# Kazhdan-Lusztig Polynomials and Cells for Affine Weyl Groups and Unequal Parameters 

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## Department of Mathematics

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#### Abstract

Let $W_{a}$ be an affine Weyl group with a set $S_{a}$ of simple reflections and parameter set $\left\{c_{s} \mid s \in S_{a}\right\}$. We study Kazhdan-Lusztig polynomials and cells for the cases that not all parameters are equal.

Denote by $\mathcal{H}$ the generic Hecke algebra corresponding to $\left(W_{a}, S_{a}\right)$ and $\left\{c_{s} \mid s \in S_{a}\right\}$. We show the existence of a canonical basis for a certain $\mathcal{H}$-module $\mathcal{M}^{0}$. The coefficients of the basis elements are generically inverses of the Kazhdan-Lusztig polynomials. We establish a formula for Kazhdan-Lusztig polynomials in terms of certain alcove polynomials. We also obtain a formula involving an analogue of Kostant's partition function.

We explicitly describe the lowest generalized two-sided cell. We find reduced expressions for its elements, provide a geometric interpretation of this cell, and we give a description in terms of a numerical function $\boldsymbol{a}$ on $W_{a}$. We also prove that the lowest generalized two-sided cell consists of at most $\left|W_{0}\right|$ generalized left cells where $W_{0}$ denotes the finite Weyl group corresponding to $W_{a}$. For parameters coming from graph automorphisms, we show that this bound is exact. For these parameters we also characterize all generalized left cells.


Thesis Supervisor: George Lusztig
Title: Professor of Mathematics

To my parents

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## Introduction

The concept of cells for an arbitrary Coxeter system ( $W, S$ ) was introduced by Kazhdan and Lusztig in [10]. They defined left, right, and two-sided cells. The definition involves a canonical basis of the generic Hecke algebra $\mathcal{H}$ of $W$. The coefficients of the canonical basis elements with respect to a standard basis of $\mathcal{H}$ are the Kazhdan-Lusztig polynomials.

Cells in Coxeter groups have been intensively studied. They are important because they give rise to representations of $\mathcal{H}$. For example, if $W$ is the finite Weyl group of type $A_{r}, r \in \mathbb{N}$, the representations attached to the left cells of $W$ are all the irreducible representations of $\mathcal{H}$, and similarly for right cells.

In [16], Lusztig extends the concept of cells to Coxeter systems ( $W, S$ ) equipped with a parameter set $\left\{c_{s} \in \mathbb{N} \mid s \in S\right\}$ such that $c_{s}=c_{s^{\prime}}$ if $s$ and $s^{\prime}$ are conjugate. We refer to these cells as generalized left, right and two-sided cells. If all parameters $c_{s}$ are equal, we get the same cells as before.

Generalized cells give rise to representations of the generic Hecke algebra corresponding to $(W, S)$ and parameter set $\left\{c_{s} \mid s \in S\right\}$. If $W$ is a finite resp. affine Weyl group the representation theory of corresponding Hecke algebras is very relevant for the representation theory of reductive groups over finite resp. $p$-adic fields.

In this thesis, we are mostly concerned with Hecke algebras corresponding to affine Weyl groups. Let $W_{a}$ be the affine Weyl group with a set $S_{a}$ of simple reflections and parameter set $\left\{c_{s} \mid s \in S_{a}\right\}$. Denote by $\mathcal{H}$ the corresponding Hecke algebra over the ring $\mathcal{A}=\mathbb{Z}\left[v, v^{-1}\right]$ where $v$ is an indeterminate. Let $W_{0}$ be the finite Weyl group corresponding to $W_{a}$.

For equal parameters, the cells in $W_{a}$ have been explicitly described for type $\tilde{A}_{r}, r \in$ $\mathbb{N}($ see $[27],[17])$, rank 2,3 (see [18], [2], [6]), and types $\tilde{B}_{4}, \tilde{C}_{4}, \tilde{D}_{4}$ (see [29], [30], [5]), and a lot of important results are obtained in [18]-[21]. In particular, Lusztig shows that the two-sided cells of $W_{a}$ are in bijective correspondence with the unipotent conjugacy classes in a simple algebraic group over $\mathbb{C}$ of type dual to that of $W_{0}$.

Several problems regarding Kazhdan-Lusztig polynomials and cells arise for arbitrary parameters but have only been solved for equal parameters. In this thesis, we solve some of these problems for unequal parameters.

The following is an outline of the contents of this thesis.
Chapter 1 contains background material on affine Weyl groups and affine Hecke algebras. In particular, the affine Weyl group $W_{a}$ is realized as a group of affine motions of a Euclidean space $V$ as well as in terms of certain alcoves in $V$. Let $X$ be this set of alcoves.

In Chapter 2, we consider a certain $\mathcal{H}$-module $\mathcal{M}^{0}$. Generalizing [14], we prove the existence of an $\mathcal{A}$-basis $\left\{D_{C} \mid C \in X\right\}$ whose characterization is similar to the characterization of the canonical basis for $\mathcal{H}$. We show that the coefficients of the elements $D_{C}, C \in X$, expressed with respect to $X$ are generically inverses of the KazhdanLusztig polynomials and that they can be calculated by a finite induction. We obtain patterns which generalize Jantzen's generic decomposition patterns for Weyl modules of a simply connected almost simple algebraic group over an algebraically closed field of prime characteristic $p$ (for $p$ large).

In Chapter 3, we prove formulas involving Kazhdan-Lusztig polynomials, their inverses, and alcove polynomials. (Alcove polynomials generically equal KazhdanLusztig polynomials.) This generalizes work by Kato and Andersen ([9] and [1]). The main statement is a formula which expresses Kazhdan-Lusztig polynomials for dominant elements in terms of alcove polynomials. We also find a way to express certain Kazhdan-Lusztig polynomials in terms of an analogue of Kostant's partition function.

Chapter 4 deals with generalized cells. We first consider the lowest generalized twosided cell. This cell contains nearly all elements of $W_{a}$. We give different descriptions, and show that this cell consists of at most $\left|W_{0}\right|$ generalized left cells. (References for the equal parameter case are given in Chapter 4.) We then study the case that the parameters $c_{s}, s \in S_{a}$, come from a graph automorphism. In this situation, the coefficients of the Kazhdan-Lusztig polynomials and of the structure constants of the canonical basis of $\mathcal{H}$ can be interpreted in terms of intersection cohomology sheaves. We derive a characterization of all generalized left cells. This characterization implies that there are only finitely many generalized left cells. We also show that for parameters coming from a graph automorphism, the lowest generalized two-sided cell consists of exactly $\left|W_{0}\right|$ generalized left cells. The results in this chapter partially answer questions raised in [31].

## 1. Preliminaries

In this chapter, we collect some basic material about affine Weyl groups and affine Hecke algebras which will be needed later on. The exposition follows [9], [14] and [31]. For more details and proofs, we refer to these publications.
1.1. The affine Weyl group. Let $V$ be a Euclidean space of finite dimension $r \geq 1$. Let $\Phi \subset V$ be an irreducible root system of $\operatorname{rank} r$ and $\check{\Phi} \subset V^{*}$ the dual root system. We denote the coroot corresponding to $\alpha \in \Phi$ by $\check{\alpha}$, and we write $\langle x, y\rangle$ for the value of $y \in V^{*}$ at $x \in V$.

Let $Q$ be the root lattice and $P$ the weight lattice. The Weyl group $W_{0}$ of $\Phi$ acts on $Q$ and $P$ (on the left), so we can form the semidirect products

$$
W_{a}=W_{0} \propto Q
$$

and

$$
\widetilde{W}_{a}=W_{0} \ltimes P .
$$

The group $W_{a}$ is the affine Weyl group of type $\check{\Phi}$, and $\widetilde{W}_{a}$ is the extended affine Weyl group associated to a simply connected, simple algebraic group $G$ (over $\mathbb{C}$ ) of root system $\Phi$. When $\lambda \in P$ is regarded as an element of $\widetilde{W}_{a}$, we will also write $p_{\lambda}$ instead of $\lambda$. The reflection in $W_{0}$ along the hyperplane orthogonal to $\alpha \in \Phi$ will be denoted by $s_{\alpha}$.

Geometrically, $W_{a}$ can be described as follows. (We will not distinguish between $V$ and the underlying affine space.) Fix a set of positive roots $\Phi^{+} \subset \Phi$, and let $\Pi \subseteq \Phi^{+}$ be the set of simple roots. For $\alpha \in \Phi^{+}$and $n \in \mathbb{Z}$, we define a hyperplane

$$
H_{\alpha, n}=\{x \in V \mid\langle x, \check{\alpha}\rangle=n\}
$$

and write $\sigma_{\alpha, n}=\sigma_{H_{\alpha, n}}$ for the reflection along $H_{\alpha, n}$. Let $\alpha_{0} \in \Phi^{+}$be such that $\check{\alpha}_{0}$ is the highest coroot in $\check{\Phi}$. Mapping $s_{\alpha}, \alpha \in \Pi$, to $\sigma_{\alpha, 0}$ and $s_{0} \stackrel{\text { def }}{=} p_{\alpha_{0}} s_{\alpha_{0}}$ to $\sigma_{\alpha_{0}, 1}$ establishes an isomorphism from $W_{a}$ to the group $\Omega$ generated by $\sigma_{\alpha, 0}, \alpha \in \Pi$, and $\sigma_{\alpha_{0}, 1}$.

Denote the set of simple reflections in $W_{0}$ by $S_{0}$. The group $W_{a}$ is a Coxeter group with generating set $S_{a}=S_{0} \cup\left\{s_{0}\right\}$.

The extended affine Weyl group $\widetilde{W}_{a}$ can be written as the semidirect product $N \ltimes W_{a}$ where $N$ is the normalizer of $S_{a}$ in $\widetilde{W}_{a}$.

We also need the following realization of $W_{a}$ (cf. [14]). Let

$$
\mathcal{F}=\left\{H_{\alpha, n} \mid \alpha \in \Phi^{+}, n \in \mathbb{Z}\right\}
$$

and let $X$ be the set of connected components of $V-\bigcup_{H \in \mathcal{F}} H$. The elements of $X$ are called alcoves. The group $\Omega$ acts on the set of faces of alcoves, and we denote the set of $\Omega$-orbits by $S_{a}^{\prime}$. If $f$ is a face contained in the orbit $t \in S_{a}^{\prime}$, we say $f$ is of type $t$. For $A \in X$ and $t \in S_{a}^{\prime}$, there is a unique alcove $t A \in X, t A \neq A$, such that $t A$ shares with $A$ its face of type $t$. The involutions $\sigma_{t}: A \mapsto t A$ on $X$ for $t \in S_{a}^{\prime}$ generate a group $W_{a}^{\prime}$. There is an isomorphism from $W_{a}$ to $W_{a}^{\prime}$, which can be described as follows. Let

$$
A^{+}=\left\{x \in V \mid\langle x, \check{\alpha}\rangle>0 \text { for all } \alpha \in \Pi,\left\langle x, \check{\alpha}_{0}\right\rangle<1\right\} .
$$

For $s=s_{\alpha}, \alpha \in \Pi$, resp. $s=s_{0}$, the hyperplane $H_{\alpha, 0}$ resp. $H_{\alpha_{0}, 1}$ contains a unique face of $A^{+}$, whose orbit in $S_{a}^{\prime}$ we denote by $t_{s}$. The isomorphism sends $s \in S_{a}$ to the involution $\sigma_{t_{s}}$.

