}
\end{equation*}
$$

We make the convention that $X^{m, n}$ without a subscript means $X_{\theta}^{m, n}$ with $\theta \in(-\pi, \pi]$. For comparison, the implied convention in [2] is $\theta \in[-\pi, \pi)$. The subalgebra $\Lambda_{F}\left(X^{m, n}\right)=\Lambda_{F}\left(X_{\theta}^{m, n}\right)$ only depends on $(m, n)$ and so does not depend on the choice of branch for the angle $\theta$. When we refer to a subalgebra $\Lambda_{\mathbf{k}}\left(X^{m, n}\right)$, which does depend on the branch, the convention $\theta \in(-\pi, \pi]$ applies.

The $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ action is now given by $\rho \cdot f\left(X_{\theta}^{m, n}\right)=f\left(X_{\theta^{\prime}}^{m^{\prime}, n^{\prime}}\right)$ for $f(X) \in \Lambda_{\mathbf{k}}(X)$ where $\rho \in \widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ acts on the indexing data in (21) by $\rho \cdot(m, n, \theta)=\left(m^{\prime}, n^{\prime}, \theta^{\prime}\right)$. Note that if $m, n$ are coprime, then so are $m^{\prime}, n^{\prime}$. The action on $F$ factors through the action of $\mathrm{SL}_{2}(\mathbb{Z})$ on the group algebra $\mathbf{k} \cdot \mathbb{Z}^{2} \cong F$.

For instance, the 'rotation by $2 \pi$ ' element $\rho \in \widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ fixes $F$, and has $\rho \cdot f\left(X_{\theta}^{m, n}\right)=$ $f\left(X_{\theta+2 \pi}^{m, n}\right)=f\left[c_{1}^{m} c_{2}^{n} X_{\theta}^{m, n}\right]$. Thus $\rho$ coincides with multiplication by $c_{1}^{r} c_{2}^{s}$ in degree $(r, s)$, and automatically preserves all relations that respect the $\mathbb{Z}^{2}$ grading.

We now turn to the defining relations of $\mathcal{E}$. Apart from the relations implicit in $F=$ $\mathbf{k}\left[c_{1}^{ \pm 1}, c_{2}^{ \pm 1}\right]$ being central and each $\Lambda_{F}\left(X^{m, n}\right)$ being isomorphic to $\Lambda_{F}(X)$, these fall into three families: Heisenberg relations, internal action relations and axis-crossing relations.

### 3.1.2. Heisenberg relations

Each pair of subalgebras $\Lambda_{F}\left(X^{m, n}\right)$ and $\Lambda_{F}\left(X^{-m,-n}\right)$ in degrees along opposite rays in $\mathbb{Z}^{2}$ satisfy Heisenberg relations

$$
\begin{equation*}
\left[p_{k}\left(X_{\theta}^{-m,-n}\right), p_{l}\left(X_{\theta+\pi}^{m, n}\right)\right]=\delta_{k, l} k p_{k}\left[\left(c_{1}^{m} c_{2}^{n}-1\right) / \widehat{M}\right] \tag{23}
\end{equation*}
$$

where $\widehat{M}$ is given by 20 . As an exercise, the reader can check, using $\sqrt[22]{ }$, that the relations in (23) are consistent with swapping the roles of $\Lambda_{F}\left(X^{m, n}\right)$ and $\Lambda_{F}\left(X^{-m,-n}\right)$.

### 3.1.3. Internal action relations

The internal action relations describe the adjoint action of each $\Lambda_{F}\left(X^{m, n}\right)$ on $\mathcal{E}$. For simplicity, we write these relations, and also the axis-crossing relations below, with $\Lambda_{F}\left(X^{1,0}\right)$ distinguished. The full set of relations is understood to be given by closing the stated relations under the $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ action.

Bearing in mind that $X^{m, n}$ means $X_{\theta}^{m, n}$ with $\theta \in(-\pi, \pi]$, the relations for the internal action of $\Lambda_{F}\left(X^{1,0}\right)$ are:

$$
\begin{align*}
\left(\operatorname{Ad} f\left(X^{1,0}\right)\right) p_{1}\left(X^{m, 1}\right) & =(\omega f)[z] \mid z^{k} \mapsto p_{1}\left(X^{m+k, 1}\right) \\
\left(\operatorname{Ad} f\left(X^{1,0}\right)\right) p_{1}\left(X^{m,-1}\right) & =(\omega f)[-z] \mid z^{k} \mapsto p_{1}\left(X^{m+k,-1}\right) \tag{24}
\end{align*}
$$

### 3.1.4. Axis-crossing relations

Again distinguishing $\Lambda_{F}\left(X^{1,0}\right)$ and taking angles on the branch $\theta \in(-\pi, \pi]$, the final set of relations is the closure under the $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ action of

$$
\begin{equation*}
\left[p_{1}\left(X^{b,-1}\right), p_{1}\left(X^{a, 1}\right)\right]=-\frac{e_{a+b}\left[-\widehat{M} X^{1,0}\right]}{\widehat{M}} \text { for } a+b>0 . \tag{25}
\end{equation*}
$$

More generally, rotating this relation by $\pi$ determines $\left[p_{1}\left(X^{b,-1}\right), p_{1}\left(X^{a, 1}\right)\right]$ for $a+b<0$, and the Heisenberg relations determine it when $a+b=0$. Combining these gives

$$
\left[p_{1}\left(X^{b,-1}\right), p_{1}\left(X^{a, 1}\right)\right]=-\frac{1}{\widehat{M}} \begin{cases}e_{a+b}\left[-\widehat{M} X^{1,0}\right] & a+b>0  \tag{26}\\ 1-c_{1}^{-b} c_{2} & a+b=0 \\ -c_{1}^{-b} c_{2} e_{-(a+b)}\left[-\widehat{M} X^{-1,0}\right] & a+b<0\end{cases}
$$

### 3.1.5. Further remarks

Define upper and lower half subalgebras $\mathcal{E}^{*,>0}, \mathcal{E}^{*,<0} \subseteq \mathcal{E}$ to be generated by the $\Lambda_{F}\left(X^{m, n}\right)$ with $n>0$ or $n<0$, respectively. Using the $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ image of the relations in 25 , one can express any $e_{k}\left[-\widehat{M} X^{m, n}\right]$ for $n>0$ in terms of iterated commutators of the elements $p_{1}\left(X^{a, 1}\right)$. This shows that $\left\{p_{1}\left(X^{a, 1}\right) \mid a \in \mathbb{Z}\right\}$ generates $\mathcal{E}^{*,>0}$ as an $F$-algebra. Similarly, $\left\{p_{1}\left(X^{a,-1}\right) \mid a \in \mathbb{Z}\right\}$ generates $\mathcal{E}^{*,<0}$.

The internal action relations give the adjoint action of $\Lambda_{F}\left(X^{1,0}\right)$ on the space spanned by $\left\{p_{1}\left(X^{a, \pm 1}\right) \mid a \in \mathbb{Z}\right\}$. Using the formula $(\operatorname{Ad} f)\left(\zeta_{1} \zeta_{2}\right)=\sum\left(\left(\operatorname{Ad} f_{(1)}\right) \zeta_{1}\right)\left(\left(\operatorname{Ad} f_{(2)}\right) \zeta_{2}\right)$, where $\Delta f=\sum f_{(1)} \otimes f_{(2)}$ in Sweedler notation, this determines the adjoint action of $\Lambda_{F}\left(X^{1,0}\right)$ on $\mathcal{E}^{*,>0}$ and $\mathcal{E}^{*,<0}$. The Heisenberg relations give the adjoint action of $\Lambda_{F}\left(X^{1,0}\right)$ on $\Lambda_{F}\left(X^{-1,0}\right)$, while $\Lambda_{F}\left(X^{1,0}\right)$ acts trivially on itself, with $(\operatorname{Ad} f) g=f[1] g$.

Together these determine the adjoint action of $\Lambda_{F}\left(X^{1,0}\right)$ on the whole algebra $\mathcal{E}$. By symmetry, the same holds for the adjoint action of any $\Lambda_{F}\left(X^{m, n}\right)$.

### 3.1.6. Anti-involution

One can check from the defining relations above that $\mathcal{E}$ has a further symmetry given by an involutory anti-automorphism (product reversing automorphism)

$$
\begin{gather*}
\Phi: \mathcal{E} \rightarrow \mathcal{E} \\
\Phi\left(g\left(c_{1}, c_{2}\right)\right)=g\left(c_{2}^{-1}, c_{1}^{-1}\right), \quad \Phi\left(f\left(X_{\theta}^{m, n}\right)\right)=f\left(X_{\pi / 2-\theta}^{n, m}\right) . \tag{27}
\end{gather*}
$$

Note that $\Phi$ is compatible with reflecting degrees in $\mathbb{Z}^{2}$ about the line $x=y$. Together with $\widehat{\mathrm{SL}_{2}(\mathbb{Z})}$ it generates a $\widehat{\mathrm{GL}_{2}(\mathbb{Z})}$ action on $\mathcal{E}$ for which $\rho \in \widehat{\mathrm{GL}_{2}(\mathbb{Z})}$ is an anti-automorphism if $\widehat{\mathrm{GL}_{2}(\mathbb{Z})} \rightarrow \mathrm{GL}_{2}(\mathbb{Z}) \xrightarrow{\text { det }}\{ \pm 1\}$ sends $\rho$ to -1 .

### 3.2. Action of $\mathcal{E}$ on $\Lambda$

We write $f$ • for the operator of multiplication by a function $f$ to better distinguish between operator expressions such as $(\omega f)^{\bullet}$ and $\omega \cdot f^{\bullet}$. For $f$ a symmetric function, $f^{\perp}$ denotes the $\langle-,-\rangle$ adjoint of $f^{\bullet}$.

Here and again later on, we use an overbar to indicate inverting the variables in any expression; for example

$$
\begin{equation*}
\bar{M}=\left(1-q^{-1}\right)\left(1-t^{-1}\right) . \tag{28}
\end{equation*}
$$

We extend the notation in (6) accordingly, setting

$$
\begin{equation*}
f[\bar{B}] \tilde{H}_{\mu}=f\left[B_{\mu}\left(q^{-1}, t^{-1}\right)\right] \tilde{H}_{\mu} . \tag{29}
\end{equation*}
$$

Proposition 3.2.1 ([2, Prop 3.3.1]). There is an action of $\mathcal{E}$ on $\Lambda$ characterized uniquely by the following properties.
(i) The central parameters $c_{1}, c_{2}$ act as scalars

$$
\begin{equation*}
c_{1} \mapsto 1, \quad c_{2} \mapsto(q t)^{-1} \tag{30}
\end{equation*}
$$

(ii) The subalgebras $\Lambda_{\mathbf{k}}\left(X^{ \pm 1,0}\right)$ act as

$$
\begin{equation*}
f\left(X^{1,0}\right) \mapsto(\omega f)[B-1 / M], \quad f\left(X^{-1,0}\right) \mapsto(\omega f)[\overline{1 / M-B}] . \tag{31}
\end{equation*}
$$

(iii) The subalgebras $\Lambda_{\mathbf{k}}\left(X^{0, \pm 1}\right)$ act as

$$
\begin{equation*}
f\left(X^{0,1}\right) \mapsto f[-X / M]^{\bullet}, \quad f\left(X^{0,-1}\right) \mapsto f(X)^{\perp} \tag{32}
\end{equation*}
$$

We will make particular use of operators representing the action on $\Lambda$ of elements $p_{1}\left(X^{a, 1}\right)$ and $p_{1}\left(X^{1, a}\right)$ in $\mathcal{E}$. For the first we need the operator $\nabla$, defined in [1] as an eigenoperator on the modified Macdonald basis by

$$
\begin{equation*}
\nabla \tilde{H}_{\mu}=t^{n(\mu)} q^{n\left(\mu^{*}\right)} \tilde{H}_{\mu} \tag{33}
\end{equation*}
$$

where $n(\mu)$ is given by (5) and $\mu^{*}$ denotes the transpose partition.
For the second, we introduce the doubly infinite generating series

$$
\begin{equation*}
D(z)=\omega \Omega\left[z^{-1} X\right]^{\bullet}(\omega \Omega[-z M X])^{\perp} \tag{34}
\end{equation*}
$$

where $\Omega(X)$ is given by 19$)$.

Definition 3.2.2. For $a \in \mathbb{Z}$, we define operators on $\Lambda=\Lambda_{\mathbf{k}}(X)$ :

$$
\begin{align*}
E_{a} & =\nabla^{a} e_{1}(X)^{\bullet} \nabla^{-a}  \tag{35}\\
D_{a} & =\left\langle z^{-a}\right\rangle D(z) \tag{36}
\end{align*}
$$

The operators $D_{a}$ are the same as in [2] and differ by a sign $(-1)^{a}$ from those in [1, 10].
Proposition 3.2.3. In the action of $\mathcal{E}$ on $\Lambda$ given by Proposition 3.2.1:
(i) the element $p_{1}\left[-M X^{1, a}\right]=-M p_{1}\left(X^{1, a}\right) \in \mathcal{E}$ acts as the operator $D_{a}$;
(ii) the element $p_{1}\left[-M X^{a, 1}\right]=-M p_{1}\left(X^{a, 1}\right) \in \mathcal{E}$ acts as the operator $E_{a}$.

Proof. Part (i) is proven in [2, Prop 3.3.4].
By (32), $p_{1}\left[-M X^{0,1}\right]$ acts on $\Lambda$ as multiplication by $p_{1}[X]=e_{1}(X)$. It was shown in [2, Lemma 3.4.1] that the action of $\mathcal{E}$ on $\Lambda$ satisfies the symmetry $\nabla f\left(X^{m, n}\right) \nabla^{-1}=f\left(X^{m+n, n}\right)$. More generally, this implies $\nabla^{a} f\left(X^{m, n}\right) \nabla^{-a}=f\left(X^{m+a n, n}\right)$ for every integer $a$. Hence, $p_{1}\left[-M X^{a, 1}\right]$ acts as $\nabla^{a} p_{1}\left[-M X^{0,1}\right] \nabla^{-a}=\nabla^{a} e_{1}(X)^{\bullet} \nabla^{-a}$.

## 3.3. $\mathrm{GL}_{l}$ characters and the shuffle algebra

As usual, the weight lattice of $\mathrm{GL}_{l}$ is $\mathbb{Z}^{l}$, with Weyl group $W=S_{l}$ permuting the coordinates. A weight $\lambda$ is dominant if $\lambda_{1} \geq \cdots \geq \lambda_{l}$. A polynomial weight is a dominant weight $\lambda$ such that $\lambda_{l} \geq 0$. In other words, polynomial weights of $\mathrm{GL}_{l}$ are integer partitions of length at most $l$.

