# An algebra associated with a flag in a subspace lattice over a finite field and the quantum affine algebra 

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#### Abstract

In this paper, we introduce an algebra $\mathcal{H}$ from a subspace lattice with respect to a fixed flag which contains its incidence algebra as a proper subalgebra. We then establish a relation between the algebra $\mathcal{H}$ and the quantum affine algebra $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$, where $q$ denotes the cardinality of the base field. It is an extension of the well-known relation between the incidence algebra of a subspace lattice and the quantum algebra $U_{q^{1 / 2}}\left(\mathfrak{s l}_{2}\right)$. We show that there exists an algebra homomorphism from $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$ to $\mathcal{H}$ and that any irreducible module for $\mathcal{H}$ is irreducible as an $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module.


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## 1 Introduction

By a subspace lattice, also known as a projective geometry, we mean the partially ordered set (poset) of all subspaces of a finite-dimensional vector space over a finite field, where the ordering is given by inclusion. In the field of combinatorics, subspace lattices are regarded as $q$-analogs of Boolean lattices and therefore they have been studied from many combinatorial points of view, such as Grassmann codes and Grassmann graphs. On the other hand, the quantum affine algebras $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$ are Hopf algebras that are $q$-deformations of the universal enveloping algebra of the affine Lie algebra $\widehat{\mathfrak{s l}}_{2}$ and their representations are developed in [1, Section 5] as trigonometric solutions of the quantum Yang-Baxter equation. Recently, the author succeeded in [7] in establishing a relation between an

[^0]algebra associated with a subspace lattice and the quantum affine algebras $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$ as an extension of the well-known relation between the incidence algebra of a subspace lattice and the quantum algebras $U_{q}\left(\mathfrak{s l}_{2}\right)$. In this paper, we introduce another algebra and establish its relation to the quantum affine algebra $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$ which is in some sense the opposite extreme to that obtained in [7].

Here we briefly recall the known facts. See [5], [6] and [7] for more detail. Let $H$ denote an $N$-dimensional vector space over a finite field $\mathbb{F}_{q}$ of $q$ elements and let $P$ denote the subspace lattice consisting of all subspaces of $H$. From the poset structure of $P$, we define the lowering matrix $L$ indexed by $P$ whose $(x, y)$-entry is 1 if $y$ covers $x$ and 0 otherwise for $x, y \in P$. Similarly, we define the raising matrix $R$ indexed by $P$ whose $(x, y)$-entry is 1 if $x$ covers $y$ and 0 otherwise for $x, y \in P$. The poset $P$ has the grading which is a partition of $P$ into nonempty sets

$$
P_{i}=\{y \in P \mid \operatorname{dim} y=i\} \quad(0 \leqslant i \leqslant N)
$$

From this grading structure, for $0 \leqslant i \leqslant N$, we define the $i$-th projection matrix $E_{i}^{\star}$ by the diagonal matrix indexed by $P$ whose $(x, x)$-entry is 1 if $x \in P_{i}$ and 0 otherwise for $x \in P$. By the incidence algebra, we mean the complex matrix algebra generated by the above three kinds of matrices $L, R$ and $E_{i}^{\star}$, where $0 \leqslant i \leqslant N$. It is known that there exists a surjective algebra homomorphism from the quantum algebra $U_{q^{1 / 2}}\left(\mathfrak{s l}_{2}\right)$ to the incidence algebra. Moreover, it is also known that any irreducible module for the incidence algebra induces an irreducible $U_{q^{1 / 2}}\left(\mathfrak{s l}_{2}\right)$-module of type 1 .

In our previous paper [7], we extended the algebra homomorphism as follows. Let us fix one subspace $x \in P$ with $0<\operatorname{dim} x<N$ and consider the following new "rectangle" partition of $P$ with respect to $x$ :

$$
\begin{equation*}
P_{i, j}=\{y \in P \mid \operatorname{dim} y=i+j, \operatorname{dim}(y \cap x)=i\} \tag{1}
\end{equation*}
$$

for $0 \leqslant i \leqslant \operatorname{dim} x$ and for $0 \leqslant j \leqslant N-\operatorname{dim} x$. Remark that this is a refinement of the grading. Then define the new projection matrices with respect to this partition and define the complex matrix algebra generated by the lowering, raising matrices and these new projection matrices. By the construction, this new algebra contains the incidence algebra as its subalgebra. Then it is shown in [7] that there exists an algebra homomorphism from the quantum affine algebra $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$ to the new algebra, which extends the above algebra homomorphism from $U_{q^{1 / 2}}\left(\mathfrak{s l}_{2}\right)$ to the incidence algebra. Moreover it is also shown in [7] that any irreducible module for the new algebra induces an irreducible $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module of type $(1,1)$ which is more precisely a tensor product of two evaluation modules.

Now we summarize the main results of this paper. We fix a (full) flag $\left\{x_{i}\right\}_{i=0}^{N}$ on $H$ instead of the subspace $x \in P$, and consider the following new "hyper-cubic" partition of $P$ with respect to $\left\{x_{i}\right\}_{i=0}^{N}$ :

$$
\begin{equation*}
P_{\mu}=\left\{y \in P \mid \operatorname{dim}\left(y \cap x_{i}\right)=\mu_{1}+\mu_{2}+\cdots+\mu_{i}(1 \leqslant i \leqslant N)\right\} \tag{2}
\end{equation*}
$$

for $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$. Then for $\mu \in\{0,1\}^{N}$, we define the projection matrix $E_{\mu}^{*}$ by the diagonal matrix indexed by $P$ whose $(y, y)$-entry is 1 if $y \in P_{\mu}$ and 0 otherwise
for $y \in P$. We next define the complex matrix algebra $\mathcal{H}$ generated by the lowering, raising matrices and these new projection matrices $E_{\mu}^{*}$, where $\mu \in\{0,1\}^{N}$. By the construction, the algebra $\mathcal{H}$ contains the incidence algebra as its subalgebra. We prove that there exists an algebra homomorphism from the quantum affine algebra $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s}}_{2}\right)$ to the algebra $\mathcal{H}$, which again extends the above algebra homomorphism from $U_{q^{1 / 2}}\left(\mathfrak{s l}_{2}\right)$ to the incidence algebra. Moreover it is also proved that any irreducible module for the algebra $\mathcal{H}$ induces an irreducible $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module of type $(1,1)$ which is more precisely a tensor product of evaluation modules of dimension 2. Our main results are Theorems 62 and 66. To prove the main theorems, we classify all the $\mathcal{H}$-modules up to isomorphism and determine the multiplicities appearing in the standard module.

Seen from the viewpoint of the action of the general linear group GL $(N, q)$ on the subspace lattice $P$, we may say the results of this paper are "opposite" to those obtained in our previous paper [7]. (In this paper, however, we will not take this point of view in any essential way. We refer the reader to [3] for this viewpoint.) Indeed, the partitions (1) and (2) turn out to be the orbits of maximal and minimal parabolic subgroups of $\mathrm{GL}(N, q)$, respectively. More precisely, the corresponding subgroups stabilize the fixed subspace $x$ and the fixed flag $\left\{x_{i}\right\}_{i=0}^{N}$, respectively.

It is worth pointing out that our proofs involve a natural and intrinsic combinatorial characterization of the subspace lattice, while the method used in our previous paper [7] is rather oriented towards Lie theory and the representation theory of quantum groups. In this paper, we fix a basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ such that $x_{i}$ is spanned by $v_{1}, v_{2}, \ldots, v_{i}$ for $1 \leqslant i \leqslant N$. With respect to the basis, we identify each subspace in $P$ with a certain matrix whose entries are in the base field $\mathbb{F}_{q}$.

Then, we relate these matrices to classical combinatorial objects, such as Ferrers boards, rook placements and inversion numbers, and interpret algebraic properties of subspaces in terms of these matrices (and moreover, of other combinatorial objects above). Almost all the problems which we concern in this paper arrive at problems in such classical combinatorial fields. This type of argument is motivated by Delsarte [2] and the technique used in this paper is a kind of a generalized version of that in [2].

Comparing the partitions (1) and (2) again, one may ask whether same kinds of results can still be obtained if we take a more general partition, which is defined by replacing a subspace or a full flag by a general flag. We will not develop this point here because the required computation is expected to be far more complicated. However we emphasize that we have done for the two extremal and the most essential cases, and conjecture that similar results still hold in the general case.

We organize this paper as follows. In Section 2, we recall the basic notation and introduce a hyper-cubic structure in a subspace lattice. In Section 3, we recall some notation on Ferrers boards, rook placements and inversion numbers which is used in this paper. In Sections 4 and 5, we introduce a matrix representation of $P$ and interpret some properties of matrices in terms of rook placements and inversion numbers. In Sections 6 and 7 , we introduce the main object of this paper, the algebra $\mathcal{H}$, and discuss the structure of it. In Sections $8,9,10$ and 11, we discuss the $\mathcal{H}$-action on the standard module and classify all the irreducible $\mathcal{H}$-modules up to isomorphism. In Section 12, for
the convenience of the reader, we repeat the relevant material, including the definition of the quantum affine algebra $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$, from [1] without proofs, thus making our exposition self-contained. In Section 13, our main results are stated and proved.

## 2 A subspace lattice and its hyper-cubic structure

We now begin our formal argument. Recall the integers $\mathbb{Z}=\{0, \pm 1, \pm 2, \ldots\}$ and the natural numbers $\mathbb{N}=\{0,1,2, \ldots\}$ and let $\mathbb{C}$ denote the complex field. The Kronecker delta is denoted by $\delta$. Throughout the paper except Section 12, we fix $N \in \mathbb{N} \backslash\{0\}$. Throughout the paper except Sections 3,10 and 12 , we fix a prime power $q$. Let $\mathbb{F}_{q}$ denote a finite field of $q$ elements and let $H$ denote a vector space over $\mathbb{F}_{q}$ with dimension $N$. Let $P$ denote the set of all subspaces of $H$. We view $P$ as a poset with the partial order given by inclusion. The poset $P$ is a graded lattice of rank $N$ where the rank function is defined by its dimension and called the subspace lattice. For two subspaces $y, z \in P$, we say $y$ covers $z$ whenever $z \subseteq y$ and $\operatorname{dim} z=\operatorname{dim} y-1$. By a (full) flag on $H$ we mean a sequence $\left\{x_{i}\right\}_{i=0}^{N}$ of subspaces in $P$ such that $\operatorname{dim} x_{i}=i$ for $0 \leqslant i \leqslant N$ and $x_{i-1} \subsetneq x_{i}$ for $1 \leqslant i \leqslant N$. For the rest of this paper, we fix a flag $\left\{x_{i}\right\}_{i=0}^{N}$ on $H$. A basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ is said to be adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$ whenever each $x_{i}$ is spanned by $v_{1}, v_{2}, \ldots, v_{i}$ for $1 \leqslant i \leqslant N$.

By the $N$-cube we mean the poset consisting of all $N$-tuples in $\{0,1\}^{N}$ with the partial order $\mu \leqslant \nu$ defined by $\mu_{m} \leqslant \nu_{m}$ for all $1 \leqslant m \leqslant N$, where $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right), \nu=$ $\left(\nu_{1}, \nu_{2}, \ldots, \nu_{N}\right) \in\{0,1\}^{N}$. (We note that it is isomorphic to the Boolean lattice of all subsets of an $N$-set.) The $N$-cube is a graded lattice of rank $N$ with the rank function defined by

$$
|\mu|=\mu_{1}+\mu_{2}+\cdots+\mu_{N}
$$

for $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$.
Proposition 1. There exists an order-preserving map from the subspace lattice $P$ to the $N$-cube which sends $y \in P$ to $\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$ where

$$
\operatorname{dim}\left(y \cap x_{m}\right)=\mu_{1}+\mu_{2}+\cdots+\mu_{m}
$$

for $1 \leqslant m \leqslant N$. Moreover this map is surjective.
Proof. Let $y \in P$ and $1 \leqslant m \leqslant N$. We have $\operatorname{dim}\left(y \cap x_{m-1}\right) \leqslant \operatorname{dim}\left(y \cap x_{m}\right)$ since $x_{m-1} \subseteq x_{m}$. We also have $\operatorname{dim}\left(y \cap x_{m}\right)-\operatorname{dim}\left(y \cap x_{m-1}\right) \leqslant 1$ since $\operatorname{dim} x_{m}-\operatorname{dim} x_{m-1}=1$. Thus $\mu_{m}=\operatorname{dim}\left(y \cap x_{m}\right)-\operatorname{dim}\left(y \cap x_{m-1}\right)$ is either 0 or 1 . Therefore this correspondence becomes a map from $P$ to the $N$-cube. Let $y, z \in P$ satisfy $y \subseteq z$ and let $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right), \lambda=$ $\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{N}\right) \in\{0,1\}^{N}$ be the images of $y, z$ under the map, respectively. If there exists $1 \leqslant m \leqslant N$ such that $z \cap x_{m-1}=z \cap x_{m}$, then

$$
\left(y \cap x_{m}\right) \backslash\left(y \cap x_{m-1}\right)=\left(y \cap x_{m}\right) \backslash x_{m-1} \subseteq\left(z \cap x_{m}\right) \backslash x_{m-1}=\left(z \cap x_{m}\right) \backslash\left(z \cap x_{m-1}\right)=\emptyset
$$

and so we have $y \cap x_{m-1}=y \cap x_{m}$. Therefore, $\lambda_{m}=0$ implies $\mu_{m}=0$ for any $1 \leqslant m \leqslant N$, which is equivalent to $\mu \leqslant \lambda$. We have now proved that the map preserves the ordering.

To show its surjectivity, let $v_{1}, v_{2}, \ldots, v_{N}$ denote a basis for $H$ adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$. For any $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$, consider the subspace $y \in P$ spanned by the vectors $\left\{v_{i} \mid 1 \leqslant i \leqslant N, \mu_{i}=1\right\}$. For each $1 \leqslant m \leqslant N$, the intersection $y \cap x_{m}$ is spanned by the vectors $\left\{v_{i} \mid 1 \leqslant i \leqslant m, \mu_{i}=1\right\}$. Therefore, $\operatorname{dim}\left(y \cap x_{m}\right)-\operatorname{dim}\left(y \cap x_{m-1}\right)=\mu_{m}$ for $1 \leqslant m \leqslant N$ and so $y$ is mapped to $\mu$. This proves the map is surjective.

Definition 2. If $\mu \in\{0,1\}^{N}$ is the image of $y \in P$ by the map in Proposition 1, we call $\mu$ the location of $y$. For $\mu \in\{0,1\}^{N}$, let $P_{\mu}$ denote the set of all subspaces at location $\mu$. For notational convenience, for $\mu \in \mathbb{Z}^{N}$ we set $P_{\mu}=\emptyset$ unless $\mu \in\{0,1\}^{N}$.

Note that $P$ is the disjoint union of $P_{\mu}$, where $\mu \in\{0,1\}^{N}$. Observe that $\operatorname{dim} y=|\mu|$ for $y \in P_{\mu}$.

Definition 3. Let $1 \leqslant m \leqslant N$. For $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right), \nu=\left(\nu_{1}, \nu_{2}, \ldots, \nu_{N}\right) \in\{0,1\}^{N}$, we say $\mu$-covers $\nu$ whenever $\nu_{m}<\mu_{m}$ and $\nu_{n}=\mu_{n}$ for $1 \leqslant n \leqslant N$ with $n \neq m$. Similarly, for $y, z \in P$, we say $y m$-covers $z$ whenever $y$ covers $z$ and the location of $y$ $m$-covers the location of $z$.

For each $1 \leqslant m \leqslant N$, let $\widehat{m}$ denote the $N$-tuple in $\{0,1\}^{N}$ with a 1 in $m$-th coordinate and 0 elsewhere. To simplify the notation, we consider the coordinate-wise addition in $\mathbb{Z}^{N}$ so that $\mu m$-covers $\nu$ if and only if $\mu=\nu+\widehat{m}$ for $\mu, \nu \in\{0,1\}^{N}$.

Lemma 4. For $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$ and for $1 \leqslant m \leqslant N$, the following (i), (ii) hold.
(i) Given $y \in P_{\mu}$, the number of subspaces $m$-covered by $y$ is

$$
\delta_{\mu_{m}, 1} q^{\mu_{m+1}+\mu_{m+2}+\cdots+\mu_{N}}
$$

(ii) Given $y \in P_{\mu}$, the number of subspaces which $m$-cover $y$ is

$$
\delta_{\mu_{m}, 0} q^{(m-1)-\left(\mu_{1}+\mu_{2}+\cdots+\mu_{m-1}\right)} .
$$

Proof. (i) Let $\widetilde{P}$ be the set of subspaces in $P$ which are $m$-covered by $y$. Then $\widetilde{P} \subseteq P_{\mu-\widehat{m}}$. If $\mu_{m}=0$, then $\mu-\widehat{m} \notin\{0,1\}^{N}$ and so $\widetilde{P}=\emptyset$. We may assume $\mu_{m}=1$. For $z \in \widetilde{P}$, we have $y \cap x_{m-1}=z \cap x_{m-1}=z \cap x_{m}$ since $z \in P_{\mu-\widehat{m}}, y \in P_{\mu}$ and $z \subseteq y$.

Set $n=\operatorname{dim} y-\operatorname{dim}\left(y \cap x_{m}\right)=\mu_{m+1}+\mu_{m+2}+\cdots+\mu_{N}$. Let $U_{n}$ denote the set of $n$-sets of linearly independent vectors $\mathbf{u}=\left\{u_{1}, u_{2}, \ldots, u_{n}\right\} \subseteq y \backslash\left(y \cap x_{m}\right)$ such that $(\operatorname{Span} \mathbf{u}) \cap\left(y \cap x_{m}\right)=0$. Since $\operatorname{dim} y=|\mu|$ and $\operatorname{dim}\left(y \cap x_{m}\right)=\mu_{1}+\cdots+\mu_{m}$, we have

$$
\left|U_{n}\right|=\prod_{k=1}^{n}\left(q^{|\mu|}-q^{\mu_{1}+\cdots+\mu_{m}+k-1}\right)
$$

For $\mathbf{u} \in U_{n}$ and $1 \leqslant k \leqslant m$, we have $(\operatorname{Span} \mathbf{u}) \cap x_{k}=0$. For $\mathbf{u} \in U_{n}$ and $m \leqslant k \leqslant N$, since $y=\operatorname{Span} \mathbf{u}+\left(y \cap x_{m}\right)$ and $y \cap x_{m} \subseteq x_{k}$, we have $y \cap x_{k}=\left(\operatorname{Span} \mathbf{u}+\left(y \cap x_{m}\right)\right) \cap x_{k}=$ $(\operatorname{Span} \mathbf{u}) \cap x_{k}+\left(y \cap x_{m}\right)$.

We count the cardinality of the following set $S$ in two ways.

$$
S=\left\{(\mathbf{u}, z) \mid \mathbf{u} \in U_{n}, z \in \widetilde{P}, z=\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)\right\}
$$

Let $\mathbf{u} \in U_{n}$, and set $z=\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)$. Since $\operatorname{Span} \mathbf{u} \subseteq y$, we have $z \subseteq y$. For $1 \leqslant k \leqslant m-1$, we have $z \cap x_{k}=\left(\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)\right) \cap x_{k} \supseteq y \cap x_{k}$, and moreover, equality must hold since $z \subseteq y$. So, we have $\operatorname{dim}\left(z \cap x_{k}\right)=\operatorname{dim}\left(y \cap x_{k}\right)$.