Identifying $W_{a}$ with $\Omega$ yields an action of $W_{a}$ on $V$ and thereby on $X$, which we consider as a right action. We also identify $W_{a}$ with $W_{a}^{\prime}$ from now on and write the action of $W_{a}$ on $X$ resulting from this identification on the left. The two actions of $W_{a}$ on $X$ can be seen to commute and to be simply transitive.

We fix parameters $c_{s} \in \mathbb{N}$ for $s \in S_{a}$ such that $c_{s}=c_{t}$ whenever $s$ and $t$ are conjugate in $W_{a}$.

Lemma 1.1.1. Let $H$ be a hyperplane in $\mathcal{F}$, and suppose $H$ supports faces of types $s, t \in S_{a}$. Then $s$ and $t$ are conjugate in $W_{a}$.

Proof. The assumptions imply that there are alcoves $A, A^{\prime} \in X$ such that $s A=$ $A \sigma_{H}$ and $t A^{\prime}=A^{\prime} \sigma_{H}$. Because of the transitivity of the left action of $W_{a}$ on $X$, we can find an element $w \in W_{a}$ such that $A^{\prime}=w A$. We have

$$
t w A=t A^{\prime}=A^{\prime} \sigma_{H}=w A \sigma_{H}=w s A
$$

and hence $t w=w s$, i.e. $s$ and $t$ are conjugate via $w$.

As a consequence of this lemma, we can associate a parameter $c_{H} \in \mathbb{N}$ to $H \in \mathcal{F}$ where $c_{H}=c_{s}$ if $H$ supports a face of type $s$.

For a 0 -dimensional facet $\lambda$ of an alcove, we define

$$
m(\lambda)=\sum_{H \in \mathcal{F}, H \ni \lambda} c_{H}
$$

and we call $\lambda$ a special point if $m(\lambda)$ is maximal. Note that, in general, the set of 0 -dimensional facets of alcoves contains $P$ as a proper subset.

Let $T \subset V$ be the set of all special points. If all parameters are equal, the notion of special points coincides with the notion in [14], so $T=P$ and $m(\lambda)=\left|\Phi^{+}\right|$for $\lambda \in T$. The next lemma will enable us to determine $T$ in all cases.

Let $\Gamma$ be the Coxeter graph of $\left(W_{a}, S_{a}\right)$, and identify the set of vertices of $\Gamma$ with $S_{a}$. If $\Gamma$ is of type $\tilde{A}_{1}$ or $\tilde{C}_{r}, r \geq 2$, there is a unique nontrivial automorphism ${ }^{\sim}$ on $\Gamma$.

Lemma 1.1.2. Let $H, H^{\prime}$ be parallel hyperplanes in $\mathcal{F}$ and let $s, s^{\prime} \in S_{a}$. If $H$ supports a face of type $s$ and $H^{\prime}$ supports a face of type $s^{\prime}$, we either have
(i) $\Gamma$ is of type $\tilde{A}_{1}$ or $\tilde{C}_{r}, r \geq 2$, and $\left\{s, s^{\prime}\right\}=\left\{s_{0}, \tilde{s}_{0}\right\}$ or
(ii) $s$ and $s^{\prime}$ are conjugate in $W_{a}$.

Proof. Suppose (i) does not hold (and $s \neq s^{\prime}$ ). W.l.o.g. we can assume that there exists an element $t \in S_{a}$ such that $(s t)^{3}=\mathrm{id}$.

Let $A$ be an alcove having its face of type $s$ on $H$, and let $H^{\prime \prime}$ be the hyperplane containing the face of $A$ of type $t$. Then $H^{\prime \prime}$ intersects $H$ and hence $H^{\prime}$ at an angle $\pm \frac{\pi}{3}$, which implies $H^{\prime \prime} \sigma_{H^{\prime} \sigma_{H^{\prime \prime}}}=H^{\prime}$. Therefore, $H^{\prime}$ supports a face of type $t$, and $s^{\prime} \sim t$ by Lemma 1.1.1. Since $(s t s) s(s t s)=t$, we have $s \sim t$. Thus $s^{\prime} \sim s$.

Throughout this paper, we refer to the situation in which the Coxeter graph $\Gamma$ is of type $\tilde{A}_{1}$ or $\tilde{C}_{r}, r \geq 2$, and $c_{s_{0}} \neq c_{\tilde{s}_{0}}$, as Case 1 and all other situations as Case 2.

In Case 1 , let $\lambda_{r}$ be the fundamental weight such that $P$ is generated by $Q$ and $\lambda_{r}$.
Claim 1.1.3. We have $T=Q$ or $T=\lambda_{r}+Q$ in Case 1 and $T=P$ in Case 2.
We first notice that if $\lambda$ is a special point and $\mu \in Q$, the point $\lambda+\mu=\lambda p_{\mu}$ is a special point as well.

Next, since according to the definition of the weight lattice, $P$ consists of all points $\lambda \in V$ that lie in the intersection of $\left|\Phi^{+}\right|$hyperplanes in $\mathcal{F}$, we have $T \subseteq P$.

Now suppose we are in Case 1 and $\Gamma$ is of type $\tilde{C}_{r}$. (For the following data about roots and weights see e.g. [4], Ch. VI.) Take an orthonormal basis $\left\{e_{1}, \ldots, e_{r}\right\}$ of $V$ and write $\alpha_{i}=e_{i}-e_{i+1}, 1 \leq i \leq r-1$, and $\alpha_{r}=e_{r}$ for the simple roots in II. Then $\lambda_{r}=\frac{1}{2}\left(\alpha_{1}+2 \alpha_{2}+\cdots+r \alpha_{r}\right)=\frac{1}{2}\left(e_{1}+\cdots+e_{r}\right)$ and $\alpha_{0}=\alpha_{1}+\cdots+\alpha_{r}=e_{1}$. Hence $\left\langle\lambda_{r}, \check{\alpha}_{0}\right\rangle=1$, i.e. $\lambda_{r} \in H_{\alpha_{0}, 1}$. More generally, we have $\lambda_{r} \in H_{\alpha, 1}$ for all short roots $\alpha \in \Phi^{+}$. Since each $H_{\alpha, 1}, \alpha \in \Phi^{+}, \alpha$ short, supports a face of type $s_{0}$, we conclude $T=Q$ if $c_{s_{0}}<c_{\tilde{s}_{0}}$ and $T=\lambda_{r}+Q$ if $c_{s_{0}}>c_{\tilde{s}_{0}}$. (The case $\Gamma$ of type $\tilde{A}_{1}$ follows by a simpler computation.)

In Case 2, parallel hyperplanes have the same parameter, so the special points are the same as in [14].

For the remainder of this paper, we assume that in Case 1 we have $c_{s_{0}}<c_{\tilde{s}_{0}}$, so $T=Q$. (We can always make this true by labeling the simple reflections accordingly.)

We define

$$
\widehat{W}_{a}=W_{0} \ltimes T
$$

The right action of $W_{a}$ on $X$ naturally extends to a right action of $\widehat{W}_{a}$ on $X$.
We need some more notation. If $w \in W_{a}$ has a reduced expression $w=s_{1} s_{2} \cdots s_{n}$, $s_{i} \in S_{a}$ for $1 \leq i \leq n$, we set

$$
m(w)=\sum_{i=1}^{n} c_{s_{i}}
$$

which is independent of the chosen reduced expression. (This follows from the fact that any two reduced expressions for $w$ can be transformed into each other by a sequence of braid relations. See e.g. [4], Ch. IV.) We remark that the length function $l$ on $W_{a}$ and the function $m$ on $W_{a}$ extend to functions on $\widetilde{W}_{a}$ via $l(n w)=l(w)$ and $m(n w)=m(w)$ for $n \in N$ and $w \in W_{a}$.

Let $W_{\lambda}, \lambda \in T$, be the stabilizer of the set of alcoves containing $\lambda$ in their closure with respect to the left action of $W_{a}$. (The definition of $W_{0}$ is consistent.) It can be shown that this group is a maximal parabolic subgroup of $W_{a}$ and that $m\left(w_{\lambda}\right)=m(\lambda)$ for the longest element $w_{\lambda} \in W_{\lambda}$.

A hyperplane $H=H_{\alpha, n} \in \mathcal{F}$ divides $V-H$ into the two parts

$$
V_{H}^{+}=\{x \in V \mid\langle x, \check{\alpha}\rangle>n\}
$$

and

$$
V_{H}^{-}=\{x \in V \mid\langle x, \check{\alpha}\rangle<n\} .
$$

For $\lambda \in T$, a quarter with vertex $\lambda$ is a connected component of

$$
V-\bigcup_{H \in \mathcal{F}, H \ni \lambda} H
$$

Hyperplanes which are adjacent to a quarter $\mathcal{C}$ are called walls of $\mathcal{C}$. The quarter

$$
\bigcap_{H \in \mathcal{F}, H \ni \lambda} V_{H}^{+}
$$

will be denoted by $\mathcal{C}_{\lambda}^{+}$, and $A_{\lambda}^{+}$is the unique alcove in $\mathcal{C}_{\lambda}^{+}$such that $\lambda$ lies in the closure $\overline{A_{\lambda}^{+}}$. Set $A_{\lambda}^{-}=w_{\lambda} A_{\lambda}^{+}$, and let $\mathcal{C}_{\lambda}^{-}$be the quarter with vertex $\lambda$ containing $A_{\lambda}^{-}$. For $\lambda=0$, we also write $\mathcal{C}^{+}, A^{+}$etc. Let $v_{\lambda} \in W_{a}$ be such that $A_{\lambda}^{+} v_{\lambda}=A_{\lambda}^{-}$.

Let $\mathcal{F}^{*}$ be the set of hyperplanes $H \in \mathcal{F}$ such that $H$ is a wall of $\mathcal{C}_{\lambda}^{+}$for some $\lambda \in T$. The connected components of $V-\bigcup_{H \in \mathcal{F} *} H$ will be called boxes. We denote by $\Pi_{\lambda}$ the box containing $A_{\lambda}^{+}$for $\lambda \in T$.

To an alcove $A \in X$, we associate the subset $\mathcal{L}(A)$ of $S_{a}$ containing all $s \in S_{a}$ such that, if $H \in \mathcal{F}$ is the hyperplane supporting the face of type $s$ of $A$, then $A \subset V_{H}^{+}$, $s A \subset V_{H}^{-}$.

We define two integers $d(A, B)$ and $c(A, B)$ for $A, B \in X$. Consider the set of hyperplanes $H \in \mathcal{F}$ separating $A$ from $B$. For each such hyperplane, we set $\varepsilon_{H}=1$ if $A \subset V_{H}^{-}, B \subset V_{H}^{+}$and $\varepsilon_{H}=-1$ if $A \subset V_{H}^{+}, B \subset V_{H}^{-}$. The integer $d(A, B)$ is the sum of all $\varepsilon_{H}$, and $c(A, B)$ is the sum of all $\varepsilon_{H} c_{H}$.