As in [2], $\S 2.3$, we identify the algebra of virtual $\mathrm{GL}_{l}$ characters over $\mathbf{k}$ with the algebra of symmetric Laurent polynomials $\mathbf{k}\left[x_{1}^{ \pm 1}, \ldots, x_{l}^{ \pm 1}\right]^{S_{l}}$. If $\lambda$ is a polynomial weight, the irreducible character $\chi_{\lambda}$ is equal to the Schur function $s_{\lambda}\left(x_{1}, \ldots, x_{l}\right)$. Given a virtual $\mathrm{GL}_{l}$ character $f(x)=$ $f\left(x_{1}, \ldots, x_{l}\right)=\sum_{\lambda} c_{\lambda} \chi_{\lambda}$, the partial sum over polynomial weights $\lambda$ is a symmetric polynomial in $l$ variables, which we denote by $f(x)_{\text {pol }}$ (this is different from the polynomial terms of $f(x)$ considered as a Laurent polynomial). We use the same notation for infinite formal sums $f(x)$ of irreducible $\mathrm{GL}_{l}$ characters, in which case $f(x)_{\text {pol }}$ is a symmetric formal power series.

The Weyl symmetrization operator for $\mathrm{GL}_{l}$ is

$$
\begin{equation*}
\boldsymbol{\sigma}\left(\phi\left(x_{1}, \ldots, x_{l}\right)\right)=\sum_{w \in S_{l}} w\left(\frac{\phi(x)}{\prod_{i<j}\left(1-x_{j} / x_{i}\right)}\right) \tag{37}
\end{equation*}
$$

For dominant weights $\lambda$, the Weyl character formula can be written $\chi_{\lambda}=\boldsymbol{\sigma}\left(x^{\lambda}\right)$. More generally, if $\phi(x)=\phi\left(x_{1}, \ldots, x_{l}\right)$ is a Laurent polynomial over any field $\mathbf{k}$, then $\boldsymbol{\sigma}(\phi(x))$ is a virtual $\mathrm{GL}_{l}$ character over $\mathbf{k}$.

The Hall-Littlewood symmetrization operator is defined by

$$
\begin{equation*}
\mathbf{H}_{q}^{l}(\phi(x))=\boldsymbol{\sigma}\left(\frac{\phi(x)}{\prod_{i<j}\left(1-q x_{i} / x_{j}\right)}\right) \tag{38}
\end{equation*}
$$

If $\phi(x)=\phi\left(x_{1}, \ldots, x_{l}\right)$ is a rational function over a field $\mathbf{k}$ containing $\mathbb{Q}(q)$, then $\mathbf{H}_{q}^{l}(\phi(x))$ is a symmetric rational function over $\mathbf{k}$. If $\phi(x)$ is a Laurent polynomial, we can also regard $\mathbf{H}_{q}^{l}(\phi(x))$ as an infinite formal sum of $\mathrm{GL}_{l}$ characters with coefficients in $\mathbf{k}$, by interpreting the factors $1 /\left(1-q x_{i} / x_{j}\right)$ as geometric series $1+q x_{i} / x_{j}+\left(q x_{i} / x_{j}\right)^{2}+\cdots$. We always understand $\mathbf{H}_{q}^{l}(\phi(x))$ in this sense when taking the polynomial part $\mathbf{H}_{q}^{l}(\phi(x))_{\text {pol }}$.

We also use the two-parameter symmetrization operator

$$
\begin{equation*}
\mathbf{H}_{q, t}^{l}(\phi(x))=\mathbf{H}_{q}^{l}\left(\phi(x) \prod_{i<j} \frac{\left(1-q t x_{i} / x_{j}\right)}{\left(1-t x_{i} / x_{j}\right)}\right)=\boldsymbol{\sigma}\left(\frac{\phi(x) \prod_{i<j}\left(1-q t x_{i} / x_{j}\right)}{\prod_{i<j}\left(\left(1-q x_{i} / x_{j}\right)\left(1-t x_{i} / x_{j}\right)\right)}\right) . \tag{39}
\end{equation*}
$$

Again, if $\phi(x)$ is a rational function over $\mathbf{k}=\mathbb{Q}(q, t)$, then $\mathbf{H}_{q, t}^{l}(\phi(x))$ is a symmetric rational function over $\mathbf{k}$, while if $\phi(x)$ is a Laurent polynomial, or more generally a Laurent polynomial times a rational function which has a power series expansion in the $x_{i} / x_{j}$ for $i<j$, we can also interpret $\mathbf{H}_{q, t}^{l}(\phi(x))$ as an infinite formal sum of $\mathrm{GL}_{l}$ characters, similarly to (39). This series interpretation always applies when taking $\mathbf{H}_{q, t}^{l}(\phi(x))_{\text {pol }}$.

Fixing $\mathbf{k}=\mathbb{Q}(q, t)$ once again, let $T=T\left(\mathbf{k}\left[z^{ \pm 1}\right]\right)$ be the tensor algebra on the Laurent polynomial ring in one variable, that is, the non-commutative polynomial algebra with generators corresponding to the basis elements $z^{a}$ of $\mathbf{k}\left[z^{ \pm 1}\right]$ as a vector space. Identifying $T^{m}=T^{m}\left(\mathbf{k}\left[z^{ \pm 1}\right]\right)$ with $\mathbf{k}\left[z_{1}^{ \pm 1}, \ldots, z_{m}^{ \pm 1}\right]$, the product in $T$ is given by 'concatenation,'

$$
\begin{equation*}
f \cdot g=f\left(z_{1}, \ldots, z_{k}\right) g\left(z_{k+1}, \ldots, z_{k+l}\right), \quad \text { for } f \in T^{k}, g \in T^{l} . \tag{40}
\end{equation*}
$$

The Feigin-Tsymbaliuk shuffle algebra [8] is the quotient $S=T / I$, where $I$ is the graded twosided ideal whose degree $l$ component $I^{l} \subseteq T^{l}$ is the kernel of the symmetrization operator $\mathbf{H}_{q, t}^{l}$ in variables $z_{1}, \ldots, z_{l}$, as explained further in [2, §3.5].

Let $\mathcal{E}^{+} \subseteq \mathcal{E}$ be the subalgebra generated by the $\Lambda_{\mathbf{k}}\left(X^{m, n}\right)$ for $m>0$. We leave out the central subalgebra $F$, since the relations of $\mathcal{E}^{+}$(as we will see in a moment) do not depend on the central parameters.

The image of $\mathcal{E}^{+}$under the anti-automorphism $\Phi$ in 3 is the subalgebra $\Phi\left(\mathcal{E}^{+}\right)$generated by the $\Lambda_{\mathbf{k}}\left(X^{m, n}\right)$ for $n>0$. Note that our convention $\theta \in(-\pi, \pi]$ when the subscript is omitted yields $\Phi\left(f\left(X^{m, n}\right)\right)=f\left(X^{n, m}\right)$ for $\Lambda_{\mathbf{k}}\left(X^{m, n}\right) \subseteq \mathcal{E}^{+}$, since the branch cut is in the third quadrant.

Schiffmann and Vasserot [25] proved the following result. See [2, §3.5] for more details on the translation of their theorem into our notation.

Proposition 3.3.1 ([25, Theorem 10.1]). There is an algebra isomorphism $\psi: S \rightarrow \mathcal{E}^{+}$and an anti-isomorphism $\psi^{\mathrm{op}}=\Phi \circ \psi: S \rightarrow \Phi\left(\mathcal{E}^{+}\right)$, given on the generators by $\psi\left(z^{a}\right)=p_{1}\left[-M X^{1, a}\right]$ and $\psi^{\mathrm{op}}\left(z^{a}\right)=p_{1}\left[-M X^{a, 1}\right]$.

To be clear, on monomials in $m$ variables, representing elements of tensor degree $m$ in $S$, the maps in Proposition 3.3.1 are given by

$$
\begin{align*}
\psi\left(z_{1}^{a_{1}} \cdots z_{m}^{a_{m}}\right) & =p_{1}\left[-M X^{1, a_{1}}\right] \cdots p_{1}\left[-M X^{1, a_{m}}\right]  \tag{41}\\
\psi^{\mathrm{op}}\left(z_{1}^{a_{1}} \cdots z_{m}^{a_{m}}\right) & =p_{1}\left[-M X^{a_{m}, 1}\right] \cdots p_{1}\left[-M X^{a_{1}, 1}\right] \tag{42}
\end{align*}
$$

Later we will need the following formula for the action of $\psi(\phi(z))$ on $\Lambda(X)$.
Proposition 3.3.2 ([2] Proposition 3.5.2]). Let $\phi(z)=\phi\left(z_{1}, \ldots, z_{l}\right)$ be a Laurent polynomial representing an element of tensor degree $l$ in $S$, and let $\zeta=\psi(\phi(z)) \in \mathcal{E}^{+}$be its image under the map in 41). With $\mathcal{E}$ acting on $\Lambda$ as in Proposition 3.2.1, we have

$$
\begin{equation*}
\omega(\zeta \cdot 1)\left(x_{1}, \ldots, x_{l}\right)=\mathbf{H}_{q, t}^{l}(\phi(x))_{\mathrm{pol}} . \tag{43}
\end{equation*}
$$

## 4. Schiffmann algebra reformulation of the symmetric function side

### 4.1. Distinguished elements $D_{\mathbf{b}}$ and $E_{\mathbf{a}}$

Negut [19] defined a family of distinguished elements $D_{\mathbf{b}} \in \mathcal{E}^{+}$, indexed by $\mathbf{b} \in \mathbb{Z}^{l}$, which in the case $l=1$ reduce to the elements in Proposition 3.2.3(i). Here a remarkable symmetry between these elements and their images $E_{\mathbf{a}}$ under the anti-involution $\Phi$ will play a crucial role. After defining the Negut elements, we derive this symmetry in Proposition 4.3 .3 with the help of a commutator formula of Negut [20].
Definition 4.1.1 (see also [2, §3.6]). Given $\mathbf{b}=\left(b_{1}, \ldots, b_{l}\right) \in \mathbb{Z}^{l}$, set

$$
\begin{equation*}
\phi(z)=\frac{z_{1}^{b_{1}} \cdots z_{l}^{b_{l}}}{\prod_{i=1}^{l-1}\left(1-q t z_{i} / z_{i+1}\right)} . \tag{44}
\end{equation*}
$$

and let $\nu(z)=\nu\left(z_{1}, \ldots, z_{l}\right)$ be a Laurent polynomial satisfying $\mathbf{H}_{q, t}^{l}(\nu(z))=\mathbf{H}_{q, t}^{l}(\phi(z))$. Such a $\nu(z)$ exists by [19, Proposition 6.1], and represents a well-defined element of the shuffle algebra $S$. The Negut element $D_{\mathbf{b}}$ and the transposed Negut element $E_{\mathbf{a}}$, where $\mathbf{a}=\left(b_{l}, \ldots, b_{1}\right)$ is the reversed sequence of indices, are defined by

$$
\begin{align*}
& D_{\mathbf{b}}=D_{b_{1}, \ldots, b_{l}}=\psi(\nu(z)) \in \mathcal{E}^{+}  \tag{45}\\
& E_{\mathbf{a}}=E_{b_{l}, \ldots, b_{1}}=\Phi\left(D_{\mathbf{b}}\right)=\psi^{\mathrm{op}}(\nu(z)) \in \Phi\left(\mathcal{E}^{+}\right) . \tag{46}
\end{align*}
$$

We should point out that, strictly speaking, the Negut elements in the case $l=1$ are defined to be elements $D_{a}=p_{1}\left[-M X^{1, a}\right]$ and $E_{a}=p_{1}\left[-M X^{a, 1}\right]$ of $\mathcal{E}$, while in Definition 3.2.2, we used the notation $D_{a}$ and $E_{a}$ for operators on $\Lambda$. However, by Proposition 3.2.3, these Negut elements act as the operators with the same name, so no confusion should ensue.

Later we will use the following product formulas, which are immediate from Definition 4.1.1.

$$
\begin{align*}
D_{b_{1}, \ldots, b_{l}} D_{b_{l+1}, \ldots, b_{n}} & =D_{b_{1}, \ldots, b_{n}}-q t D_{b_{1}, \ldots, b_{l}+1, b_{l+1}-1, \ldots, b_{n}}  \tag{47}\\
E_{a_{n}, \ldots, a_{l+1}} E_{a_{l}, \ldots, a_{1}} & =E_{a_{n}, \ldots, a_{1}}-q t E_{a_{n}, \ldots, a_{l+1}-1, a_{l}+1, \ldots, a_{1}} \tag{48}
\end{align*}
$$

As noted in $\$ 3.1 .5$, the internal action relations determine the action of $\Lambda_{\mathbf{k}}\left(X^{0,1}\right)$ on $\Phi\left(\mathcal{E}^{+}\right)$. Using the anti-isomorphism between $\Phi\left(\mathcal{E}^{+}\right)$and the shuffle algebra we can make this more explicit.
Lemma 4.1.2. Let $\phi(z)=\phi\left(z_{1}, \ldots, z_{n}\right)$ be a Laurent polynomial representing an element of tensor degree $n$ in $S$. Then

$$
\begin{equation*}
\left(\operatorname{Ad} f\left(X^{1,0}\right)\right) \psi^{\mathrm{op}}(\phi(z))=\psi^{\mathrm{op}}\left((\omega f)\left(z_{1}, \ldots, z_{n}\right) \cdot \phi(z)\right) . \tag{49}
\end{equation*}
$$

As a particular consequence, we have

$$
\begin{equation*}
\left(\operatorname{Ad} f\left(X^{1,0}\right)\right) E_{a_{n}, \ldots, a_{1}}=\psi^{\mathrm{op}}\left(\frac{(\omega f)\left(z_{1}, \ldots, z_{n}\right) \cdot z_{1}^{a_{1}} \cdots z_{n}^{a_{n}}}{\prod_{i=1}^{n-1}\left(1-q t z_{i} / z_{i+1}\right)}\right) \tag{50}
\end{equation*}
$$

Proof. This follows immediately from the rule in $\$ 3.1 .5$ for $\operatorname{Ad} f$ acting on a product.