For $m \leqslant k \leqslant N$, since $y \cap x_{m-1} \subseteq x_{k}$, we have $z \cap x_{k}=\left(\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)\right) \cap x_{k}=$ (Span $\left.\mathbf{u} \cap x_{k}\right)+\left(y \cap x_{m-1}\right)$. Recall that $y \cap x_{k}=\left(\operatorname{Span} \mathbf{u} \cap x_{k}\right)+\left(y \cap x_{m}\right)$. Since the sums in these two equations are direct, we have $\operatorname{dim}\left(z \cap x_{k}\right)-\operatorname{dim}\left(y \cap x_{k}\right)=\operatorname{dim}\left(y \cap x_{m}\right)-$ $\operatorname{dim}\left(y \cap x_{m-1}\right)=1$. Thus, $z \in \widetilde{P}$. By these comments, we have

$$
\begin{equation*}
|S|=\left|U_{n}\right| \tag{3}
\end{equation*}
$$

Conversely, let $z \in \widetilde{P}$. We have $\operatorname{dim} z-\operatorname{dim}\left(z \cap x_{m}\right)=n$ since $z \in P_{\mu-\widehat{m}}$. So, there exists an $n$-set of linearly independent vectors $\mathbf{u}=\left\{u_{1}, u_{2}, \ldots, u_{n}\right\} \subseteq z \backslash\left(z \cap x_{m}\right)$ such that $z=\operatorname{Span} \mathbf{u}+\left(z \cap x_{m}\right)$ and $(\operatorname{Span} \mathbf{u}) \cap\left(z \cap x_{m}\right)=0$. Let $U_{n}(z)$ denote the set of such $n$-sets. Since $\operatorname{dim} z=|\mu|-1$ and $\operatorname{dim}\left(z \cap x_{m}\right)=\mu_{1}+\cdots+\mu_{m}-1$, we have

$$
\left|U_{n}(z)\right|=\prod_{k=1}^{n}\left(q^{|\mu|-1}-q^{\mu_{1}+\cdots+\mu_{m}+k-2}\right)=q^{-n}\left|U_{n}\right| .
$$

For $\mathbf{u} \in U_{n}(z)$, we have $\mathbf{u} \subseteq y \backslash\left(y \cap x_{m}\right)$ since $z \subseteq y$, and we also have (Span $\left.\mathbf{u}\right) \cap(y \cap$ $\left.x_{m}\right)=(\operatorname{Span} \mathbf{u}) \cap\left(z \cap x_{m}\right)=0$ since $\operatorname{Span} \mathbf{u} \subseteq z \subseteq y$. Thus, $U_{n}(z) \subseteq U_{n}$. We may write $z=\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)$ since $y \cap x_{m-1}=z \cap x_{m}$. Moreover, if $\mathbf{u} \in U_{n}$ satisfies $z=\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)$, then $\mathbf{u} \subseteq z$ and $(\operatorname{Span} \mathbf{u}) \cap\left(z \cap x_{m}\right)=0$, which imply $\mathbf{u} \in U_{n}(z)$. By these comments, we have $U_{n}(z)=\left\{\mathbf{u} \in U_{n} \mid z=\operatorname{Span} \mathbf{u}+\left(y \cap x_{m-1}\right)\right\}$ for $z \in \widetilde{P}$, and so

$$
\begin{equation*}
|S|=\sum_{z \in \widetilde{P}}\left|U_{n}(z)\right|=|\widetilde{P}| \times q^{-n}\left|U_{n}\right| . \tag{4}
\end{equation*}
$$

Thus, by (3) and (4), we have $|\widetilde{P}|=q^{n}$. The result follows.
(ii) Let $\widetilde{P}$ be the set of subspaces in $P$ which $m$-cover $y$. Then $\widetilde{P} \subseteq P_{\mu+\widehat{m}}$. If $\mu_{m}=1$, then $\mu+\widehat{m} \notin\{0,1\}^{N}$ and so $\widetilde{P}=\emptyset$. We may assume $\mu_{m}=0$. Let $\bar{U}=x_{m} \backslash x_{m-1}$. We have $|U|=q^{m}-q^{m-1}$. We also have $U \cap y \subseteq\left(y \cap x_{m}\right) \backslash\left(y \cap x_{m-1}\right)=\emptyset$.

We count the cardinality of the following set $S$ in two ways.

$$
S=\{(u, z) \mid u \in U, z \in \widetilde{P}, z=(\operatorname{Span} u)+y\} .
$$

Let $u \in U$, and set $z=(\operatorname{Span} u)+y$. Then $y \subseteq z$. For $m \leqslant k \leqslant N$, since $\operatorname{Span} u \subseteq x_{k}$, we have $z \cap x_{k}=(\operatorname{Span} u+y) \cap x_{k}=(\operatorname{Span} u)+\left(y \cap x_{k}\right)$. Since the sum is direct, we have $\operatorname{dim}\left(z \cap x_{k}\right)=\operatorname{dim}\left(y \cap x_{k}\right)+1$. For $1 \leqslant k \leqslant m-1$, since $y \cap x_{m} \subseteq x_{m-1}$ and Span $u \cap x_{m-1}=0$, we have

$$
\begin{aligned}
z \cap x_{k} & =\left(z \cap x_{m}\right) \cap x_{m-1} \cap x_{k} \\
& =\left(\operatorname{Span} u+y \cap x_{m}\right) \cap x_{m-1} \cap x_{k} \\
& =\left(\operatorname{Span} u \cap x_{m-1}+y \cap x_{m-1}\right) \cap x_{k} \\
& =y \cap x_{k} .
\end{aligned}
$$

In particular, $\operatorname{dim}\left(z \cap x_{k}\right)=\operatorname{dim}\left(y \cap x_{k}\right)$. Thus $z \in \widetilde{P}$. By these comments, we have

$$
\begin{equation*}
|S|=|U| . \tag{5}
\end{equation*}
$$

Conversely, let $z \in \widetilde{P}$. We have $\operatorname{dim}\left(z \cap x_{m}\right)-\operatorname{dim}\left(z \cap x_{m-1}\right)=1$ since $z \in P_{\mu+\widehat{m}}$. Denote $U(z)=\left(z \cap x_{m}\right) \backslash\left(z \cap x_{m-1}\right)$. For any $u \in U(z)$, we have $z \cap x_{m}=\operatorname{Span} u+\left(z \cap x_{m-1}\right)$. Since $\operatorname{dim}\left(z \cap x_{m}\right)=\mu_{1}+\cdots+\mu_{m-1}+1$ and $\operatorname{dim}\left(z \cap x_{m-1}\right)=\mu_{1}+\cdots+\mu_{m-1}$, we have

$$
|U(z)|=q^{\mu_{1}+\cdots+\mu_{m-1}+1}-q^{\mu_{1}+\cdots+\mu_{m-1}}=q^{\left(\mu_{1}+\cdots+\mu_{m-1}\right)-(m-1)}|U| .
$$

By definition, we have $U(z) \subseteq U$. Let $u \in U$ satisfy $z=\operatorname{Span} u+y$. Then $u \in z$, and so $u \in U(z)$. By these comments, we have $\{u \in U \mid z=(\operatorname{Span} u)+y\}=U(z)$, and so

$$
\begin{equation*}
|S|=\sum_{z \in \widetilde{P}}|U(z)|=|\widetilde{P}| \times q^{\left(\mu_{1}+\cdots+\mu_{m-1}\right)-(m-1)}|U| . \tag{6}
\end{equation*}
$$

Thus, by (5) and (6), we have $|\widetilde{P}|=q^{(m-1)-\left(\mu_{1}+\cdots+\mu_{m-1}\right)}$. The result follows.
Lemma 5. Let $1 \leqslant m<n \leqslant N$. For $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$ with $\mu_{m}=\mu_{n}=1$, the following hold.
(i) Given $z \in P_{\mu}$ and $y \in P_{\mu-\widehat{m}-\widehat{n}}$ with $y \subseteq z$, there exists a unique element in $P_{\mu-\widehat{n}}$ which $m$-covers $y$ and which is $n$-covered by $z$.
(ii) Given $z \in P_{\mu}$ and $y \in P_{\mu-\widehat{m}-\widehat{n}}$ with $y \subseteq z$, there exist exactly $q$ elements in $P_{\mu-\widehat{m}}$ which $n$-cover $y$ and which are $m$-covered by $z$.
(iii) Given $y \in P_{\mu-\widehat{m}}$ and $z \in P_{\mu-\widehat{n}}$, if there exists an element that is covered by both $y$ and $z$, then there exists a unique element that covers both $y$ and $z$.
(iv) Given $y \in P_{\mu-\widehat{m}}$ and $z \in P_{\mu-\widehat{n}}$, if there exists an element that covers both $y$ and $z$, then there exists a unique element that is covered by both $y$ and $z$.

Proof. (i) We first show the existence of such element. Set $w=y+\left(z \cap x_{n-1}\right)$ and let $\mu^{\prime}$ denote the location of $w$. We have $y \subseteq w \subseteq z$ and $\mu-\widehat{m}-\widehat{n} \leqslant \mu^{\prime} \leqslant \mu$. Since $m \leqslant n-1$, we have $w \cap x_{m}=\left(y+z \cap x_{n-1}\right) \cap x_{m} \supseteq z \cap x_{m}$, and moreover equality must hold since $w \subseteq z$. So, $\operatorname{dim}\left(w \cap x_{m}\right)=\operatorname{dim}\left(z \cap x_{m}\right)$. Since $z \cap x_{n-1} \subseteq x_{n}$, we have $w \cap x_{n}=\left(y+z \cap x_{n-1}\right) \cap x_{n}=y \cap x_{n}+z \cap x_{n-1}$. Since $z \in P_{\mu}$ and $y \in P_{\mu-\widehat{m}-\widehat{n}}$, we have $\operatorname{dim}\left(w \cap x_{n}\right)=\operatorname{dim}\left(y \cap x_{n}\right)+\operatorname{dim}\left(z \cap x_{n-1}\right)-\operatorname{dim}\left(y \cap x_{n-1}\right)=\operatorname{dim}\left(y \cap x_{n}\right)-1$. Thus the location $\mu^{\prime}$ must be $\mu-\widehat{n}$, i.e., $w \in P_{\mu-\widehat{n}}$. Since $y \subseteq w \subseteq z$, the element $w$ must $m$-cover $y$ and be $n$-covered by $z$.

We next show the uniqueness of such element. Take any $w^{\prime} \in P_{\mu-\widehat{n}}$ which covers $y$ and which is covered by $z$. Then $w^{\prime}$ must contain both $y$ and $z \cap x_{n-1}$. So $w \subseteq w^{\prime}$. By computing dimensions, $w$ and $w^{\prime}$ must coincide. The result follows.
(ii) Let $\widetilde{P}$ be the set of subspaces in $P$ which cover $y$ and which are covered by $z$. Since $\operatorname{dim}(z / y)=2$, we have $|\widetilde{P}|=\left(q^{2}-1\right) /(q-1)=q+1$. Let $w \in \widetilde{P}$ and let $\mu^{\prime}$ be
the location of $w$. Then, we have $\mu^{\prime} \in\{\mu-\widehat{n}, \mu-\widehat{m}\}$. So, the $q+1$ elements in $\widetilde{P}$ must belong to either $P_{\mu-\widehat{n}}$ or $P_{\mu-\widehat{m}}$. Therefore the result follows from (i).
(iii) Let $w$ be an element that is covered by both $y$ and $z$. Then we have $w \subseteq y \cap z$. Since $y$ and $z$ are distinct, we have $\operatorname{dim} y-1=\operatorname{dim} w \leqslant \operatorname{dim}(y \cap z) \leqslant \operatorname{dim} y-1$, and so $w=y \cap z$. Set $w^{\prime}=y+z$. Then $\operatorname{dim} w^{\prime}=\operatorname{dim} y+1=\operatorname{dim} z+1$. This means $w^{\prime}$ is an element that covers both $y$ and $z$. The uniqueness is clear.
(iv) Similar to (iii).

## 3 Ferrers boards

We introduce the notion of Ferrers boards. For the general theory on this topic, we refer the reader to [4, Chapters 1 and 2]. Note that we modify the notations of [4] to fit our setting.

Let $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$. Then $\mu$ has a natural correspondence with a bipartition of $\{1,2, \ldots, N\}$, which is defined by

$$
\begin{equation*}
S_{\mu}=\left\{s \in \mathbb{N} \mid 1 \leqslant s \leqslant N, \mu_{s}=0\right\}, \quad T_{\mu}=\left\{t \in \mathbb{N} \mid 1 \leqslant t \leqslant N, \mu_{t}=1\right\} . \tag{7}
\end{equation*}
$$

We remark that $S_{\mu}$ and $T_{\mu}$ are empty if and only if $\mu=\mathbf{1}=(1,1, \ldots, 1)$ and $\mu=\mathbf{0}=$ $(0,0, \ldots, 0)$, respectively. The Ferrers board of shape $\mu$ is defined by

$$
\begin{equation*}
B_{\mu}=\left\{(s, t) \in S_{\mu} \times T_{\mu} \mid s<t\right\} . \tag{8}
\end{equation*}
$$

If both $S_{\mu}$ and $T_{\mu}$ are not empty, i.e. if $\mu \neq \mathbf{0}, \mathbf{1}$, we can draw a Ferrers board as a two-dimensional subarray of a matrix whose rows indexed by $S_{\mu}$ and columns indexed by $T_{\mu}$, whose $(s, t)$-entry has a box for all $(s, t) \in B_{\mu}$.

This subarray is also known as a Young diagram of shape $\mu$.
Example $6(\boldsymbol{N}=\mathbf{1 3})$. Let $\mu=(0,1,1,0,1,1,0,1,1,0,0,1,0) \in\{0,1\}^{13}$. Then the corresponding Ferrers board $B_{\mu}$ has the following subarray form:


13

Take a nonempty Ferrers board $B_{\mu}$ of shape $\mu$. For $\left(s_{0}, t_{0}\right) \in B_{\mu}$, the rectangle in $B_{\mu}$ with respect to $\left(s_{0}, t_{0}\right)$, denoted by $B_{\mu}\left(s_{0}, t_{0}\right)$, is defined by

$$
\begin{equation*}
B_{\mu}\left(s_{0}, t_{0}\right)=\left\{(s, t) \in B_{\mu} \mid s \leqslant s_{0}, t \geqslant t_{0}\right\} . \tag{9}
\end{equation*}
$$

It is actually the rectangle in the corresponding Young diagram which includes the topright corner and the $\left(s_{0}, t_{0}\right)$-th box as its bottom-left corner. We remark that such a rectangle is called the Durfee square if it is the largest square in $B_{\mu}$. To see the rectangle structure, we use the following notation:

$$
\begin{equation*}
S_{\mu}(m)=\left\{s \in S_{\mu} \mid s \leqslant m\right\}, \quad T_{\mu}(m)=\left\{t \in T_{\mu} \mid t \geqslant m\right\} \tag{10}
\end{equation*}
$$

for $1 \leqslant m \leqslant N$ so that we can write $B_{\mu}\left(s_{0}, t_{0}\right)=S_{\mu}\left(s_{0}\right) \times T_{\mu}\left(t_{0}\right)$.
Example $7(\boldsymbol{N}=\mathbf{1 3})$. Take $\mu \in\{0,1\}^{13}$ as in Example 6. Then $(4,6) \in B_{\mu}$ and the rectangle $B_{\mu}(4,6)$ is the set of the following eight elements:
$(1,6)$,
$(1,8)$,
$(1,9)$,
$(1,12)$,
$(4,6)$,
$(4,8), \quad(4,9)$,

In the corresponding Young diagram, $B_{\mu}(4,6)$ is the following gray rectangle.


13
Take a nonempty Ferrers board $B_{\mu}$ of shape $\mu$. A subset of $B_{\mu}$ such that no two elements have a common entry is called a rook placement on $B_{\mu}$. Let $\sigma$ denote a rook placement on $B_{\mu}$. The row index set $\pi_{1}(\sigma)$ and the column index set $\pi_{2}(\sigma)$ of $\sigma$ are defined by

$$
\begin{equation*}
\pi_{1}(\sigma)=\left\{s \in S_{\mu} \mid(s, t) \in \sigma \text { for some } t\right\}, \quad \pi_{2}(\sigma)=\left\{t \in T_{\mu} \mid(s, t) \in \sigma \text { for some } s\right\}, \tag{11}
\end{equation*}
$$

respectively. Remark that $\left|\pi_{1}(\sigma)\right|=\left|\pi_{2}(\sigma)\right|=|\sigma|$. Assume $\sigma \neq \emptyset$. For $1 \leqslant i \leqslant|\sigma|$, we denote by $s_{i}$ and by $t_{i}$ the $i$-th smallest element in $\pi_{1}(\sigma)$ and in $\pi_{2}(\sigma)$, respectively. Then $\sigma$ gives rise to a permutation of $\{1,2, \ldots,|\sigma|\}$ which sends $i$ to $j$ where $\left(s_{i}, t_{j}\right) \in \sigma$.
Lemma 8. Let $\mu \in\{0,1\}^{N}$ and $\sigma$ be a rook placement on $B_{\mu}$ with the row/column index sets $\pi_{1}=\pi_{1}(\sigma), \pi_{2}=\pi_{2}(\sigma)$, respectively. Then the pair $\left(\pi_{1}, \pi_{2}\right)$ satisfies the following.
(i) $\left|\pi_{1}\right|=\left|\pi_{2}\right|$.
(ii) Let $n$ denote the common value in (i). For $1 \leqslant i \leqslant n$, the $i$-th smallest element in $\pi_{1}$ is strictly smaller than the $i$-th smallest element in $\pi_{2}$.
Proof. (i) It is clear.
(ii) We may assume $\sigma \neq \emptyset$ since otherwise the assertion is clear. Let $\widetilde{\sigma}$ denote the permutation of $\{1,2, \ldots, n\}$ corresponding to $\sigma$. For $1 \leqslant i \leqslant n$, we write $s_{i}$, $t_{i}$ for the $i$-th smallest element in $\pi_{1}, \pi_{2}$, respectively. Fix $1 \leqslant i \leqslant n$. Since $\widetilde{\sigma}$ is a permutation, there exists $i \leqslant k \leqslant n$ such that $\widetilde{\sigma}(k) \leqslant i$. So we have $\left(s_{k}, t_{\tilde{\sigma}(k)}\right) \in \sigma$. Therefore $s_{i} \leqslant s_{k}<t_{\widetilde{\sigma}(k)} \leqslant t_{i}$ as desired.

Proposition 9. Let $\mu \in\{0,1\}^{N}$. For a pair $\left(\pi_{1}, \pi_{2}\right)$ such that $\pi_{1} \subseteq S_{\mu}$ and $\pi_{2} \subseteq T_{\mu}$, the following are equivalent:
(i) there exists a rook placement $\sigma$ on $B_{\mu}$ such that $\pi_{1}=\pi_{1}(\sigma)$ and $\pi_{2}=\pi_{2}(\sigma)$;
(ii) it satisfies (i), (ii) in Lemma 8.

Proof. We have shown in Lemma 8 that (i) implies (ii).
Suppose we are given $\pi_{1} \subseteq S_{\mu}$ and $\pi_{2} \subseteq T_{\mu}$ satisfying (i), (ii) in Lemma 8. By the condition (i) in Lemma 8, we set $n=\left|\pi_{1}\right|=\left|\pi_{2}\right|$. Let $\sigma=\left\{\left(s_{i}, t_{i}\right) \mid 1 \leqslant i \leqslant n\right\}$, where each $s_{i}, t_{i}$ is the $i$-th smallest element in $\pi_{1}, \pi_{2}$, respectively. By the condition (ii) in Lemma 8, we have $\sigma \subseteq B_{\mu}$ and so $\sigma$ is a rook placement on $B_{\mu}$. By construction, it is clear that $\pi_{1}=\pi_{1}(\sigma)$ and $\pi_{2}=\pi_{2}(\sigma)$. So (ii) implies (i).

Definition 10. Let $\mu \in\{0,1\}^{N}$ and consider the Ferrers board $B_{\mu}$ of shape $\mu$. Then the type of a rook placement $\sigma$ on $B_{\mu}$ is defined by the disjoint union

$$
\pi_{1}(\sigma) \cup \pi_{2}(\sigma) \subseteq\{1,2, \ldots, N\}
$$

where $\pi_{1}(\sigma), \pi_{2}(\sigma)$ are the row/column index sets of $\sigma$ defined in (11).
Lemma 11. Let $\mu \in\{0,1\}^{N}$. For $\lambda \subseteq\{1,2, \ldots, N\}$, the following are equivalent:
(i) there exists a rook placement on $B_{\mu}$ of type $\lambda$;
(ii) the pair $\left(\lambda \cap S_{\mu}, \lambda \cap T_{\mu}\right)$ satisfies (i), (ii) in Lemma 8.