A length function on $X$ is a function $\delta: X \rightarrow \mathbb{Z}$ such that

$$
d(A, B)=\delta(B)-\delta(A)
$$

for all $A, B \in X$. Similarly, we call a function $\gamma: X \rightarrow \mathbb{Z}$ a weighted length function on $X$ if it satisfies

$$
c(A, B)=\gamma(B)-\gamma(A)
$$

for all $A, B \in X$.
We have the following partial order on $X$ (cf. [14]). For $A, B \in X$, we say $A \leq B$ if there exists a sequence $A=A_{0}, A_{1}, \ldots, A_{n}=B$ of alcoves such that $d\left(A_{i-1}, A_{i}\right)=1$ and $A_{i}=A_{i-1} \sigma_{H_{i}}$ for some $H_{i} \in \mathcal{F}$ for all $1 \leq i \leq n$.
1.2. The affine Hecke algebra. Let $\mathcal{A}=\mathbb{Z}\left[v, v^{-1}\right]$ be the ring of Laurent polynomials in an indeterminate $v$, and set $\mathcal{A}^{+}=\mathbb{Z}[v]$. The generic Hecke algebra $\widehat{\mathcal{H}}$ of $\widehat{W}_{a}$ with parameters $c_{s}, s \in S_{a}$, can be defined as a free $\mathcal{A}$-module with basis $\left\{T_{w} \mid w \in \widehat{W}_{a}\right\}$ and relations

$$
\left(T_{s}-v^{2 c_{s}}\right)\left(T_{s}+1\right)=0
$$

for $s \in S_{a}$ and

$$
T_{w} T_{w^{\prime}}=T_{w w^{\prime}}
$$

for $w, w^{\prime} \in \widehat{W}_{a}$ with $l(w)+l\left(w^{\prime}\right)=l\left(w w^{\prime}\right)$. We denote by $\mathcal{H}$ the subalgebra of $\widehat{\mathcal{H}}$ generated by $T_{s}$ for $s \in S_{a}$. Note that in Case 1 , the algebras $\mathcal{H}$ and $\widehat{\mathcal{H}}$ coincide.

We have

$$
T_{s}^{-1}=v^{-2 c_{s}} T_{s}+v^{-2 c_{s}}-1
$$

for $s \in S_{a}$ and

$$
T_{n}^{-1}=T_{n^{-1}}
$$

for $n \in N$, which implies the invertibility of all elements $T_{w}, w \in W_{a}$ resp. $w \in \widehat{W}_{a}$, in $\mathcal{H}$ resp. $\widehat{\mathcal{H}}$.

Let ${ }^{-}: \widehat{\mathcal{H}} \rightarrow \widehat{\mathcal{H}}$ be the unique involutive automorphism sending $v$ to $v^{-1}$ and $T_{w}$ to $T_{w^{-1}}^{-1}, w \in \widehat{W}_{a}$.

For $w \in W_{a}$, there exists a unique element $C_{w}^{*} \in \mathcal{H}$ such that $\overline{C_{w}^{*}}=v^{-2 m(w)} C_{w}^{*}$ and

$$
C_{w}^{*}=\sum_{y \leq w} P_{y, w} T_{y}
$$

where the degree of $P_{y, w} \in \mathcal{A}^{+}$is smaller than $m(w)-m(y)$ for $y<w$ and $P_{w, w}=1$. (This is proven in [16] for $w \in W_{a}$. The extension to $\widehat{W}_{a}$ is straightforward.)

The polynomials $P_{y, w}$ for $y, w \in \widehat{W}_{a}$ are the Kazhdan-Lusztig polynomials. It is easily seen that the elements $C_{w}^{*}$ for $w \in \widehat{W}_{a}$ form an $\mathcal{A}$-basis for $\widehat{\mathcal{H}}$.

Let $\mathcal{M}$ be the free $\mathcal{A}$-module on $X$. There is a unique $\mathcal{H}$-module structure on $\mathcal{M}$ such that

$$
T_{s} A=\left\{\begin{array}{lcc}
s A & \text { if } & s \notin \mathcal{L}(A) \\
v^{2 c_{s}} s A+\left(v^{2 c_{s}}-1\right) A & \text { if } & s \in \mathcal{L}(A)
\end{array}\right.
$$

for $A \in X$ and $s \in S_{a}$ (comp. [31]).
For $\lambda \in T$, we set

$$
e_{\lambda}=\sum_{A \in X, \bar{A} \ni \lambda} A \in \mathcal{M}
$$

and denote by $\mathcal{M}_{\lambda}$ the $\mathcal{H}$-submodule of $\mathcal{M}$ generated by $e_{\lambda}$. (See Remark 2.3 .2 for a motivation of this definition.) The $\mathcal{H}$-submodule of $\mathcal{M}$ generated by all elements $e_{\lambda}$ will be called $\mathcal{M}^{0}$.

Let $T^{+}=T \cap \overline{\mathcal{C}^{+}}$be the set of dominant weights in $T$. For $\lambda \in T$, we choose an element $\mu \in T^{+}$such that $\lambda+\mu \in T^{+}$and set

$$
\check{T}_{\lambda}=T_{\lambda+\mu} T_{\mu}^{-1} .
$$

This is a well-defined element in $\widehat{\mathcal{H}}$ (cf. [22]). For $w \in \widehat{W}_{a}$ of the form $w=x p_{\lambda}$, $x \in W_{0}, \lambda \in T$, we set

$$
\check{T}_{w}=T_{x} \check{T}_{\lambda}
$$

Proposition 1.2.1. (comp. [9], Proposition 1.10)
(i) The set $\left\{\check{T}_{w} \mid w \in W_{a}\right\}$ is an $\mathcal{A}$-basis for $\mathcal{H}$.
(ii) The $\operatorname{map} \phi: \mathcal{H} \rightarrow \mathcal{M}$ sending $f \in \mathcal{H}$ to $f\left(A^{-}\right) \in \mathcal{M}$ is an $\mathcal{H}$-module isomorphism. (Consider $\mathcal{H}$ as a left $\mathcal{H}$-module.)

For the proof, we refer the reader to the proof of Proposition 1.10 in loc. cit. The generalization of the arguments used there is straightforward.

We remark that $\phi^{-1}$ sends $A^{-} w$ to $\check{T}_{w}$ for $w \in W_{a}$.

In loc. cit., Kato introduces the generic length function $g: \widetilde{W}_{a} \rightarrow \mathbb{Z}$ as follows. Let ht $: Q \rightarrow \mathbb{Z}$ be the linear function satisfying ht $(\alpha)=1$ for all $\alpha \in \Pi$. This function uniquely extends to a linear function ht : P $\rightarrow \frac{1}{2} \mathbb{Z}$. The generic length function is given by

$$
g\left(x p_{\lambda}\right)=l(x)+2 \operatorname{ht}(\lambda)
$$

for $x \in W_{0}, \lambda \in P$. Kato shows that this function satisfies

$$
g(w)=d\left(A^{-}, A^{-} w\right)
$$

for $w \in \widetilde{W}_{a}$.
Similarly, let $\tilde{\mathrm{ht}}: T \rightarrow \frac{1}{2} \mathbb{Z}$ be the linear function satisfying $\tilde{\mathrm{ht}}(\alpha)=c_{s_{\alpha}}$ for $\alpha \in \Pi$. We set

$$
h\left(x p_{\lambda}\right)=m(x)+2 \tilde{h t}(\lambda)
$$

for $x \in W_{0}, \lambda \in T$. Arguments analogous to the ones in loc. cit. then show that

$$
h(w)=c\left(A^{-}, A^{-} w\right)
$$

for $w \in \widehat{W}_{a}$. We call the function $h: \widehat{W}_{a} \rightarrow \mathbb{Z}$ the generic weighted length function. The restriction of $h$ to $T$ will again be denoted by $h$. Note that $h: T \rightarrow \mathbb{Z}$ is linear.

For $\lambda \in T$, we set

$$
\hat{T}_{\lambda}=v^{-h(\lambda)} \check{T}_{\lambda}
$$

(As shown in [22], the sets $\left\{T_{w} \hat{T}_{\lambda} \mid w \in W_{0}, \lambda \in T\right\}$ and $\left\{\hat{T}_{\lambda} T_{w} \mid w \in W_{0}, \lambda \in T\right\}$ are $\mathcal{A}$-bases of $\widehat{\mathcal{H}}$; for equal parameters, the elements $\hat{T}_{\lambda}$ were introduced by Bernstein.)

## 2. A CANONICAL BASIS FOR $\mathcal{M}^{0}$

The main result of this chapter is Theorem 2.3.1, which generalizes the main theorem of [14] to unequal parameters $c_{s}, s \in S_{a}$. We follow the proof in loc.cit. Sections 2.1 and 2.2 deal with those parts of the proof whose generalization is not straightforward. Section 2.3 also contains the generalization of several other statements in loc. cit.
2.1. Intertwining operators. Our main goal in this section is to define an element $\theta_{w} \in \operatorname{End}_{\mathcal{H}}\left(\mathcal{M}^{0}\right)$ for $w \in W_{a}$. We do this in Proposition 2.1.3.

For $\lambda \in T$, we set

$$
d_{\lambda}=T_{n\left(p_{\lambda}\right)^{-1}}\left(\sum_{x \in W_{0}} T_{x}\right) \check{T}_{\lambda}
$$

where $n(w)$ for $w \in \widehat{W_{a}}$ denotes the element in $N$ such that $n(w)^{-1} w \in W_{a}$. The elements $d_{\lambda}$ lie in $\mathcal{H}$ (cf. [9] for equal parameters; the general case is completely analogous). Also, the arguments in loc.cit. show that $\phi\left(d_{\lambda}\right)=e_{\lambda}$.

The $\mathcal{H}$-submodule of $\mathcal{H}$ generated by all elements $d_{\lambda}$ will be called $\mathcal{H}^{0}$. We first define an element $\Theta_{w} \in \operatorname{End}_{\mathcal{H}}\left(\mathcal{H}^{0}\right)$ for $w \in W_{a}$. For equal parameters, this has been done in loc. cit.

We extend the definition of the automorphism ~ by saying that ${ }^{\sim}$ is the identity on the Coxeter graph $\Gamma$ (and on $S_{0}$ ) whenever $\Gamma$ is not of type $\widetilde{A}_{1}$ or $\widetilde{C}_{r}$. For $\alpha \in \Pi$ and $s=s_{\alpha}$ the corresponding simple reflection, we define

$$
\tilde{I}_{s}=T_{s}\left(1-\hat{T}_{-2 \alpha}\right)-\hat{T}_{-\alpha}\left(v^{c_{s}+c_{\tilde{s}}}-v^{c_{s}-c_{\tilde{s}}}\right)+1-v^{2 c_{s}} \in \mathcal{H} .
$$

This relates to the element $I_{s}$ in loc.cit. as follows. Suppose $s \in S_{0}$ with $c_{s}=c_{\tilde{s}}$. Then

$$
\tilde{I}_{s}=I_{s}\left(1+\hat{T}_{-\alpha}\right)
$$

where

$$
I_{s}=T_{s}\left(1-\hat{T}_{-\alpha}\right)+1-v^{2 c_{s}} .
$$

Proposition 2.1.1. For $\lambda \in T$ and $s \in S_{0}$, we have

$$
\hat{T}_{\lambda} \tilde{I}_{s}=\tilde{I}_{s} \hat{T}_{\lambda s}
$$

Proof. Let $U$ be the set of coroots in Case 1 and the set of coweights in Case 2. Since $(T, U, \Phi, \check{\Phi}, \Pi$ ) is a root system, as defined in [22], Proposition 3.6 in loc. cit. implies the stated identity.