### 4.2. Commutator identity

We use a formula for the commutator of elements $D_{a}$ and $D_{\mathbf{b}}$, and a similar identity for $E_{a}$ and $E_{\mathbf{b}}$. This commutation relation was proved geometrically by Negut in [20], but to keep
things self-contained, we provide an elementary algebraic proof. It is convenient to express the formula using the notation

$$
\sum_{i=a}^{b} \# f_{i}= \begin{cases}\sum_{i=a}^{b} f_{i} & \text { for } a \leq b+1  \tag{51}\\ -\sum_{i=b+1}^{a-1} f_{i} & \text { for } a \geq b+1\end{cases}
$$

As a mnemonic device, note that both cases can be interpreted as $\sum_{i=a}^{\infty} f_{i}-\sum_{i=b+1}^{\infty} f_{i}$.
Proposition 4.2.1 ([20, Proposition 4.7]). For any $a \in \mathbb{Z}$ and $\mathbf{b}=\left(b_{1}, \ldots, b_{l}\right) \in \mathbb{Z}^{l}$, we have

$$
\begin{align*}
{\left[D_{a}, D_{b_{1}, b_{2}, \ldots, b_{l}}\right] } & =-M \sum_{i=1}^{l} \sum_{k=a+1}^{b_{i}} \# D_{b_{1}, \ldots, b_{i-1}, k, a+b_{i}-k, b_{i+1}, \ldots, b_{l}}  \tag{52}\\
{\left[E_{b_{l}, \ldots, b_{2}, b_{1}}, E_{a}\right] } & =-M \sum_{i=1}^{l} \sum_{k=a+1}^{b_{i}} \# E_{b_{l}, \ldots, b_{i+1}, a+b_{i}-k, k, b_{i-1}, \ldots, b_{1}} . \tag{53}
\end{align*}
$$

We will need the following lemma for the proof. The notation $\Omega(X)$ is defined in 19). Since plethystic substitution into $\Omega(X)$ is characterized by

$$
\begin{equation*}
\Omega\left[a_{1}+a_{2}+\cdots-b_{1}-b_{2}-\cdots\right]=\frac{\prod_{i}\left(1-b_{i}\right)}{\prod_{i}\left(1-a_{i}\right)} \tag{54}
\end{equation*}
$$

we have

$$
\begin{equation*}
\Omega[M z]=\frac{(1-q z)(1-t z)}{(1-z)(1-q t z)} \quad \text { and } \quad \Omega[-M z]=\frac{(1-z)(1-q t z)}{(1-q z)(1-t z)} \tag{55}
\end{equation*}
$$

Lemma 4.2.2. For any $f(z)=f\left(z_{1}, \ldots, z_{m}\right)$ antisymmetric in $z_{i}$ and $z_{i+1}$, we have

$$
\begin{equation*}
\mathbf{H}_{q, t}^{m}\left(\Omega\left[M z_{i} / z_{i+1}\right] f(z)\right)=0 \tag{56}
\end{equation*}
$$

Proof. The definition of $\mathbf{H}_{q, t}^{m}$ and (55) imply that

$$
\begin{equation*}
\mathbf{H}_{q, t}^{m}\left(\Omega\left[M z_{i} / z_{i+1}\right] f(z)\right)=\sum_{w \in S_{m}} w\left(f(z) \prod_{j \neq k} \frac{1}{1-z_{j} / z_{k}} \prod_{\substack{j<k \\(j, k) \neq(i, i+1)}} \Omega\left[-M z_{j} / z_{k}\right]\right) \tag{57}
\end{equation*}
$$

which vanishes since $f(z)$ is antisymmetric in $z_{i}$ and $z_{i+1}$.
Proof of Proposition 4.2.1. Identity (53) for $\left[E_{b_{l}, \ldots, b_{1}}, E_{a}\right]$ follows from (52) by applying the anti-homomorphism $\Phi$, so we only prove 52, which can be written

$$
\begin{equation*}
D_{a} D_{\mathbf{b}}-D_{\mathbf{b}} D_{a}+M \sum_{i=1}^{l} \sum_{k=a+1}^{b_{i}} \# D_{b_{1}, \ldots, b_{i-1}, k, a+b_{i}-k, b_{i+1}, \ldots, b_{l}}=0 \tag{58}
\end{equation*}
$$

Using Definition 4.1.1 and the isomorphism $\psi: S \rightarrow \mathcal{E}^{+}$, we can prove 58) by showing that a rational function representing the left hand side is in the kernel of the symmetrization operator $\mathbf{H}_{q, t}^{l+1}$. For this we can work directly with the rational functions $\phi(z)$ in 44 ; there is no need to replace them explicitly with Laurent polynomials having the same symmetrization.

Let $\phi(z)$ be the function in 44 for $D_{\mathbf{b}}$, and set

$$
\begin{equation*}
\phi\left(\hat{z}_{i}\right)=\phi\left(z_{1}, \ldots, z_{i-1}, z_{i+1}, \ldots, z_{l+1}\right)=\frac{z_{1}^{b_{1}} \cdots z_{i-1}^{b_{i-1}} z_{i+1}^{b_{i}} \cdots z_{l+1}^{b_{l}}}{\left(1-q t z_{i-1} / z_{i+1}\right) \prod_{\substack{1 \leq j \leq l \\ j \neq i, i-1}}\left(1-q t z_{j} / z_{j+1}\right)} . \tag{59}
\end{equation*}
$$

To prove (58), we want to show

$$
\begin{equation*}
\mathbf{H}_{q, t}^{l+1}\left(z_{1}^{a} \phi\left(\hat{z_{1}}\right)-\phi\left(z_{l+1}\right) z_{l+1}^{a}+M \frac{\sum_{i=1}^{l} \sum_{k=a+1}^{b_{i}} \# z_{1}^{b_{1}} \cdots z_{i-1}^{b_{i-1}} z_{i}^{k} z_{i+1}^{a+b_{i}-k} z_{i+2}^{b_{i+1}} \cdots z_{l+1}^{b_{l}}}{\prod_{j=1}^{l}\left(1-q t z_{j} / z_{j+1}\right)}\right)=0 . \tag{60}
\end{equation*}
$$

Since $z_{i}^{a} \phi\left(\hat{z}_{i}\right)-\phi\left(\hat{z}_{i+1}\right) z_{i+1}^{a}$ is antisymmetric in $z_{i}$ and $z_{i+1}$, Lemma 4.2.2 implies

$$
\begin{equation*}
\sum_{i=1}^{l} \mathbf{H}_{q, t}^{l+1}\left(\Omega\left[M z_{i} / z_{i+1}\right]\left(z_{i}^{a} \phi\left(\hat{z}_{i}\right)-\phi\left(\hat{z}_{i+1}\right) z_{i+1}^{a}\right)\right)=0 \tag{61}
\end{equation*}
$$

The first formula in (55) is algebraically the same as

$$
\Omega[M z]=1-\frac{M}{\left(1-z^{-1}\right)(1-q t z)} .
$$

After substituting this into (61), the linearity of $\mathbf{H}_{q, t}^{l+1}$ gives

$$
\begin{equation*}
\mathbf{H}_{q, t}^{l+1}\left(\sum_{i=1}^{l}\left(z_{i}^{a} \phi\left(\hat{z}_{i}\right)-\phi\left(z_{i+1}\right) z_{i+1}^{a}-M \frac{z_{i}^{a} \phi\left(\hat{z_{i}}\right)-\phi\left(z_{i+1}\right) z_{i+1}^{a}}{\left(1-z_{i+1} / z_{i}\right)\left(1-q t z_{i} / z_{i+1}\right)}\right)\right)=0 \tag{62}
\end{equation*}
$$

The terms $z_{i}^{a} \phi\left(\hat{z_{i}}\right)-\phi\left(z_{i+1}\right) z_{i+1}^{a}$ telescope, reducing this to

$$
\begin{equation*}
\mathbf{H}_{q, t}^{l+1}\left(z_{1}^{a} \phi\left(\hat{z}_{1}\right)-\phi\left(z_{l+1}\right) z_{l+1}^{a}-M \sum_{i=1}^{l} \frac{z_{i}^{a} \phi\left(\hat{z_{i}}\right)-\phi\left(z_{i+1}\right) z_{i+1}^{a}}{\left(1-z_{i+1} / z_{i}\right)\left(1-q t z_{i} / z_{i+1}\right)}\right)=0 . \tag{63}
\end{equation*}
$$

If we use the convention $z_{0}=0$ and $z_{l+2}=\infty$, collecting terms in $z_{i}^{a} \phi\left(\hat{z}_{i}\right)$ and some further algebra manipulations give

$$
\begin{aligned}
& \sum_{i=1}^{l} \frac{z_{i}^{a} \phi\left(\hat{z}_{i}\right)-\phi\left(z_{i+1}\right) z_{i+1}^{a}}{\left(1-\frac{z_{i+1}}{z_{i}}\right)\left(1-q t \frac{z_{i}}{z_{i+1}}\right)}=\sum_{i=1}^{l+1}\left[\frac{1}{\left(1-\frac{z_{i+1}}{z_{i}}\right)\left(1-q t \frac{z_{i}}{z_{i+1}}\right)}-\frac{1}{\left(1-\frac{z_{i}}{z_{i-1}}\right)\left(1-q t \frac{z_{i-1}}{z_{i}}\right)}\right] z_{i}^{a} \phi\left(\hat{z_{i}}\right) \\
&=\sum_{i=1}^{l+1} \frac{z_{i}^{a} \phi\left(\hat{z}_{i}\right)\left(1-q t \frac{z_{i-1}}{z_{i+1}}\right)}{\left(1-q t \frac{z_{i-1}}{z_{i}}\right)\left(1-q t \frac{z_{i}}{z_{i+1}}\right)}\left(\frac{1}{\left.1-\frac{z_{i+1}}{z_{i}}-\frac{1}{1-\frac{z_{i}}{z_{i-1}}}\right)}\right. \\
&=\sum_{i=1}^{l+1} \frac{z_{i}^{a} \phi\left(\hat{z_{i}}\right)\left(1-q t \frac{z_{i-1}}{z_{i+1}}\right)}{\left(1-q t \frac{z_{i-1}}{z_{i}}\right)\left(1-q t \frac{z_{i}}{z_{i+1}}\right)}-\frac{z_{i+1}^{a} \phi\left(z_{i+1}\right)\left(1-q t \frac{z_{i}}{z_{i+2}}\right)}{\left(1-q t \frac{z_{i}}{z_{i+1}}\right)\left(1-q t \frac{z_{i+1}}{z_{i+2}}\right)} \\
& 1-\frac{z_{i+1}}{z_{i}}
\end{aligned}
$$

Expanding the definition 59) of $\phi\left(\hat{z}_{i}\right)$ for each $i$ yields

$$
\frac{z_{i}^{a} \phi\left(\hat{z_{i}}\right)\left(1-q t z_{i-1} / z_{i+1}\right)}{\left(1-q t z_{i-1} / z_{i}\right)\left(1-q t z_{i} / z_{i+1}\right)}=\frac{z_{1}^{b_{1}} \cdots z_{i-1}^{b_{i-1}} z_{i}^{a} z_{i+1}^{b_{i}} \cdots z_{l+1}^{b_{l}}}{\prod_{j=1}^{l}\left(1-q t z_{j} / z_{j+1}\right)},
$$

so that

$$
\begin{aligned}
\sum_{i=1}^{l} \frac{z_{i}^{a} \phi\left(\hat{z_{i}}\right)-\phi\left(z_{i+1}\right) z_{i+1}^{a}}{\left(1-z_{i+1} / z_{i}\right)\left(1-q t z_{i} / z_{i+1}\right)} & =\frac{\sum_{i=1}^{l} z_{1}^{b_{1}} \cdots z_{i-1}^{b_{i-1}} \cdot \frac{z_{i}^{a} z_{i+1}^{b_{i}}-z_{i}^{b_{i}} z_{i+1}^{a}}{1-z_{i+1} / z_{i}} \cdot z_{i+2}^{b_{i+1}} \cdots z_{l+1}^{b_{l}}}{\prod_{j=1}^{l}\left(1-q t z_{j} / z_{j+1}\right)} \\
& =\frac{-\sum_{i=1}^{l} z_{1}^{b_{1}} \cdots z_{i-1}^{b_{i-1}} \cdot\left(\sum_{k=a+1}^{b_{i}} \# z_{i}^{k} z_{i+1}^{a+b_{i}-k}\right) \cdot z_{i+2}^{b_{i+1}} \cdots z_{l+1}^{b_{l}}}{\prod_{j=1}^{l}\left(1-q t z_{j} / z_{j+1}\right)}
\end{aligned}
$$

Identity (60) follows by substituting this back into (63).

### 4.3. Symmetry identity for $D_{\mathbf{b}}$ and $E_{\mathbf{a}}$

Next we will prove an identity between certain instances of the Negut elements $D_{\mathbf{b}} \in \mathcal{E}^{+}$and transposed Negut elements $E_{\mathbf{a}} \in \Phi\left(\mathcal{E}^{+}\right)$. Before stating the identity we need to describe how the indices $\mathbf{a}$ and $\mathbf{b}$ will correspond.
Definition 4.3.1. A south-east lattice path $\gamma$ from $(0, n)$ to $(m, 0)$, for positive integers $m, n$, is admissible if it starts with a south step and ends with an east step; that is, $\gamma$ has a step from $(0, n)$ to $(0, n-1)$ and one from $(m-1,0)$ to $(m, 0)$. Define $\mathbf{b}(\gamma)=\left(b_{1}, \ldots, b_{m}\right)$ by taking $b_{i}=$ (vertical run of $\gamma$ at $x=i-1$ ) for $i=1, \ldots, m$ and $\mathbf{a}(\gamma)=\left(a_{n}, \ldots, a_{1}\right)$ with $a_{j}=($ horizontal run of $\gamma$ at $y=j-1)$ for $j=1, \ldots, n$. Set $D_{\gamma}=D_{\mathbf{b}(\gamma)}$ and $E_{\gamma}=E_{\mathbf{a}(\gamma)}$.

Note that if $\gamma^{*}$ is the transpose of an admissible path $\gamma$ with $\mathbf{b}(\gamma)=\left(b_{1}, \ldots, b_{m}\right)$ and $\mathbf{a}(\gamma)=$ $\left(a_{n}, \ldots, a_{1}\right)$, as above, then $\mathbf{a}\left(\gamma^{*}\right)=\left(b_{m}, \ldots, b_{1}\right)$ and $\mathbf{b}\left(\gamma^{*}\right)=\left(a_{1}, \ldots, a_{n}\right)$, and $E_{\gamma}=\Phi\left(D_{\gamma^{*}}\right)$. Example 4.3.2. Paths $\gamma$ and $\gamma^{*}$ below are both admissible. $\gamma$ is from $(0,8)$ to $(4,0)$ with $\mathbf{b}(\gamma)=(2,1,3,2)$ and $\mathbf{a}(\gamma)=(0,1,1,0,0,1,0,1)$, whereas $\gamma^{*}$ is from $(0,4)$ to $(8,0)$ and has $\mathbf{a}\left(\gamma^{*}\right)=(2,3,1,2)$ and $\mathbf{b}\left(\gamma^{*}\right)=(1,0,1,0,0,1,1,0)$.


Proposition 4.3.3. For every admissible path $\gamma$ we have $D_{\gamma}=E_{\gamma}$.
Proof. Let $\gamma$ be an admissible path $\gamma$ from $(0, n)$ to $(m, 0)$, where $m, n$ are positive integers.
We first establish the case when $n=1$. In this case, $E_{\gamma}=E_{m}=p_{1}\left[-M X^{m, 1}\right]$ and $D_{\gamma}=D_{10^{m-1}}$. If $m=1$, these are $E_{1}=D_{1}=p_{1}\left[-M X^{1,1}\right]$. In general, 24) implies $E_{m}=$ $p_{1}\left[-M X^{m, 1}\right]=\left(\operatorname{Ad} p_{1}\left(X^{1,0}\right)\right)^{m-1} p_{1}\left[-M X^{1,1}\right]=\left(\operatorname{Ad} p_{1}\left(X^{1,0}\right)\right)^{m-1} D_{1}$, while 17) and the commutator identity (52) imply $\left(\operatorname{Ad} p_{1}\left(X^{1,0}\right)\right) D_{10^{k}}=\left[p_{1}\left(X^{1,0}\right), D_{10^{k}}\right]=-(1 / M)\left[D_{0}, D_{10^{k}}\right]=$ $D_{10^{k+1}}$, and therefore $\left(\operatorname{Ad} p_{1}\left(X^{1,0}\right)\right)^{m-1} D_{1}=D_{10^{m-1}}$.