Proof. Immediate from Proposition 9.
Lemma 12. For $\lambda \subseteq\{1,2, \ldots, N\}$, the following are equivalent:
(i) there exists a rook placement on $B_{\mu}$ of type $\lambda$ for some $\mu \in\{0,1\}^{N}$;
(ii) the cardinality of $\lambda$ is even.

Proof. Fix $\lambda \subseteq\{1,2, \ldots, N\}$. Suppose there exists a rook placement $\sigma$ on $B_{\mu}$ of type $\lambda$ for some $\mu \in\{0,1\}^{N}$. Then by Lemma 11, the pair ( $\lambda \cap S_{\mu}, \lambda \cap T_{\mu}$ ) satisfies (i), (ii) in Lemma 8. In particular, $|\lambda|=\left|\lambda \cap S_{\mu}\right|+\left|\lambda \cap T_{\mu}\right|$ is even. So (ii) holds.

Conversely, we suppose $|\lambda|=2 n$ for some $n \in \mathbb{N}$ and show (i) holds. Let ( $\pi_{1}, \pi_{2}$ ) denote the bipartition of $\lambda$ where $\pi_{1}$ contains the first $n$ smallest elements in $\lambda$ and $\pi_{2}$ contains the remaining $n$ elements in $\lambda$. Take any $\mu \in\{0,1\}^{N}$ such that $\pi_{1} \subseteq S_{\mu}$ and $\pi_{2} \subseteq T_{\mu}$. Then we have $\pi_{1}=\lambda \cap S_{\mu}$ and $\pi_{2}=\lambda \cap T_{\mu}$. Observe that the pair ( $\pi_{1}, \pi_{2}$ ) satisfies (i), (ii) in Lemma 8. So by Lemma 11, there exists a rook placement on $B_{\mu}$ of type $\lambda$. In particular, (i) holds.

Since rook placements can be seen as permutations, we define the concept of inversions. Let $\sigma$ be a nonempty rook placement on a Ferrers board $B_{\mu}$ of shape $\mu$. For $\left(s_{0}, t_{0}\right) \in \sigma$, the local inversion number of $\sigma$ at $\left(s_{0}, t_{0}\right)$, denoted by $\operatorname{inv}\left(\sigma, s_{0}, t_{0}\right)$, is defined by

$$
\begin{equation*}
\operatorname{inv}\left(\sigma, s_{0}, t_{0}\right)=\left|\left\{(s, t) \in \sigma \mid s<s_{0}, t>t_{0}\right\}\right|=\left|\sigma \cap B_{\mu}\left(s_{0}, t_{0}\right)\right|-1 \tag{12}
\end{equation*}
$$

For a rook placement $\sigma$, the (total) inversion number of $\sigma$, denoted by $\operatorname{inv}(\sigma)$, is defined by

$$
\operatorname{inv}(\sigma)=\sum_{(s, t) \in \sigma} \operatorname{inv}(\sigma, s, t)
$$

Example $13(\boldsymbol{N}=\mathbf{1 3})$. Take $\mu \in\{0,1\}^{13}$ as in Example 6. Consider the following rook placement $\sigma$ on $B_{\mu}$ :

$$
\sigma=\{(1,9),(4,6),(10,12)\} .
$$

Then we have $\operatorname{inv}(\sigma, 1,9)=\operatorname{inv}(\sigma, 10,12)=0$ and $\operatorname{inv}(\sigma, 4,6)=1$. Thus $\operatorname{inv}(\sigma)=1$.


Lemma 14. Let $\mu \in\{0,1\}^{N}$ and let $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. For $1 \leqslant m \leqslant N$ and for a rook placement $\sigma$ on $B_{\mu}$ of type $\lambda$, we have

$$
\begin{equation*}
\left|\sigma \cap\left(S_{\mu}(m) \times T_{\mu}(m)\right)\right|=\left|\lambda \cap S_{\mu}(m)\right|+\left|\lambda \cap T_{\mu}(m)\right|-\frac{|\lambda|}{2} . \tag{13}
\end{equation*}
$$

In particular, this number is independent of the choice of $\sigma$.
Proof. Since $\pi_{1}(\sigma)=\lambda \cap S_{\mu}, \pi_{2}(\sigma)=\lambda \cap T_{\mu}$ and $\sigma$ is a rook placement, we have

$$
\begin{aligned}
\left|\lambda \cap S_{\mu}(m)\right| & =\left|\pi_{1}(\sigma) \cap S_{\mu}(m)\right|=\left|\sigma \cap\left(S_{\mu}(m) \times T_{\mu}\right)\right|, \\
\left|\lambda \cap T_{\mu}(m)\right| & =\left|\pi_{2}(\sigma) \cap T_{\mu}(m)\right|=\left|\sigma \cap\left(S_{\mu} \times T_{\mu}(m)\right)\right|, \\
|\lambda| & =\left|\pi_{1}(\sigma)\right|+\left|\pi_{2}(\sigma)\right|=2|\sigma| .
\end{aligned}
$$

Set $\overline{S_{\mu}(m)}=S_{\mu} \backslash S_{\mu}(m)$ and $\overline{T_{\mu}(m)}=T_{\mu} \backslash T_{\mu}(m)$. Then we have

$$
\begin{aligned}
& \left|\lambda \cap S_{\mu}(m)\right|+\left|\lambda \cap T_{\mu}(m)\right|-\frac{|\lambda|}{2} \\
& =\left|\sigma \cap\left(S_{\mu}(m) \times T_{\mu}\right)\right|+\left|\sigma \cap\left(S_{\mu} \times T_{\mu}(m)\right)\right|-|\sigma| \\
& =\left|\sigma \cap\left(S_{\mu}(m) \times T_{\mu}(m)\right)\right|-\left|\sigma \cap\left(\overline{S_{\mu}(m)} \times \overline{T_{\mu}(m)}\right)\right| .
\end{aligned}
$$

So, it remains to show that $\sigma \cap\left(\overline{S_{\mu}(m)} \times \overline{T_{\mu}(m)}\right)=\emptyset$. Observe that

$$
B_{\mu} \cap\left(\overline{S_{\mu}(m)} \times \overline{T_{\mu}(m)}\right)=\left\{(s, t) \in B_{\mu} \mid t<m<s\right\}=\emptyset .
$$

Therefore, from $\sigma \subseteq B_{\mu}$, the result follows.
The next lemma is a generalization of [4, Corollary 1.3.10] and the proof of the next lemma is motivated by that of [4, Corollary 1.3.10].

Lemma 15. Let $\mu \in\{0,1\}^{N}$ and let $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. For $1 \leqslant m \leqslant N$, let $\rho(m, \mu, \lambda)$ denote the left-hand side of (13). Then for $q \in \mathbb{C}$ with $q \neq 0,1$, we have

$$
\sum_{\sigma} q^{\operatorname{inv}(\sigma)}=\prod_{s \in \lambda \cap S_{\mu}} \frac{q^{\rho(s, \mu, \lambda)}-1}{q-1},
$$

where the sum is taken over all rook placements $\sigma$ on $B_{\mu}$ of type $\lambda$.
Proof. If $\lambda=\emptyset$, the assertion is clear. (Note that $\operatorname{inv}(\emptyset)=0$.) We assume $\lambda \neq \emptyset$. We claim that there exists a bijection between the following two sets:
(i) rook placements $\sigma$ on $B_{\mu}$ of type $\lambda$,
(ii) integer sequences $\left(a_{s}\right)_{s \in \lambda \cap S_{\mu}}$ such that $0 \leqslant a_{s} \leqslant \rho(s, \mu, \lambda)-1$ for $s \in \lambda \cap S_{\mu}$,
such that $\operatorname{inv}(\sigma)=\sum_{s \in \lambda \cap S_{\mu}} a_{s}$. Suppose for the moment that the claim is true. Then we have

$$
\sum_{\sigma} q^{\operatorname{inv}(\sigma)}=\prod_{s \in \lambda \cap S_{\mu}}\left(\sum_{a_{s}=0}^{\rho(s, \mu, \lambda)-1} q^{a_{s}}\right)=\prod_{s \in \lambda \cap S_{\mu}} \frac{q^{\rho(s, \mu, \lambda)}-1}{q-1} .
$$

So the result follows.
Therefore, it remains to prove the claim. For a given rook placement $\sigma$ on $B_{\mu}$ of type $\lambda$ and for $s \in \lambda \cap S_{\mu}$, there exists a unique $t(s) \in \lambda \cap T_{\mu}$ such that $(s, t(s)) \in \sigma$. Thus, we consider the map $\iota$ that sends $\sigma$ to $\left(a_{s}\right)_{s \in \lambda \cap S_{\mu}}$, where $a_{s}=\operatorname{inv}(\sigma, s, t(s))$. Then for $s \in \lambda \cap S_{\mu}$, we have

$$
\begin{aligned}
0 \leqslant a_{s} & =\left|\sigma \cap\left(S_{\mu}(s) \times T_{\mu}(t(s))\right)\right|-1 \\
& \leqslant\left|\sigma \cap\left(S_{\mu}(s) \times T_{\mu}(s)\right)\right|-1 \\
& =\rho(s, \mu, \lambda)-1,
\end{aligned}
$$

where the second inequality follows from the fact that $s \leqslant t(s)$. This implies that the map $\iota$ is from (i) to (ii). To show the bijectivity of $\iota$, take a sequence $\left(a_{s}\right)_{s \in \lambda \cap S_{\mu}}$ in the set (ii). Set $r=\left|\lambda \cap S_{\mu}\right|$ and for $1 \leqslant i \leqslant r$, we write $s_{i}$ the $i$-th smallest element in $\lambda \cap S_{\mu}$. By definition, observe that

$$
0 \leqslant a_{s_{i}} \leqslant \rho\left(s_{i}, \mu, \lambda\right)-1 \leqslant\left|\lambda \cap S_{\mu}\left(s_{i}\right)\right|-1=i-1
$$

where the third inequality follows from $|\lambda| / 2-\left|\lambda \cap T_{\mu}\left(s_{i}\right)\right|=\left|\left\{t \in \lambda \cap T_{\mu} \mid t<s_{i}\right\}\right| \geqslant 0$. Then, there exists a unique permutation $\widetilde{\sigma}$ of $\{1,2, \ldots, r\}$ such that

$$
\begin{equation*}
a_{s_{i}}=|\{j \mid 1 \leqslant j<i, \widetilde{\sigma}(i)<\widetilde{\sigma}(j)\}| . \tag{14}
\end{equation*}
$$

Then consider the set $\sigma=\left\{\left(s_{i}, t_{\widetilde{\sigma}(i)}\right) \mid 1 \leqslant i \leqslant r\right\}$, where $t_{i}$ is the $i$-th smallest element in $\lambda \cap T_{\mu}$. Fix $1 \leqslant i \leqslant r$. By (14), we have $\widetilde{\sigma}(i) \geqslant i-a_{s_{i}}$ and so we have

$$
\widetilde{\sigma}(i) \geqslant i-a_{s_{i}} \geqslant i-\rho\left(s_{i}, \mu, \lambda\right)+1=\left|\left\{t \in \lambda \cap T_{\mu} \mid t<s_{i}\right\}\right|+1 .
$$

This implies that $s_{i}<t_{\widetilde{\sigma}(i)}$. This holds for any $1 \leqslant i \leqslant r$ and so $\sigma$ becomes a rook placement on $B_{\mu}$. It is clear that $\sigma$ is of type $\lambda$. By construction, the map which sends $\left(a_{s}\right)_{s \in \lambda \cap S_{\mu}}$ to $\sigma$ becomes the inverse of $\iota$. Therefore, our claim holds.

## 4 The matrix representation of $P$

For a field $\mathbb{K}$ and for two finite nonempty sets $S$ and $T$, let $\mathrm{Mat}_{S, T}(\mathbb{K})$ denote the set of all matrices with rows indexed by $S$ and columns indexed by $T$ whose entries are in $\mathbb{K}$. If $S=T$, we write it $\operatorname{Mat}_{S}(\mathbb{K})$ for short. For $M \in \operatorname{Mat}_{S, T}(\mathbb{K})$, the support of $M$, denoted by $\operatorname{Supp}(M)$, is the set of indices containing nonzero entries:

$$
\operatorname{Supp}(M)=\left\{(s, t) \in S \times T \mid M_{s, t} \neq 0\right\} .
$$

For $\mu \in\{0,1\}^{N}$, recall the corresponding bipartition $S_{\mu}, T_{\mu}$ from (7) and the Ferrers board $B_{\mu}$ of shape $\mu$ from (8). We will assume $\mu \neq \mathbf{0}, \mathbf{1}$ in this section so that both $S_{\mu}$ and $T_{\mu}$ are nonempty.

Definition 16. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}$, 1. Let $\mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ denote the set of matrices in $\operatorname{Mat}_{S_{\mu}, T_{\mu}}\left(\mathbb{F}_{q}\right)$ such that $\operatorname{Supp}(M) \subseteq B_{\mu}$.

Recall the set $P_{\mu}$ of subspaces at location $\mu \in\{0,1\}^{N}$ from Definition 2.
Proposition 17. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}$, 1. Fix a basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$. There exists a bijection from $P_{\mu}$ to the set $\mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ in Definition 16 that sends $y \in P_{\mu}$ to $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$, where $y$ has a basis

$$
\sum_{s \in S_{\mu}} Y_{s, t} v_{s}+v_{t}, \quad t \in T_{\mu}
$$

Proof. For $y \in P_{\mu}$, there exists a basis $w_{t}\left(t \in T_{\mu}\right)$ for $y$ such that $w_{t} \in x_{t} \backslash x_{t-1}$ for each $t \in T_{\mu}$. Write each vector $w_{t}$ as a linear combination of the fixed basis $v_{1}, v_{2}, \ldots, v_{t}$ for $x_{t}$. Without loss of generality, we may assume the coefficient of $v_{t}$ is 1 . Use linear operations on the basis $w_{t}\left(t \in T_{\mu}\right)$ to make the coefficient of $v_{t^{\prime}} 0$ for any $t^{\prime} \in T_{\mu}$ with $t \neq t^{\prime}$. Observe that the resulting basis $w_{t}^{\prime}\left(t \in T_{\mu}\right)$ is uniquely determined by $y$. Then from the basis $w_{t}^{\prime}$ $\left(t \in T_{\mu}\right)$, we construct the matrix $Y \in \operatorname{Mat}_{S_{\mu}, T_{\mu}}\left(\mathbb{F}_{q}\right)$ such that $Y_{s, t}$ is the coefficient of $v_{s}$ in $w_{t}^{\prime}$. Then we have $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ since $w_{t}^{\prime} \in x_{t}$. On the other hand, let $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$. For $t \in T_{\mu}$, we write $w_{t}=\sum_{s \in S_{\mu}} Y_{s, t} v_{s}+v_{t}$. Since $\operatorname{Supp}(Y) \subseteq B_{\mu}$, the vector $w_{t}$ is a linear combination of $v_{1}, v_{2}, \ldots, v_{t}$, that means $w_{t} \in x_{t} \backslash x_{t-1}$. Therefore the subspace $y$ spanned by the vectors $w_{t}\left(t \in T_{\mu}\right)$ must belong to $P_{\mu}$.

Definition 18. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}, \mathbf{1}$. Take $y \in P_{\mu}$. By the matrix form of $y$, we mean the matrix $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ which is the image of $y$ under the bijection in Proposition 17. We note that the matrix form of $y$ depends on the basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$.

Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}, \mathbf{1}$. For $s \in S_{\mu}$, we denote by $s^{-}$the one smaller element in $S_{\mu}$. If there is no such element, we set $s^{-}=0$. For $t \in T_{\mu}$, we denote by $t^{+}$the one larger element in $T_{\mu}$. If there is no such element, we set $t^{+}=N+1$. Observe that for $(s, t) \in B_{\mu}$, we have $\left(s^{-}, t\right) \in B_{\mu}$ if $s^{-} \neq 0$ and we have $\left(s, t^{+}\right) \in B_{\mu}$ if $t^{+} \neq N+1$. For $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ and for $(s, t) \in B_{\mu}$, let $M(s, t)$ denote the submatrix of $M$ indexed by the rectangle with respect to $(s, t)$ in (9). Moreover, we set

$$
\begin{align*}
r^{-}(M, s, t) & = \begin{cases}\operatorname{rank}\left(M\left(s^{-}, t\right)\right) & \text { if } s^{-} \neq 0, \\
0 & \text { if } s^{-}=0,\end{cases}  \tag{15}\\
r^{+}(M, s, t) & = \begin{cases}\operatorname{rank}\left(M\left(s, t^{+}\right)\right) & \text {if } t^{+} \neq N+1, \\
0 & \text { if } t^{+}=N+1,\end{cases}  \tag{16}\\
r^{-+}(M, s, t) & = \begin{cases}\operatorname{rank}\left(M\left(s^{-}, t^{+}\right)\right) & \text {if } s^{-} \neq 0 \text { and } t^{+} \neq N+1, \\
0 & \text { if } s^{-}=0 \text { or } t^{+}=N+1 .\end{cases} \tag{17}
\end{align*}
$$

Definition 19. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}$, 1. For $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$, we define the set $\sigma(M)$ consisting of all indices $(s, t) \in B_{\mu}$ such that

$$
r^{\epsilon}(M, s, t)=\operatorname{rank}(M(s, t))-1
$$

for all $\epsilon \in\{-,+,-+\}$.
Lemma 20. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}$, 1. For $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$, the set $\sigma(M)$ in Definition 19 is a rook placement on $B_{\mu}$.

Proof. Fix $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$. Since $\sigma(M)$ is a subset of $B_{\mu}$, it suffices to show that no two elements in $\sigma(M)$ have a common entry. To do this, we take $\left(s_{1}, t\right),\left(s_{2}, t\right) \in \sigma(M)$ and assume $s_{1}<s_{2}$. Observe that $s_{2}^{-} \neq 0$. Since $\left(s_{1}, t\right) \in \sigma(M)$, we have

$$
\begin{equation*}
r^{+}\left(M, s_{1}, t\right)=\operatorname{rank}\left(M\left(s_{1}, t\right)\right)-1 \tag{18}
\end{equation*}
$$

Since $\left(s_{2}, t\right) \in \sigma(M)$, we have $r^{-+}\left(M, s_{2}, t\right)=r^{-}\left(M, s_{2}, t\right)$. By definition, $r^{-+}\left(M, s_{2}, t\right)=$ $r^{+}\left(M, s_{2}^{-}, t\right), r^{-}\left(M, s_{2}, t\right)=\operatorname{rank}\left(M\left(s_{2}^{-}, t\right)\right)$ and so we obtain

$$
\begin{equation*}
r^{+}\left(M, s_{2}^{-}, t\right)=\operatorname{rank}\left(M\left(s_{2}^{-}, t\right)\right) . \tag{19}
\end{equation*}
$$

By (18), the $t$-th column of $M\left(s_{1}, t\right)$ can't be expressed as a linear combination of other columns of $M\left(s_{1}, t\right)$. By (19), the $t$-th column of $M\left(s_{2}^{-}, t\right)$ can be expressed as a linear combination of other columns of $M\left(s_{2}^{-}, t\right)$. This implies $s_{2}^{-}<s_{1}$, which contradicts to $s_{1}<s_{2}$. Therefore we must have $s_{1}=s_{2}$. Similarly, if we take $\left(s, t_{1}\right),\left(s, t_{2}\right) \in \sigma(M)$, then one can show that $t_{1}=t_{2}$. So the result follows.