Let $\alpha \in \Pi, s=s_{\alpha}$. In loc. cit., Lusztig defines an element

$$
\mathcal{G}(\alpha)=\frac{\left(\hat{T}_{\alpha} v^{c_{s}+c_{\bar{z}}}-1\right)\left(\hat{T}_{\alpha} v^{c_{s}-c_{\tilde{s}}}+1\right)}{\hat{T}_{2 \alpha}-1}
$$

which lies in the quotient field of the $\mathcal{A}$-submodule of $\mathcal{H}$ generated by the elements $\hat{T}_{\lambda}$ for $\lambda \in Q$. (Note the similarity between $\mathcal{G}(\alpha)$ and the function $c_{\alpha}$ in [26], p.98, which is used in order to construct intertwining operators between the principal series representations of $\mathcal{H}$; also, see [25], p. 51.)

In terms of $\mathcal{G}(\alpha)$, we can express $\tilde{I}_{s}$ as

$$
\tilde{I}_{s}=\left(T_{s}+1-\mathcal{G}(\alpha)\right)\left(1-\hat{T}_{-2 \alpha}\right)
$$

We define

$$
\tilde{\mathcal{G}}(\alpha)=\left(\hat{T}_{-\alpha} v^{c_{s}+c_{\tilde{s}}}-1\right)\left(\hat{T}_{-\alpha} v^{c_{s}-c_{\tilde{s}}}+1\right)
$$

and denote by $\tilde{G}(\alpha) \in \operatorname{End}_{\mathcal{H}}(\mathcal{H})$ right multiplication with $\tilde{\mathcal{G}}(\alpha)$.
To $s \in S_{0}$, we associate an $\mathcal{H}$-endomorphism $\tilde{\Theta}_{s}$ of $\mathcal{H}$, which maps $f \in \mathcal{H}$ to

$$
\tilde{\Theta}_{s}(f)=f \tilde{I}_{s}
$$

Lemma 2.1.2. Let $\alpha \in \Pi, s=s_{\alpha}$, and $\lambda \in T$. Then

$$
\tilde{\Theta}_{s}\left(\left(T_{s}+1\right) \check{T}_{\lambda}\right)=-v^{h(\lambda-\lambda s)} \tilde{G}(\alpha)\left(\left(T_{s}+1\right) \check{T}_{\lambda s}\right)
$$

Proof. Using Proposition 2.1.1, we calculate

$$
\begin{aligned}
& \tilde{\Theta}_{s}\left(\left(T_{s}+1\right) \check{T}_{\lambda}\right) \\
& \quad=\left(T_{s}+1\right) \check{T}_{\lambda} \tilde{I}_{s} \\
& \quad=v^{h(\lambda-\lambda s)}\left(T_{s}+1\right) \tilde{I}_{s} \check{T}_{\lambda s} \\
& \quad=v^{h(\lambda-\lambda s)}\left(T_{s}+1\right)\left(v^{2 c_{s}}\left(1-\hat{T}_{-2 \alpha}\right)-\hat{T}_{-\alpha}\left(v^{c_{s}+c_{\bar{s}}}-v^{c_{s}-c_{\tilde{s}}}\right)+1-v^{2 c_{s}}\right) \check{T}_{\lambda s} \\
& \quad=-v^{h(\lambda-\lambda s)}\left(T_{s}+1\right)\left(\hat{T}_{-2 \alpha} v^{2 c_{s}}+\hat{T}_{-\alpha}\left(v^{c_{s}+c_{s}}-v^{c_{s}-c_{\tilde{s}}}\right)-1\right) \check{T}_{\lambda s} \\
& \quad=-v^{h(\lambda-\lambda s)}\left(T_{s}+1\right) \tilde{\mathcal{G}}(\alpha) \check{T}_{\lambda s} \\
& \quad=-v^{h(\lambda-\lambda s)} \tilde{G}(\alpha)\left(\left(T_{s}+1\right) \check{T}_{\lambda s}\right)
\end{aligned}
$$

which is what we claimed.

We deduce from this lemma that

$$
\tilde{\Theta}_{s}\left(d_{\lambda}\right)=-v^{h(\lambda-\lambda s)} \tilde{G}(\alpha) d_{\lambda s} .
$$

Thus $\tilde{\Theta}_{s}$ induces an $\mathcal{H}$-endomorphism on $\mathcal{H}^{0}$, for which we write $\tilde{\Theta}_{s}$ as well. Since $\tilde{G}(\alpha)$ is an injective endomorphism on $\mathcal{H}$ (and on $\mathcal{H}^{0}$ ), we can define

$$
\Theta_{s}=-\tilde{G}(\alpha)^{-1} \tilde{\Theta}_{s} \in \operatorname{End}_{\mathcal{H}}\left(\mathcal{H}^{0}\right)
$$

If $w=s_{1} \cdots s_{n} \in W_{0}$ with $s_{i} \in S_{0}, 1 \leq i \leq n$, we set $\Theta_{w}=\Theta_{s_{n}} \cdots \Theta_{s_{1}} \in \operatorname{End}_{\mathcal{H}}\left(\mathcal{H}^{0}\right)$.
Now take an arbitrary element $w=x p_{\lambda}, x \in W_{0}, \lambda \in Q$, in $W_{a}$. Right multiplication with $\hat{T}_{\lambda}$ yields an element $\Theta_{\lambda}$ in $\operatorname{End}_{\mathcal{H}}\left(\mathcal{H}^{0}\right)$ (and in $\operatorname{End}_{\mathcal{H}}(\mathcal{H})$ ). We let $\Theta_{w}=\Theta_{\lambda} \cdot \Theta_{x}$ and thereby obtain an element in $\operatorname{End}_{\mathcal{H}}\left(\mathcal{H}^{0}\right)$, which satisfies

$$
\Theta_{w}\left(d_{\lambda}\right)=v^{h(\lambda-\lambda w)} d_{\lambda w}
$$

for $\lambda \in T$.
Proposition 2.1.3. (comp. [14], Proposition 2.8) For any $w \in W_{a}$, there is a unique $\mathcal{H}$-linear isomorphism $\theta_{w}$ of $\mathcal{M}^{0}$ such that

$$
\theta_{w}\left(e_{\lambda}\right)=v^{c\left(A_{\lambda w}^{+}, A_{\lambda}^{+}\right)} e_{\lambda w}
$$

for any $\lambda \in T$.
Proof. We have

$$
\begin{aligned}
\phi \Theta_{w} \phi^{-1}\left(e_{\lambda}\right) & =\phi \Theta_{w}\left(d_{\lambda}\right) \\
& =v^{h(\lambda-\lambda w)} \phi\left(d_{\lambda w}\right) \\
& =v^{c\left(A_{\lambda w}^{-}, A_{\lambda}^{-}\right)} e_{\lambda w} \\
& =v^{c\left(A_{\lambda w}^{+}, A_{\lambda}^{+}\right)} e_{\lambda w}
\end{aligned}
$$

and the map $\theta_{w} \stackrel{\text { def }}{=} \phi \Theta_{w} \phi^{-1}$ has the desired properties. The uniqueness is clear.
Note the equality $\theta_{w} \theta_{w^{\prime}}=\theta_{w w^{\prime}}$ for $w, w^{\prime} \in W_{a}$.

There are further maps introduced in [14], which we will need later on. Let $q=v^{2}$. Replacing $q$ by the appropriate $q^{c_{s}}, s \in S_{a}$, and $\delta$, a length function, by $\gamma$, a weighted length function, will prove the following statements in this section.

Lemma 2.1.4. (comp. [14], Lemma 2.10) For $\lambda \in T$, the map $\varphi_{\lambda}: \mathcal{M} \rightarrow \mathcal{M}$ defined by

$$
\varphi_{\lambda}\left(\sum_{A \in X} c_{A} A\right)=\sum_{A \in X} \bar{c}_{A} A v_{\lambda}
$$

$c_{A} \in \mathcal{A}$ for $A \in X$, is $\mathcal{H}$-antilinear.

Theorem 2.1.5. (comp. [14], Theorem 2.12) Let $\gamma$ be a weighted length function on $X$.
(i) There is a unique $\mathcal{H}$-antilinear map $\Phi_{\gamma}: \mathcal{M}^{0} \rightarrow \mathcal{M}^{0}$ such that

$$
\Phi_{\gamma}\left(e_{\lambda}\right)=v^{-2 \gamma\left(A_{\lambda}^{+}\right)} e_{\lambda}
$$

for all $\lambda \in T$.
(ii) The map $\Phi_{\gamma}$ is an involution.
(iii) If $\gamma^{\prime}=\gamma+n, n \in \mathbb{Z}$, is another weighted length function on $X$ we have $\Phi_{\gamma^{\prime}}=$ $v^{-2 n} \Phi_{\gamma}$.

Corollary 2.1.6. (comp. [14], Corollary 2.13) Let $\lambda \in T$ and $\gamma_{\lambda}$ be the weighted length function satisfying $\gamma_{\lambda}\left(A_{\lambda}^{+}\right)=0$. Then $\Phi_{\gamma_{\lambda}}(m)=\varphi_{\lambda}(m)$ for all $m \in \mathcal{M}_{\lambda}$.
2.2. Some estimates. We now generalize the degree estimates in [14], Section 4. We achieve this by means of a more detailed analysis of the arguments in loc. cit.

Fix a special point $\lambda$. If $A$ is an alcove such that $\lambda \in \bar{A}$ and if $y \in W_{a}$ is such that $y\left(A_{\lambda}^{+}\right) \subset \mathcal{C}_{\lambda}^{+}$, we write

$$
T_{y}(A)=\sum_{B \in X} \pi_{B, y}^{A} B(\in \mathcal{M})
$$

Proposition 2.2.1. (comp. [14], Proposition 4.2) Let $A$ and $y$ be as above.
(i) The coefficient $\pi_{B, y}^{A}$ of $B \in X$ is zero unless $B \leq y(A)$, in which case $\pi_{B, y}^{A}$ is a polynomial in $q$ of degree at most $\frac{1}{2} c\left(B, y\left(A_{\lambda}^{+}\right)\right)$.
(ii) Let $c^{*}=\min _{s \in S_{a}} c_{s}$, and suppose $y\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$. If $B<y(A)$, the polynomial $\pi_{B, y}^{A}$ has degree at most $\frac{1}{2}\left(c\left(B, y\left(A_{\lambda}^{+}\right)\right)-c^{*}\right)$. If $B=y(A)$, we have $\pi_{B, y}^{A}=1$.
(iii) If $c\left(B, y\left(A_{\lambda}^{+}\right)\right)$is even, the coefficient of $q^{\frac{1}{2} c\left(B, y\left(A_{\lambda}^{+}\right)\right)}$in $\pi_{B, y}^{A}$ is non-negative.
(iv) In part (ii), if $c\left(B, y\left(A_{\lambda}^{+}\right)\right)-c^{*}$ is even, the coefficient of $q^{\frac{1}{2}\left(c\left(B, y\left(A_{\lambda}^{+}\right)\right)-c^{*}\right)}$ in $\pi_{B, y}^{A}$ is non-negative.

In order to prove this proposition, we need the following result, which is Corollary 3.4 in loc. cit.

Lemma 2.2.2. Let $A \in X$ and $s_{1}, s_{2}, \ldots, s_{k} \in S_{a}$ be such that $d\left(A, s_{k} \cdots s_{2} s_{1}(A)\right)=$ $k$. For any sequence $1 \leq i_{1}<\cdots<i_{p} \leq k$, we have $s_{i_{p}} \cdots s_{i_{1}}(A) \leq s_{k} \cdots s_{2} s_{1}(A)$.