Using the involution $\Phi$, we can deduce the $m=1$ case from the $n=1$ case:

$$
\begin{equation*}
D_{\gamma}=D_{n}=\Phi\left(E_{n}\right)=\Phi\left(D_{1,0^{n-1}}\right)=E_{0^{n-1}, 1}=E_{\gamma} \tag{64}
\end{equation*}
$$

For $m, n>1$, we proceed by induction, assuming that the result holds for paths from $\left(0, n^{\prime}\right)$ to $\left(m^{\prime}, 0\right)$ when $m^{\prime} \leq m$ and $n^{\prime} \leq n$ and $\left(m^{\prime}, n^{\prime}\right) \neq(m, n)$.

For a given $m, n$, there are finitely many admissible paths $\gamma$, and thus a finite dimensional space $V$ of linear combinations $\sum_{\gamma} c_{\gamma} D_{\gamma}$ involving these paths. Let $V^{\prime} \subseteq V$ denote the subspace consisting of linear combinations which form the left hand side of a valid instance of the identity

$$
\begin{equation*}
\sum_{\gamma} c_{\gamma} D_{\gamma}=\sum_{\gamma} c_{\gamma} E_{\gamma} \tag{65}
\end{equation*}
$$

Note that $V^{\prime}=V$ if and only if $D_{\gamma}=E_{\gamma}$ for all the paths $\gamma$ in question.
We will use the induction hypothesis to construct enough instances of (65) to reduce each $D_{\gamma}$ modulo $V^{\prime}$ to a scalar multiple of $D_{\gamma_{0}}$, where $\gamma_{0}$ is the path with a south run from $(0, n)$ to $(0,0)$ followed by an east run to $(m, 0)$. We will then prove one more instance of (65) for which the left hand side reduces to a non-zero scalar multiple of $D_{\gamma_{0}}$, showing that $V^{\prime}=V$.

Suppose now that $\gamma \neq \gamma_{0}$. Then $\gamma$ contains an east step from $\left(m_{1}-1, n_{2}\right)$ to $\left(m_{1}, n_{2}\right)$ and a south step from $\left(m_{1}, n_{2}\right)$ to $\left(m_{1}, n_{2}-1\right)$ for some $m_{1}+m_{2}=m$ and $n_{1}+n_{2}=n$. In particular, $\gamma=\nu \cdot \eta$ for shorter admissible paths $\nu$ and $\eta$, where $\nu \cdot \eta$ is defined to be the lattice path obtained by placing $\nu$ and $\eta$ end to end; thus $\nu \cdot \eta$ traces a copy of $\nu$ from $\left(0, n_{1}+n_{2}\right)$ to $\left(m_{1}, n_{2}\right)$ and then traces a copy of $\eta$ from $\left(m_{1}, n_{2}\right)$ to $\left(m_{1}+m_{2}, 0\right)$.

Define $\gamma^{\prime}=\nu!^{\prime} \eta$ to be the admissible path obtained from $\nu \cdot \eta$ by replacing the east-south corner at $\left(m_{1}, n_{2}\right)$ with a south-east corner at $\left(m_{1}-1, n_{2}-1\right) ; \gamma^{\prime}$ contains a south step from $\left(m_{1}-1, n_{2}\right)$ to $\left(m_{1}-1, n_{2}-1\right)$ and an east step from $\left(m_{1}-1, n_{2}-1\right)$ to $\left(m_{1}, n_{2}-1\right)$.

The product formulas (47) and (48) imply that the elements corresponding to the paths constructed in this way satisfy

$$
\begin{equation*}
D_{\nu} D_{\eta}=D_{\nu \cdot \eta}-q t D_{\nu^{\prime} \eta} \quad \text { and } \quad E_{\nu} E_{\eta}=E_{\nu \cdot \eta}-q t E_{\nu^{\prime} \prime \eta} . \tag{66}
\end{equation*}
$$

By induction, $D_{\nu}=E_{\nu}$ and $D_{\eta}=E_{\eta}$, so (66) implies $D_{\gamma}-q t D_{\gamma^{\prime}}=E_{\gamma}-q t E_{\gamma^{\prime}}$. In other words, in terms of the space $V^{\prime}$ defined above, we have $D_{\gamma} \equiv q t D_{\gamma^{\prime}}\left(\bmod V^{\prime}\right)$. Using this repeatedly, we obtain $D_{\gamma} \equiv(q t)^{h(\gamma)} D_{\gamma_{0}}\left(\bmod V^{\prime}\right)$ for every path $\gamma$, where $h(\gamma)$ is the area enclosed by the path $\gamma$ and the $x$ and $y$ axes.

To complete the proof it suffices to establish one more identity of the form (65), for which the congruences $D_{\gamma} \equiv(q t)^{h(\gamma)} D_{\gamma_{0}}\left(\bmod V^{\prime}\right)$ reduce the left hand side to a non-zero scalar multiple of $D_{\gamma_{0}}$.

We can assume by induction that $D_{n, 0^{m-2}}=E_{0^{n-1}, m-1}$, since this case has the same $n$ and a smaller $m$. Taking the commutator with $p_{1}\left(X^{1,0}\right)$ on both sides gives

$$
\begin{equation*}
-\frac{1}{M}\left[D_{0}, D_{n, 0^{m-2}}\right]=\left[p_{1}\left(X^{1,0}\right), D_{n, 0^{m-2}}\right]=\left(\operatorname{Ad} p_{1}\left(X^{1,0}\right)\right) E_{0^{n-1}, m-1} \tag{67}
\end{equation*}
$$

Using (52) on the left hand side and (50) on the right hand side gives

$$
\begin{equation*}
\sum_{k=0}^{n-1} D_{\left(n-k, k, 0^{m-2}\right)}=\sum_{k=0}^{n-1} E_{\left(0^{n-1}, m-1\right)+\varepsilon_{n-k}} \tag{68}
\end{equation*}
$$

Now, for $1 \leq k \leq n-1$, we have $D_{\left(n-k, k, 0^{m-2}\right)}=D_{\gamma}$ and $E_{\left(0^{n-1}, m-1\right)+\varepsilon_{n-k}}=E_{\gamma}$ for an admissible path with $h(\gamma)=k$, as displayed below.


This shows that (68) is an instance of (65). The previous congruences reduce the left hand side of 68 to $\left(1+q t+\cdots+(q t)^{n-1}\right) D_{\gamma_{0}}$. Since the coefficient is non-zero, we have now established a set of instances of (65) whose left hand sides span $V$.

Corollary 4.3.4. For any indices $a_{1}, \ldots, a_{l}$, we have

$$
\begin{equation*}
E_{a_{l}, \ldots, a_{2}, a_{1}} \cdot 1=E_{a_{l}, \ldots, a_{2}, 0} \cdot 1 \tag{69}
\end{equation*}
$$

Proof. To rephrase, we are to show that $E_{a_{l}, \ldots, a_{2}, a_{1}} \cdot 1$ does not depend on $a_{1}$. The symmetry $f\left(X^{m, n}\right) \mapsto f\left(X^{m+r n, n}\right)$ of $\Phi\left(\mathcal{E}^{+}\right)$sends $E_{a_{l}, \ldots, a_{1}}$ to $E_{a_{l}+r, \ldots, a_{1}+r}$. By [2, Lemma 3.4.1], the action of $\mathcal{E}$ on $\Lambda$ satisfies $\nabla^{r} f\left(X^{m, n}\right) \nabla^{-r}=f\left(X^{m+r n, n}\right)$, and since $\nabla(1)=1$, this gives $\nabla^{r} E_{a_{l}, \ldots, a_{2}, a_{1}} \cdot 1=E_{a_{l}+r, \ldots, a_{2}+r, a_{1}+r} \cdot 1$. Hence, we can reduce to the case that $a_{i}>0$ for all $i$.

By [2, Lemma 3.6.2], we have that $D_{b_{1}, \ldots, b_{n}, 0, \ldots, 0} \cdot 1$ is independent of the number of trailing zeroes. In the case that $b_{i} \geq 0$ for all $i$ and $b_{1}>0$, this and Proposition 4.3.3 imply that $E_{a_{l}, \ldots, a_{1}} \cdot 1$ is independent of $a_{1}$, provided that $a_{i} \geq 0$ for all $i$ and $a_{1}>0$. However, we already saw that this suffices.

### 4.4. Shuffling the symmetric function side of the Extended Delta Conjecture

We can now give the promised reformulation of (7).
Theorem 4.4.1. For $0 \leq l<m \leq N$, we have

$$
\begin{equation*}
\left(\omega\left(h_{l}[B] e_{m-l-1}[B-1] e_{N-l}\right)\right)\left(x_{1}, \ldots, x_{m}\right)=\mathbf{H}_{q, t}^{m}(\phi(x))_{\mathrm{pol}}, \tag{70}
\end{equation*}
$$

where

$$
\begin{equation*}
\phi(x)=\frac{x_{1} \cdots x_{m}}{\prod_{i}\left(1-q t x_{i} / x_{i+1}\right)} h_{N-m}\left(x_{1}, \ldots, x_{m}\right) \overline{e_{l}\left(x_{2}, \ldots, x_{m}\right)}, \tag{71}
\end{equation*}
$$

and $\overline{e_{l}\left(x_{2}, \ldots, x_{m}\right)}=e_{l}\left(x_{2}^{-1}, \ldots, x_{m}^{-1}\right)$ by our convention on the use of the overbar.
Proof. For any symmetric function $f$ set $g(X)=(\omega f)[X+1 / M]$; then (31) gives an identity in $\Lambda$ for every $\zeta \in \mathcal{E}$

$$
\begin{equation*}
f[B] \zeta \cdot 1=g\left(X^{1,0}\right) \zeta \cdot 1=\sum\left(\left(\operatorname{Ad} g_{(1)}\left(X^{1,0}\right)\right) \zeta\right) g_{(2)}\left(X^{1,0}\right) \cdot 1 \tag{72}
\end{equation*}
$$

where $g[X+Y]=\sum g_{(1)}(X) g_{(2)}(Y)$ in Sweedler notation and we used the general formula $g \zeta=\sum\left(\left(\operatorname{Ad} g_{(1)}\right) \zeta\right) g_{(2)}$. Since $g[X+Y]=(\omega f)[X+Y+1 / M]$, and $h[B] \cdot 1=h[0] \cdot 1$ for any
$h(X)$, the right hand side of 72 is equal to

$$
\begin{align*}
& \sum\left(\left(\operatorname{Ad}(\omega f)_{(1)}\left(X^{1,0}\right)\right) \zeta\right)(\omega f)_{(2)}\left[X^{1,0}+1 / M\right] \cdot 1 \\
& \quad=\sum\left(\left(\operatorname{Ad}(\omega f)_{(1)}\left(X^{1,0}\right)\right) \zeta\right)(\omega f)_{(2)}[0] \cdot 1=\left(\left(\operatorname{Ad}(\omega f)\left(X^{1,0}\right)\right) \zeta\right) \cdot 1 \tag{73}
\end{align*}
$$

Let $n=N-l$. Taking $\zeta=E_{a_{n}, \ldots, a_{1}}$ and using (50), this gives

$$
\begin{equation*}
f[B] E_{a_{n}, \ldots, a_{1}} \cdot 1=f\left(z_{n}, \ldots, z_{1}\right) \mid z_{n}^{r_{n}} \cdots z_{1}^{r_{1}} \mapsto E_{a_{n}+r_{n}, \ldots, a_{2}+r_{2}, a_{1}+r_{1}} \cdot 1 \tag{74}
\end{equation*}
$$

By Corollary 4.3.4 the right hand side is a function of $f\left(z_{n}, \ldots, z_{2}, 1\right)$, since the substitution for the monomial $z^{\mathbf{r}}$ does not depend on the exponent $r_{1}$. Expressing $f\left(z_{n}, \ldots, z_{2}, 1\right)$ as $f\left[z_{n}+\right.$ $\left.\cdots+z_{2}+1\right]$ and then substituting $f[X-1]$ for $f(X)$ yields

$$
\begin{equation*}
f[B-1] E_{a_{n}, \ldots, a_{1}} \cdot 1=f\left[z_{n}+\cdots+z_{2}\right] \mid z_{n}^{r_{n}} \cdots z_{2}^{r_{2}} \mapsto E_{a_{n}+r_{n}, \ldots, a_{2}+r_{2}, a_{1}} \cdot 1 \tag{75}
\end{equation*}
$$

By [19, Proposition 6.7], $E_{0^{n}}=\Phi\left(D_{0^{n}}\right)=\Phi\left(e_{n}\left[-M X^{1,0}\right]\right)=e_{n}\left[-M X^{0,1}\right]$ (see also [2, Proposition 3.6.1]).

Using (75), we therefore obtain

$$
\begin{align*}
e_{k-1}[B-1] e_{n}=e_{k-1}\left[z_{n}+\cdots+z_{2}\right] \mid z_{n}^{r_{n}} \cdots z_{2}^{r_{2}} & \mapsto
\end{align*} E_{r_{n}, \ldots, r_{2}, 0} \cdot 1 .
$$

where the sum is over subsets $I \subseteq[n-1]$ and $\varepsilon_{I}=\sum_{i \in I} \varepsilon_{i}$. The terms in the last sum are just $E_{\mathbf{a}(\nu)} \cdot 1$ for paths $\nu$ from $(0, n)$ to $(k, 0)$ with single east steps on any $k-1$ chosen lines $y=j$ for $j \in[n-1]$, and a final east step at $y=0$. Denote the set of these admissible paths by $\mathcal{P}_{k, n}$. For instance, with $n=8$ and $k=4$, the path $\gamma$ in Example 4.3 .2 corresponds to $E_{\gamma}=E_{0,1,1,0,0,1,0,1}$.

By (74), applying $h_{l}[B]$ to (76) gives

$$
\begin{equation*}
h_{l}[B] e_{k-1}[B-1] e_{n}=\sum_{\nu \in \mathcal{P}_{k, n}} \sum_{\substack{\mathbf{r} \in \mathbb{N}^{n} \\|\mathbf{r}|=l}} E_{\mathbf{r}+\mathbf{a}(\nu)} \cdot 1 . \tag{77}
\end{equation*}
$$

This last expression is the sum of $E_{\gamma} \cdot 1$ over admissible paths $\gamma$ from $(0, n)$ to $(k+l, 0)$, together with a choice of $k-1$ indices $j \in[n-1]$ for which $\gamma$ has at least one east step on the line $y=j$. We can consider these indices as distinguishing $k-1$ east-south corners in $\gamma$. However, we can also distinguish these corners by their $x$ coordinates, that is, by a set of $k-1$ indices $i \in[k+l-1]$ for which $\gamma$ has at least one south step on the line $x=i$. Setting $m=k+l$ and using Proposition 4.3.3. this yields the identity

$$
\begin{equation*}
h_{l}[B] e_{m-l-1}[B-1] e_{n}=\sum_{\substack{\mathbf{s} \in \mathbb{N}^{m}:|\mathbf{s}|=n-k \\ I \subseteq[2, m],|I|=l}} D_{\mathbf{s}+\left(1^{m}\right)-\varepsilon_{I}} \cdot 1 . \tag{78}
\end{equation*}
$$

Now, since

$$
\begin{equation*}
\sum_{\substack{\mathbf{s} \in \mathbb{N}^{m}:|\mathbf{s}|=n-k \\ I \subseteq[2, m],|I|=l}} x^{\mathbf{s}+\left(1^{m}\right)-\varepsilon_{I}}=x_{1} x_{2} \cdots x_{m} h_{n-k}\left(x_{1}, \ldots, x_{m}\right) \overline{e_{l}\left(x_{2}, \ldots, x_{m}\right)}, \tag{79}
\end{equation*}
$$

the definition of $D_{\mathbf{b}}$ and Proposition 3.3.2 imply that

$$
\begin{equation*}
\omega\left(\sum_{\substack{\mathbf{s} \in \mathbb{N}^{m}:|\mathbf{s}|=n-k \\ I \subseteq[2, m],|I|=l}} D_{\mathbf{s}+\left(1^{m}\right)-\varepsilon_{I}} \cdot 1\right)\left(x_{1}, \ldots, x_{m}\right)=\mathbf{H}_{q, t}^{m}(\phi(x))_{\mathrm{pol}} \tag{80}
\end{equation*}
$$

with $\phi(x)$ given by 71).
Remark 4.4.2. For any $\mathbf{b} \in \mathbb{Z}^{m}$, [2, Corollary 3.7.2] gives that the Schur expansion of $\omega\left(D_{\mathbf{b}} \cdot 1\right)$ involves only $s_{\lambda}(X)$ with $\ell(\lambda) \leq m$. Hence, although Theorem 4.4.1 is a statement in $m$ variables, it determines $\omega\left(h_{l}[B] e_{m-l-1}[B-1] e_{N-l}\right)$ by 78$]$.