Recall the local inversion numbers of a rook placement from (12).
Lemma 21. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}$, 1. For $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$, we have

$$
\operatorname{rank}(M(s, t))=\operatorname{inv}(\sigma(M), s, t)+1
$$

for $(s, t) \in \sigma(M)$.
Proof. Fix $(s, t) \in \sigma(M)$. Observe that $\operatorname{rank}(M(s, t))$ can be computed as follows:

$$
\sum_{\left(s^{\prime}, t^{\prime}\right) \in B_{\mu}(s, t)}\left(\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)-r^{-}\left(M, s^{\prime}, t^{\prime}\right)-r^{+}\left(M, s^{\prime}, t^{\prime}\right)+r^{-+}\left(M, s^{\prime}, t^{\prime}\right)\right) .
$$

Then by the definition of $\sigma(M)$, each summand is 1 if $\left(s^{\prime}, t^{\prime}\right) \in \sigma(M)$. We claim that each summand is 0 if $\left(s^{\prime}, t^{\prime}\right) \notin \sigma(M)$. Suppose for the moment that the claim is true. Then $\operatorname{rank}(M(s, t))$ is equal to the cardinality of $\sigma(M) \cap B_{\mu}(s, t)$. The result follows from the definition of local inversion numbers.

Therefore, it remains to prove the claim. If $\left(s^{\prime}, t^{\prime}\right) \notin \sigma(M)$, then there exists $\epsilon \in$ $\{-,+,-+\}$ such that $r^{\epsilon}\left(M, s^{\prime}, t^{\prime}\right) \neq \operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)-1$. If $\epsilon=+$, then $r^{+}\left(M, s^{\prime}, t^{\prime}\right)=$ $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)$. In this case, the $t^{\prime}$-th column of $M\left(s^{\prime}, t^{\prime}\right)$ can be expressed as a linear combination of other columns of $M\left(s^{\prime}, t^{\prime}\right)$. In particular, if $s^{\prime-} \neq 0$, the $t^{\prime}$-th column of $M\left(s^{\prime-}, t^{\prime}\right)$ can be expressed as a linear combination of other columns of $M\left(s^{\prime-}, t^{\prime}\right)$. This implies $r^{-+}\left(M, s^{\prime}, t^{\prime}\right)=r^{-}\left(M, s^{\prime}, t^{\prime}\right)$, which is also true if $s^{\prime-}=0$. Therefore, the summand is 0 . Similarly, if $\epsilon=-$, the summand is 0 . If $\epsilon=-+$, then we have two possibilities: $r^{-+}\left(M, s^{\prime}, t^{\prime}\right)=\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)$ or $r^{-+}\left(M, s^{\prime}, t^{\prime}\right)=\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)-2$. For the first case, we have $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)=r^{-}\left(M, s^{\prime}, t^{\prime}\right)=r^{+}\left(M, s^{\prime}, t^{\prime}\right)=r^{-+}\left(M, s^{\prime}, t^{\prime}\right)$ since we have $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right) \geqslant r^{-}\left(M, s^{\prime}, t^{\prime}\right) \geqslant r^{-+}\left(M, s^{\prime}, t^{\prime}\right)$ and $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right) \geqslant r^{+}\left(M, s^{\prime}, t^{\prime}\right) \geqslant$ $r^{-+}\left(M, s^{\prime}, t^{\prime}\right)$ by definition. This also implies the summand is 0 . For the second case, we have $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right)=r^{-}\left(M, s^{\prime}, t^{\prime}\right)+1=r^{+}\left(M, s^{\prime}, t^{\prime}\right)+1=r^{-+}\left(M, s^{\prime}, t^{\prime}\right)+2$ since we have $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right) \leqslant r^{-}\left(M, s^{\prime}, t^{\prime}\right)+1 \leqslant r^{-+}\left(M, s^{\prime}, t^{\prime}\right)+2$ and $\operatorname{rank}\left(M\left(s^{\prime}, t^{\prime}\right)\right) \leqslant$ $r^{+}\left(M, s^{\prime}, t^{\prime}\right)+1 \leqslant r^{-+}\left(M, s^{\prime}, t^{\prime}\right)+2$ by definition. This also implies the summand is 0 . Hence the claim holds.

Lemma 22. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}$, 1. For a subset $\sigma \subseteq B_{\mu}$, the following are equivalent:
(i) there exists $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ such that $\sigma(M)=\sigma$.
(ii) it is a rook placement on $B_{\mu}$.

Proof. Lemma 20 shows that (i) implies (ii).
Assume we are given a rook placement $\sigma$ on $B_{\mu}$. Consider the matrix $M_{\sigma} \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ defined by

$$
\left(M_{\sigma}\right)_{s, t}= \begin{cases}1 & \text { if }(s, t) \in \sigma \\ 0 & \text { otherwise }\end{cases}
$$

for $s \in S_{\mu}, t \in T_{\mu}$. Then it is easy to check that $\sigma(M)=\sigma$. So (ii) implies (i).

## 5 The number of matrices with given parameter

Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}, \mathbf{1}$. Recall from Lemma 22 that each matrix $\mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ corresponds to a rook placement on the Ferrers board $B_{\mu}$ of shape $\mu$. Recall the sets from (7) and (10). To simplify the notation, we set

$$
\begin{equation*}
n\left(\pi_{1}\right)=\sum_{s \in \pi_{1}}\left|S_{\mu}(s)\right| \tag{20}
\end{equation*}
$$

for a subset $\pi_{1} \subseteq S_{\mu}$.
Definition 23. Let $\mu \in\{0,1\}^{N}$. A subset $\lambda \subseteq\{1,2, \ldots, N\}$ is said to be column-full with respect to $\mu$ whenever $T_{\mu} \subseteq \lambda$. Moreover, a rook placement $\sigma$ on $B_{\mu}$ is said to be column-full whenever the type of $\sigma$ is column-full.

Let $\mu \in\{0,1\}^{N}$. We remark that a rook placement $\sigma$ on $B_{\mu}$ is column-full if and only if the column index set $\pi_{2}(\sigma)$, defined in (11), is maximal.

Proposition 24. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}, \mathbf{1}$ and let $\sigma$ denote a rook placement on $B_{\mu}$. Assume $\sigma$ is column-full in Definition 23. Then the number of matrices $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ such that $\sigma=\sigma(M)$ in Definition 19 is given by

$$
(q-1)^{|\mu|} q^{\operatorname{inv}(\sigma)+\left|B_{\mu}\right|-n\left(\pi_{1}(\sigma)\right)} .
$$

Proof. Let $t \in T_{\mu}$. We count the number of possibilities for the $t$-th column of $M$ with $\sigma=\sigma(M)$. Since $\sigma$ is a column-full rook placement, there uniquely exists $s \in S_{\mu}$ such that $(s, t) \in \sigma$. Since $(s, t) \in \sigma$, we have

$$
\begin{equation*}
r^{-}(M, s, t)=r^{-+}(M, s, t)=\operatorname{rank}(M(s, t))-1 . \tag{21}
\end{equation*}
$$

This means that the $t$-th column of the submatrix $M\left(s^{-}, t\right)$ is a linear combination of other columns. Therefore, the number of possibilities for the $t$-th column of $M\left(s^{-}, t\right)$ is $q^{r(M, s, t)-1}$. For a given such column of $M\left(s^{-}, t\right)$, the number of possibilities for the $t$-th column of $M(s, t)$ is at most $q$ since $M(s, t)$ has one more row than $M\left(s^{-}, t\right)$. In other words, the number of possibilities for the $t$-th column of $M(s, t)$ is at most $q^{r(M, s, t)-1} \times q=$ $q^{r(M, s, t)}$. Similarly, since $(s, t) \in \sigma$, we have

$$
\begin{equation*}
r(M, s, t)-1=r^{+}(M, s, t) . \tag{22}
\end{equation*}
$$

This means that the $t$-th column of the submatrix $M(s, t)$ is not a linear combination of other columns. Since there are $q^{r(M, s, t)-1}$ columns which are linear combinations of columns of $M\left(s, t^{+}\right)$, the number of possibilities for the $t$-th column of $M(s, t)$ is

$$
q^{r(M, s, t)}-q^{r(M, s, t)-1}=(q-1) q^{r(M, s, t)-1}=(q-1) q^{\operatorname{inv}(\sigma, s, t)} .
$$

The second equality follows from Lemma 21 . Since $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$, or equivalently $\operatorname{Supp}(M) \subseteq B_{\mu}$, the $\left(s^{\prime}, t\right)$-entries are 0 if $s^{\prime}>t$. Therefore, for a given $t$-th column of $M(s, t)$, the number of possibilities for the $t$-th column of $M$ is at most $q^{l}$, where

$$
l=\left|\left\{s^{\prime} \in S_{\mu} \mid s<s^{\prime} \leqslant t\right\}\right|=\left|S_{\mu}(t)\right|-\left|S_{\mu}(s)\right| .
$$

Observe that any choices of the $t$-th column among the $q^{l}$ possibilities satisfy both (21) and (22) by construction. Since the conditions (21) and (22) are equivalent to $(s, t) \in \sigma$, the number is exactly $q^{l}$. We have shown that the number of possibilities for the $t$-th column of $M$ is

$$
(q-1) q^{\operatorname{inv}(\sigma, s, t)} \times q^{\left|S_{\mu}(t)\right|-\left|S_{\mu}(s)\right|},
$$

which is independent of the choice of other columns of $M$. Therefore the number of $M$ is obtained by taking the product of the values for all $t \in T_{\mu}$ since $\sigma$ is column-full. The result follows from the definition of $\operatorname{inv}(\sigma)$ and the column-full property and

$$
\sum_{t \in T_{\mu}}\left|S_{\mu}(t)\right|=\left|\left\{(s, t) \in S_{\mu} \times T_{\mu} \mid s<t\right\}\right|=\left|B_{\mu}\right|
$$

Corollary 25. Let $\mu \in\{0,1\}^{N}$ with $\mu \neq \mathbf{0}, \mathbf{1}$ and let $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. Assume $\lambda$ is column-full with respect to $\mu$ in Definition 23. Then the number of matrices $M \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ such that $\sigma(M)$ is of type $\lambda$ in Definitions 10 and 19 is given by

$$
q^{\left|B_{\mu}\right|-n\left(\lambda \cap S_{\mu}\right)} \prod_{s \in \lambda \cap S_{\mu}}\left(q^{\rho(s, \mu, \lambda)}-1\right)
$$

where $\rho(s, \mu, \lambda)$ is defined in Lemma 15.
Proof. Use Lemma 15 and Proposition 24.

## 6 The algebra $\mathcal{H}$

Recall $\operatorname{Mat}_{P}(\mathbb{C})$, the set of all matrices whose rows and columns are indexed by $P$ and whose entries are in $\mathbb{C}$. We see it as a $\mathbb{C}$-algebra. We write $I \in \operatorname{Mat}_{P}(\mathbb{C})$ for the identity matrix and $O \in \operatorname{Mat}_{P}(\mathbb{C})$ for the zero matrix. In this section, we introduce a subalgebra $\mathcal{H}$ of $\operatorname{Mat}_{P}(\mathbb{C})$ which represents the $N$-cube structure in $P$.

Let $V=\mathbb{C} P$ denote the vector space over $\mathbb{C}$ consisting of the column vectors whose coordinates are indexed by $P$ and whose entries are in $\mathbb{C}$. Observe that $\operatorname{Mat}_{P}(\mathbb{C})$ acts on $V$ by left multiplication. We call $V$ the standard module for $\operatorname{Mat}_{P}(\mathbb{C})$. We equip $V$ with the standard Hermitian inner product defined by $\langle u, v\rangle=u^{T} \bar{v}$ for $u, v \in V$, where ${ }^{T}$ denotes transpose and ${ }^{-}$denotes complex conjugate.

Recall from Definition 2 that we have partitioned $P$ into the sets $P_{\mu}$ of all subspaces at location $\mu$ for $\mu \in\{0,1\}^{N}$. For $\mu \in \mathbb{Z}^{N}$, define a diagonal matrix $E_{\mu}^{*} \in \operatorname{Mat}_{P}(\mathbb{C})$ by

$$
\left(E_{\mu}^{*}\right)_{y, y}=\left\{\begin{array}{ll}
1 & \text { if } y \in P_{\mu}, \\
0 & \text { if } y \notin P_{\mu},
\end{array} \quad y \in P\right.
$$

Observe that $E_{\mu}^{*}=O$ unless $\mu \in\{0,1\}^{N}$. By construction, we have

$$
E_{\mu}^{*} E_{\nu}^{*}=\delta_{\mu, \nu} E_{\mu}^{*}, \quad \mu, \nu \in\{0,1\}^{N}
$$

$$
I=\sum_{\mu \in\{0,1\}^{N}} E_{\mu}^{*}
$$

Moreover, we have a decomposition of $V$ :

$$
V=\sum_{\mu \in\{0,1\}^{N}} E_{\mu}^{*} V, \quad \quad(\text { direct sum })
$$

where $E_{\mu}^{*} V$ is the subspace of $V$ consisting of the vectors whose nonzero entries are indexed by elements in $P_{\mu}$. Thus, the matrix $E_{\mu}^{*}$ is the projection from $V$ onto $E_{\mu}^{*} V$ and we call it the projection matrix.

Definition 26. By the above comments, the matrices $E_{\mu}^{*}$, where $\mu \in\{0,1\}^{N}$ form a basis for a commutative subalgebra of $\operatorname{Mat}_{P}(\mathbb{C})$. We denote this subalgebra by $\mathcal{K}$.

We now introduce matrices that generate $\mathcal{K}$. For $1 \leqslant m \leqslant N$, we define diagonal matrices $K_{m} \in \operatorname{Mat}_{P}(\mathbb{C})$ by

$$
\left(K_{m}\right)_{y, y}=q^{1 / 2-\mu_{m}}, \quad y \in P_{\mu},
$$

where $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right)$.
Lemma 27. For $1 \leqslant m \leqslant N$, we have

$$
K_{m}=\sum_{\mu \in\{0,1\}^{N}} q^{1 / 2-\mu_{m}} E_{\mu}^{*},
$$

where $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right)$.
Proof. Immediate from the construction.
Proposition 28. The algebra $\mathcal{K}$ in Definition 26 is generated by $K_{m}$ for $1 \leqslant m \leqslant N$.
Proof. By Lemma 27, the matrices $K_{m}(1 \leqslant m \leqslant N)$ generate a subalgebra $\mathcal{K}^{\prime}$ of $\mathcal{K}$. By Lemma 27 and since $E_{\mu}^{*}$ are idempotent, for $\nu=\left(\nu_{1}, \nu_{2}, \ldots, \nu_{N}\right) \in\{0,1\}^{N}$, we have

$$
K_{1}^{\nu_{1}} K_{2}^{\nu_{2}} \cdots K_{N}^{\nu_{N}}=\sum_{\mu \in\{0,1\}^{N}} q^{\sum_{m=1}^{N}\left(\nu_{m} / 2-\mu_{m} \nu_{m}\right)} E_{\mu}^{*}
$$

By linear algebra, if the coefficient matrix $Q_{N}$ indexed by $\{0,1\}^{N}$, whose $(\nu, \mu)$-entry is $q^{\sum_{m=1}^{N}\left(\nu_{m} / 2-\mu_{m} \nu_{m}\right)}$, is invertible, then each $E_{\mu}^{*}$ is a linear combination of $K_{1}^{\nu_{1}} K_{2}^{\nu_{2}} \cdots K_{N}^{\nu_{N}}$ $\left(\nu=\left(\nu_{1}, \nu_{2}, \ldots, \nu_{N}\right) \in\{0,1\}^{N}\right)$. In particular, in this case, $E_{\mu}^{*}$ is a polynomial in $K_{m}$ $(1 \leqslant m \leqslant N)$ for every $\mu \in\{0,1\}^{N}$ and consequently, $\mathcal{K}^{\prime}=\mathcal{K}$. So, it remains to show that the determinant of $Q_{N}$ is nonzero. First, observe that if $N=1$,

$$
\operatorname{det} Q_{1}=\operatorname{det}\left(\begin{array}{cc}
1 & 1 \\
q^{1 / 2} & q^{-1 / 2}
\end{array}\right)=q^{-1 / 2}-q^{1 / 2} \neq 0
$$

since $q \neq 1$. We next consider the matrix $Q_{1}^{\otimes N}$ indexed by $\{0,1\}^{N}$. The $(\nu, \mu)$-entry of $Q_{1}^{\otimes N}$ is given by

$$
q^{\sum_{m=1}^{N}\left(\nu_{m} / 2-\mu_{m} \nu_{m}\right)},
$$

which is same as that of $Q_{N}$. This means $Q_{N}=Q_{1}^{\otimes N}$. By $Q_{N}=Q_{1}^{\otimes N}$ and $\operatorname{det} Q_{1} \neq 0$, we conclude that $\operatorname{det} Q_{N} \neq 0$ as desired.

Next we introduce two kinds of matrices from covering relations in Definition 3. For $1 \leqslant m \leqslant N$, the matrices $L_{m}, R_{m} \in \operatorname{Mat}_{P}(\mathbb{C})$ are defined by

$$
\left(L_{m}\right)_{y, z}=\left\{\begin{array}{ll}
1 & \text { if } z m \text {-covers } y, \\
0 & \text { otherwise },
\end{array} \quad\left(R_{m}\right)_{y, z}= \begin{cases}1 & \text { if } y m \text {-covers } z, \\
0 & \text { otherwise }\end{cases}\right.
$$

for $y, z \in P$. We remark that for each $1 \leqslant m \leqslant N$, the matrices $L_{m}$ and $R_{m}$ are transposes of one another. Recall the comment in the above of Lemma 4.

Lemma 29. For $1 \leqslant m \leqslant N$ and $\mu \in\{0,1\}^{N}$, we have the following.
(i) $L_{m} E_{\mu}^{*}=E_{\mu-\widehat{m}}^{*} L_{m}$ and $R_{m} E_{\mu}^{*}=E_{\mu+\widehat{m}}^{*} R_{m}$.
(ii) $L_{m} E_{\mu}^{*} V \subseteq E_{\mu-\widehat{m}}^{*} V$ and $R_{m} E_{\mu}^{*} V \subseteq E_{\mu+\widehat{m}}^{*} V$.

Proof. Immediate from the construction.
Because of Lemma 29 (ii), we call $L_{m}$ the lowering matrices and $R_{m}$ the raising matrices.

Definition 30. Let $\mathcal{H}$ denote the subalgebra of $\operatorname{Mat}_{P}(\mathbb{C})$ generated by $L_{m}, R_{m}(1 \leqslant m \leqslant$ $N$ ) and the algebra $\mathcal{K}$ in Definition 26.

Proposition 31. The algebra $\mathcal{H}$ in Definition 30 is semisimple.
Proof. This follows since $\mathcal{H}$ is closed under the conjugate-transpose map.
We recall the incidence algebra, which is generated by $L, R$ and $E_{i}^{\star}(0 \leqslant i \leqslant N)$ from the second paragraph in Section 1. We remark that $\mathcal{H}$ contains the incidence algebra as its subalgebra because $L=\sum_{m=1}^{N} L_{m}, R=\sum_{m=1}^{N} R_{m}$ and $E_{i}^{\star}=\sum_{\mu \in\{0,1\}^{N},|\mu|=i} E_{\mu}^{*}$. Moreover, if $N \geqslant 2$, the incidence algebra is a proper subalgebra of $\mathcal{H}$.

## 7 The structure of the algebra $\mathcal{H}$

In this section, we discuss the relations among the generators $L_{m}, R_{m}, K_{m}$ of the algebra $\mathcal{H}$.

Proposition 32. For $1 \leqslant m, n \leqslant N$ with $m \neq n$, the following hold.
(i) $L_{m} K_{n}=K_{n} L_{m}$.
(ii) $R_{m} K_{n}=K_{n} R_{m}$.
(iii) $q L_{m} K_{m}=K_{m} L_{m}$.
(iv) $R_{m} K_{m}=q K_{m} R_{m}$.