Proof of Proposition 2.2.1. Let $s_{1}, \ldots, s_{k} \in S_{a}$. Let $\mathcal{I}$ be the collection of all $I=\left\{i_{1}, \ldots, i_{p_{I}}\right\}$ such that $1 \leq i_{1}<\cdots<i_{p_{I}} \leq k$ and

$$
s_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A)<\hat{s}_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A)
$$

for all $t=1, \ldots, p_{I}$. For $I \in \mathcal{I}$, define

$$
\begin{gathered}
\tilde{p}_{I}=\sum_{j=1}^{p_{I}} c_{s_{i_{j}}} \\
\widetilde{\Pi}_{I}=\prod_{j=1}^{p_{I}}\left(q^{c_{s_{i_{j}}}}-1\right) \text { and } \\
\tilde{m}_{I}=\sum_{j} c_{s_{j}}
\end{gathered}
$$

where the last sum runs over all $j, 1 \leq j \leq k$, such that $j \notin I$ and

$$
s_{j} s_{j-1} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A)<s_{j-1} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A) .
$$

(In the last expression we omit all $s_{i}$ such that $i \in I$ and $i<j$.)
One verifies by induction on $k$ that

$$
\begin{equation*}
T_{s_{k}} \cdots T_{s_{1}}(A)=\sum_{I \in \mathcal{I}} q^{\tilde{m}_{I}} \widetilde{\Pi}_{I} s_{k} \cdots \hat{s}_{i_{p_{I}}} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A) . \tag{*}
\end{equation*}
$$

We now take a reduced expression $s_{k} \cdots s_{1}=y$. If $s_{1}^{\prime}, \ldots, s_{n}^{\prime} \in S_{a} \cap W_{\lambda}$ are such that $s_{n}^{\prime} \cdots s_{1}^{\prime}(A)=A_{\lambda}^{+}$and $d\left(A, A_{\lambda}^{+}\right)=n$, we have

$$
d\left(A, y s_{n}^{\prime} \cdots s_{1}^{\prime}(A)\right)=d\left(A, y\left(A_{\lambda}^{+}\right)\right)=n+k=l\left(y s_{n}^{\prime} \cdots s_{1}^{\prime}\right)
$$

We conclude with Lemma 2.2.2 that

$$
s_{k} \cdots \hat{s}_{i_{p_{I}}} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A) \leq y\left(A_{\lambda}^{+}\right)
$$

for all $I=\left\{i_{1}, \ldots, i_{p_{I}}\right\}, 1 \leq i_{1}<\cdots<i_{p_{I}} \leq k$. Suppose $I \in \mathcal{I}$. We set

$$
B=s_{k} \cdots \hat{s}_{i_{p_{I}}} \cdots \hat{s}_{i_{2}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A)
$$

Then

$$
d(A, B)=k-p_{I}-2 m_{I}
$$

and

$$
c(A, B)=\sum_{i=1}^{k} c_{s_{i}}-\tilde{p}_{I}-2 \tilde{m}_{I}
$$

We deduce

$$
\begin{aligned}
c\left(B, y\left(A_{\lambda}^{+}\right)\right) & =c\left(A, y\left(A_{\lambda}^{+}\right)\right)-c(A, B) \\
& =\sum_{i=1}^{k} c_{s_{i}}+\sum_{i=1}^{n} c_{s_{i}^{\prime}}-\sum_{i=1}^{k} c_{s_{i}}+\tilde{p}_{I}+2 \tilde{m}_{I} \\
& =\sum_{i=1}^{n} c_{s_{i}^{\prime}}+\tilde{p}_{I}+2 \tilde{m}_{I}
\end{aligned}
$$

In view of $(*)$, it remains to prove

$$
\sum_{i=1}^{n} c_{s_{i}^{\prime}} \geq \tilde{p}_{I}
$$

for part (i), and

$$
\sum_{i=1}^{n} c_{s_{i}^{\prime}}-c^{*} \geq \tilde{p}_{I}
$$

(under the additional assumption that $y\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$ ) for part (ii). Parts (iii) and (iv) then follow.

Let $\overline{\mathcal{F}}$ be the set of directions of hyperplanes in $\mathcal{F}$. We denote the direction of a hyperplane $H \in \mathcal{F}$ by $i(H)$. If $\mathcal{C}$ is a quarter, the set $\mathcal{I}(\mathcal{C})$ is said to contain all directions $i \in \overline{\mathcal{F}}$ such that $\mathcal{C} \subseteq V_{H}^{-}$for some $H \in \mathcal{F}$ with $i(H)=i$.

The proof in loc. cit. makes use of the following two facts.
(2.2.3) Let $\mathcal{C}$ be the quarter with vertex $\lambda$ containing $A$. Then

$$
n=|\mathcal{I}(\mathcal{C})|
$$

(2.2.4) If $\mathcal{C}$ is any quarter and $H \in \mathcal{F}$ with direction $i(H) \in \mathcal{I}(\mathcal{C})$, then

$$
\left|\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right|<|\mathcal{I}(\mathcal{C})|
$$

We will need stronger versions. For $i \in \overline{\mathcal{F}}$, set

$$
c_{i}=\max _{H \in \mathcal{F}, i(H)=i} c_{H}
$$

and if $\mathcal{J} \subseteq \overline{\mathcal{F}}$, we write

$$
m(\mathcal{J})=\sum_{i \in \mathcal{J}} c_{i}
$$

Claim 2.2.5. Let $\mathcal{C}$ be the quarter with vertex $\lambda$ containing $A$. We have

$$
\sum_{j=1}^{n} c_{s_{j}^{\prime}}=m(\mathcal{I}(\mathcal{C}))
$$

Indeed, let $H_{j}, 1 \leq j \leq n$, be the hyperplane in $\mathcal{F}$ separating $s_{j-1}^{\prime} \cdots s_{1}^{\prime}(A)$ from $s_{j}^{\prime} \cdots s_{1}^{\prime}(A)$. Since $s_{n}^{\prime} \cdots s_{1}^{\prime}(A)=A_{\lambda}^{+} \subset V_{H_{j}}^{+}$for all $H_{j} \in \mathcal{F}, 1 \leq j \leq n$, we have $\mathcal{C} \subseteq V_{H_{j}}^{-}$for all $1 \leq j \leq n$, i.e. $\mathcal{I}(\mathcal{C})=\left\{i\left(H_{j}\right) \mid 1 \leq j \leq n\right\}$. Besides, for each $j$, $1 \leq j \leq n$, the hyperplane $H_{j}$ passes through $\lambda \in T$, so $c_{i\left(H_{j}\right)}=c_{H_{j}}$, and since $H_{j}$ supports a face of type $s_{j}^{\prime} \in W_{0}$, the assertion follows.

Claim 2.2.6. If $\mathcal{C}$ is a quarter with vertex in $T$ and if $H \in \mathcal{F}$ has direction $i=$ $i(H) \in \mathcal{I}(\mathcal{C})$, there is an injective map $\mathcal{I}\left(\mathcal{C} \sigma_{H}\right) \rightarrow \mathcal{I}(\mathcal{C})-\{i\}$ which preserves the parameter.
W.l.o.g., let $\mathcal{C}$ be the quarter with vertex $\lambda$ containing $A$. We distinguish two cases.

First, suppose $H$ passes through $\lambda$, hence $H$ equals one of the above hyperplanes $H_{j}$ for some $1 \leq j \leq n$. We have $A \sigma_{H} \subset \mathcal{C} \sigma_{H}$ and $\lambda \in \overline{A \sigma_{H}}$. Since $s_{j}^{\prime} \cdots s_{1}^{\prime}(A)=$ $s_{j-1}^{\prime} \cdots s_{1}^{\prime}\left(A \sigma_{H}\right)$, we get

$$
A_{\lambda}^{+}=s_{n}^{\prime} \cdots s_{1}^{\prime}(A)=s_{n}^{\prime} \cdots s_{j+1}^{\prime} s_{j-1}^{\prime} \cdots s_{1}^{\prime}\left(A \sigma_{H}\right)
$$

Therefore, any hyperplane $H^{\prime}$ separating $A_{\lambda}^{+}$and $A \sigma_{H}$ supports a face of type $s_{k}^{\prime}$ for some $1 \leq k \leq n, k \neq j$, and we can map $i^{\prime}=i\left(H^{\prime}\right)$ to $i\left(H_{k}\right)$. (Note that the expression $s_{n}^{\prime} \cdots s_{j+1}^{\prime} s_{j-1}^{\prime} \cdots s_{1}^{\prime}$ does not have to be reduced.)

Next, suppose $H$ does not pass through $\lambda$, i.e. $H=H_{j} p_{\mu}$ for some $1 \leq j \leq n$ and some $\mu \in P$. Then

$$
\mathcal{C} \sigma_{H}=\mathcal{C} \sigma_{H_{j} p_{\mu}}=\left(\mathcal{C} \sigma_{H_{j}}\right) p_{2 \mu}
$$

Since $\lambda p_{2 \mu}$ is again a special point and $\mathcal{I}(\mathcal{C})=\mathcal{I}\left(\mathcal{C} p_{2 \mu}\right)$, we can replace $\mathcal{C}$ by $\mathcal{C} p_{2 \mu}$ and argue as before.

By Claim 2.2.5, the proposition becomes a consequence of the following statement.
Lemma 2.2.7. (comp. [14], Lemma 4.3) Let $\lambda$ be a special point, $A$ an alcove containing $\lambda$ in its closure and $\mathcal{C}$ the quarter with vertex $\lambda$ containing $A$. Let $s_{1}, \ldots, s_{k} \in$ $S_{a}$ be such that $d\left(A_{\lambda}^{+}, s_{k} \cdots s_{1}\left(A_{\lambda}^{+}\right)\right)=k$, and let $1 \leq i_{1}<\cdots<i_{p} \leq k$ be such that

$$
s_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A)<\hat{s}_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{1}} \cdots s_{1}(A)
$$

for all $t=1, \ldots, p$. Write $\tilde{p}=\sum_{j=1}^{p} c_{s_{i}}$.
(i) We have $\tilde{p} \leq m(\mathcal{I}(\mathcal{C}))$.
(ii) If $s_{k} \cdots s_{1}\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$ and $A \neq A_{\lambda}^{+}$, then $\tilde{p} \leq m(\mathcal{I}(\mathcal{C}))-c^{*}$.

Proof. The following facts can be found in loc. cit.
Let $w \in W_{a}$ be such that $\lambda w=\lambda$ and $A_{\lambda}^{+} w=A$. We denote by $H_{j}, 1 \leq j \leq k$, the hyperplane in $\mathcal{F}$ separating $s_{j-1} \cdots s_{1}\left(A_{\lambda}^{+}\right)$from $s_{j} \cdots s_{1}\left(A_{\lambda}^{+}\right)$. Let $1 \leq j \leq k$.

- If $s_{j} s_{j-1} \cdots s_{1}(A)<s_{j-1} \cdots s_{1}(A)$ we have $i\left(H_{j} w\right) \in \mathcal{I}(\mathcal{C})$.
- If $s_{k} \cdots s_{1}\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$ the hyperplane $H_{j} w$ is not parallel to any of the walls of $\mathcal{C}$.

The proof proceeds by induction on $|\mathcal{I}(\mathcal{C})|$. If $\mathcal{I}(\mathcal{C})=\emptyset$ we have $A=A_{\lambda}^{+}$and $p=\tilde{p}=0$. There is nothing to prove for part (ii).