## 5. Reformulation of the combinatorial side

### 5.1. Statement of the reformulation

We reformulate (14) by explicitly extracting the coefficient of $z^{N-m}$. The most natural form of the resulting expression involves a generating function $N_{\beta / \alpha}$ for $q$-weighted tableaux rather than partially labelled paths. For now, we work only with the tableau description of $N_{\beta / \alpha}$, but in $\$ 6.2$ we will see that $N_{\beta / \alpha}$ is a truncation of an LLT series introduced by Grojnowski and Haiman in [12].

The $q$-weight in our reformulation involves two auxiliary statistics: for $\eta, \tau \in \mathbb{N}^{m}$, define

$$
\begin{equation*}
d(\eta, \tau)=\sum_{1 \leq j<r \leq m}\left|\left[\eta_{j}, \eta_{j}+\tau_{j}\right] \cap\left[\eta_{r}, \eta_{r}+\tau_{r}-1\right]\right|, \tag{81}
\end{equation*}
$$

with $[a, b]=\{a, \ldots, b\}$ and $[b]=[1, b]$, and for a vector $\eta$ of length $n$ and $I \subseteq[n]$, define

$$
\begin{equation*}
h_{I}(\eta)=\left|\left\{(r<s): r \in I, s \notin I, \eta_{s}=\eta_{r}+1\right\}\right|, \tag{82}
\end{equation*}
$$

where $(r<s)$ denotes a pair of positions $(r, s)$ in $\eta$ with $1 \leq r<s \leq n$.
Our reformulation of (14) is stated in the following theorem, proven at the end of this section.
Theorem 5.1.1. For $0 \leq l<m \leq N$, we have

$$
\begin{align*}
&\left\langle z^{N-m}\right\rangle \sum_{\substack{\lambda \in \mathbf{D}_{N} \\
P \in \in \mathbf{L}_{N, l}(\lambda)}} t^{|\delta / \lambda|} \prod_{\substack{1<i \leq N \\
c_{i}(\lambda)=c_{i-1}(\lambda)+1}}\left(1+z t^{-c_{i}(\lambda)}\right) q^{\operatorname{dinv}(P)} x^{\mathrm{wt}+(P)} \\
&=\sum_{\substack{J \subseteq[m-1] \\
|J|=l}} \sum_{\substack{\tau,(0, \mathbf{a}) \in \mathbb{N}^{m} \\
|\tau|=N-m}} t^{|\mathbf{a}|} q^{d((0, \mathbf{a}), \tau)+h_{J}(\mathbf{a})} N_{\left((0, \mathbf{a})+\left(1^{m}\right)+\tau\right) /\left((\mathbf{a}, 0)+\varepsilon_{J}\right)}(X ; q), \tag{83}
\end{align*}
$$

where $N_{\beta / \alpha}$ is given by Definition 5.2.1, below.

### 5.2. Definition of $N_{\beta / \alpha}$

For $\alpha, \beta \in \mathbb{Z}^{l}$ such that $\alpha_{j} \leq \beta_{j}$ for all $j$, define $\beta / \alpha$ to be the tuple of single row skew shapes $\left(\beta_{j}\right) /\left(\alpha_{j}\right)$ such that the $x$ coordinates of the right edges of boxes $a$ in the $j$-th row are the integers $\alpha_{j}+1, \ldots, \beta_{j}$. The boxes just outside the $j$-th row, adjacent to the left and right ends of the row, then have $x$ coordinates $\alpha_{j}$ and $\beta_{j}+1$. We consider these two boxes to be adjacent to the ends of an empty row, with $\alpha_{j}=\beta_{j}$, as well.

Given a tuple of skew row shapes $\beta / \alpha$, three boxes $(u, v, w)$ form a $w_{0}$-triple when box $v$ is in row $r$ of $\beta / \alpha$, boxes $u$ and $w$ are in or adjacent to a row $j$ with $j>r$, and the $x$-coordinates

$$
\begin{gathered}
-\infty \begin{array}{|c|c|c}
\hline 1|3| 4 \\
& -\infty \boxed{0} \infty \\
& -\infty \mid \infty \\
-\infty \boxed{3|5|} \infty \\
-\infty \mid \infty \\
S= & -\infty \mid \infty \\
-\infty \boxed{2} \infty \\
-\infty \boxed{1} \infty \\
-\infty & 0 \mid 6 \infty \\
-\infty \boxed{4} \infty \\
-\infty \mid \infty
\end{array}
\end{gathered}
$$

Figure 2. For $\beta=(12211123233)$, $\alpha=(11000121220)$, there are $h_{w_{0}}(\beta / \alpha)=29 w_{0}$-triples in $\beta / \alpha$. The row strict tableau $S$ of shape $\beta / \alpha$ has $h_{w_{0}}(S)=15$ increasing $w_{0}$-triples, $x^{\mathrm{wt}+(S)}=$ $x_{1}^{2} x_{2} x_{3}^{2} x_{4}^{2} x_{5} x_{6}$, and $x^{\mathrm{wt}(S)}=x_{0}^{2} x_{1}^{2} x_{2} x_{3}^{2} x_{4}^{2} x_{5} x_{6}$.
$i_{u}, i_{v}, i_{w}$ of these boxes satisfy $i_{u}=i_{v}$ and $i_{w}=i_{v}+1$. These triples are a special case of $\sigma$ triples defined for any $\sigma \in S_{l}$ in [2]. We denote the number of $w_{0}$-triples in $\beta / \alpha$ by $h_{w_{0}}(\beta / \alpha)$. The reader can verify that

$$
\begin{equation*}
h_{w_{0}}(\beta / \alpha)=\sum_{r<j}\left|\left[\alpha_{r}+1, \beta_{r}\right] \cap\left[\alpha_{j}, \beta_{j}\right]\right| . \tag{84}
\end{equation*}
$$

For a totally ordered alphabet $\mathcal{A}$, a row strict tableau of shape $\beta / \alpha$ is a map $S: \beta / \alpha \rightarrow \mathcal{A}$ that is strictly increasing on each row. The set of these maps is denoted by $\operatorname{RST}(\beta / \alpha, \mathcal{A})$. For convenience, given $\alpha, \beta \in \mathbb{Z}^{l}$ with some $\alpha_{j}>\beta_{j}$, we set $\operatorname{RST}(\beta / \alpha, \mathcal{A})=\varnothing$.

A $w_{0}$-triple $(u, v, w)$ is an increasing $w_{0}$-triple in $S$ if $S(u)<S(v)<S(w)$, with the convention that $S(u)=-\infty$ if $u$ is adjacent to the left end of a row of $\beta / \alpha$, and $S(w)=\infty$ if $w$ is adjacent to the right end of a row. Let $h_{w_{0}}(S)$ be the number of increasing $w_{0}$-triples in $S$.

For $S \in \operatorname{RST}(\beta / \alpha, \mathbb{N})$, define

$$
\begin{equation*}
x^{\mathrm{wt}+(S)}=\prod_{u \in \beta / \alpha, S(u) \neq 0} x_{S(u)} \quad \text { and } \quad x^{\mathrm{wt}(S)}=\prod_{u \in \beta / \alpha} x_{S(u)} . \tag{85}
\end{equation*}
$$

Definition 5.2.1. For $\alpha, \beta \in \mathbb{N}^{m}$, define

$$
\begin{equation*}
N_{\beta / \alpha}=N_{\beta / \alpha}(X ; q)=\sum_{S \in \operatorname{RST}\left(\beta / \alpha, \mathbb{Z}_{>0}\right)} q^{h_{w_{0}}(S)} x^{\mathrm{wt}(S)} . \tag{86}
\end{equation*}
$$

Note that if $\alpha_{j}>\beta_{j}$ for any $j$ then $N_{\beta / \alpha}=0$ by our convention that $\operatorname{RST}(\beta / \alpha, \mathcal{A})=\varnothing$.
Remark 5.2.2. It is shown in [2, Proposition 4.5.2] and its proof that, for $\alpha, \beta \in \mathbb{N}^{m}, \omega N_{\beta / \alpha}$ is a symmetric function whose Schur expansion involves only $s_{\lambda}$ where $\ell(\lambda) \leq m$.

### 5.3. Transforming the combinatorial side

To prove (83), we first associate each Dyck path with a tuple of row shapes recording vertical runs.

Definition 5.3.1. The LLT data associated to a path $\lambda \in \mathbf{D}_{N}$ is

$$
\beta=\left(1, c_{2}(\lambda)+1, \ldots, c_{N}(\lambda)+1\right) \text { and } \alpha=\left(c_{2}(\lambda), \ldots, c_{N}(\lambda), 0\right)
$$

where $c_{i}(\lambda)$ counts lattice squares between $\lambda$ and the line segment connecting $(0, N)$ to $(N, 0)$ in column $i$, numbered from right to left, as in Lemma 2.2.4.

Figure 2 shows the LLT data $\beta, \alpha$ associated to the path $\lambda$ in Figure 1 . Note that $\beta_{i}$ (resp. $\alpha_{i}$ ) is the furthest (resp. closest) distance from the diagonal to the path $\lambda$ on the line $x=N-i$, so that $\beta_{i}-\alpha_{i}$ is the number of south steps of $\lambda$ on that line.

This association allows us to relate $q$-weighted sums over partial labellings to the $N_{\beta / \alpha}$.
Lemma 5.3.2. For $\lambda \in \mathbf{D}_{N}$ and its associated LLT data $\alpha, \beta$, we have

$$
\begin{equation*}
\sum_{P \in \mathbf{L}_{N, l}(\lambda)} q^{\operatorname{dinv}(P)} x^{\mathrm{wt}(P)}=\sum_{\substack{\subseteq[N-1] \\|I|=l}} q^{h_{I}(\alpha)} N_{\beta /\left(\alpha+\varepsilon_{I}\right)}(X ; q) . \tag{87}
\end{equation*}
$$

Proof. There is a natural weight-preserving bijection mapping $P \in \mathbf{L}_{N}(\lambda)$ to $S \in \operatorname{RST}(\beta / \alpha, \mathbb{N})$, where the labels of column $x=i$ of $P$, read north to south, are placed into row $N-i$ of $\beta / \alpha$, west to east. See Figures 1 and 2. Moreover, $\operatorname{dinv}(P)=h_{w_{0}}(S)$. To see this, let $\hat{P}$ be the same labelling as $P$ but with the ordering on letters taken to be $0>1>2 \cdots$. It is proven in 2, Proposition 6.1.1] that $\operatorname{dinv}_{1}(\hat{P})=h_{w_{0}}(S)$, where $\operatorname{dinv}_{1}(\hat{P})$ was introduced in [13] and matches $\operatorname{dinv}(P)$ as discussed in Remark 2.2.3. The bijection restricts to a bijection from $\mathbf{L}_{N, l}(\lambda)$ to the subset of tableaux $S \in \operatorname{RST}(\beta / \alpha, \mathbb{N})$ with exactly $l 0$ 's, none in row $N$. This gives

$$
\begin{equation*}
\sum_{P \in \mathbf{L}_{N, l}(\lambda)} q^{\operatorname{dinv}(P)} x^{\mathrm{wt}_{+}(P)}=\sum_{\substack{I \subseteq[N-1] \\|I|=l}} \sum_{\substack{S \in \operatorname{RST}(\beta / \alpha, \mathbb{N}) \\ 0 \text { in rows } i \in I}} q^{h_{w_{0}}(S)} x^{\mathrm{wt}_{+}(S)} . \tag{88}
\end{equation*}
$$

The claim then follows from Definition 5.2.1 and the following Lemma.
Lemma 5.3.3. For $\alpha, \beta \in \mathbb{N}^{N}$ and $S \in \operatorname{RST}(\beta / \alpha, \mathbb{N})$, let $I \subseteq[N]$ be the rows of $S$ containing a zero and let $T$ be the tableau in $\operatorname{RST}\left(\beta /\left(\alpha+\varepsilon_{I}\right), \mathbb{Z}_{>0}\right)$ obtained by deleting all zeros from $S$. Then

$$
\begin{equation*}
h_{w_{0}}(T)=h_{w_{0}}(S)-h_{I}(\alpha), \tag{89}
\end{equation*}
$$

where $h_{I}(\alpha)$ is defined in 82).
Proof. Consider an increasing $w_{0}$-triple $(u, v, w)$ of $S$; the entries satisfy $S(u)<S(v)<S(w)$, $v$ lies in some row $r$, and both $u$ and $w$ lie in a row $j>r$. When $r \notin I$, either $j \notin I$ so that $(u, v, w)$ is an increasing $w_{0}$-triple of $T$ with the same entries as $S$, or $j \in I$ and $S(u)=0$ changes to $T(u)=-\infty$ where still $(u, v, w)$ is an increasing $w_{0}$-triple of $T$. However, if $r \in I$, $S(v)=0$ changes to $T(v)=-\infty$ and thus $(u, v, w)$ is not an increasing $w_{0}$-triple of $T$. Note the increasing condition implies that this happens only when $j \notin I$ and $\alpha_{r}=\alpha_{j}-1$ since $S(u)<0<S(w)$. Thus 89) follows.

Definition 5.3.4. Given $\mathbf{a}=\left(a_{1}, \ldots, a_{m-1}\right) \in \mathbb{N}^{m-1}$ and $\tau=\left(\tau_{1}, \ldots, \tau_{m}\right) \in \mathbb{N}^{m}$, we define two sequences $\beta_{\mathbf{a} \tau}$ and $\alpha_{\mathbf{a} \tau}$ of length $|\tau|+m$ as follows.