Proof. This lemma follows by combining Lemmas 27 and 29 (i).
Proposition 33. For $1 \leqslant m, n \leqslant N$, we have the following.
(i) $L_{m}^{2}=R_{m}^{2}=0$.
(ii) $q L_{m} L_{n}=L_{n} L_{m}$ if $m<n$.
(iii) $R_{m} R_{n}=q R_{n} R_{m}$ if $m<n$.
(iv) $L_{m} R_{n}=R_{n} L_{m}$ if $m \neq n$.

Proof. (i) It follows from the definition of $L_{m}$ and $R_{m}$. (ii), (iii) These are matrix reformulations of Lemma 5 (i), (ii). (iv) This is a matrix reformulation of Lemma 5 (iii), (iv).

## 8 The $L_{m^{-}}$and $R_{m^{-}}$-actions on $V$

We now describe a basis for $V$, which is the key in this paper. In this section, we fix a basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$ and assume that the matrix forms in Definition 18 are always taken with respect to this basis $v_{1}, v_{2}, \ldots, v_{N}$.

Definition 34. Let $\chi$ denote a nontrivial character of the additive group $\mathbb{F}_{q}$ and let $\mu \in\{0,1\}^{N}$. For $y \in P_{\mu}$, define a vector $\chi_{y} \in V$ as follows.
(i) If $\mu=\mathbf{0}$ or $\mathbf{1}$, then for $z \in P$, the $z$-th entry of $\chi_{y}$ is 1 if $y=z$ and 0 otherwise.
(ii) If $\mu \neq \mathbf{0}, \mathbf{1}$, then for $z \in P$, the $z$-th entry of $\chi_{y}$ is defined by

$$
\begin{cases}\chi\left(\operatorname{tr}\left(Y Z^{T}\right)\right) & \text { if } z \in P_{\mu} \\ 0 & \text { if } z \notin P_{\mu}\end{cases}
$$

where $Y, Z \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ are the matrix forms of $y, z$, respectively in Definition 18. Here ${ }^{T}$ denotes transpose and $\operatorname{tr}$ denotes the trace map of matrices.

For the rest of this section, we fix a nontrivial character $\chi$ of the additive group $\mathbb{F}_{q}$.
Lemma 35. For $\mu \in\{0,1\}^{N}$, the set of vectors $\chi_{y} \in V$ for $y \in P_{\mu}$ in Definition 34 forms an orthogonal basis for the vector space $E_{\mu}^{*} V$.

Proof. Let $\mu \in\{0,1\}^{N}$. For $y \in P_{\mu}$, observe that $\chi_{y} \in E_{\mu}^{*} V$ from the construction. If $\mu=\mathbf{0}$ or $\mathbf{1}$, then the assertion is trivial, since $\operatorname{dim} E_{\mu}^{*} V=1$. Assume $\mu \neq \mathbf{0}, \mathbf{1}$ and take $y, y^{\prime} \in P_{\mu}$. Consider the Hermitian inner product

$$
\left\langle\chi_{y}, \chi_{y^{\prime}}\right\rangle=\sum_{z \in P} \chi_{y}(z) \overline{\chi_{y^{\prime}}(z)}
$$

where $\chi_{y}(z), \chi_{y^{\prime}}(z)$ denote the $z$-th entries of $\chi_{y}, \chi_{y^{\prime}}$, respectively. By the definitions of $\chi_{y}(z), \chi_{y^{\prime}}(z)$, we have

$$
\left\langle\chi_{y}, \chi_{y^{\prime}}\right\rangle=\sum_{Z \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)} \chi\left(\operatorname{tr}\left(Y-Y^{\prime}\right) Z^{T}\right),
$$

where $Y, Y^{\prime}$ are the matrix forms of $y, y^{\prime}$, respectively. Assume $y \neq y^{\prime}$ and equivalently $Y \neq Y^{\prime}$. Observe that for $g \in \mathbb{F}_{q}$, the number of $Z \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ such that $\operatorname{tr}\left(Y-Y^{\prime}\right) Z^{T}=g$ does not depend on $g$, and so the number is $\left|\mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)\right| /\left|\mathbb{F}_{q}\right|=q^{\left|B_{\mu}\right|-1}$. Therefore, we have

$$
\left\langle\chi_{y}, \chi_{y^{\prime}}\right\rangle=q^{\left|B_{\mu}\right|-1} \sum_{g \in \mathbb{F}_{q}} \chi(g)=0 .
$$

The last equality follows from the orthogonality of the character $\chi$ and the trivial character. Therefore the set of vectors $\chi_{y}$ for $y \in P_{\mu}$ becomes an orthogonal basis for a subspace $V_{\mu}$ of $E_{\mu}^{*} V$. By comparing their dimensions, we have $V_{\mu}=E_{\mu}^{*} V$ and the result follows.

Recall the $m$-covering relation from Definition 3 .
Lemma 36. Let $1 \leqslant m \leqslant N$ and let $\mu, \nu \in\{0,1\}^{N}$ with $\mu, \nu \neq \mathbf{0}, \mathbf{1}$ such that $\mu$-covers $\nu$. Take $y \in P_{\mu}, z \in P_{\nu}$ and let $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ and $Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ denote the matrix forms of $y, z$, respectively in Definition 18. Then $y$ m-covers $z$ if and only if

$$
Z_{s, t}=Y_{s, t}+Y_{s, m} Z_{m, t}
$$

for $s \in S_{\mu}$ and for $t \in T_{\nu}$.
Proof. Recalling the bijection of Proposition 17, for $t \in T_{\mu}$ and $t^{\prime} \in T_{\nu}$, we write

$$
w_{t}(Y)=\sum_{s \in S_{\mu}} Y_{s, t} v_{s}+v_{t}, \quad w_{t^{\prime}}(Z)=\sum_{s^{\prime} \in S_{\nu}} Z_{s^{\prime}, t^{\prime}} v_{s^{\prime}}+v_{t^{\prime}}
$$

Assume $y$ covers $z$. For each $t^{\prime} \in T_{\nu}$, since $z \subseteq y$, the vector $w_{t^{\prime}}(Z)$ is a linear combination of $w_{t}(Y)$, where $t \in T_{\mu}$. Comparing the coefficients of $v_{t}$ for $t \in T_{\mu}$, we have $w_{t^{\prime}}(Z)=$ $Z_{m, t^{\prime}} w_{m}(Y)+w_{t^{\prime}}(Y)$. Then comparing the coefficients of $v_{s}$ for $s \in S_{\mu}$, we obtain the desired equality. On the other hand, assume the equality $Z_{s, t^{\prime}}=Y_{s, t^{\prime}}+Z_{m, t^{\prime}} Y_{s, m}$ for $s \in S_{\mu}$ and $t^{\prime} \in T_{\nu}$. By the same argument above, we have $w_{t^{\prime}}(Z) \in y$ for all $t^{\prime} \in T_{\nu}$. This implies $y$ covers $z$, as desired.

Lemma 37. Let $1 \leqslant m \leqslant N$ and let $\mu, \nu \in\{0,1\}^{N}$ with $\mu, \nu \neq \mathbf{0}, \mathbf{1}$ such that $\mu$-covers $\nu$. Take $y \in P_{\mu}, z \in P_{\nu}$ and let $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right), Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ denote the matrix forms of $y, z$, respectively in Definition 18. Then the $z$-th entry of $L_{m} \chi_{y}$ is given by

$$
L_{m} \chi_{y}(z)=q^{\left|S_{\mu}(m-1)\right|} \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\nu}} Y_{s, t} Z_{s, t}\right)
$$

if $Y_{s, m}=\sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}$ for all $s \in S_{\mu}$ with $s<m$ and 0 otherwise.
Proof. By the definition of $L_{m}$, the $z$-th entry of $L_{m} \chi_{y}$ is defined by

$$
L_{m} \chi_{y}(z)=\sum_{y^{\prime}} \chi_{y}\left(y^{\prime}\right),
$$

where the sum is taken over all $y^{\prime} \in P_{\mu}$ such that $y^{\prime} m$-covers $z$. Then by Definition 34 and Lemma 36, we have

$$
L_{m} \chi_{y}(z)=\sum_{Y^{\prime}} \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\mu}} Y_{s, t} Y_{s, t}^{\prime}\right)
$$

where the sum is taken over all $Y^{\prime} \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ such that $Z_{s, t}=Y_{s, t}^{\prime}+Y_{s, m}^{\prime} Z_{m, t}$ for $s \in S_{\mu}$ and for $t \in T_{\nu}$. Observe that $T_{\mu} \backslash T_{\nu}=\{m\}$ and so we have

$$
\begin{aligned}
L_{m} \chi_{y}(z) & =\sum \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\nu}} Y_{s, t}\left(Z_{s, t}-Y_{s, m}^{\prime} Z_{m, t}\right)+\sum_{s \in S_{\mu}} Y_{s, m} Y_{s, m}^{\prime}\right) \\
& =\sum \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\nu}} Y_{s, t} Z_{s, t}\right) \chi\left(\sum_{s \in S_{\mu}}\left(Y_{s, m}-\sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}\right) Y_{s, m}^{\prime}\right),
\end{aligned}
$$

where the first sum in each line is taken over all $Y_{s, m}^{\prime} \in \mathbb{F}_{q}$ such that $Y_{s, m}^{\prime}=0$ if $s>m$. If $Y_{s, m} \neq \sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}$ for some $s \in S_{\mu}$ with $s<m$, then by the same argument as in the proof of Lemma 35, the sum is 0 . If $Y_{s, m}=\sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}$ for all $s \in S_{\mu}$ with $s<m$, then

$$
L_{m} \chi_{y}(z)=q^{\left|S_{\mu}(m-1)\right|} \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\nu}} Y_{s, t} Z_{s, t}\right)
$$

Here the coefficient $q^{\left|S_{\mu}(m-1)\right|}$ is the number of choices for $Y_{s, m}^{\prime} \in \mathbb{F}_{q}$ for $s \in S_{\mu}$ with $s<m$. The result follows.
Lemma 38. Let $1 \leqslant m \leqslant N$ and let $\mu, \nu \in\{0,1\}^{N}$ with $\mu, \nu \neq \mathbf{0}, \mathbf{1}$ such that $\mu$-covers $\nu$. Take $y \in P_{\mu}, z \in P_{\nu}$ and let $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right), Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ denote the matrix forms of $y, z$, respectively in Definition 18. Then the $y$-th entry of $R_{m} \chi_{z}$ is given by

$$
R_{m} \chi_{z}(y)=q^{\left|T_{\nu}(m+1)\right|} \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\nu}} Y_{s, t} Z_{s, t}\right)
$$

if $Z_{m, t}=-\sum_{s \in S_{\mu}} Z_{s, t} Y_{s, m}$ for all $t \in T_{\nu}$ with $t>m$ and 0 otherwise.

Proof. Similar to the proof of Lemma 37.
Lemma 39. Referring to Lemma 37, let $\lambda$ denote the type of $\sigma(Y)$ in Definitions 10 and 19. Then the number of $Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ such that $Y_{s, m}=\sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}$ for all $s \in S_{\mu}$ with $s<m$ is given by $q^{l}$ where

$$
l=\left|B_{\nu}\right|-\left|\lambda \cap S_{\mu}(m-1)\right|-\left|\lambda \cap T_{\mu}(m+1)\right|+|\lambda| / 2
$$

if $m \notin \lambda$, and 0 otherwise.
Proof. We count the number of possibilities for $Z_{s, t} \in \mathbb{F}_{q}$ for $s \in S_{\nu}$ and $t \in T_{\nu}$. If $s>t$, then $Z_{s, t}=0$ since $\operatorname{Supp}(Z) \subseteq B_{\nu}$. If $s \neq m$ and $s<t$, then $Z_{s, t}$ is arbitrary and therefore the number of possibilities is $q$. The number of such pairs $(s, t)$ is given by

$$
\left|\left\{(s, t) \in B_{\nu} \mid s \neq m\right\}\right|=\left|B_{\nu}\right|-\left|T_{\nu}(m+1)\right| .
$$

For the case $s=m$ and $m<t$, by the constraint, the sequence $\left(Z_{m, t}\right)_{t \in T_{\nu}, t>m}$ must be a solution of the system of linear equations over $\mathbb{F}_{q}$ :

$$
C \mathbf{u}=\mathbf{c}
$$

where $C=\left(Y_{s, t}\right)_{s \in S_{\mu}, s<m, t \in T_{\nu}, t>m}$ is the coefficient matrix, $\mathbf{u}=\left(u_{t}\right)_{t \in T_{\nu}, t>m}$ is the unknown vector and $\mathbf{c}=\left(Y_{s, m}\right)_{s \in S_{\mu}, s<m}$ is the constant vector. By linear algebra, the system $C \mathbf{u}=\mathbf{c}$ has a solution if and only if the rank of the augmented matrix $[C, \mathbf{c}]$ is equal to the rank of the coefficient matrix $C$. By Definition 19, it is also equivalent to $(s, m) \notin \sigma(Y)$ for all $s \in S_{\mu}$ with $s<m$, which means $m \notin \lambda$. Moreover, suppose there is a solution of the system $C \mathbf{u}=\mathbf{c}$. Since there are $\left|T_{\nu}(m+1)\right|$ columns in $C$, the number of solutions is given by

$$
q^{\left|T_{\nu}(m+1)\right|-\mathrm{rank} C} .
$$

By the proof of Lemma 21, the rank of $C$ is computed as follows:

$$
\begin{aligned}
\operatorname{rank} C & =|\{(s, t) \in \sigma(Y) \mid s \leqslant m-1, t \geqslant m+1\}| \\
& =|\{(s, t) \in \sigma(Y) \mid s \leqslant m-1\}|+|\{(s, t) \in \sigma(Y) \mid t \geqslant m+1\}|-|\sigma(Y)| \\
& =\left|\lambda \cap S_{\mu}(m-1)\right|+\left|\lambda \cap T_{\mu}(m+1)\right|-|\lambda| / 2 .
\end{aligned}
$$

Therefore the result follows.
Lemma 40. Referring to Lemma 38, let $\lambda$ denote the type of $\sigma(Z)$ in Definitions 10 and 19. Then the number of $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ such that $Z_{m, t}=-\sum_{s \in S_{\mu}} Z_{s, t} Y_{s, m}$ for all $t \in T_{\nu}$ with $t>m$ is given by $q^{l}$ where

$$
l=\left|B_{\mu}\right|-\left|\lambda \cap S_{\mu}(m-1)\right|-\left|\lambda \cap T_{\mu}(m+1)\right|+|\lambda| / 2
$$

if $m \notin \lambda$, and 0 otherwise.
Proof. Similar to the proof of Lemma 39.

Definition 41. Let $\mu \in\{0,1\}^{N}$ and take $y \in P_{\mu}$. If $\mu \neq \mathbf{0}, \mathbf{1}$, then let $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ denote the matrix form of $y$ in Definition 18. Then the type of $y$ is defined to be the type of $\sigma(Y)$ in Definitions 10 and 19. If $\mu=\mathbf{0}$ or $\mathbf{1}$, then the type of $y$ is defined to be the empty set. We note that the type of $y$ depends on the basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ since the matrix form does.

Lemma 42. Let $\mu \in\{0,1\}^{N}$ and let $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. For $1 \leqslant m \leqslant N$, the following are equivalent:
(i) for any $y \in P_{\mu}$ of type $\lambda$, we have $L_{m} \chi_{y}=0$;
(ii) $m \in S_{\mu}$ or $m \in \lambda$.

Proof. Set $\nu=\mu-\widehat{m}$ so that $\mu$-covers $\nu$. Then $\nu \neq 1$. For $y \in P_{\mu}$, observe that $L_{m} \chi_{y} \in E_{\nu}^{*} V$ by Lemma 29 (ii).
(i) $\Rightarrow$ (ii) Suppose $L_{m} \chi_{y}=0$ for any $y \in P_{\mu}$ of type $\lambda$. If $\mu=\mathbf{0}$, then $m \in S_{\mu}$ and so (ii) holds. If $\mu=1$, then $P_{\mu}=\{y=H\}$ and any subspaces $z \in P_{\nu}$ are $m$-covered by $y$, and so the $z$-th entry of $L_{m} \chi_{y}$ is

$$
L_{m} \chi_{y}(z)=\chi_{y}(y)=1
$$

by Definition 34. This is a contradiction to $L_{m} \chi_{y}=0$. If $\nu=\mathbf{0}$, then $P_{\nu}=\{0\}$ and any subspaces $y^{\prime} \in P_{\mu} m$-cover 0 , and so the 0 -th entry of $L_{m} \chi_{y}$ is

$$
L_{m} \chi_{y}(0)=\sum_{Y^{\prime} \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)} \chi\left(\operatorname{tr}\left(Y Y^{\prime T}\right)\right)
$$

where $Y$ is the matrix form of $y$. By the same argument as in the proof of Lemmas 35 and 37 , the sum vanishes (if and) only if $Y$ is not the zero matrix from the orthogonality of the characters $\chi$ and the trivial character. Since $y \in P_{\widehat{m}}$, we must have $m \in \lambda$. If $\mu, \nu \neq \mathbf{0}, \mathbf{1}$, then by Lemma $37, L_{m} \chi_{y}=0$ implies that there is no $Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ such that $Y_{s, m}=\sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}$ for all $s \in S_{\mu}$ with $s<m$, where $Y \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ denote the matrix form of $y$ in Definition 18. In this case, by Lemma 39, we have $m \in \lambda$, where $\lambda$ is the type of $y$.
(ii) $\Rightarrow$ (i) Suppose $m \in S_{\mu}$ or $m \in \lambda$. If $m \in S_{\mu}$, then $\nu \notin\{0,1\}^{N}$ and so $E_{\mu}^{*} V=0$. This implies $L_{m} \chi_{y}=0$ since $L_{m} \chi_{y} \in E_{\nu}^{*} V$. We now assume $m \in T_{\mu}$ and $m \in \lambda$. Observe that $\mu \neq \mathbf{0}$. If $\mu=\mathbf{1}$, then $\lambda=\emptyset$ by Definition 41. This contradicts to $m \in \lambda$. If $\nu=\mathbf{0}$, then by the similar argument above, $m \in \lambda$ implies the matrix form of $y$ is not the zero matrix. Then this implies the 0 -th entry of $L_{m} \chi_{y}$ is 0 , which means $L_{m} \chi_{y}=0$. If $\mu, \nu \neq \mathbf{0}, \mathbf{1}$, then the result follows from Lemmas 37 and 39 .

Lemma 43. Let $\nu \in\{0,1\}^{N}$ and let $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11 with $\mu$ replaced by $\nu$. For $1 \leqslant m \leqslant N$, the following are equivalent:
(i) for any $z \in P_{\nu}$ of type $\lambda$, we have $R_{m} \chi_{z}=0$;
(ii) $m \in T_{\nu}$ or $m \in \lambda$.

Proof. Similar to the proof of Lemma 42.
Recall from Lemma 12 , a subset $\lambda \subseteq\{1,2, \ldots, N\}$ becomes a type if and only if it has even cardinality. For $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality, let $V_{\lambda}$ denote the subspace of $V$ spanned by the vectors $\chi_{y} \in V$ for all $y \in P$ of type $\lambda$ in Definitions 34 and 41 . Then for $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality, we define a matrix $E_{\lambda} \in \operatorname{Mat}_{P}(\mathbb{C})$ such that

$$
\begin{aligned}
&\left(E_{\lambda}-I\right) V_{\lambda}=0 \\
& E_{\lambda} V_{\lambda^{\prime}}=0 \text { if } \lambda \neq \lambda^{\prime}
\end{aligned}
$$

where $\lambda^{\prime} \subseteq\{1,2, \ldots, N\}$ with even cardinality. In other words, $E_{\lambda}$ is the projection from $V$ onto $V_{\lambda}$. Observe that $E_{\mu}^{*}$ and $E_{\lambda}$ commute for all $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality.