Now let $|\mathcal{I}(\mathcal{C})| \geq 1$. The case $p=0$ is clear, so we assume $p \geq 1$. Let $H$ be the unique hyperplane in $\mathcal{F}$ separating $s_{i_{1}} \cdots s_{2} s_{1}(A)$ from $\hat{s}_{i_{1}} \cdots s_{2} s_{1}(A)$. We set
$\lambda^{\prime}=\lambda \sigma_{H}, A^{\prime}=A \sigma_{H}, \mathcal{C}^{\prime}=\mathcal{C} \sigma_{H}$. Then $d\left(A_{\lambda^{\prime}}^{+}, s_{k} \cdots s_{1}\left(A_{\lambda^{\prime}}^{+}\right)\right)=k$, and $\hat{s}_{i_{1}} \cdots s_{2} s_{1}(A)=$ $s_{i_{1}} \cdots s_{2} s_{1}\left(A^{\prime}\right)$ implies

$$
s_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{2}} \cdots s_{i_{1}} \cdots s_{1}\left(A^{\prime}\right)<\hat{s}_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{2}} \cdots s_{i_{1}} \cdots s_{1}\left(A^{\prime}\right)
$$

for $t=2,3, \ldots, p$.
We have $\left|\mathcal{I}\left(\mathcal{C}^{\prime}\right)\right|<|\mathcal{I}(\mathcal{C})|$. Applying the inductive assumption, part (i), to $\lambda^{\prime}, A^{\prime}$, $s_{1}, \ldots, s_{k}$ and $i_{2}, \ldots, i_{p}$ yields

$$
\sum_{j=2}^{p} c_{s_{i_{j}}} \leq m\left(\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right)
$$

and therefore

$$
\begin{aligned}
\sum_{j=1}^{p} c_{s_{i_{j}}} & \leq m\left(\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right)+c_{s_{i_{1}}} \\
& =m\left(\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right)+c_{H} \\
& \leq m(\mathcal{I}(\mathcal{C}))
\end{aligned}
$$

We used Claim 2.2.6 for the last inequality. (Note that it is feasible that

$$
\left.m\left(\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right)+c_{H}<m\left(\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right)+c_{i(H)}<m(\mathcal{I}(\mathcal{C})) .\right)
$$

Under the assumptions of part (ii), if $\mathcal{C}^{\prime} \neq \mathcal{C}_{\lambda^{\prime}}^{+}$it follows by induction that

$$
\sum_{j=2}^{p} c_{s_{i_{j}}} \leq m\left(\mathcal{I}\left(\mathcal{C} \sigma_{H}\right)\right)-c^{*}
$$

and therefore

$$
\sum_{j=1}^{p} c_{s_{i_{j}}} \leq m(\mathcal{I}(\mathcal{C}))-c^{*}
$$

Otherwise, we have

$$
\sum_{j=2}^{p} c_{s_{i_{j}}}=0 .
$$

But then, according to the two facts mentioned at the beginning of the proof, $\mathcal{I}(\mathcal{C})$ contains the direction of $H w$, which is not parallel to any wall of $\mathcal{C}$, and also the direction of some wall of $\mathcal{C}$. Consequently,

$$
\sum_{j=1}^{p} c_{s_{i_{j}}}=c_{s_{i_{1}}}=c_{H} \leq m(\mathcal{I}(\mathcal{C}))-c^{*}
$$

which completes the proof of the lemma.
Note that $\tilde{p}=m(\mathcal{I}(\mathcal{C}))$ implies $p=|\mathcal{I}(\mathcal{C})|$. We therefore obtain the following corollary of the proof of the above lemma in the same way as for equal parameters.

Corollary 2.2.8. (comp. [14], Corollary 4.4) Let $A$ and $y$ be as in Proposition 2.2.1 (i), and let $B$ be an alcove such that $B \leq y\left(A_{\lambda}^{+}\right)$and such that $\pi_{B, y}^{A}$ has degree equal to $\frac{1}{2} c\left(B, y\left(A_{\lambda}^{+}\right)\right)$. Then $B=y\left(A_{\lambda}^{+}\right) \tau$ for some translation $\tau \in W_{a}$.
2.3. The basis elements. The following results all have their counterpart in [14]. In view of the two previous sections, most proofs in loc. cit. can easily be generalized, and we refer the reader to loc. cit. whenever a statement in this section is presented without proof.

Theorem 2.3.1. (comp. [14], Theorem 2.15) Let $\lambda$ be a special point and $C$ an alcove in $\Pi_{\lambda}$. There exists a unique element $D_{C} \in \mathcal{M}_{\lambda}$ such that
(i) $D_{C}=\sum_{A \leq C} Q_{A, C} A$ where $Q_{A, C}$ is a polynomial in $v$ with integer coefficients of degree less than $c(A, C)$ if $A<C$ and $Q_{C, C}=1$ and
(ii) $v^{2 c\left(A_{\lambda}^{+}, C\right)} \overline{Q_{A, C}}=Q_{A v_{\lambda}, C}$.
(We set $Q_{A, C}=0$ for $A, C \in X, A \not \leq C$.)
Remark 2.3.2. With $\lambda$ and $C$ as above, let $B=C v_{\lambda}$ and $D^{B}=\varphi_{\lambda}\left(D_{C}\right)$. The element $D^{B}$ is characterized by the conditions $D^{B} \in \mathcal{M}_{\lambda}$ and
(i) $D^{B}=\sum_{A \geq B} \overline{Q^{B, A}} A$ where $Q^{B, A}$ is a polynomial in $v$ with integer coefficients of degree less than $c(B, A)$ if $B<A$ and $Q^{B, B}=1$ and
(ii) $v^{2 c\left(B, A_{\lambda}^{-}\right)} \overline{Q^{B, A}}=Q^{B, A v_{\lambda}}$.

Let $G$ be a simply connected almost simple algebraic group of type $\Phi$ over an algebraically closed field of characteristic $p>1$. Assume that $p$ is sufficiently large.

In loc. cit., Lusztig conjectured that the integer $Q^{B, A}(1), A \geq B$, equals the multiplicity of a certain irreducible $G$-module in a Jordan-Hölder series of a certain Weyl module of $G$ associated to $A$. He pointed out that for $v=1$ the condition $D^{B} \in \mathcal{M}_{\lambda}$ becomes the condition that

$$
\sum_{A \geq B} Q^{B, A}(1) A
$$

is invariant under the stabilizer $\Omega_{\lambda}$ of $\lambda$ in $W_{a}$ acting on $X$ on the right. Assuming Lusztig's conjecture, this $\Omega_{\lambda}$-invariance is due to Jantzen (cf. [7]).

The $\Omega_{\lambda}$-invariance of $\sum_{A \geq B} Q^{B, A}(1) A$ also implies the symmetry condition (ii) for $v=1$.
(The above conjecture by Lusztig is equivalent to his conjecture on the irreducible characters of rational $G$-modules in [13]. It is now known to be true - due to work involving affine Kac-Moody Lie algebras as well as quantum groups.)

For $\lambda \in Q$, define

$$
Y_{\lambda}=\left\{y p_{\mu} \mid y \in W_{0}, \mu \in Q, \mu \preceq \lambda\right\}
$$

where $\preceq$ denotes the usual partial order on $Q$. A subset $K \subset W_{a}$ is called $x$-bounded, $x \in W_{0}$, if

$$
K \subseteq Y_{\lambda} x
$$

for some $\mu \in Q$. Instead of $e$-bounded, $e$ being the identity in $W_{0}$, we also say bounded. For $x \in W_{0}$, let $\hat{\mathcal{H}}_{x}$ be the set of formal sums $f=\sum_{w \in W_{a}} a_{w} \check{T}_{w}, a_{w} \in \mathcal{A}$, such that Supp $f \stackrel{\text { def }}{=}\left\{w \in W_{a} \mid a_{w} \neq 0\right\}$ is $x$-bounded. The sets $\hat{\mathcal{H}}_{x}$ can naturally be regarded as $\mathcal{H}$-modules. Similarly, let $\hat{\mathcal{M}}_{x}$ be the $\mathcal{H}$-module which contains all formal sums $f=\sum_{A \in X} b_{A} A, b_{A} \in \mathcal{A}$, such that $\left\{w \in W_{a} \mid b_{A^{-} w} \neq 0\right\}$ is $x$-bounded. We set $\hat{\mathcal{H}}=\hat{\mathcal{H}}_{e}$ and $\hat{\mathcal{M}}=\hat{\mathcal{M}}_{e}$. In this section we will only need $\hat{\mathcal{M}}$.

Choose a weighted length function $\gamma$ on $X$.
As in loc.cit., we can extend $\Phi_{\gamma}: \mathcal{M}^{0} \rightarrow \mathcal{M}^{0}$ to a $\operatorname{map} \hat{\Phi}_{\gamma}: \hat{\mathcal{M}} \rightarrow \hat{\mathcal{M}}$, and we define elements $\mathcal{R}_{B, A} \in \mathcal{A}$ for $A, B \in X$ by

$$
\hat{\Phi}_{\gamma}(A)=v^{-2 \gamma(A)} \sum_{B \in X}(-1)^{d(A, B)} \mathcal{R}_{B, A} B .
$$

We remark that $\mathcal{R}_{B, A}=0$ for all $B \not \leq A$ and that $\mathcal{R}_{A, A}=1$.
The following statement strengthens the uniqueness part of Theorem 2.3.1.
Theorem 2.3.3. (comp. [14], Theorem 7.3) For any $C \in X$, there is a unique element $D_{C} \in \hat{\mathcal{M}}$ such that $\hat{\Phi}_{\gamma}\left(D_{C}\right)=v^{-2 \gamma(C)} D_{C}$ and

$$
D_{C}=\sum_{A \leq C} Q_{A, C} A
$$

where $Q_{A, C} \in \mathcal{A}^{+}$has degree less than $c(A, C)$ if $A<C$ and $Q_{C, C}=1$.
(It follows from Theorem 2.1.5 and Corollary 2.1.6 that the conditions $\hat{\Phi}_{\gamma}\left(D_{C}\right)=$ $v^{-2 \gamma(C)} D_{C}$ and $v^{2 c\left(A_{\lambda}^{+}, C\right)} \overline{Q_{A, C}}=Q_{A v_{\lambda}, C}$ are equivalent.)

For $C \in X$, let $\lambda \in T$ and $w \in W_{a}$ be such that $C \subseteq \Pi_{\lambda}$ and $w\left(A_{\lambda}^{+}\right)=C$. The element $D_{C}$ can be shown to satisfy

$$
\begin{aligned}
D_{C} & =\sum_{y \leq w w_{\lambda}, l\left(y w_{\lambda}\right)=l(y)+l\left(w_{\lambda}\right)} P_{y w_{\lambda}, w w_{\lambda}} T_{y}\left(\sum_{z \in W_{\lambda}} T_{z}\left(A_{\lambda}^{-}\right)\right) \\
& =\sum_{y^{\prime} \leq w w_{\lambda}} P_{y^{\prime}, w w_{\lambda}} T_{y^{\prime}}\left(A_{\lambda}^{-}\right) \\
& =C_{w w_{\lambda}}^{*}\left(A_{\lambda}^{-}\right) .
\end{aligned}
$$

This representation of $D_{C}$ is used to prove the next two corollaries. (We preserve the notation.)