The sequence $\beta_{\mathbf{a} \tau}$ is the concatenation of sequences $\left(1,2, \ldots, \tau_{1}+1\right)$ and $\left(a_{i-1}+1, a_{i-1}+\right.$ $\left.2, \ldots, a_{i-1}+\tau_{i}+1\right)$ for $i=2, \ldots, m$. The sequence $\alpha_{\mathbf{a} \tau}$ is the same as $\beta_{\mathbf{a} \tau}$ except in the positions
corresponding to the ends of the concatenated subsequences. In these positions, we change the entries $\tau_{1}+1, a_{1}+\tau_{2}+1, \ldots, a_{m-1}+\tau_{m}+1$ in $\beta_{\mathbf{a} \tau}$ to $a_{1}, a_{2}, \ldots, a_{m-1}, 0$. Equivalently, $\alpha_{\mathbf{a} \tau}$ is the same as the sequence obtained by subtracting 1 from all entries of $\beta_{\mathbf{a} \tau}$ and shifting one place to the left, deleting the first entry and adding a zero at the end.

Example 5.3.5. For $\mathbf{a}=(130012)$ and $\tau=(2311022)$,

$$
\begin{align*}
& (0, \mathbf{a})+\left(1^{m}\right)+\tau=\left(\begin{array}{lllllllll}
3 & 5 & 5 & 2 & 1 & 4 & 5
\end{array}\right) \\
& \beta_{\mathbf{a} \tau}=\left(\begin{array}{llllllllll}
1 & 2 & 3 & 2 & 3 & 45 & 45 & 12 & 1 & 234 \\
\hline
\end{array}\right) \\
& \alpha_{\mathbf{a} \tau}=\left(\begin{array}{lllllllllll}
1 & 2 & 1 & 2 & 3 & 4 & 3 & 4 & 0 & 1 & 0 \\
1 & 2 & 3 & 3 & 4
\end{array}\right)  \tag{90}\\
& (\mathbf{a}, 0)=\left(\begin{array}{llllllll}
1 & 3 & 0 & 0 & 1 & 2 & 0
\end{array}\right)
\end{align*}
$$

The wider spaces show the division into blocks of size $\tau_{i}+1$. The last entry of $\alpha_{\mathbf{a} \tau}$ in each block is $a_{i}$, and the next block in $\alpha_{\mathbf{a} \tau}$ and $\beta_{\mathbf{a} \tau}$ starts with $a_{i}+1$.

Lemma 5.3.6. For $0 \leq l<m \leq N$,

$$
\begin{align*}
&\left\langle z^{N-m}\right\rangle \sum_{\substack{\lambda \in \mathbf{D}_{N} \\
P \in \mathbf{L}_{N, l}(\lambda)}} t^{|\delta / \lambda|} \prod_{\substack{1<i \leq N \\
c_{i}(\lambda)=c_{i-1}(\lambda)+1}}\left(1+z t^{-c_{i}(\lambda)}\right) q^{\operatorname{dinv}(P)} x^{\mathrm{wt}+(P)} \\
&=\sum_{\substack{I \subseteq[N-1] \\
|I|=l}} \sum_{\substack{\left(,(0, \mathbf{a}) \in \mathbb{N}^{m} \\
|\tau|=N-m\right.}} t^{|\mathbf{a}|} q^{h_{I}\left(\alpha_{\mathbf{a} \tau}\right)} N_{\beta_{\mathbf{a} \tau} /\left(\alpha_{\mathbf{a} \tau}+\varepsilon_{I}\right)}(X ; q) . \tag{91}
\end{align*}
$$

Proof. Use Lemma 5.3.2 to rewrite the left hand side of (91) as

$$
\begin{equation*}
\left\langle z^{N-m}\right\rangle \sum_{\lambda \in \mathbf{D}_{N}} t^{|\delta / \lambda|} \prod_{\substack{1<i \leq N \\ c_{i}(\lambda)=c_{i-1}(\lambda)+1}}\left(1+z t^{-c_{i}(\lambda)}\right) \sum_{\substack{I \subseteq[N-1] \\|I|=l}} q^{h_{I}(\alpha)} N_{\beta /\left(\alpha+\varepsilon_{I}\right)} \tag{92}
\end{equation*}
$$

where $\beta=\left(1^{N}\right)+\left(0, c_{2}(\lambda), \ldots, c_{N}(\lambda)\right), \alpha=\left(c_{2}(\lambda), \ldots, c_{N}(\lambda), 0\right)$ are the LLT data for $\lambda$. Note that a tuple $\mathbf{c}=\left(c_{1}, c_{2}, \ldots, c_{N}\right) \in \mathbb{N}^{N}$ is the sequence of column heights $c_{i}(\lambda)$ of a path $\lambda \in \mathbf{D}_{N}$ if and only if $c_{s} \leq c_{s-1}+1$ for all $s>1$ and $c_{1}=0$; in this case, $|\delta / \lambda|=|\mathbf{c}|$. Replace $\mathbf{D}_{N}$ in (92) by these tuples, and expand the product to obtain

$$
\begin{align*}
\left\langle z^{N-m}\right\rangle & \sum_{\substack{ \\
A \subseteq[N] \backslash\{1\}}} \sum_{\substack{c_{i} \leq c_{i-1}+1 \\
c_{i}=c_{i-1}+1 \\
\forall i \in A}} t^{|\mathbf{c}|-\sum_{i \in A} c_{i}} z^{|A|} \sum_{\substack{I \subseteq[N-1]}} q^{h_{I}(\alpha)} N_{\beta /\left(\alpha+\varepsilon_{I}\right)} \\
& =\sum_{\substack{|I|=l}} \sum_{\substack{\subseteq \subseteq J \subseteq[N] \\
|J|=m}} t^{\sum_{j \in J} c_{j}} \sum_{\substack{ \\
c_{j}-1+1}} \sum_{\substack{ \\
I \subseteq[N-1]}} q^{h_{I}(\alpha)} N_{\beta /\left(\alpha+\varepsilon_{I}\right)}, \tag{93}
\end{align*}
$$

where the equality comes from re-indexing with $J=[N] \backslash A$ and noting that we can drop the condition $c_{j} \leq c_{j-1}+1 \forall j \in J$ because $N_{\beta /\left(\alpha+\varepsilon_{I}\right)}=0$ if any $\left(\alpha+\varepsilon_{I}\right)_{j} \geq \alpha_{j}>\beta_{j}$.

If we replace the sum over $J$ by a sum over $\left\{\tau \in \mathbb{N}^{m}:|\tau|=N-m\right\}$ using $J=\left\{1, \tau_{1}+\right.$ $\left.2, \tau_{1}+\tau_{2}+3, \ldots, \tau_{1}+\cdots+\tau_{m-1}+m\right\}$, then, for fixed $J$ (or fixed $\tau$ ), the sum over can be replaced by a sum over

$$
\begin{equation*}
\mathbf{c}=\left(0,1,2, \ldots, \tau_{1}, a_{1}, a_{1}+1, \ldots, a_{1}+\tau_{2}, a_{2}, \ldots, a_{m-1}+\tau_{m}\right) \tag{94}
\end{equation*}
$$

for a ranging over $\mathbb{N}^{m-1}$. Note that $\sum_{i \in J} c_{j}=|\mathbf{a}|$. With this encoding of $\mathbf{c}$, we have $\beta / \alpha=$ $\beta_{\mathbf{a} \tau} / \alpha_{\mathbf{a} \tau}$ in the notation of Definition 5.3.4, and (93) becomes the right hand side of (91).


Figure 3. Comparing the tuples of rows $\beta_{\mathbf{a} \tau} / \alpha_{\mathbf{a} \tau}$ and $\left((0, \mathbf{a})+\left(1^{m}\right)+\tau\right) /(\mathbf{a}, 0)$ for $\mathbf{a} \in \mathbb{N}^{m-1}$ and $\tau \in \mathbb{N}^{m}$. Here $a_{j}=2, a_{r-1}=0, a_{r}=3$, and $\tau_{r}=5$.

We make a final adjustment to the right hand side of 91. This sum runs over tuples $\beta_{\mathbf{a} \tau} /\left(\alpha_{\mathbf{a} \tau}+\varepsilon_{I}\right)$ with $|\tau|$ necessarily empty rows which can be removed at the cost of a $q$ factor. We introduce some notation depending on a given $\mathbf{a} \in \mathbb{N}^{m-1}, \tau=\left(\tau_{1}, \ldots, \tau_{m}\right) \in \mathbb{N}^{m}$, and the associated $\beta_{\mathbf{a} \tau} / \alpha_{\mathbf{a} \tau}$ from Definition 5.3.4. First we set $j_{\uparrow}=j+\sum_{x \leq j} \tau_{x}$ for $j \in[m]$, so the entry of $\beta_{\mathbf{a} \tau}$ in position $j_{\uparrow}$ is $a_{j-1}+\tau_{j}+1$, or $\tau_{1}+1$ if $j=1$, and the entry of $\alpha_{\mathbf{a} \tau}$ in the same position is $a_{j}$, or 0 if $j=m$. For a subset $J \subseteq[m]$, we set $J_{\uparrow}=\left\{j_{\uparrow}: j \in J\right\}$. In positions $i \notin[m]_{\uparrow}$, the sequences $\beta_{\mathbf{a} \tau}$ and $\alpha_{\mathbf{a} \tau}$ agree, so row $i$ is empty in $\beta_{\mathbf{a} \tau} / \alpha_{\mathbf{a} \tau}$. The tuple of row shapes obtained by deleting these empty rows from $\beta_{\mathbf{a} \tau} / \alpha_{\mathbf{a} \tau}$ is $\left((0, \mathbf{a})+\left(1^{m}\right)+\tau\right) /(\mathbf{a}, 0)$, where row $j \in[m]$ corresponds to row $j_{\uparrow}$ of $\beta_{\mathbf{a} \tau} / \alpha_{\mathbf{a} \tau}$; note that rows $(j-1)_{\uparrow}$ and $j_{\uparrow}$ are separated by $\tau_{j}$ empty rows. See Figure 3 .
Lemma 5.3.7. For $J \subseteq[m], \mathbf{a} \in \mathbb{N}^{m-1}$ and $\tau \in \mathbb{N}^{m}$, let $I=J_{\uparrow}$. Then

$$
\begin{equation*}
N_{\beta_{\mathbf{a} \tau} /\left(\alpha_{\mathbf{a} \tau}+\varepsilon_{I}\right)}=q^{d((0, \mathbf{a}), \tau)-h_{J}^{\prime}(\mathbf{a}, \tau)} N_{\left((0, \mathbf{a})+\left(1^{m}\right)+\tau\right) /\left((\mathbf{a}, 0)+\varepsilon_{J}\right)}, \tag{95}
\end{equation*}
$$

where $h_{J}^{\prime}(\mathbf{a}, \tau)=\left|\left\{(j<r): j \in J, r \in[m], a_{j} \in\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right\}\right|$ with $a_{0}=0$, and $d((0, \mathbf{a}), \tau)$ is defined by (81).

Proof. Set $a_{0}=0$. We can assume $a_{j}+\left(\varepsilon_{J}\right)_{j} \leq a_{j-1}+\tau_{j}+1$ for all $j \in[m]$ since otherwise both sides of (95) vanish by Definition 5.2.1. Hence, each side is a $q$-generating function for row strict tableaux on tuples of single row skew shapes; rows of $\beta_{\mathbf{a} \tau} /\left(\alpha_{\mathbf{a} \tau}+\varepsilon_{I}\right)$ on the left hand side differ from the right hand side only by the removal of empty rows $r \notin[m]_{\uparrow}$. Thus, the two sides agree up to a factor $q^{d}$, where $d$ counts $w_{0}$-triples of $\beta_{\mathbf{a} \tau} /\left(\alpha_{\mathbf{a} \tau}+\varepsilon_{I}\right)$ involving one of these empty rows.

To evaluate $d$, consider such an empty row $(b) /(b)$, coming from $b \in\left\{a_{r-1}+1, \ldots, a_{r-1}+\tau_{r}\right\}$ for some $r \in[m]$. The adjacent boxes on the left and right of this empty row form a $w_{0}-$ triple, increasing in every tableau, with one box in each non-empty lower row $j_{\uparrow}$, of the form $\left(a_{j-1}+\tau_{j}+1\right) /\left(a_{j}+\left(\varepsilon_{J}\right)_{j}\right)$, such that $b \in\left[a_{j}+\left(\varepsilon_{J}\right)_{j}+1, a_{j-1}+\tau_{j}+1\right]$. Hence,

$$
\begin{aligned}
d & =\sum_{1 \leq j<r \leq m}\left|\left[a_{j}+\left(\varepsilon_{J}\right)_{j}, a_{j-1}+\tau_{j}\right] \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right| \\
& =\sum_{1 \leq j<r \leq m}\left|\left[a_{j}, a_{j-1}+\tau_{j}\right] \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|-\sum_{\substack{1 \leq j<r \leq m \\
j \in J}}\left|\left\{a_{j}\right\} \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right| .
\end{aligned}
$$

The sum after the minus sign is $h_{J}^{\prime}(\mathbf{a}, \tau)$. To prove that the remaining sum is $d((0, \mathbf{a}), \tau)$, first rewrite it as

$$
\begin{equation*}
\sum_{1 \leq j<r \leq m}\left(\left|\left[a_{j}, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|-\left|\left[a_{j-1}+\tau_{j}+1, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|\right) \tag{96}
\end{equation*}
$$

using the fact that $a_{j} \leq a_{j-1}+\tau_{j}+1$ by assumption. Next observe that since $a_{0}=0 \leq a_{r-1}$,

$$
\left|\left[a_{r-1}, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|=\left|\left[a_{0}, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|
$$

Adding $\sum_{1<j<r}\left|\left[a_{j-1}, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|$ to both sides, it follows that

$$
\sum_{1 \leq j<r}\left|\left[a_{j}, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|=\sum_{1 \leq j<r}\left|\left[a_{j-1}, \infty\right) \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right| .
$$

Hence (96] is unchanged upon replacing $\left[a_{j}, \infty\right)$ with $\left[a_{j-1}, \infty\right)$ and is thus equal to

$$
\sum_{1 \leq j<r \leq m}\left|\left[a_{j-1}, a_{j-1}+\tau_{j}\right] \cap\left[a_{r-1}, a_{r-1}+\tau_{r}-1\right]\right|=d((0, \mathbf{a}), \tau)
$$

Proof of Theorem 5.1.1. Consider a summand $t^{|\mathbf{a}|} q^{h_{I}\left(\alpha_{\mathbf{a} \tau}\right)} N_{\beta_{\mathbf{a} \tau} /\left(\alpha_{\mathbf{a} \tau}+\varepsilon_{I}\right)}$ on the right hand side of identity (91) for $I \subseteq[N-1], \mathbf{a} \in \mathbb{N}^{m-1}, \tau \in \mathbb{N}^{m}$. It vanishes unless $I=J_{\uparrow}$ for some $J \subseteq[m-$ 1], since $N_{\beta /\left(\alpha+\varepsilon_{I}\right)}=0$ when $\left(\alpha+\varepsilon_{I}\right)_{i}>\beta_{i}$ for some index $i$. For $I=J_{\uparrow}$, we can use Lemma 5.3.7 to replace this summand with $t^{|\mathbf{a}|} q^{d((0, \mathbf{a}), \tau)+h_{I}\left(\alpha_{\mathbf{a} \tau)}\right)-h_{J}^{\prime}(\mathbf{a}, \tau)} N_{\left((0, \mathbf{a})+\left(1^{m}\right)+\tau\right) /\left((\mathbf{a}, 0)+\varepsilon_{J}\right)}$.