Lemma 44. For $\mu \in\{0,1\}^{N}$ and for $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality, the following are equivalent:
(i) $E_{\mu}^{*} E_{\lambda}=E_{\lambda} E_{\mu}^{*} \neq 0$;
(ii) the pair $\left(\lambda \cap S_{\mu}, \lambda \cap T_{\mu}\right)$ satisfies (i), (ii) in Lemma 8.

Proof. This is a matrix interpretation of Lemma 11.

## 9 The $L_{m} R_{m^{-}}$and $R_{m} L_{m^{-}}$-actions on $V$

In this section, we fix a basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$ and assume that the matrix forms in Definition 18 and the types in Definition 41 are always taken with respect to this basis $v_{1}, v_{2}, \ldots, v_{N}$. We also fix a nontrivial character $\chi$ of the additive group $\mathbb{F}_{q}$. Recall from Section 8, the definition of $E_{\lambda}$ for $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality depends on the basis $v_{1}, v_{2}, \ldots, v_{N}$ and on the character $\chi$. We show in this section, that $E_{\lambda}$ is independent of the basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$ and the nontrivial character $\chi$ of the additive group $\mathbb{F}_{q}$.

Lemma 45. Let $1 \leqslant m \leqslant N$, and let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. Set

$$
\begin{equation*}
\kappa(m, \mu, \lambda)=\left|S_{\mu}(m-1) \backslash \lambda\right|+\left|T_{\mu}(m+1) \backslash \lambda\right|+|\lambda| / 2 . \tag{23}
\end{equation*}
$$

Then for $v \in E_{\mu}^{*} E_{\lambda} V$, we have the following:

$$
R_{m} L_{m} v= \begin{cases}q^{\kappa(m, \mu, \lambda)} v & \text { if } m \in T_{\mu} \text { and } m \notin \lambda, \\ 0 & \text { if } m \in S_{\mu} \text { or } m \in \lambda\end{cases}
$$

Proof. Observe that $R_{m} L_{m}$ acts on $E_{\mu}^{*} V$ by Lemma 29 (ii). Fix $y \in P_{\mu}$ of type $\lambda$ in Definition 41. We show that $\chi_{y}$ is an eigenvector for $R_{m} L_{m}$. If $m \in S_{\mu}$ or $m \in \lambda$, then by Lemma 42, we have $L_{m} \chi_{y}=0$ and so $\chi_{y}$ is an eigenvector for $R_{m} L_{m}$ with respect to the eigenvalue 0. If $\mu=\mathbf{1}$, then $P_{\mu}=\{H\}$ and $\lambda=\emptyset$. So we have $\operatorname{dim} E_{\mu}^{*} V=1$. Therefore, $\chi_{y}$ is an eigenvector of $R_{m} L_{m}$ and the corresponding eigenvalue is the number of subspaces which are $m$-covered by $y=H$, which is equal to $q^{N-m}=q^{k(m, 1,0)}$ by Lemma 4 (i). Set $\nu=\mu-\widehat{m}$ so that $\mu m$-covers $\nu$. If $m \in T_{\mu}, m \notin \lambda$ and $\nu=\mathbf{0}$, then $P_{\nu}=\{0\}$ and $\lambda=\emptyset$. In other words, the matrix form of $y$ in Definition 18 equals the zero matrix $O$, and so $y^{\prime}$-th entry $\chi_{y}\left(y^{\prime}\right)$ of $\chi_{y}$ is 1 if $y^{\prime} \in P_{\mu}$ and 0 if $y^{\prime} \notin P_{\mu}$. Since $P_{\nu}=\{0\}, \chi_{y}$ is an eigenvector of $R_{m} L_{m}$ and the corresponding eigenvalue is the number of subspaces which $m$-covers $z=0$, which is equal to $q^{m-1}=q^{\kappa(m, \widehat{m}, \emptyset)}$ by Lemma 4 (ii). If $m \in T_{\mu}, m \notin \lambda$, $\mu \neq 1$ and $\nu \neq 0$, then we have

$$
R_{m} L_{m} \chi_{y}=\frac{1}{\left|P_{\mu}\right|} \sum_{y^{\prime} \in P_{\mu}}\left\langle R_{m} L_{m} \chi_{y}, \chi_{y^{\prime}}\right\rangle \chi_{y^{\prime}} .
$$

Let $y^{\prime} \in P_{\mu}$. Since $L_{m}$ and $R_{m}$ are (conjugate-)transposes of one another, we have

$$
\begin{aligned}
\left\langle R_{m} L_{m} \chi_{y}, \chi_{y^{\prime}}\right\rangle & =\left\langle L_{m} \chi_{y}, L_{m} \chi_{y^{\prime}}\right\rangle \\
& =\sum_{z \in P_{\nu}} L_{m} \chi_{y}(z) \overline{L_{m} \chi_{y^{\prime}}(z)}
\end{aligned}
$$

Let $Y, Y^{\prime} \in \mathcal{M}_{\mu}\left(\mathbb{F}_{q}\right)$ and $Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ be the matrix forms of $y, y^{\prime}, z$, respectively in Definition 18. Then by Lemma 37, it becomes

$$
\sum_{z \in P_{\nu}} L_{m} \chi_{y}(z) \overline{L_{m} \chi_{y^{\prime}}(z)}=q^{2\left|S_{\mu}(m-1)\right|} \sum \chi\left(\sum_{s \in S_{\mu}} \sum_{t \in T_{\nu}}\left(Y_{s, t}-Y_{s, t}^{\prime}\right) Z_{s, t}\right)
$$

where the sum is taken over all $Z \in \mathcal{M}_{\nu}\left(\mathbb{F}_{q}\right)$ such that

$$
\begin{equation*}
\sum_{t \in T_{\nu}} Y_{s, t} Z_{m, t}=Y_{s, m}, \quad \sum_{t \in T_{\nu}} Y_{s, t}^{\prime} Z_{m, t}=Y_{s, m}^{\prime} \tag{24}
\end{equation*}
$$

for all $s \in S_{\mu}$ with $s<m$. Then, since $\operatorname{Supp}(Z) \subseteq B_{\nu}$, by the orthogonality of the character $\chi$ and the trivial character, the sum vanishes unless $Y_{s, t}=Y_{s, t}^{\prime}$ for all $s \in S_{\mu}$ and $t \in T_{\nu}$ with $s<t$, which by (24) and Lemma 39 implies $Y=Y^{\prime}$ and so $y=y^{\prime}$. In particular, $\chi_{y}$ is an eigenvector of $R_{m} L_{m}$. Moreover, using Lemma 39 and $\left|P_{\mu}\right|=q^{\left|B_{\mu}\right|}$, we can easily show that the corresponding eigenvalues is $q^{\kappa(m, \mu, \lambda)}$.

Lemma 46. Let $1 \leqslant m \leqslant N$, and let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. Recall $\kappa(m, \mu, \lambda)$ from (23). Then for $v \in E_{\mu}^{*} E_{\lambda} V$, we have the following:

$$
L_{m} R_{m} v= \begin{cases}q^{\kappa(m, \mu, \lambda)} v & \text { if } m \in S_{\mu} \text { and } m \notin \lambda, \\ 0 & \text { if } m \in T_{\mu} \text { or } m \in \lambda .\end{cases}
$$

Proof. Similar to the proof of Lemma 45.
Proposition 47. For $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality, the matrix $E_{\lambda}$ belongs to the algebra $\mathcal{H}$ in Definition 30.

Proof. Referring to (23), we set

$$
\theta(m, \mu, \lambda)= \begin{cases}q^{\kappa(m, \mu, \lambda)} & \text { if } m \notin \lambda, \\ 0 & \text { if } m \in \lambda\end{cases}
$$

for $1 \leqslant m \leqslant N, \mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11. Then by Lemmas 45 and 46, we have

$$
R_{m} L_{m}+L_{m} R_{m}=\sum_{\mu, \lambda} \theta(m, \mu, \lambda) E_{\mu}^{*} E_{\lambda},
$$

where the sum is taken over all pairs $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11. Pick $\mu \in\{0,1\}^{N}$ and multiply each term on the left of the above equation, by $E_{\mu}^{*}$. Then we obtain

$$
E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}=\sum_{\lambda} \theta(m, \mu, \lambda) E_{\mu}^{*} E_{\lambda},
$$

where the sum is taken over $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11. For a subset $\lambda^{\prime} \subseteq\{1,2, \ldots, N\}$, since $E_{\mu}^{*}, E_{\lambda}$ are mutually commutative and they are idempotents, we have

$$
\begin{equation*}
\prod_{m \in \lambda^{\prime}}\left(E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}\right)=\sum_{\lambda}\left(\prod_{m \in \lambda^{\prime}} \theta(m, \mu, \lambda)\right) E_{\mu}^{*} E_{\lambda}, \tag{25}
\end{equation*}
$$

where the sum is taken over $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11. Observe that the coefficient $\prod_{m \in \lambda^{\prime}} \theta(m, \mu, \lambda)$ vanishes if and only if $\lambda \cap \lambda^{\prime} \neq \emptyset$.

We show that each $E_{\mu}^{*} E_{\lambda}$ is a polynomial in $E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}(1 \leqslant m \leqslant N)$ by induction on $|\lambda|$. If we apply $\lambda^{\prime}=\{1,2, \ldots, N\}$ to the equation (25), then the right-hand side becomes a nonzero scalar multiple of $E_{\mu}^{*} E_{\emptyset}$. This means that $E_{\mu}^{*} E_{\emptyset}$ is a polynomial in $E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}(1 \leqslant m \leqslant N)$. Suppose each $E_{\mu}^{*} E_{\lambda^{\prime \prime}}$ is a polynomial in $E_{\mu}^{*} R_{m} L_{m}+$ $E_{\mu}^{*} L_{m} R_{m}(1 \leqslant m \leqslant N)$ for all $\left|\lambda^{\prime \prime}\right|<k$. Then for $\lambda \subseteq\{1,2, \ldots, N\}$ with $|\lambda|=k$ satisfying (ii) in Lemma 11, we apply $\lambda^{\prime}=\{1,2, \ldots, N\} \backslash \lambda$ to the equation (25). The right-hand side is a nonzero scalar multiple of $E_{\mu}^{*} E_{\lambda}$ plus a linear combination of $E_{\mu}^{*} E_{\lambda^{\prime \prime}}$ with $\left|\lambda^{\prime \prime}\right|<k$, which is a polynomial in $E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}(1 \leqslant m \leqslant N)$ by inductive hypothesis. This means $E_{\mu}^{*} E_{\lambda}$ is also a polynomial in $E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}(1 \leqslant m \leqslant N)$. Therefore each $E_{\mu}^{*} E_{\lambda}$ is a polynomial in $E_{\mu}^{*} R_{m} L_{m}+E_{\mu}^{*} L_{m} R_{m}(1 \leqslant m \leqslant N)$. Observe that for $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality, we have

$$
E_{\lambda}=\sum_{\mu} E_{\mu}^{*} E_{\lambda}
$$

where the sum is taken over all $\mu \in\{0,1\}^{N}$ such that the pair $\left(\lambda \cap S_{\mu}, \lambda \cap T_{\mu}\right)$ satisfies (i), (ii) in Lemma 8. Then the result follows.

We remark that the above proof of Proposition 47 also shows that the matrices $E_{\lambda}$ are independent of the basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ adapted to the flag $\left\{x_{i}\right\}_{i=0}^{N}$ and the nontrivial character $\chi$ of the additive group $\mathbb{F}_{q}$.
Lemma 48. Let $V_{\text {new }}$ denote the set of all $v \in V$ such that $L_{m} v=0$ for all $1 \leqslant m \leqslant N$. Then we have

$$
V_{\mathrm{new}}=\sum_{\mu, \lambda} E_{\mu}^{*} E_{\lambda} V \quad \quad(\text { direct sum })
$$

where the sum is taken over all pairs $(\mu, \lambda)$ with $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11 such that $\lambda$ is column-full with respect to $\mu$ in Definition 23. Proof. Take $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11. Observe that the following are equivalent:
(i) for $1 \leqslant m \leqslant N$, we have either $m \in S_{\mu}$ or $m \in \lambda$;
(ii) $\lambda$ is column-full with respect to $\mu$.

Then by Lemma 42 , if $\lambda$ is column-full with respect to $\mu$, we have $E_{\mu}^{*} E_{\lambda} V \subseteq V_{\text {new }}$. Suppose $\lambda$ is not column-full with respect to $\mu$. Then there exists $1 \leqslant m \leqslant N$ such that $m \in T_{\mu}$ and $m \notin \lambda$. By Lemma 45, for any $v \in E_{\mu}^{*} E_{\lambda} V, R_{m} L_{m} v$ is a nonzero scalar multiple of $v$. In particular, $L_{m} v \neq 0$ and so $v \notin V_{\text {new }}$. By above comments and by the fact that $V$ is the direct sum of $E_{\mu}^{*} E_{\lambda} V$, the result follows.

Recall the column-full property in Definition 23. For $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11, we say $\lambda$ is row-full with respect to $\mu$ if $S_{\mu} \subseteq \lambda$.
Lemma 49. Let $V_{\text {old }}$ denote the set of all $v \in V$ such that $R_{m} v=0$ for all $1 \leqslant m \leqslant N$. Then we have

$$
V_{\text {old }}=\sum_{\mu, \lambda} E_{\mu}^{*} E_{\lambda} V \quad \quad(\text { direct sum })
$$

where the sum is taken over all pairs $(\mu, \lambda)$ with $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11 such that $\lambda$ is row-full with respect to $\mu$.

Proof. Similar to the proof of Lemma 48.

## 10 The scalar $\kappa(m, \mu, \lambda)$

In this section, we discuss on the scalar $\kappa(m, \mu, \lambda)$ in (23).
Lemma 50. Let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. Referring to (23), we have the following.

$$
\sum_{m}(-1)^{\mu_{m}} \kappa(m, \mu, \lambda)=\frac{(N-1)(N-2|\mu|)}{2}
$$

where the sum is taken over all $1 \leqslant m \leqslant N$ with $m \notin \lambda$.

Proof. Fix $\mu \in\{0,1\}^{N}$ and we prove the assertion by induction on the cardinality of $\lambda$. Let $F(\lambda)$ denote the left-hand side of the equation. Observe that

$$
\begin{aligned}
F(\lambda)= & \left(\sum_{s \in S_{\mu} \backslash \lambda}\left|S_{\mu}(s-1) \backslash \lambda\right|+\sum_{s \in S_{\mu} \backslash \lambda}\left|T_{\mu}(s+1) \backslash \lambda\right|+\sum_{s \in S_{\mu} \backslash \lambda} \frac{|\lambda|}{2}\right) \\
& -\left(\sum_{t \in T_{\mu} \backslash \lambda}\left|S_{\mu}(t-1) \backslash \lambda\right|+\sum_{t \in T_{\mu} \backslash \lambda}\left|T_{\mu}(t+1) \backslash \lambda\right|+\sum_{t \in T_{\mu} \backslash \lambda} \frac{|\lambda|}{2}\right) .
\end{aligned}
$$

Each of the second and fourth sums counts the number of pairs $(s, t) \in S_{\mu} \times T_{\mu}$ with $s, t \notin \lambda$ and $t>s$. Thus, the second and fourth terms cancel out, i.e.,

$$
F(\lambda)=\left(\sum_{s \in S_{\mu} \backslash \lambda}\left|S_{\mu}(s-1) \backslash \lambda\right|\right)-\left(\sum_{t \in T_{\mu} \backslash \lambda}\left|T_{\mu}(t+1) \backslash \lambda\right|\right)+\frac{|\lambda|}{2}\left(\left|S_{\mu} \backslash \lambda\right|-\left|T_{\mu} \backslash \lambda\right|\right) .
$$

If $\lambda=\emptyset$, then we have

$$
\sum_{s \in S_{\mu}}\left|S_{\mu}(s-1)\right|=0+1+\cdots+(N-|\mu|-1)=\frac{(N-|\mu|)(N-|\mu|-1)}{2}
$$

and

$$
\sum_{t \in T_{\mu}}\left|T_{\mu}(t+1)\right|=0+1+\cdots+(|\mu|-1)=\frac{|\mu|(|\mu|-1)}{2}
$$

Therefore, we have

$$
F(\emptyset)=\frac{(N-|\mu|)(N-|\mu|-1)}{2}-\frac{|\mu|(|\mu|-1)}{2}=\frac{(N-1)(N-2|\mu|)}{2}
$$

and the result follows.
If $|\lambda| \geqslant 1$, there exist $s=\max \left(\lambda \cap S_{\mu}\right)$ and $t=\max \left(\lambda \cap T_{\mu}\right)$ since the pair $\left(\lambda \cap S_{\mu}, \lambda \cap T_{\mu}\right)$ satisfies (i) in Lemma 8. Set $\lambda^{\prime}=\lambda \backslash\{s, t\}$ and observe that $\lambda^{\prime}$ satisfies (ii) in Lemma 11 and we have

$$
\sum_{s^{\prime} \in S_{\mu} \backslash \lambda}\left|S_{\mu}\left(s^{\prime}-1\right) \backslash \lambda\right|=\left(\sum_{s^{\prime} \in S_{\mu} \backslash \lambda^{\prime}}\left|S_{\mu}\left(s^{\prime}-1\right) \backslash \lambda^{\prime}\right|\right)-\left|S_{\mu} \backslash \lambda\right|,
$$

and

$$
\sum_{t^{\prime} \in T_{\mu} \backslash \lambda}\left|T_{\mu}\left(t^{\prime}+1\right) \backslash \lambda\right|=\left(\sum_{t^{\prime} \in T_{\mu} \backslash \lambda^{\prime}}\left|T_{\mu}\left(t^{\prime}+1\right) \backslash \lambda^{\prime}\right|\right)-\left|T_{\mu} \backslash \lambda\right| .
$$

Therefore, since $|\lambda|=\left|\lambda^{\prime}\right|+2$, we have

$$
F(\lambda)=F\left(\lambda^{\prime}\right)
$$

and by the inductive hypothesis, the result follows.

In the next lemma, we do not assume $q$ to be a prime power.
Lemma 51. Let $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. Referring to (23), for $q \in \mathbb{C}$ with $q \neq 0$, 1 , we have the following.

$$
\sum_{m}(-1)^{\mu_{m}} q^{\kappa(m, \mu, \lambda)}=\frac{q^{N-|\mu|}-q^{|\mu|}}{q-1}
$$

where the sum is taken over all $1 \leqslant m \leqslant N$ with $m \notin \lambda$.
Proof. For notational convenience, in this proof we use the following notation. Take $n \in \mathbb{N} \backslash\{0\}$. For $\nu=\left(\nu_{1}, \nu_{2}, \ldots, \nu_{n}\right) \in\{0,1\}^{n}$, a sequence $\mathfrak{a}=\left(\mathfrak{a}_{1}, \mathfrak{a}_{2}, \ldots, \mathfrak{a}_{n}\right) \in \mathbb{Z}^{n}$ is called a $\kappa$-sequence with respect to $\nu$ whenever it satisfies

$$
\mathfrak{a}_{i}= \begin{cases}\mathfrak{a}_{i-1}+1 & \text { if } \nu_{i-1}=\nu_{i}, \\ -\mathfrak{a}_{i-1} & \text { if } \nu_{i-1} \neq \nu_{i}\end{cases}
$$

for $2 \leqslant i \leqslant n$. We call $\nu \in\{0,1\}^{n}$ reduced if $n \leqslant 2$ or $\nu$ is either $\mathbf{0}$ or $\mathbf{1}$. Let $\mathfrak{a}=$ $\left(\mathfrak{a}_{1}, \mathfrak{a}_{2}, \ldots, \mathfrak{a}_{n}\right) \in \mathbb{Z}^{n}$ be a $\kappa$-sequence with respect to a non-reduced $\nu=\left(\nu_{1}, \nu_{2}, \ldots, \nu_{n}\right) \in$ $\{0,1\}^{n}$. Then we have $\nu_{i-1} \neq \nu_{i}$ for some $2 \leqslant i \leqslant n$. Let $\nu^{\prime} \in\{0,1\}^{n-2}$ be the sequence obtained from $\nu$ by removing the coordinates $i-1$ and $i$, and let $\mathfrak{a}^{\prime} \in \mathbb{Z}^{n-2}$ denote the sequence obtained from $\mathfrak{a}$ by removing the same pair of coordinates. Then it is easy to show that the sequence $\mathfrak{a}^{\prime}$ is again a $\kappa$-sequence with respect to $\nu^{\prime}$. Moreover, by continuing this process, any $\kappa$-sequence reaches a $\kappa$-sequence with respect to a reduced tuple $\nu$. More precisely, a $\kappa$-sequence $\mathfrak{a}$ with respect to $\nu \in\{0,1\}^{n}$ becomes
(i) a $\kappa$-sequence of length 2 with respect to $(0,1)$ or $(1,0)$ if $2|\nu|=n$,
(ii) a $\kappa$-sequence of length $n-2|\nu|$ with respect to $\mathbf{0} \in\{0,1\}^{n-2|\nu|}$ if $2|\nu|<n$,
(iii) a $\kappa$-sequence of length $2|\nu|-n$ with respect to $1 \in\{0,1\}^{2|\nu|-n}$ if $2|\nu|>n$.