Corollary 2.3.4. (comp. [14], Corollary 5.3) With $y \in W_{a}$ such that $A=y\left(A_{\lambda}^{-}\right)$, we have

$$
Q_{A, C}(1)=P_{y, w w_{\lambda}}(1)
$$

Corollary 2.3.5. (comp. [14], Corollary 5.4) Let $\mathcal{O} \subset X$ be an orbit under the translation subgroup $Q$ of $W_{a}$. For $C \in X$, the sum

$$
\sum_{A \in \mathcal{O}} Q_{A, C}
$$

is independent of the choice of $\mathcal{O}$, and it is equal to

$$
\sum_{y \leq w w_{\lambda}, l\left(y w_{\lambda}\right)=l(y)+l\left(w_{\lambda}\right)} v^{2 m(y)} P_{y w_{\lambda}, w w_{\lambda}} .
$$

Theorem 2.3.3 is the main ingredient of the following statement.

Corollary 2.3.6. (comp. [14], Corollary 7.4) Let $\tau \in T$ be a translation on $X$. For any $A \leq C$ in $X$, we have $Q_{A, C}=Q_{A \tau, C \tau}$.

Finally note that for $A, C \in X$ and $\lambda \in T$ such that $C \subseteq \Pi_{\lambda}$, the polynomial $Q_{A, C}$ is non-zero only if $C v_{\lambda} \leq A$.

In analogy with loc.cit., we now want to further examine the elements $D_{C}$. The aim is to find an inductive formula.

For $A, C \in X$ and $s \in S_{a}$ such that $s A<A<C<s C$, we define elements $M_{A, C}^{s} \in \mathcal{A}$ by the conditions
(i) $\sum_{A \leq B<C, s B<B} v^{c(B, C)} Q_{A, B} M_{B, C}^{s}-v^{c_{s}} Q_{A, C}$ has degree less than $c(A, C)$ and (ii) $\bar{M}_{A, C}^{s}=\overline{M_{A, C}^{s}}$.
(We set $M_{A, C}^{s}=0$ in all other cases.) This definition is similar to the definition of $M_{y, w}^{s} \in \mathcal{A}$ for $y, w \in W_{a}$ and $s \in S_{a}$ with $s y<y<w<s w$ in [16], Section 3.

Suppose $C$ is an alcove in $\Pi_{\lambda}, \lambda \in T$. We claim that $M_{A, C}^{s}=0$ unless $C v_{\lambda} \leq A \leq C$. In particular, for any $C \in X$, there are only finitely many $M_{A, C}^{s}, A \in X, s \in S_{a}$, which are non-zero.

Indeed, assume $A \not \leq C$. Then $Q_{A, C}=0$, and of course there is no $B \in X$ satisfying $A<B<C$. Thus $M_{A, C}^{s}=0$ for any $s \in S_{a}$.

Next assume $C v_{\lambda} \not \leq A$ or equivalently $A v_{\lambda} \not \leq C$. Again, we know that $Q_{A, C}=0$, and if $B \in X, A<B<C$, is such that $Q_{A, B} \neq 0$, then $A v_{\lambda} \leq B$, hence $A v_{\lambda} \leq C$, a contradiction. The claim follows.

Theorem 2.3.7. (comp. [14], Theorem 8.2) Let $C \in X$ and $s \in S_{a}$. We have

$$
T_{s} D_{C}= \begin{cases}-D_{C}+D_{s C}+\sum_{A<C, s \in \mathcal{L}(A)} v^{c(A, C)+c_{s}} M_{A, C}^{s} D_{A} & \text { if } s \notin \mathcal{L}(C) \\ v^{2 c_{s}} D_{C} & \text { if } s \in \mathcal{L}(\mathcal{C})\end{cases}
$$

Note that the sum is finite.
Proof. The case $s \in \mathcal{L}(C)$ is dealt with as in loc. cit. Consider the case $s \notin \mathcal{L}(C)$. We set

$$
D_{s C}^{\prime}=\left(T_{s}+1\right) D_{C}-\sum_{A<C, s \in \mathcal{L}(A)} v^{c(A, C)+c_{s}} M_{A, C}^{s} D_{A}
$$

As in loc. cit., we deduce that

$$
\Phi_{\gamma}\left(D_{s C}^{\prime}\right)=v^{-2 \gamma(s C)} D_{s C}^{\prime}
$$

The coefficient of $B \in X$ in $D_{s C}^{\prime}$ equals

$$
\tilde{Q}_{B}=\left\{\begin{array}{ll}
v^{2 c_{s}} Q_{s B, C}+Q_{B, C}, & \text { if } s B>B \\
Q_{s B, C}+v^{2 c_{s}} Q_{B, C}, & \text { if } s B<B
\end{array}\right\}-\sum_{A<C, s \in \mathcal{L}(A)} v^{c(A, C)+c_{s}} Q_{B, A} M_{A, C}^{s} .
$$

We have $\tilde{Q}_{s C}=Q_{C, C}=1$, and since the condition $B \leq s C$ is equivalent to $s B \leq s C$, the coefficient $\tilde{Q}_{B}$ is zero if $B \not \leq s C$.

Suppose $B<s C$. If $s B>B$, the definition of $M_{s B, C}^{s}$ implies that

$$
\sum_{A<C, s \in \mathcal{L}(A)} v^{c(A, C)} Q_{s B, A} M_{A, C}^{s}-v^{c_{s}} Q_{s B, C}
$$

has degree less than $c(s B, C)$. From $T_{s} D_{A}=v^{2 c_{s}} D_{A}$ for $A \in X$ with $s A<A$, we derive $Q_{B, A}=Q_{s B, A}$ for any $B \leq A$. Consequently, the degree of

$$
\sum_{A<C, s \in \mathcal{L}(A)} v^{c(A, C)} Q_{B, A} M_{A, C}^{s}-v^{c_{s}} Q_{B, C}
$$

is less than $c(B, C)-c_{s}$, and therefore $\tilde{Q}_{B}$ has degree less than $c(B, C)<c(B, s C)$.
If $s B<B$ the degree of

$$
\sum_{A<C, s \in \mathcal{L}(A)} v^{c(A, C)} Q_{B, A} M_{A, C}^{s}-v^{c_{s}} Q_{B, C}
$$

is less than $c(B, C)$. Thus $\tilde{Q}_{B}$ has again degree less than $c(B, s C)\left(=c(B, C)+c_{s}\right)$.
Now we apply the uniqueness statement in Theorem 2.3.3 to conclude that $D_{s C}^{\prime}=$ $D_{s C}\left(\right.$ and $\left.\tilde{Q}_{B}=Q_{B, s C}\right)$, hence the assertion.

Notice one equality which occurred in the above proof.
Corollary 2.3.8. (comp. [14], Corollary 8.4 (a)) For $A, C \in X$ with $A \leq C$ and $s \in \mathcal{L}(C)$, we have $Q_{A, C}=Q_{s A, C}$.

The following statement can be deduced from Theorem 2.3.7.

Corollary 2.3.9. (comp. [14], Corollary 8.3) The elements $D_{C}, C \in X$, form an $\mathcal{A}$-basis of $\mathcal{M}^{0}$.

For $A, C \in X$, we write $A \triangleleft C$ if for some (or equivalently any) box $\Pi$ the alcoves $A^{\prime}, C^{\prime}$ in $\Pi$ obtained from $A, C$ by translation under $T$ satisfy $A^{\prime}<C^{\prime}$.

Proposition 2.3.10. (comp. [14], Corollaries 10.5 and 10.6) Let $C \in X$ and $s \in$ $S_{a}-\mathcal{L}(C)$ such that $s C$ and $C$ lie in the same box. Then $\left(T_{s}+1\right) D_{C}-D_{s C}$ is an $\mathcal{A}^{+}$-linear combination of elements $D_{A}$ such that $A<C, A \triangleleft C$. In particular, if $A \in X$ is such that $s A<A<C$ and $M_{A, C}^{s} \neq 0$, then $A \triangleleft C$.

This result enables us to give an algorithm to compute the elements $D_{C}$ by a finite induction.

Let $C \in X$. There exists a unique $\lambda \in T$ such that $C \subseteq \Pi_{\lambda}$, and we can find elements $s_{1}, s_{2}, \ldots, s_{n_{C}}$ in $S_{a}$ such that $C=s_{1} s_{2} \cdots s_{n_{C}} A_{\lambda}^{+}$and $d\left(A_{\lambda}^{+}, C\right)=n_{C}$. We proceed by induction on $n_{C}$.

If $n_{C}=0$ we have $C=A_{\lambda}^{+}$, hence $D_{C}=e_{\lambda}$ and $M_{A, C}^{s}=0$ for all $A \in X, s \in S_{a}$. Now assume $n_{C} \geq 1$. Let $C^{\prime}=s_{2} \cdots s_{n_{C}}\left(A_{\lambda}^{+}\right)$, which lies in $\Pi_{\lambda}$. By Propositions 2.3.7 and 2.3.10, we can write

$$
D_{C}=\left(T_{s_{1}}+1\right) D_{C^{\prime}}-\sum_{A \triangleleft C^{\prime}, s_{1} \in \mathcal{L}(A)} v^{c\left(A, C^{\prime}\right)+c_{s_{1}}} M_{A, C^{\prime}}^{s_{1}} D_{A}
$$

Since $n_{C^{\prime}}<n_{C}$ and $A \triangleleft C^{\prime}$ implies $n_{A}<n_{C^{\prime}}$, the elements $D_{C^{\prime}}$ and $D_{A}$ for $A \triangleleft C^{\prime}$, as well as the polynomials $M_{A, C^{\prime}}^{s}$ are known by induction. Thus $D_{C}$ and by induction on $d(A, C)$ the polynomials $M_{A, C}^{s}$ for $A \in X, s \in S_{a}$, can be calculated.

Example 2.3.11. Suppose $W_{a}$ is of type $\tilde{C}_{2}$ with parameters

and

represents a box $\Pi_{\lambda}, \lambda \in T$. (So if $\lambda$ lies in the same $\Omega$-orbit as 0 then $a=A_{\lambda}^{+}$, $b=s_{0} A_{\lambda}^{+}, c=s_{1} s_{0} A_{\lambda}^{+}$, and $d=s_{2} s_{1} s_{0} A_{\lambda}^{+}$.) Since $c_{s_{0}}=c_{s_{2}}$, we can assume that $\lambda$ and 0 lie in the same $\Omega$-orbit.

Set $x=v^{2 a}$ and $y=v^{2 b}$. We obtain

$$
\begin{aligned}
& D_{b}=\left(T_{s_{0}}+1\right) D_{a}, \\
& D_{c}=\left(T_{s_{1}}+1\right) D_{b}-\left\{\begin{array}{ll}
x+y & \text { if } a<b \\
x & \text { if } a=b \\
0 & \text { if } a>b
\end{array}\right\} D_{a}, \\
& D_{d}=\left(T_{s_{2}}+1\right) D_{c}-\left\{\begin{array}{ll}
0 & \text { if } a<b \\
x & \text { if } a=b \\
x+y & \text { if } a>b
\end{array}\right\} D_{b}-\left\{\begin{array}{ll}
0 & \text { if } 2 a<b \\
-x^{2} & \text { if } 2 a=b \\
-x^{2}-y & \text { if } a<b<2 a \\
0 & \text { if } a=b \\
x y+y & \text { if } a>b
\end{array}\right\} D_{a} .
\end{aligned}
$$

The following patterns describe the elements $D_{C}$ for $C \in\{a, b, c, d\}$. Each pattern has center $\lambda$. The alcove $C$ is singled out, and the entry of an alcove $A$ is the polynomial $Q_{A, C}$.