It now suffices to prove that for $\alpha=\alpha_{\mathbf{a} \tau}$,

$$
\begin{equation*}
h_{I}(\alpha)=h_{J}^{\prime}(\mathbf{a}, \tau)+h_{J}(\mathbf{a}) . \tag{97}
\end{equation*}
$$

We recall that $N=m_{\uparrow}$ and note that $[N] \backslash I=\left([N] \backslash[m]_{\uparrow}\right) \sqcup\left([m]_{\uparrow} \backslash I\right)=\left([N] \backslash[m]_{\uparrow}\right) \sqcup([m] \backslash J)_{\uparrow}$. Hence, $h_{I}(\alpha)=\left|\left\{(x<y): x \in I, y \in[N] \backslash I, \alpha_{y}=\alpha_{x}+1\right\}\right|=\left|S_{1}\right|+\left|S_{2}\right|$ for

$$
\begin{aligned}
& S_{1}=\left\{(x<y): x \in J_{\uparrow}, y \in[N] \backslash[m]_{\uparrow}, \alpha_{y}=\alpha_{x}+1\right\}, \\
& S_{2}=\left\{(x<y): x \in J_{\uparrow}, y \in([m] \backslash J)_{\uparrow}, \alpha_{y}=\alpha_{x}+1\right\} .
\end{aligned}
$$

Since $\alpha_{m_{\uparrow}}=0$ implies $\left(x<m_{\uparrow}\right) \notin S_{2}$ for all $x<m_{\uparrow}$, we use that $a_{u}=\alpha_{u_{\uparrow}}$ for every $u \in[m-1]$ to see that

$$
\begin{equation*}
h_{J}(\mathbf{a})=\left|S_{2}\right|=\left|\left\{(j<r): j \in J, r \in[m-1] \backslash J, a_{r}=a_{j}+1\right\}\right| . \tag{98}
\end{equation*}
$$

Furthermore, $\left\{(j<r): j \in J, r \in[m], a_{r-1}+1 \leq a_{j}+1 \leq a_{r-1}+\tau_{r}\right\}$ and $S_{1}$ are equinumerous, as we can see by letting a pair $(j<r)$ in the first set correspond to the pair $\left(j_{\uparrow}<y\right)$ in $S_{1}$, where $y$ is the unique row index in the range $(r-1)_{\uparrow}<y<r_{\uparrow}$ such that $\alpha_{y}=\alpha_{j_{\uparrow}}+1=a_{j}+1$, as illustrated in Figure 3 .

## 6. Stable Unstraightened Extended Delta theorem

### 6.1. Overview

By Theorems 4.4.1 and 5.1.1, the Extended Delta Conjecture is equivalent to

$$
\begin{gather*}
\mathbf{H}_{q}^{m}\left(\frac{\prod_{i+1<j \leq m}\left(1-q t x_{i} / x_{j}\right)}{\prod_{i<j \leq m}\left(1-t x_{i} / x_{j}\right)} x_{1} \cdots x_{m} h_{N-m}\left(x_{1}, \ldots, x_{m}\right) \overline{e_{l}\left(x_{2}, \ldots, x_{m}\right)}\right)_{\mathrm{pol}} \\
=\sum_{\substack{J \subseteq[m-1] \\
|J|=l}} \sum_{\substack{0, \mathbf{a}), \tau \in \mathbb{N}^{m} \\
|\tau|=N-m}} t^{|\mathbf{a}|} q^{d((0, \mathbf{a}), \tau)+h_{J}(\mathbf{a})}\left(\omega N_{\beta / \alpha}\right)\left(x_{1}, \ldots, x_{m} ; q\right), \tag{99}
\end{gather*}
$$

where $\beta=(0, \mathbf{a})+\left(1^{m}\right)+\tau, \alpha=(\mathbf{a}, 0)+\varepsilon_{J}$, and $\left(\omega N_{\beta / \alpha}\right)\left(x_{1}, \ldots, x_{m} ; q\right)$ is $\omega N_{\beta / \alpha}(X ; q)$ evaluated in $m$ variables.

Although this is an identity in only $m$ variables, it does amount to the Extended Delta Conjecture by Remarks 4.4 .2 and 5.2 .2 both $\omega\left(h_{l}[B] e_{m-l-1}[B-1] e_{N-l}\right)$ and $\omega N_{\beta / \alpha}(X ; q)$ for the $\alpha, \beta$ arising in (99) are linear combinations of Schur functions $s_{\lambda}$ with $\ell(\lambda) \leq m$.

By Proposition 6.2.2 (below, proven in [2]), the functions $\omega N_{\beta / \alpha}$ on the right hand side of (99) are the polynomial parts of 'LLT series' introduced in [12], making each side of (99) the polynomial part of an infinite series of $\mathrm{GL}_{m}$ characters. We then prove 99 as a consequence of a stronger identity between these infinite series.

Hereafter, we use $x$ to abbreviate the alphabet $x_{1}, \ldots, x_{m}$.

### 6.2. LLT series

We will work with the (twisted) non-symmetric Hall-Littlewood polynomials as in [2]. For a $\mathrm{GL}_{m}$ weight $\lambda \in \mathbb{Z}^{m}$ and $\sigma \in S_{m}$, the twisted non-symmetric Hall-Littlewood polynomial $E_{\lambda}^{\sigma}(x ; q)$ is an element of $\mathbb{Z}\left[q^{ \pm 1}\right]\left[x_{1}^{ \pm 1}, \ldots, x_{m}^{ \pm 1}\right]$ defined using an action of the Hecke algebra on this ring; we refer the reader to [2, §4.3] for the precise definition, citing properties as needed. We also have their variants

$$
\begin{equation*}
F_{\lambda}^{\sigma}(x ; q)=\overline{E_{-\lambda}^{\sigma w_{0}}(x ; q)} \tag{100}
\end{equation*}
$$

recalling that $\overline{f\left(x_{1}, \ldots, x_{m} ; q\right)}=f\left(x_{1}^{-1}, \ldots, x_{m}^{-1} ; q^{-1}\right)$.
For any weights $\alpha, \beta \in \mathbb{Z}^{m}$ and a permutation $\sigma \in S_{m}$, the LLT series $\mathcal{L}_{\beta / \alpha}^{\sigma}(x ; q)=$ $\mathcal{L}_{\beta / \alpha}^{\sigma}\left(x_{1}, \ldots, x_{m} ; q\right)$ is defined in [2, §4.4] by

$$
\begin{equation*}
\left\langle\chi_{\lambda}\right\rangle \mathcal{L}_{\beta / \alpha}^{\sigma^{-1}}\left(x ; q^{-1}\right)=\left\langle E_{\beta}^{\sigma}\right\rangle \chi_{\lambda} \cdot E_{\alpha}^{\sigma} . \tag{101}
\end{equation*}
$$

Alternatively, [2, Proposition 4.4.2] gives the following expression in terms of the HallLittlewood symmetrization operator in (38):

$$
\begin{equation*}
\mathcal{L}_{\beta / \alpha}^{\sigma}(x ; q)=\mathbf{H}_{q}^{m}\left(w_{0}\left(F_{\beta}^{\sigma^{-1}}(x ; q) \overline{E_{\alpha}^{\sigma^{-1}}(x ; q)}\right)\right), \tag{102}
\end{equation*}
$$

where $w_{0}$ denotes the permutation of maximum length here and after. We will only need the LLT series for $\sigma=w_{0}$ and $\sigma=i d$, although most of what follows can be generalized to any $\sigma$.

In addition to the above formulas, we have the following combinatorial expressions for the polynomial truncations of LLT series as tableau generating functions with $q$ weights that count triples. As usual, a semistandard tableau on a tuple of skew row shapes $\nu=\beta / \alpha$ is a map $T: \nu \rightarrow[m]$ which is weakly increasing on rows. Let $\operatorname{SSYT}(\nu)$ denote the set of these, and define $x^{\mathrm{wt}(T)}=\prod_{b \in \nu} x_{T(b)}$.

Proposition 6.2.1 ([2, Remark 4.5.5 and Corollary 4.5.7]). If $\alpha_{i} \leq \beta_{i}$ for all $i$, then

$$
\begin{equation*}
\mathcal{L}_{\beta / \alpha}^{w_{0}}(x ; q)_{\mathrm{pol}}=\sum_{T \in \operatorname{SSYT}(\beta / \alpha)} q^{h_{w_{0}}^{\prime}(T)} x^{\mathrm{wt}(T)} \tag{103}
\end{equation*}
$$

where $h_{w_{0}}^{\prime}(T)$ is the number of $w_{0}$-triples $(u, v, w)$ of $\beta / \alpha$ such that $T(u) \leq T(v) \leq T(w)$.
Proposition 6.2.2 ([2, Proposition 4.5.2]). For any $\alpha, \beta \in \mathbb{Z}^{m}$,

$$
\begin{equation*}
\mathcal{L}_{\beta / \alpha}^{w_{0}}(x ; q)_{\mathrm{pol}}=\left(\omega N_{\beta / \alpha}\right)(x ; q) . \tag{104}
\end{equation*}
$$

### 6.3. Extended Delta Theorem

We now give several lemmas on non-symmetric Hall-Littlewood polynomials, then conclude by using the Cauchy formula for these polynomials to prove Theorem 6.3.6, below, yielding the stronger series identity that implies (99).

Lemma 6.3.1. For $\mathbf{a} \in \mathbb{N}^{m-1}$ and $w_{0} \in S_{m}$ and $\tilde{w}_{0} \in S_{m-1}$ the permutations of maximum length, we have

$$
\begin{align*}
& E_{(\mathbf{a}, 0)}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right)=E_{\mathbf{a}}^{\tilde{\mathbf{w}}_{0}}\left(x_{1}, \ldots, x_{m-1} ; q\right)  \tag{105}\\
& F_{(0, \mathbf{a})}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right)=F_{\mathbf{a}}^{\tilde{w}_{0}}\left(x_{2}, \ldots, x_{m} ; q\right) . \tag{106}
\end{align*}
$$

Proof. By [2, Lemma 4.3.4], we have $E_{(\mathbf{a}, 0)}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right)=E_{\mathbf{a}}^{\tilde{w}_{0}}\left(x_{1}, \ldots, x_{m-1} ; q\right) E_{(0)}^{i d}\left(x_{m} ; q\right)$ and $E_{(0,-\mathbf{a})}^{i d}\left(x_{1}, \ldots, x_{m} ; q\right)=E_{(0)}^{i d}\left(x_{1} ; q\right) E_{-\mathbf{a}}^{i d}\left(x_{2}, \ldots, x_{m} ; q\right)$. The claim then follows from the definition $F_{\mathbf{a}}^{\sigma}=\overline{E_{-\mathbf{a}}^{w_{0} \sigma}}$ and noting that $E_{(0)}^{i d}\left(x_{m} ; q\right)=1=F_{(0)}^{i d}\left(x_{1} ; q\right)$.

Inverting all variables and specializing $\sigma=w_{0}$ in [2, Lemma 4.5.1] yields the following lemma.

Lemma 6.3.2. For $l \leq m$, $\mathbf{a} \in \mathbb{Z}^{m}$, we have

$$
\begin{equation*}
\overline{e_{l}(x)} \overline{E_{\mathbf{a}}^{w_{0}}(x ; q)}=\sum_{I \subseteq[m]:|I|=l} q^{h_{I}(\mathbf{a})} \overline{E_{\mathbf{a}+\varepsilon_{I}}^{w_{0}}(x ; q)}, \tag{107}
\end{equation*}
$$

where $h_{I}(\mathbf{a})=\left|\left\{(i<j) \mid a_{j}=a_{i}+1, i \in I, j \notin I\right\}\right|$, as defined in 82).
Lemma 6.3.3. For every $\lambda \in \mathbb{Z}^{m}$ and $\sigma \in S_{m}$, we have

$$
\begin{equation*}
F_{\lambda}^{\sigma}(x ; q)=w_{0} E_{w_{0} \lambda}^{w_{0} \sigma}\left(x ; q^{-1}\right) \tag{108}
\end{equation*}
$$

Proof. The desired identity follows from

$$
\begin{equation*}
w_{0} E_{\lambda}^{\sigma}\left(x_{1}^{-1}, \ldots, x_{m}^{-1} ; q\right)=E_{-w_{0} \lambda}^{w_{0} \sigma w_{0}}(x ; q) \tag{109}
\end{equation*}
$$

by applying $w_{0}$ to both sides, substituting $\sigma \mapsto \sigma w_{0}, \lambda \mapsto-\lambda$, and $q \mapsto q^{-1}$, and using the definition of $F_{\lambda}^{\sigma}$.

To prove 109, we use the characterization of $E_{\lambda}^{\sigma}(x ; q)$ by the recurrence [2, (77)] and initial condition $E_{\lambda}^{\sigma}=x^{\lambda}$ for $\lambda$ dominant. The change of variables $x^{\mu} \mapsto x^{-w_{0}(\mu)}$ replaces the Hecke algebra operator $T_{i}=T_{s_{i}}$ in the recurrence with $T_{w_{0} s_{i} w_{0}}$, giving a modified recurrence satisfied by the left hand side of (109). It is straightforward to verify that the right hand side of 109 , satisfies the same modified recurrence. Since both sides reduce to $x^{-w_{0}(\lambda)}$ for $\lambda$ dominant, 109 holds.

Lemma 6.3.4. Given $\alpha, \beta \in \mathbb{Z}^{m}$ and a symmetric Laurent polynomial $f\left(x_{1}, \ldots, x_{m}\right)$, we have, for any $\sigma \in S_{m}$,

$$
\begin{equation*}
\left\langle E_{w_{0} \beta}^{w_{0} \sigma w_{0}}\left(x ; q^{-1}\right)\right\rangle f(x) \cdot E_{w_{0} \alpha}^{w_{0} \sigma w_{0}}\left(x ; q^{-1}\right)=\left\langle F_{-\alpha}^{\sigma}(x ; q)\right\rangle f(x) \cdot F_{-\beta}^{\sigma}(x ; q) . \tag{110}
\end{equation*}
$$

Proof. In fact, we will show that

$$
\begin{equation*}
\left\langle E_{w_{0} \beta}^{w_{0} \sigma w_{0}}\left(x ; q^{-1}\right)\right\rangle f(x) \cdot E_{w_{0} \alpha}^{w_{0} \sigma w_{0}}\left(x ; q^{-1}\right)=\left\langle F_{-\alpha}^{\sigma}(x ; q)\right\rangle w_{0}(f(x)) \cdot F_{-\beta}^{\sigma}(x ; q), \tag{111}
\end{equation*}
$$

even if we do not assume that $f(x)$ is symmetric. By Lemma 6.3.3, the right hand side of 111 is equal to

$$
\begin{equation*}
\left\langle E_{-w_{0} \alpha}^{w_{0} \sigma}\left(x ; q^{-1}\right)\right\rangle f(x) \cdot E_{-w_{0} \beta}^{w_{0} \sigma}\left(x ; q^{-1}\right) . \tag{112}
\end{equation*}
$$

By [2, Proposition 4.3.2], the functions $E_{\lambda}^{\sigma}(x ; q)$ and $E_{-\lambda}^{\sigma w_{0}}(x ; q)$ are dual bases with respect to to the inner product $\langle-,-\rangle_{q}$ defined there. Moreover, it is immediate from the construction of the inner product that multiplication by any $f(x)$ is self-adjoint. This gives

$$
\begin{equation*}
\left\langle f(x) E_{w_{0} \alpha}^{w_{0} \sigma w_{0}}\left(x ; q^{-1}\right), E_{-w_{0} \beta}^{w_{0} \sigma}\left(x ; q^{-1}\right)\right\rangle_{q^{-1}}=\left\langle E_{w_{0} \alpha}^{w_{0} \sigma w_{0}}\left(x ; q^{-1}\right), f(x) E_{-w_{0} \beta}^{w_{0} \sigma}\left(x ; q^{-1}\right)\right\rangle_{q^{-1}} \tag{113}
\end{equation*}
$$

in which the left hand side is equal to the left hand side of 111), and the right hand side is equal to 112 .