We call this a reduced $\kappa$-sequence from $\mathfrak{a}$. For a $\kappa$-sequence $\mathfrak{a}=\left(\mathfrak{a}_{1}, \mathfrak{a}_{2}, \ldots, \mathfrak{a}_{n}\right) \in \mathbb{Z}^{n}$ with respect to $\nu=\left(\nu_{1}, \nu_{2}, \ldots, \nu_{n}\right) \in\{0,1\}^{n}$, we define

$$
f(\nu, \mathfrak{a} ; q)=\sum_{i=1}^{n}(-1)^{\nu_{i}} q^{(-1)^{\nu_{i}} \mathfrak{a}_{i}} .
$$

Observe that the value $f(\nu, \mathfrak{a} ; q)$ is invariant under the reducing process above. In particular, if $\mathfrak{a}^{\prime}$ is a reduced $\kappa$-sequence with respect to $\nu^{\prime}$ obtained from a $\kappa$-sequence $\mathfrak{a}$ with respect to $\nu$, then we have $f(\nu, \mathfrak{a} ; q)=f\left(\nu^{\prime}, \mathfrak{a}^{\prime} ; q\right)$.

Set $n=N-|\lambda|$. Let $\nu=\nu(\mu, \lambda) \in\{0,1\}^{n}$ be the sequence obtained from $\mu$ by removing all the coordinates indexed by $\lambda$. Consider the sequence $\mathfrak{a} \in \mathbb{Z}^{n}$ defined by

$$
\mathfrak{a}=\left((-1)^{\mu_{m}} \kappa(m, \mu, \lambda)\right)_{m \in\{1,2, \ldots, N\} \backslash \lambda},
$$

where the index $m$ increases from left to right. For $1 \leqslant m<m^{\prime} \leqslant N$ with $m, m^{\prime} \notin \lambda$, observe that

$$
\kappa(m, \mu, \lambda)-\kappa\left(m^{\prime}, \mu, \lambda\right)=\left|\left\{t \in T_{\mu} \backslash \lambda \mid m<t \leqslant m^{\prime}\right\}\right|-\left|\left\{s \in S_{\mu} \backslash \lambda \mid m \leqslant s<m^{\prime}\right\}\right| .
$$

Therefore, the sequence $\mathfrak{a}$ is a $\kappa$-sequence with respect to $\nu$. Let $\mathfrak{a}^{\prime}$ be a reduced $\kappa$ sequence with respect to $\nu^{\prime}$ from $\mathfrak{a}$. Then the left-hand side of the desired identity becomes $f\left(\nu^{\prime}, \mathfrak{a}^{\prime} ; q\right)$.

We first consider the case $2|\mu|=N$. Then we have $\left|S_{\mu}\right|=\left|T_{\mu}\right|$ and so $2|\nu|=n$ since the pair $\left(\lambda \cap S_{\mu}, \lambda \cap T_{\mu}\right)$ satisfies (i) in Lemma 8. Thus, $\mathfrak{a}^{\prime}$ is a $\kappa$-sequence of length 2 with respect to $(0,1)$ or $(1,0)$ and so $f\left(\nu^{\prime}, \mathfrak{a}^{\prime} ; q\right)=0$ and the result follows. We next consider the case $2|\mu|<N$. Then by the similar argument above, we have $2|\nu|<n$. Thus, $\mathfrak{a}^{\prime}$ is a $\kappa$-sequence of length $n-2|\nu|=N-2|\mu|$ with respect to $\mathbf{0} \in\{0,1\}^{n-2|\nu|}$. By the definition of $\kappa$-sequence, $\mathfrak{a}^{\prime}$ is an arithmetic sequence with common difference 1 . We claim that

$$
\mathfrak{a}^{\prime}=(|\mu|,|\mu|+1, \ldots, N-|\mu|-1) .
$$

To show this, since it is an arithmetic sequence, it suffices to show that

$$
\sum_{a^{\prime} \in \mathfrak{a}^{\prime}} a^{\prime}=\frac{(N-1)(N-2|\mu|)}{2} .
$$

This follows from Lemma 50 since $\sum_{a^{\prime} \in \mathfrak{a}^{\prime}} a^{\prime}=\sum_{a \in \mathfrak{a}} a$. For the case $2|\mu|>N$, the proof is similar to that for the case $2|\mu|<N$. Hence the result follows.

## 11 The $\mathcal{H}$-modules

Recall from Proposition 31 that the algebra $\mathcal{H}$ is semisimple. Thus the standard module $V$ is a direct sum of irreducible $\mathcal{H}$-modules, and every irreducible $\mathcal{H}$-module appears in $V$ up to isomorphism. We now discuss the $\mathcal{H}$-submodules of $V$, which from now on we call $\mathcal{H}$-modules for short.

Proposition 52. Any irreducible $\mathcal{H}$-module is generated by a nonzero vector $v \in V$ such that $L_{m} v=0$ for all $1 \leqslant m \leqslant N$.

Proof. Set $\Phi(v)=\left\{m \mid 1 \leqslant m \leqslant N, L_{m} v \neq 0\right\}$ for $v \in V$. Let $W$ denote an irreducible $\mathcal{H}$ module and take a nonzero vector $w \in W$. If $\Phi(w)=\emptyset$, then $L_{m} w=0$ for all $1 \leqslant m \leqslant N$ and by the irreducibility of $W$, the module $W$ is generated by $w$ and so the result follows. Suppose $\Phi(w) \neq \emptyset$. Let $m=\min \Phi(w)$ and set $w^{\prime}=L_{m} w \in W$. By Proposition 33 (i) and (ii), we have $\Phi\left(w^{\prime}\right) \subsetneq \Phi(w)$. By continuing this process at most $|\Phi(w)|$ times, we get a nonzero vector $v \in W$ such that $\Phi(v)=\emptyset$. By the same argument above, the assertion holds.

Recall from Sections 8 and 9 , that there are the matrices $E_{\lambda}$ in $\mathcal{H}$ and that they turn out to be independent of the basis $v_{1}, v_{2}, \ldots, v_{N}$ for $H$ and the nontrivial character $\chi$ of
the additive group $\mathbb{F}_{q}$. By Lemma 48 and Proposition 52, it suffices to consider the module $\mathcal{H} v$ for $v \in \sum_{\mu, \lambda} E_{\mu}^{*} E_{\lambda} V$, where the sum is taken over all pairs $(\mu, \lambda)$ with $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11 such that $\lambda$ is column-full with respect to $\mu$ in Definition 23.

Proposition 53. Let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11, and assume that $\lambda$ is column-full with respect to $\mu$ in Definition 23. Recall $\kappa(m, \mu, \lambda)$ in (23). For a nonzero vector $v \in E_{\mu}^{*} E_{\lambda} V$, the $\mathcal{H}$-module $\mathcal{H} v$ has a basis

$$
w(\varepsilon) \in E_{\mu+\varepsilon}^{*} V, \quad \varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}, \ldots, \varepsilon_{N}\right), \quad \varepsilon_{m}= \begin{cases}0 & \text { if } m \in \lambda,  \tag{26}\\ 0 \text { or } 1 & \text { if } m \notin \lambda,\end{cases}
$$

on which the generators $L_{m}, R_{m}(1 \leqslant m \leqslant N)$ act as follows:

$$
\begin{aligned}
& L_{m} w(\varepsilon)=q^{\kappa(m, \mu, \lambda)-\left(\varepsilon_{1}+\cdots+\varepsilon_{m-1}\right)} w(\varepsilon-\widehat{m}), \\
& R_{m} w(\varepsilon)=q^{\varepsilon_{m+1}+\cdots+\varepsilon_{N}} w(\varepsilon+\widehat{m}),
\end{aligned}
$$

where we set $w(\varepsilon)=0$ if $\varepsilon$ is not of the form in (26).
Proof. Let $\mathcal{H}^{+}$denote the subalgebra of $\mathcal{H}$ generated by $R_{1}, R_{2}, \ldots, R_{N}$. Consider $\mathcal{H}^{+} v$, the $\mathcal{H}^{+}$-module generated by $v$. We show that $\mathcal{H}^{+} v$ is an $\mathcal{H}$-module. Let $1 \leqslant m \leqslant N$. Then $\mathcal{H}^{+} v$ is $R_{m}$-invariant by the construction and $K_{m}$-invariant by Proposition 32 (ii), (iv). In addition, $\mathcal{H}^{+} v$ is $L_{m}$-invariant by Proposition 33 (i), (iii), (iv), Lemma 46 and since $L_{m} v=0$ by Lemma 48. Since $\mathcal{H}$ is generated by $R_{m}, L_{m}$ and $K_{m}$, for $1 \leqslant m \leqslant N$, $\mathcal{H}^{+} v$ is an $\mathcal{H}$-module. Thus we have $\mathcal{H}^{+} v=\mathcal{H} v$. By Proposition 33 (i), (iii), $\mathcal{H}^{+} v$ is spanned by

$$
w(\varepsilon)=R_{N}^{\varepsilon_{N}} R_{N-1}^{\varepsilon_{N-1}} \cdots R_{1}^{\varepsilon_{1}} v
$$

for $\varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}, \ldots, \varepsilon_{N}\right) \in\{0,1\}^{N}$. By Lemma 29 (ii), $w(\varepsilon) \in E_{\mu+\varepsilon}^{*} V$. By Lemma 43, $w(\varepsilon) \neq 0$ if and only if $m \in S_{\mu}$ and $m \notin \lambda$ for all $1 \leqslant m \leqslant N$ with $\varepsilon_{m}=1$. Thus (26) forms a basis for $\mathcal{H} v$. For $1 \leqslant m \leqslant N$, the $L_{m}$-actions on $w(\varepsilon)$ follow from Proposition 33 (iii), (iv), Lemma 46 and $L_{m} v=0$. Similarly, for $1 \leqslant m \leqslant N$, the $R_{m}$-actions on $w(\varepsilon)$ follow from Proposition 33 (i), (iii). The result follows.

Proposition 54. Referring to Proposition 53, the basis (26) for $\mathcal{H} v$ satisfies the following.

$$
K_{m} w(\varepsilon)=q^{1 / 2-\left(\mu_{m}+\varepsilon_{m}\right)} w(\varepsilon),
$$

for $1 \leqslant m \leqslant N$, where $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right)$ and $\varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}, \ldots, \varepsilon_{N}\right)$.
Proof. By Proposition 53, we have $w(\varepsilon) \in E_{\mu+\varepsilon}^{*} V$. The result follows from the definition of $K_{m}$.

Theorem 55. For any irreducible $\mathcal{H}$-module $W$, there uniquely exist $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11 where $\lambda$ is column-full with respect to $\mu$, such that $W$ is generated by a nonzero vector in $E_{\mu}^{*} E_{\lambda} V$. Moreover, $W$ is determined up to isomorphism by $\mu$ and $\lambda$.

Proof. By Proposition 52, there exists a nonzero vector $v \in W$ with $L_{m} v=0$ for all $1 \leqslant m \leqslant N$ such that $W=\mathcal{H} v$. According to the direct sum decomposition in Lemma 48, we write

$$
v=\sum_{\mu, \lambda} E_{\mu}^{*} E_{\lambda} v .
$$

Since $v$ is nonzero, there exists a pair $(\mu, \lambda)$ such that $E_{\mu}^{*} E_{\lambda} v \neq 0$. By Proposition 47, $E_{\mu}^{*} E_{\lambda} v$ belongs to $W$ and so by the irreducibility of $W, E_{\mu}^{*} E_{\lambda} v$ generates $W$. Suppose there exist another pair $\left(\mu^{\prime}, \lambda^{\prime}\right)$ and a vector $v^{\prime} \in V$ such that $E_{\mu^{\prime}}^{*} E_{\lambda^{\prime}} v^{\prime}$ also generates $W$. Thus we have the two bases (26) for $W$. However, by comparing them, we obtain $\left(\mu^{\prime}, \lambda^{\prime}\right)=(\mu, \lambda)$ and the result follows.
Definition 56. Referring to Theorem 55, we call $\mu \in\{0,1\}^{N}$ the endpoint of $W$ and $\lambda \subseteq\{1,2, \ldots, N\}$ the shape of $W$.
Corollary 57. Let $\lambda \subseteq\{1,2, \ldots, N\}$ with even cardinality. For an irreducible $\mathcal{H}$-module $W$ of shape $\lambda$, we have

$$
\operatorname{dim} W=2^{N-|\lambda|} .
$$

Proof. Count the vectors in the basis (26) for $W$.
Theorem 58. For $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfying (ii) in Lemma 11 where $\lambda$ is column-full with respect to $\mu$, there exists an irreducible $\mathcal{H}$-module of endpoint $\mu$ and shape $\lambda$. Moreover, the multiplicity in $V$ is given by

$$
q^{\left|B_{\mu}\right|-n\left(\lambda \cap S_{\mu}\right)} \prod_{s \in \lambda \cap S_{\mu}}\left(q^{\rho(s, \mu, \lambda)}-1\right),
$$

where $n\left(\lambda \cap S_{\mu}\right)$ is defined in (20) and $\rho(s, \mu, \lambda)$ is defined in Lemma 15.
Proof. Take a nonzero vector $v \in E_{\mu}^{*} E_{\lambda} V$. We show that $W=\mathcal{H} v$ is irreducible. Consider an irreducible $\mathcal{H}$-module decomposition of $W$ as follows.

$$
W=W_{1}+W_{2}+\cdots+W_{r}, \quad \text { (direct sum) }
$$

for some positive integer $r \geqslant 1$. According to this decomposition, we write $v=w_{1}+w_{2}+$ $\cdots+w_{r}$ such that $w_{n} \in W_{n}(1 \leqslant n \leqslant r)$. Since this sum is direct and $v \in E_{\mu}^{*} E_{\lambda} W$, we find $w_{n} \in E_{\mu}^{*} E_{\lambda} W$ for $1 \leqslant n \leqslant r$. However, by Proposition 53, we have $\operatorname{dim} E_{\mu}^{*} E_{\lambda} W=1$. Thus, all the vectors $w_{n}(1 \leqslant n \leqslant r)$ are scalar multiples of $v$. This forces $r=1$, i.e., $W$ is irreducible.

The multiplicity of $W$ in $V$ is $\operatorname{dim} E_{\mu}^{*} E_{\lambda} V$, which is determined in Corollary 25.

## 12 The quantum affine algebra $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$

In this section, we fix a nonzero scalar $q \in \mathbb{C}$ which is not a root of unity. For $n \in \mathbb{N}$, we define

$$
[n]_{q}=\frac{q^{n}-q^{-n}}{q-q^{-1}}
$$

We recall the definition of $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$ from [1] in terms of Chevalley generators.

Definition 59 ([1, Section 2]). The quantum affine algebra $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$ is the associative $\mathbb{C}$-algebra generated by $e_{i}^{ \pm}, k_{i}, k_{i}^{-1}(i=0,1)$ with the relations

$$
\begin{array}{ll}
k_{i} k_{i}^{-1}=k_{i}^{-1} k_{i}=1, & k_{0} k_{1}=k_{1} k_{0}, \\
k_{i} e_{i}^{ \pm}=q^{ \pm 2} e_{i}^{ \pm} k_{i}, & k_{i} e_{j}^{ \pm}=q^{\mp 2} e_{j}^{ \pm} k_{i}, \quad i \neq j, \\
e_{i}^{+} e_{i}^{-}-e_{i}^{-} e_{i}^{+}=\frac{k_{i}-k_{i}^{-1}}{q-q^{-1}}, & e_{0}^{ \pm} e_{1}^{\mp}-e_{1}^{\mp} e_{0}^{ \pm}=0, \\
\left(e_{i}^{ \pm}\right)^{3} e_{j}^{ \pm}-[3]_{q}\left(e_{i}^{ \pm}\right)^{2} e_{j}^{ \pm} e_{i}^{ \pm}+[3]_{q} e_{i}^{ \pm} e_{j}^{ \pm}\left(e_{i}^{ \pm}\right)^{2}-e_{j}^{ \pm}\left(e_{i}^{ \pm}\right)^{3}=0, \quad i \neq j . \tag{30}
\end{array}
$$

We call $e_{i}^{ \pm}, k_{i}, k_{i}^{-1}(i=0,1)$ the Chevalley generators for $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$.
It is known that the quantum affine algebra $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$ has the following Hopf algebra structure. The comultiplication $\Delta$ satisfies

$$
\Delta\left(e_{i}^{+}\right)=e_{i}^{+} \otimes k_{i}+1 \otimes e_{i}^{+}, \quad \Delta\left(e_{i}^{-}\right)=e_{i}^{-} \otimes 1+k_{i}^{-1} \otimes e_{i}^{-}, \quad \Delta\left(k_{i}\right)=k_{i} \otimes k_{i} .
$$

It is also known that there exists a family of finite-dimensional irreducible $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$-modules $V_{d}(\alpha)$ for $d \in \mathbb{N}, \alpha \in \mathbb{C} \backslash\{0\}$, where $V_{d}(\alpha)$ has a basis $\left\{u_{i}\right\}_{i=0}^{d}$ satisfying

$$
\begin{array}{lrr}
e_{0}^{+} u_{i}=\alpha[i+1]_{q} u_{i+1} & (0 \leqslant i \leqslant d-1), & e_{0}^{+} u_{d}=0, \\
e_{1}^{+} u_{i}=[d-i+1]_{q} u_{i-1} & (1 \leqslant i \leqslant d), & e_{1}^{+} u_{0}=0, \\
e_{0}^{-} u_{i}=\alpha^{-1}[d-i+1]_{q} u_{i-1} & (1 \leqslant \varepsilon \leqslant d), & e_{0}^{-} u_{0}=0, \\
e_{1}^{-} u_{i}=[i+1]_{q} u_{i+1} & (0 \leqslant i \leqslant d-1), & e_{1}^{-} u_{d}=0, \\
k_{0} u_{i}=q^{2 i-d} u_{i} & (0 \leqslant i \leqslant d), & \\
k_{1} u_{i}=q^{d-2 i} u_{i} & (0 \leqslant i \leqslant d) . &
\end{array}
$$

We call $V_{d}(\alpha)$ the evaluation module for $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$ with the evaluation parameter $\alpha$. We recurrently define the algebra homomorphism $\Delta^{(N)}: U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right) \rightarrow \underbrace{U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right) \otimes \cdots \otimes U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)}_{(N+1) \text { times }}$ for $N \in \mathbb{N}$ by

$$
\begin{align*}
\Delta^{(0)} & =\mathrm{id} \\
\Delta^{(1)} & =\Delta \\
\Delta^{(N)} & =(\underbrace{\operatorname{id} \otimes \cdots \otimes \operatorname{id}}_{(N-2) \text { times }} \otimes \Delta) \circ \Delta^{(N-1)}
\end{align*}
$$

This algebra homomorphism $\Delta^{(N)}$ is called the $N$-fold comultiplication. For each $N \geqslant 1$, by the $(N-1)$-fold comultiplication $\Delta^{(N-1)}$, a tensor product of $N$ evaluation modules again becomes a $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module. More precisely, a tensor product $V_{d_{1}}\left(\alpha_{1}\right) \otimes \cdots \otimes V_{d_{N}}\left(\alpha_{N}\right)$ has a basis

$$
\begin{equation*}
u(\varepsilon)=u_{\varepsilon_{1}} \otimes \cdots \otimes u_{\varepsilon_{N}}, \quad 0 \leqslant \varepsilon_{1} \leqslant d_{1}, \quad \ldots, \quad 0 \leqslant \varepsilon_{N} \leqslant d_{N} \tag{31}
\end{equation*}
$$

on which the Chevalley generators act as follows:

$$
\begin{align*}
& e_{0}^{+} u(\varepsilon)=\sum_{m=1}^{N} \alpha_{m}\left[\varepsilon_{m}+1\right]_{q} q^{2\left(\varepsilon_{m+1}+\cdots+\varepsilon_{N}\right)-\left(d_{m+1}+\cdots+d_{N}\right)} u(\varepsilon+\widehat{m}),  \tag{32}\\
& e_{1}^{+} u(\varepsilon)=\sum_{m=1}^{N}\left[d_{m}-\varepsilon_{m}+1\right]_{q} q^{\left(d_{m+1}+\cdots+d_{N}\right)-2\left(\varepsilon_{m+1}+\cdots+\varepsilon_{N}\right)} u(\varepsilon-\widehat{m}),  \tag{33}\\
& e_{0}^{-} u(\varepsilon)=\sum_{m=1}^{N} \alpha_{m}^{-1}\left[d_{m}-\varepsilon_{m}+1\right]_{q} q^{\left(d_{1}+\cdots+d_{m-1}\right)-2\left(\varepsilon_{1}+\cdots+\varepsilon_{m-1}\right)} u(\varepsilon-\widehat{m}),  \tag{34}\\
& e_{1}^{-} u(\varepsilon)=\sum_{m=1}^{N}\left[\varepsilon_{m}+1\right]_{q} q^{2\left(\varepsilon_{1}+\cdots+\varepsilon_{m-1}\right)-\left(d_{1}+\cdots+d_{m-1}\right)} u(\varepsilon+\widehat{m}),  \tag{35}\\
& k_{0} u(\varepsilon)=q^{2\left(\varepsilon_{1}+\cdots+\varepsilon_{N}\right)-\left(d_{1}+\cdots+d_{N}\right)} u(\varepsilon),  \tag{36}\\
& k_{1} u(\varepsilon)=q^{\left(d_{1}+\cdots+d_{N}\right)-2\left(\varepsilon_{1}+\cdots+\varepsilon_{N}\right)} u(\varepsilon), \tag{37}
\end{align*}
$$

where $\varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}, \ldots, \varepsilon_{N}\right) \in \mathbb{Z}^{N}$ and we define $u(\varepsilon)=0$ if $\varepsilon$ is not of the form in (31).
Let $W$ denote a finite-dimensional irreducible $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module. By [1, Proposition 3.2], there exist scalars $\epsilon_{0}, \epsilon_{1} \in\{-1,1\}$ such that each eigenvalue of $k_{i}$ on $W$ is $\epsilon_{i}$ times an integral power of $q$ for $i=0,1$. The pair $\left(\epsilon_{0}, \epsilon_{1}\right)$ is called the type of $W$. For each pair $\epsilon_{0}, \epsilon_{1} \in\{-1,1\}$, there exists an algebra automorphism of $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$ that sends

$$
k_{i} \mapsto \epsilon_{i} k_{i}, \quad \quad e_{i}^{+} \mapsto \epsilon_{i} e_{i}^{+}, \quad \quad e_{i}^{-} \mapsto e_{i}^{-}, \quad(i=0,1)
$$

By this automorphism, any finite-dimensional irreducible $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$-module of type $\left(\epsilon_{0}, \epsilon_{1}\right)$ becomes that of type $(1,1)$.
Theorem 60 ([1, Theorem 4.11]). Every finite-dimensional irreducible $U_{q}\left(\widehat{\mathfrak{s}}_{2}\right)$-module of type $(1,1)$ is isomorphic to a tensor product of evaluation modules. Moreover, two such tensor products are isomorphic if and only if one is obtained from the other by permuting the factors in the tensor product.

With an evaluation module $V_{d}(\alpha)$, we associate the set of scalars

$$
S_{d}(\alpha)=\left\{\alpha q^{d-1}, \alpha q^{d-3}, \ldots, \alpha q^{-d+1}\right\} .
$$

The set $S_{d}(\alpha)$ is called a $q$-string of length $d$. Two $q$-strings $S_{d_{1}}\left(\alpha_{1}\right), S_{d_{2}}\left(\alpha_{2}\right)$ are said to be in general position if one of the following occurs:
(i) $S_{d_{1}}\left(\alpha_{1}\right) \cup S_{d_{2}}\left(\alpha_{2}\right)$ is not a $q$-string,
(ii) $S_{d_{1}}\left(\alpha_{1}\right) \subseteq S_{d_{2}}\left(\alpha_{2}\right)$ or $S_{d_{2}}\left(\alpha_{2}\right) \subseteq S_{d_{1}}\left(\alpha_{1}\right)$.

Moreover, several $q$-strings are said to be in general position if every two $q$-strings are in general position.
Theorem 61 ([1, Theorem 4.8]). A tensor product of evaluation modules for $U_{q}\left(\widehat{\mathfrak{s l}}_{2}\right)$ is irreducible if and only if the associated $q$-strings are in general position.

## 13 The algebra $\mathcal{H}$ and the quantum affine algebra $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}_{2}}\right)$

In this section, we get back to the subspace lattice $P$ over $\mathbb{F}_{q}$. Recall the matrices $E_{\lambda} \in \mathcal{H}$ in Sections 8 and 9. Let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11. For $v \in E_{\mu}^{*} E_{\lambda} V$ and $1 \leqslant m \leqslant N$, if $L_{m} v \neq 0$, then we have $m \in T_{\mu}$ and $m \notin \lambda$ by Lemma 42 and so $\left(L_{m} R_{m}\right) L_{m} v=q^{\kappa(m, \mu, \lambda)} L_{m} v$ by Lemma 45. Therefore, we define the matrix $\left(L_{m} R_{m}\right)^{-1} L_{m}$ by

$$
\left(L_{m} R_{m}\right)^{-1} L_{m} v= \begin{cases}q^{-\kappa(m, \mu, \lambda)} L_{m} v & \text { if } L_{m} v \neq 0  \tag{38}\\ 0 & \text { if } L_{m} v=0\end{cases}
$$

for $v \in V$. We remark that $\left(L_{m} R_{m}\right)^{-1} L_{m}$ does not mean the product of $\left(L_{m} R_{m}\right)^{-1}$ and $L_{m}$ since $L_{m} R_{m}$ is not invertible by Lemma 45. Similarly, we define the matrix $\left(R_{m} L_{m}\right)^{-1} R_{m}$ by

$$
\left(R_{m} L_{m}\right)^{-1} R_{m} v= \begin{cases}q^{-\kappa(m, \mu, \lambda)} R_{m} v & \text { if } R_{m} v \neq 0  \tag{39}\\ 0 & \text { if } R_{m} v=0\end{cases}
$$

for $v \in V$.
Theorem 62. Let $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{N}$ denote nonzero scalars. The standard module $V$ supports a $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module structure on which the Chevalley generators act as follows:

| generators | actions on $V$ |
| :---: | :---: |
| $e_{0}^{+}$ | $q^{(1-N) / 2} \sum_{m=1}^{N} \alpha_{m} R_{m}$ |
| $e_{1}^{+}$ | $q^{(N-1) / 2} \sum_{m=1}^{N}\left(L_{m} R_{m}\right)^{-1} L_{m}$ |
| $e_{0}^{-}$ | $\sum_{m=1}^{N} \alpha_{m}^{-1} L_{m}$ |
| $e_{1}^{-}$ | $\sum_{m=1}^{N}\left(R_{m} L_{m}\right)^{-1} R_{m}$ |
| $k_{0}$ | $\prod_{m=1}^{N} K_{m}^{-1}$ |
| $k_{0}^{-1}$ | $\prod_{m=1}^{N} K_{m}$ |
| $k_{1}$ | $\prod_{m=1}^{N} K_{m}$ |
| $k_{1}^{-1}$ | $\prod_{m=1}^{N} K_{m}^{-1}$ |

Here the matrices $\left(L_{m} R_{m}\right)^{-1} L_{m}$ and $\left(R_{m} L_{m}\right)^{-1} R_{m}$ are defined in (38) and in (39), respectively.
Proof. Referring to the above table, for $i=0,1$ let $\widehat{e}_{i}^{+}, \widehat{e}_{i}, \widehat{k}_{i}, \widehat{k}_{i}^{-1}$ denote the expressions to the right of $e_{i}^{+}, e_{i}^{-}, k_{i}, k_{i}^{-1}$ respectively. We show these elements $\widehat{e}_{i}^{+}, \widehat{e}_{i}, \widehat{k}_{i}, \widehat{k}_{i}^{-1}(i=0,1)$ satisfy the defining relations (27)-(30) of $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s}}_{2}\right)$ on $V$.

We first show $\widehat{e}_{i}^{+}, \widehat{e}_{i}, \widehat{k_{i}}, \widehat{k}_{i}^{-1}(i=0,1)$ satisfy the relations except the first relation in (29). They satisfy the relations in (27) by the definitions of $\widehat{k}_{i}, \widehat{k}_{i}^{-1}(i=0,1)$. They satisfy the first relation in (28) with $i=0$ by Proposition 32. They satisfy the second relation in $(28)$ with $(i, j)=(1,0)$ by Proposition 32 . Since the other relations involve $\widehat{e}_{1}^{+}, \widehat{e_{1}}$, we show them as follows. Fix a nonzero vector $v \in V$. Then we apply both sides of each defining relation to $v$ and check the results are the same. These elements $\widehat{e}_{i}^{+}, \widehat{e_{i}}, \widehat{k}_{i}, \widehat{k}_{i}^{-1}$ $(i=0,1)$ satisfy the first relation in (28) with $i=1$ by Proposition 32 . They satisfy the second relation in $(28)$ with $(i, j)=(0,1)$ by Proposition 32 . They satisfy the second relation in (29) and the relations in (30) by Proposition 33.

It remains to show that they satisfy the first relation in (29). Take a nonzero vector $v \in E_{\mu}^{*} E_{\lambda} V$ for some $\mu=\left(\mu_{1}, \mu_{2}, \ldots, \mu_{N}\right) \in\{0,1\}^{N}, \lambda \subseteq\{1,2, \ldots, N\}$. By Lemmas 45 and 46 , we have

$$
\left(\widehat{e}_{0}^{+} \widehat{e}_{0}-\widehat{e}_{0} \widehat{e}_{0}^{+}\right) v=-\left(q^{(1-N) / 2} \sum_{m}(-1)^{\mu_{m}} q^{\kappa(m, \mu, \lambda)}\right) v,
$$

where the sum is taken over all $1 \leqslant m \leqslant N$ with $m \notin \lambda$. On the other hand, by the definition of $K_{m}$, we have

$$
\left(\frac{\widehat{k}_{0}-\widehat{k}_{0}^{-1}}{q^{1 / 2}-q^{-1 / 2}}\right) v=\left(\frac{q^{|\mu|-N / 2}-q^{N / 2-|\mu|}}{q^{1 / 2}-q^{-1 / 2}}\right) v
$$

By Lemma 51 , it turns out that both scalars are the same and so $\widehat{e}_{0}^{+}, \widehat{e}_{0}^{-}, \widehat{k}_{0}, \widehat{k}_{0}^{-1}$ satisfy the first relation in (29). Similarly, $\widehat{e}_{1}^{+}, \widehat{e}_{1}^{-}, \widehat{k}_{1}, \widehat{k}_{1}^{-1}$ satisfy the first relation in (29).

Corollary 63. Let $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{N}$ denote nonzero scalars. There exists an algebra homomorphism from $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$ to $\mathcal{H}$ that sends

$$
\begin{array}{ll}
e_{0}^{+} \mapsto q^{(1-N) / 2} \sum_{m=1}^{N} \alpha_{m} R_{m}, & e_{1}^{+} \mapsto q^{(N-1) / 2} \sum_{m=1}^{N}\left(L_{m} R_{m}\right)^{-1} L_{m}, \\
e_{0}^{-} \mapsto \sum_{m=1}^{N} \alpha_{m}^{-1} L_{m}, & \\
k_{0} & \mapsto \prod_{m=1}^{N} K_{m}^{-1},
\end{array}
$$

Proof. Immediate from Proposition 62.
The algebra homomorphism in Corollary 63 turns an $\mathcal{H}$-module into a $U_{q^{1 / 2}}\left(\widehat{\mathfrak{F}}_{2}\right)$ module.

Lemma 64. Let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11 where $\lambda$ is column-full with respect to $\mu$ in Definition 23. Let $W_{\mu, \lambda}$ denote an irreducible $\mathcal{H}$ module with endpoint $\mu$ and shape $\lambda$. The basis (26) for $W_{\mu, \lambda}$ has the following actions of

Chevalley generators via the algebra homomorphism in Corollary 63.

$$
\begin{align*}
& e_{0}^{+} w(\varepsilon)=q^{(1-N) / 2} \sum_{m=1}^{N} \alpha_{m} q^{\varepsilon_{m+1}+\cdots+\varepsilon_{N}} w(\varepsilon+\widehat{m}),  \tag{40}\\
& e_{1}^{+} w(\varepsilon)=q^{(N-1) / 2} \sum_{m=1}^{N} q^{-\left(\varepsilon_{m+1}+\cdots+\varepsilon_{N}\right)} w(\varepsilon-\widehat{m}),  \tag{41}\\
& e_{0}^{-} w(\varepsilon)=\sum_{m=1}^{N} \alpha_{m}^{-1} \theta_{m}(\mu, \lambda) q^{-\left(\varepsilon_{1}+\cdots+\varepsilon_{m-1}\right)} w(\varepsilon-\widehat{m}),  \tag{42}\\
& e_{1}^{-} w(\varepsilon)=\sum_{m=1}^{N} \theta_{m}(\mu, \lambda)^{-1} q^{\varepsilon_{1}+\cdots+\varepsilon_{m-1}} w(\varepsilon+\widehat{m}),  \tag{43}\\
& k_{0} w(\varepsilon)=q^{-N / 2+|\mu|+|\varepsilon|} w(\varepsilon),  \tag{44}\\
& k_{1} w(\varepsilon)=q^{N / 2-|\mu|-|\varepsilon|} w(\varepsilon), \tag{45}
\end{align*}
$$

where $\varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}, \ldots, \varepsilon_{N}\right) \in\{0,1\}^{N}$. Here we define $w(\varepsilon)=0$ if $\varepsilon$ is not of the form in (26).

Proof. Use Propositions 53, 54 and Corollary 63.
Lemma 65. Let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11 where $\lambda$ is column-full with respect to $\mu$ in Definition 23. We define $d=\left(d_{1}, d_{2}, \ldots, d_{N}\right) \in\{0,1\}^{N}$ by

$$
d_{m}=\left\{\begin{array}{ll}
1 & \text { if } m \notin \lambda, \\
0 & \text { if } m \in \lambda,
\end{array} \quad(1 \leqslant m \leqslant N)\right.
$$

Then we have the following.
(i) $|d|=N-2|\mu|$.
(ii) If $m \notin \lambda$, then $\kappa(m, \mu, \lambda)=(N-1) / 2+\left(d_{1}+\cdots+d_{m-1}\right) / 2-\left(d_{m+1}+\cdots+d_{N}\right) / 2$ defined in (23).

Proof. (i) By the definition of $d$, we have $|d|=N-|\lambda|$. By the assumption, we have $|\lambda|=2|\mu|$ and so the result follows.
(ii) Assume $m \notin \lambda$. Observe that

$$
\left|S_{\mu}(m-1) \backslash \lambda\right|=d_{1}+\cdots+d_{m-1}, \quad\left|T_{\mu}(m+1) \backslash \lambda\right|=0
$$

By the definition of $d$,

$$
|\lambda| / 2=N / 2-\left(d_{1}+\cdots+d_{N}\right) / 2 .
$$

Hence the result follows from the above comments and $d_{m}=1$.

Theorem 66. Let $\mu \in\{0,1\}^{N}$ and $\lambda \subseteq\{1,2, \ldots, N\}$ satisfy (ii) in Lemma 11 where $\lambda$ is column-full with respect to $\mu$ in Definition 23. Let $W_{\mu, \lambda}$ denote an irreducible $\mathcal{H}$-module with endpoint $\mu$ and shape $\lambda$. Then by the algebra homomorphism in Corollary 63, $W_{\mu, \lambda}$ becomes a $U_{q^{1 / 2}}\left(\widehat{\mathfrak{s l}}_{2}\right)$-module and we have the following.
(i) $W_{\mu, \lambda}$ has type $(1,1)$.
(ii) $W_{\mu, \lambda}$ is isomorphic to the tensor product of $V_{1}\left(\alpha_{m}\right)$, where $1 \leqslant m \leqslant N$ such that $m \notin \lambda$.

Proof. (i) This follows from (44) and (45).
(ii) Recall $\left(d_{1}, d_{2}, \ldots, d_{N}\right) \in\{0,1\}^{N}$ from Lemma 65. It suffices to show that

$$
W_{\mu, \lambda} \simeq V_{d_{1}}\left(\alpha_{1}\right) \otimes \cdots \otimes V_{d_{N}}\left(\alpha_{N}\right)
$$

Recall the basis $w(\varepsilon)$ in (26) for $W_{\mu, \lambda}$ and the basis $u(\varepsilon)$ in (31) for $V_{d_{1}}\left(\alpha_{1}\right) \otimes \cdots \otimes V_{d_{N}}\left(\alpha_{N}\right)$, where $\varepsilon=\left(\varepsilon_{1}, \varepsilon_{2}, \ldots, \varepsilon_{N}\right) \in\{0,1\}^{N}$ such that $w(\varepsilon)=0$ and $u(\varepsilon)=0$ if $d_{m}<\varepsilon_{m}$ for some $1 \leqslant m \leqslant N$. We define a linear map $\varphi$ from $V_{d_{1}}\left(\alpha_{1}\right) \otimes \cdots \otimes V_{d_{N}}\left(\alpha_{N}\right)$ to $W_{\mu, \lambda}$ that sends $u(\varepsilon)$ to $\gamma(\varepsilon) w(\varepsilon)$, where

$$
\gamma(\varepsilon)=q^{|\varepsilon|(1-N) / 2} \prod_{m \in T_{\varepsilon}} q^{\left(d_{m+1}+\cdots+d_{N}\right) / 2}
$$

We check $\varphi$ preserves the actions of Chevalley generators. Observe that

$$
\begin{equation*}
\gamma(\varepsilon)=q^{(N-1) / 2} q^{-\left(d_{m+1}+\cdots+d_{N}\right) / 2} \gamma(\varepsilon+\widehat{m}) \tag{46}
\end{equation*}
$$

for $\varepsilon \in\{0,1\}^{N}$.
By (36) and (44) and Lemma 65 (i), $\varphi$ preserves the action of $k_{0}$. By (37) and (45) and Lemma 65 (i), $\varphi$ preserves the action of $k_{1}$. By (32), (40) and (46), the map $\varphi$ preserves the action of $e_{0}^{+}$. By (33), (41) and (46), the map $\varphi$ preserves the action of $e_{1}^{+}$. By (34), (42), (46) and Lemma 65 (ii), the map $\varphi$ preserves the action of $e_{0}^{-}$. By (35), (43), (46) and Lemma 65 (ii), the map $\varphi$ preserves the action of $e_{1}^{-}$.

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