Lemma 6.3.5. For $w_{0}$ the maximum length permutation in $S_{m}$ and $\eta \in \mathbb{N}^{m}$, we have

$$
\begin{equation*}
h_{l}(x) F_{\eta}^{w_{0}}(x ; q)=\sum_{\substack{\tau \in \mathbb{N}^{m} \\|\tau|=l}} q^{d(\eta, \tau)} F_{\eta+\tau}^{w_{0}}(x ; q), \tag{114}
\end{equation*}
$$

recalling from (81) that $d(\eta, \tau)=\sum_{j<r}\left|\left[\eta_{j}, \eta_{j}+\tau_{j}\right] \cap\left[\eta_{r}, \eta_{r}+\tau_{r}-1\right]\right|$.
Proof. Set $\alpha=-\eta-\tau$ and $\beta=-\eta$. By 101) and Lemma 6.3.4 (with $\sigma=w_{0}$ ), we have

$$
\begin{equation*}
\left\langle h_{l}(x)\right\rangle \mathcal{L}_{w_{0}(\beta / \alpha)}^{w_{0}}(x ; q)=\left\langle E_{w_{0} \beta}^{w_{0}}\left(x ; q^{-1}\right)\right\rangle h_{l}(x) E_{w_{0} \alpha}^{w_{0}}\left(x ; q^{-1}\right)=\left\langle F_{-\alpha}^{w_{0}}(x ; q)\right\rangle h_{l}(x) F_{-\beta}^{w_{0}}(x ; q) . \tag{115}
\end{equation*}
$$

By specializing all but one variable in (103) to zero, Proposition 6.2.1 implies that the coefficient of $h_{l}$ in $\mathcal{L}_{w_{0}(\beta / \alpha)}^{w_{0}}(x ; q)_{\text {pol }}$ is $q^{h_{w_{0}}^{\prime}(T)}$ for $T$ the semistandard tableau of shape $w_{0}(\beta / \alpha)$ filled with a single letter, where $h_{w_{0}}^{\prime}(T)$ is the number of $w_{0}$-triples of $w_{0}(\beta / \alpha)=w_{0}(-\eta /(-\eta-\tau))$. By (84), this number is $d(\eta, \tau)$.

Theorem 6.3.6. For $0 \leq l<m \leq N$ and $w_{0} \in S_{m}$ the maximum length permutation, we have

$$
\begin{aligned}
& \frac{\prod_{i+1<j \leq m}\left(1-q t x_{i} / x_{j}\right)}{\prod_{i<j \leq m}\left(1-t x_{i} / x_{j}\right)} x_{1} \cdots x_{m} h_{N-m}\left(x_{1}, \ldots, x_{m}\right) \overline{e_{l}\left(x_{2}, \ldots, x_{m}\right)} \\
& \quad=\sum_{\substack{(0, \mathbf{a}), \tau \in \mathbb{N}^{m} \\
I \subseteq[m-1] \\
|\tau|=N-m,|I|=l}} t^{|\mathbf{a}|} q^{d((0, \mathbf{a}), \tau)+h_{I}(\mathbf{a})} w_{0}\left(F_{(0, \mathbf{a})+\tau+\left(1^{m}\right)}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right) \overline{E_{(\mathbf{a}, 0)+\varepsilon_{I}}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right)}\right) .
\end{aligned}
$$

Proof. Our starting point is the Cauchy formula [2, Theorem 5.1.1] for the twisted nonsymmetric Hall-Littlewood polynomials associated to any $\tilde{\sigma} \in S_{m-1}$ :

$$
\begin{equation*}
\frac{\prod_{i<j<m}\left(1-q t x_{i} y_{j}\right)}{\prod_{i \leq j<m}\left(1-t x_{i} y_{j}\right)}=\sum_{\mathbf{a} \in \mathbb{N}^{m}-1} t^{|\mathbf{a}|} E_{\mathbf{a}}^{\tilde{\sigma}}\left(x_{1}, \ldots, x_{m-1} ; q^{-1}\right) F_{\mathbf{a}}^{\tilde{\sigma}}\left(y_{1}, \ldots, y_{m-1} ; q\right) \tag{116}
\end{equation*}
$$

Take $\tilde{\sigma}=\tilde{w}_{0}$ the maximum length permutation in $S_{m-1}$, replace $x_{i}$ by $x_{i}^{-1}$, and then let $y_{j}=x_{j+1}$ to get

$$
\begin{equation*}
\frac{\prod_{i+1<j \leq m}\left(1-q t x_{j} / x_{i}\right)}{\prod_{i<j \leq m}\left(1-t x_{j} / x_{i}\right)}=\sum_{\mathbf{a} \in \mathbb{N}^{m-1}} t^{|\mathbf{a}|} F_{\mathbf{a}}^{\tilde{w}_{0}}\left(x_{2}, \ldots, x_{m} ; q\right) \overline{E_{\mathbf{a}}^{\tilde{w}_{0}}\left(x_{1}, \ldots, x_{m-1} ; q\right)} \tag{117}
\end{equation*}
$$

By (106) and the definition of $F^{\sigma}$,

$$
\left(x_{1} \cdots x_{m}\right) F_{\mathbf{a}}^{\tilde{w}_{0}}\left(x_{2}, \ldots, x_{m} ; q\right)=\left(x_{1} \cdots x_{m}\right) F_{(0, \mathbf{a})}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right)=F_{(0, \mathbf{a})+\left(1^{m}\right)}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right)
$$

for $w_{0} \in S_{m}$. Hence,

$$
\begin{aligned}
& \frac{\prod_{i+1<j \leq m}\left(1-q t x_{j} / x_{i}\right)}{\prod_{i<j \leq m}\left(1-t x_{j} / x_{i}\right)}\left(x_{1} \cdots x_{m}\right) \\
&=\sum_{\mathbf{a} \in \mathbb{N}^{m}-1} t^{|\mathbf{a}|} F_{(0, \mathbf{a})+\left(1^{m}\right)}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right) \overline{E_{\mathbf{a}}^{\tilde{\tilde{w}}_{0}}\left(x_{1}, \ldots, x_{m-1} ; q\right)}
\end{aligned}
$$

Multiplying by $h_{N-m}\left(x_{1}, \ldots, x_{m}\right)$ with the help of Lemma 6.3 .5 yields

$$
\begin{aligned}
& \frac{\prod_{i+1<j \leq m}\left(1-q t x_{j} / x_{i}\right)}{\prod_{i<j \leq m}\left(1-t x_{j} / x_{i}\right)}\left(x_{1} \cdots x_{m}\right) h_{N-m}\left(x_{1}, \ldots, x_{m}\right) \\
&=\sum_{\substack{\left(0, \mathbf{a}, \tau \in \mathbb{N}^{m} \\
|\tau|=N-m\right.}} t^{|\mathbf{a}|} q^{d((0, \mathbf{a}), \tau)} F_{\eta+\tau}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right) \overline{E_{\mathbf{a}}^{\tilde{\mathbf{a}}_{0}}\left(x_{1}, \ldots, x_{m-1} ; q\right)}
\end{aligned}
$$

where $\eta=\left(1^{m}\right)+(0, \mathbf{a})$ and we have used that $d(\eta, \tau)=d((0, \mathbf{a}), \tau)$ by 81$)$. Now multiply by $\overline{e_{l}\left(x_{1}, \ldots, x_{m-1}\right)}$ and apply 107) to get

$$
\begin{align*}
& \frac{\prod_{i+1<j \leq m}\left(1-q t x_{j} / x_{i}\right)}{\prod_{i<j \leq m}\left(1-t x_{j} / x_{i}\right)}\left(x_{1} \cdots x_{m}\right) \overline{e_{l}\left(x_{1}, \ldots, x_{m-1}\right)} h_{N-m}\left(x_{1}, \ldots, x_{m}\right) \\
& \quad=\sum_{\substack{(0, \mathbf{a}), \tau \in \mathbb{N}^{m} \\
|\tau|=N-m}} \sum_{|I|=l} t^{|\mathbf{a}|} q^{d((0, \mathbf{a}), \tau)+h_{I}(\mathbf{a})} F_{\eta+\tau}^{w_{0}}\left(x_{1}, \ldots, x_{m} ; q\right) \overline{E_{\mathbf{a}+\varepsilon_{I}}^{\tilde{w}_{0}}\left(x_{1}, \ldots, x_{m-1} ; q\right)} \tag{118}
\end{align*}
$$

where $I \subseteq[m-1]$. The result then follows by using 105 on the right hand side and applying $w_{0} \in S_{m}$ to both sides, noting that $w_{0}\left(\overline{\left.e_{l}\left(x_{2}, \ldots, x_{m}\right)\right)}=\overline{e_{l}\left(x_{1}, \ldots, x_{m-1}\right)}\right.$.

Proof of the Extended Delta Conjecture. It suffices to prove the reformulation in 99p; this follows by applying $\mathbf{H}_{q}^{m}$ and $\sqrt{102}$ to the identity of Theorem 6.3.6. taking the polynomial part, and using Proposition 6.2.2.

## References

[1] F. Bergeron, A. M. Garsia, M. Haiman, and G. Tesler, Identities and positivity conjectures for some remarkable operators in the theory of symmetric functions, Methods Appl. Anal. 6 (1999), no. 3, 363-420, Dedicated to Richard A. Askey on the occasion of his 65 th birthday, Part III.
[2] Jonah Blasiak, Mark Haiman, Jennifer Morse, Anna Pun, and George H. Seelinger, A Shuffle Theorem for paths under any line, Forum of Mathematics, Pi (????), ??-??, arXiv:math.CO/2102.07931.
[3] Igor Burban and Olivier Schiffmann, On the Hall algebra of an elliptic curve, I, Duke Math. J. 161 (2012), no. 7, 1171-1231.
[4] Erik Carlsson and Anton Mellit, A proof of the shuffle conjecture, J. Amer. Math. Soc. 31 (2018), no. 3, 661-697.
[5] Michele D'Adderio, Alessandro Iraci, and Anna Vanden Wyngaerd, The Schröder case of the generalized Delta conjecture, European J. Combin. 81 (2019), 58-83.
[6] _, Theta operators, refined delta conjectures, and coinvariants, 2019, arXiv:math.CO/1906.02623.
[7] Michele D'Adderio and Anton Mellit, A proof of the compositional Delta conjecture, 2020, arXiv:math.CO/2011.11467.
[8] B. L. Feigin and A. I. Tsymbaliuk, Equivariant K-theory of Hilbert schemes via shuffle algebra, Kyoto J. Math. 51 (2011), no. 4, 831-854.
[9] A. M. Garsia and M. Haiman, Some natural bigraded $S_{n}$-modules and $q, t$-Kostka coefficients, Electron. J. Combin. 3 (1996), no. 2, Research Paper 24, approx. 60 pp. (electronic), The Foata Festschrift.
[10] A. M. Garsia, M. Haiman, and G. Tesler, Explicit plethystic formulas for Macdonald q, tKostka coefficients, Sém. Lothar. Combin. 42 (1999), Art. B42m, 45 pp. (electronic), The Andrews Festschrift (Maratea, 1998).
[11] Adriano Garsia, Jim Haglund, Jeffrey B. Remmel, and Meesue Yoo, A proof of the delta conjecture when $q=0$, Ann. Comb. 23 (2019), no. 2, 317-333.
[12] I. Grojnowski and M. Haiman, Affine Hecke algebras and positivity of LLT and Macdonald polynomials, Unpublished manuscript, 2007.
[13] J. Haglund, M. Haiman, N. Loehr, J. B. Remmel, and A. Ulyanov, A combinatorial formula for the character of the diagonal coinvariants, Duke Math. J. 126 (2005), no. 2, 195-232.
[14] J. Haglund, J. B. Remmel, and A. T. Wilson, The delta conjecture, Trans. Amer. Math. Soc. 370 (2018), no. 6, 4029-4057.
[15] James Haglund, Brendon Rhoades, and Mark Shimozono, Ordered set partitions, generalized coinvariant algebras, and the delta conjecture, Adv. Math. 329 (2018), 851-915.
[16] ——Hall-Littlewood expansions of Schur delta operators at $t=0$, Sém. Lothar. Combin. 79 (2018-2019), Art. B79c, 20.
[17] Mark Haiman, Vanishing theorems and character formulas for the Hilbert scheme of points in the plane, Invent. Math. 149 (2002), no. 2, 371-407.
[18] I. G. Macdonald, Symmetric functions and Hall polynomials, second ed., The Clarendon Press, Oxford University Press, New York, 1995, With contributions by A. Zelevinsky, Oxford Science Publications.
[19] Andrei Negut, The shuffle algebra revisited, Int. Math. Res. Not. IMRN (2014), no. 22, 6242-6275.
[20] Andrei Negut, Hecke correspondences for smooth moduli spaces of sheaves, 2018, arXiv:math.AG/1804.03645.
[21] Dun Qiu and Andrew Timothy Wilson, The valley version of the extended delta conjecture, J. Combin. Theory Ser. A 175 (2020), 105271, 31.
[22] Brendon Rhoades, Ordered set partition statistics and the delta conjecture, J. Combin. Theory Ser. A 154 (2018), 172-217.
[23] Marino Romero, The delta conjecture at $q=1$, Trans. Amer. Math. Soc. 369 (2017), no. 10, 7509-7530.
[24] Olivier Schiffmann, On the Hall algebra of an elliptic curve, II, Duke Math. J. 161 (2012), no. 9, 1711-1750.
[25] Olivier Schiffmann and Eric Vasserot, The elliptic Hall algebra and the K-theory of the Hilbert scheme of $\mathbb{A}^{2}$, Duke Math. J. 162 (2013), no. 2, 279-366.
[26] Andrew Timothy Wilson, A weighted sum over generalized Tesler matrices, J. Algebraic Combin. 45 (2017), no. 3, 825-855.
[27] Mike Zabrocki, A module for the delta conjecture, 2019, arXiv:math.CO/1902.08966.
[28] Mike Zabrocki, A proof of the 4-variable Catalan polynomial of the delta conjecture, J. Comb. 10 (2019), no. 4, 599-632.


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