$D_{b}$


$$
D_{c}
$$



where

$$
\begin{gathered}
A=\left\{\begin{array}{c}
1-x \\
1 \\
1+y
\end{array}, \quad B=\left\{\begin{array}{c}
0 \\
x \\
x+y
\end{array}, \quad C=\left\{\begin{array}{c}
x y-y \\
x^{2} \\
x y+y
\end{array},\right.\right.\right. \\
A^{\prime}=\left\{\begin{array}{c}
1+x \\
1 \\
1-y
\end{array}, \quad B^{\prime}=\left\{\begin{array}{c}
x+y \\
x \\
0
\end{array}, \quad C^{\prime}=\left\{\begin{array}{c}
x^{2}+x y \\
x^{2} \\
0
\end{array}, \quad D^{\prime}=\left\{\begin{array}{c}
x y+x^{2} y \\
x^{3} \\
x^{2} y-x^{2}
\end{array}\right.\right.\right.\right.
\end{gathered}
$$

for

$$
\left\{\begin{array}{l}
a<b \\
a=b \\
a>b
\end{array}\right.
$$

and

$$
E^{\prime}=\left\{\begin{array}{c}
x-x^{2} \\
x \\
x+y \\
x \\
x-y
\end{array}, \quad F^{\prime}=\left\{\begin{array} { c } 
{ x y - y } \\
{ x y - y + x ^ { 2 } } \\
{ x y + x ^ { 2 } } \\
{ x ^ { 2 } } \\
{ 0 }
\end{array} \quad \text { for } \left\{\begin{array}{l}
2 a<b \\
2 a=b \\
a<b<2 a \\
a=b \\
a>b
\end{array}\right.\right.\right.
$$

(For $W_{a}$ of type $\tilde{G}_{2}$, we get 21 cases for the relationship between $a$ and $b$.)
We now introduce another right action of $W_{a}$ on $X$. Let $A \in X, A=y\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$ for some $y \in W_{a}$ and some $\lambda \in T$. For $w \in W_{0}$, we define

$$
A * w=y\left(A_{\lambda w}^{+}\right) \subseteq \Pi_{\lambda w}
$$

Proposition 2.3.12. (comp. [14], Proposition 8.7) For $C \in X$ and $w \in W_{a}$, we have

$$
\theta_{w}\left(D_{C}\right)=v^{c(C * w, C)} D_{C * w} .
$$

For $C \in X$, let $\tilde{D}_{C}=v^{-\gamma(C)} D_{C}$, and define

$$
\tilde{M}_{A, C}^{s}= \begin{cases}M_{A, C}^{s} & \text { if } A<C \\ 1 & \text { if } A=s C>C \\ 0 & \text { otherwise }\end{cases}
$$

for $A, C \in X$ and $s \in S_{a}$. We rewrite some of the previous statements in terms of the $\mathcal{A}$-basis $\left\{\tilde{D}_{C} \mid C \in X\right\}$ of $\mathcal{M}^{0}$.

Corollary 2.3.13. (comp. [14], Corollary 8.9) Let $C \in X, w \in W_{a}$ and $s \in S_{a}$. We have

$$
\begin{aligned}
& \Phi_{\gamma}\left(\tilde{D}_{C}\right)=\tilde{D}_{C} \\
& \tilde{D}_{C} \theta_{w}=\tilde{D}_{C * w}
\end{aligned}
$$

and

$$
T_{s} \tilde{D}_{C}= \begin{cases}-\tilde{D}_{C}+v^{c_{s}} \sum_{A<C, s \in \mathcal{L}(A)} \tilde{M}_{A, C}^{s} \tilde{D}_{A} & \text { if } s \notin \mathcal{L}(C) \\ v^{2 c_{s}} \tilde{D}_{C} & \text { if } s \in \mathcal{L}(\mathcal{C})\end{cases}
$$

Corollary 2.3.14. (comp. [14], Corollary 8.10) For $A, C \in X, s \in S_{a}$, and $w \in W_{a}$, we have $\tilde{M}_{A, C}^{s}=\tilde{M}_{A * w, C * w}^{s}$.

Remark 2.3.15. Suppose we specialize $q$ to a prime power $p^{s}$. Let $\lambda$ be a homomorphism of the group of translations $Q$ in $W_{a}$ to $\mathbb{C}^{*}$. The $\mathbb{C}$-vector space $\mathcal{M}_{\lambda}$ spanned by all infinite formal linear combinations

$$
\tilde{D}_{C}^{*}=\sum_{\tau \in Q} \tilde{D}_{C \tau} \lambda\left(\tau^{-1}\right) .
$$

has dimension equal to the number of orbits of $Q$ on $X$, and it has a natural $\mathcal{H}$-action. For generic $\lambda$, the $\mathcal{H}$-module $\mathcal{M}_{\lambda}$ is isomorphic to the principal series representation defined by Matsumoto in [26].

We conclude this section with a few results concerning the polynomials $\mathcal{R}_{B, A}$ and $Q_{A, B}, A, B \in X$. For more details, we refer to [14].

Define polynomials $R_{y, w} \in \mathcal{A}^{+}, y, w \in W_{a}$, by

$$
T_{w^{-1}}^{-1}=v^{-2 m(w)} \sum_{y \in W_{a}}(-1)^{l(w)-l(y)} R_{y, w} T_{y}
$$

and polynomials $\tilde{Q}_{y, w} \in \mathcal{A}^{+}$for $y, w \in W_{a}, y\left(A^{+}\right) \subset \mathcal{C}^{-}, w\left(A^{+}\right) \subset \mathcal{C}^{-}$by

$$
\sum_{y \leq z \leq w, z\left(A^{+}\right) \subset \mathcal{C}^{-}}(-1)^{l(z)-l(y)} P_{y, z} \tilde{Q}_{z, w}=\delta_{y, w} .
$$

Theorem 2.3.16. (comp. [14], Theorem 11.6 and Corollary 11.9) Suppose $A, B \in X$ and $\lambda \in T$ are such that $A, B \subset \mathcal{C}_{\lambda}^{-}$and $A, B$ are sufficiently far from the walls of $\mathcal{C}_{\lambda}^{-}$. Let $y, w \in W_{a}$ be such that $y\left(A_{\lambda}^{+}\right)=B, w\left(A_{\lambda}^{+}\right)=A$. Then

$$
\mathcal{R}_{B, A}=\sum_{b \in W_{\lambda}} R_{w, y b}
$$

and

$$
Q_{A, B}=\tilde{Q}_{y, w} .
$$

We remark that as in loc.cit., the function defined by

$$
(A, B) \mapsto(-1)^{d(A, B)} Q_{A, B}(0)
$$

for $A, B \in X$ is the Möbius function of the partially ordered set $(X, \leq)$. In particular, $Q_{A, B}(0)$ does not depend on the choice of parameters.

Proposition 2.3.17. (comp. [14], Proposition 11.15) Let $\lambda \in T$ and $y, w \in W_{\lambda}$. We have

$$
\mathcal{R}_{y\left(A_{\lambda}^{-}\right), w\left(A_{\lambda}^{-}\right)}=R_{y, w}
$$

and

$$
Q_{y\left(A_{\lambda}^{-}\right), w\left(A_{\lambda}^{-}\right)}=P_{y, w} .
$$

## 3. A formula for the Kazhdan-Lusztig polynomials

In Corollary 3.3.3, we express the Kazhdan-Lusztig polynomials in terms of alcove polynomials (see Section 3.2 for the definition). For equal parameters this has been achieved by Kato in [9], and we follow his approach. Statements which are not proven here are obtained via a straightforward generalization of the proof of the corresponding statement in loc. cit.
3.1. Definitions. We need some more notation from loc. cit.

Let $\psi: \mathcal{H} \rightarrow \mathcal{H}$ be the map such that

$$
\psi(f)=\bar{f} T_{w_{0}}
$$

for $f \in \mathcal{H}$. This map is $\mathcal{H}$-antilinear, and it satisfies $\psi\left(\mathcal{H}^{0}\right)=\mathcal{H}^{0}$.
Proposition 3.1.1. (comp. [9], Proposition 2.8) There exists a unique involutive $\mathcal{H}$ antilinear map $\Psi^{0}: \mathcal{H}^{0} \rightarrow \mathcal{H}^{0}$ such that

$$
\Psi^{0}\left(d_{\lambda}\right)=v^{-2 h\left(w_{0} p_{\lambda}\right)} d_{\lambda}
$$

for any $\lambda \in T$. We have

$$
\Psi^{0}=v^{-2 m\left(w_{0}\right)} \psi \Theta_{w_{0}}
$$

Remark 3.1.2. Let $\gamma$ be the weighted length function on $X$ defined by

$$
\gamma\left(A^{-} w\right)=h(w)
$$

for $w \in W_{a}$. It turns out that under the isomorphism $\phi: \mathcal{H} \rightarrow \mathcal{M}$, the map $\psi$ corresponds to the map $\varphi \stackrel{\text { def }}{=} \varphi_{0}$ and $\Psi^{0}$ corresponds to $\Phi_{\gamma}$. From now on, we fix $\gamma$ and also the length function $\delta$ on $X$ defined by

$$
\delta\left(A^{-} w\right)=g(w)
$$

for $w \in W_{a}$.

We define $\tilde{\mathcal{G}}(\alpha)$ for arbitrary $\alpha \in \Phi$. Let $\alpha \in \Phi$. If $\alpha$ lies in the $W_{0}$-orbit of a simple root $\beta$, we set $c_{\alpha}=c_{s_{\beta}}, \tilde{c}_{\alpha}=c_{s_{\beta}}$, and

$$
\tilde{\mathcal{G}}(\alpha)=\left(\hat{T}_{-\alpha} v^{c_{\alpha}+\tilde{c}_{\alpha}}-1\right)\left(\hat{T}_{-\alpha} v^{c_{\alpha}-\tilde{c}_{\alpha}}+1\right)
$$

Besides, we let

$$
\mathbf{e}(\alpha)=-\tilde{G}(\alpha)
$$

where $\tilde{G}(\alpha)$ is again right multiplication with $\tilde{\mathcal{G}}(\alpha)$.
Let $x \in W_{0}$. For $\lambda \in Q$, we can regard $\Theta_{\lambda}$ as an element of $\operatorname{End}_{\mathcal{H}}\left(\hat{\mathcal{H}}_{x}\right)$ (and $\theta_{\lambda}$ as an element of $\operatorname{End}_{\mathcal{H}}\left(\hat{\mathcal{M}}_{x}\right)$ ). The operators of the form $1-p \Theta_{\alpha}, p \in \mathcal{A}, \alpha \in \Phi$, are invertible on $\hat{\mathcal{H}}_{x}$ with inverses as follows. Consider the formal sums

$$
\left(1-p \Theta_{\alpha}\right)_{+}^{-1}= \begin{cases}\sum_{n \geq 0} p^{n} \Theta_{n \alpha} & \text { if } \alpha>0 \\ -p \Theta_{\alpha} \sum_{n \geq 0} p^{-n} \Theta_{-n \alpha} & \text { if } \alpha<0\end{cases}
$$

and

$$
\left(1-p \Theta_{\alpha}\right)_{-}^{-1}= \begin{cases}-p \Theta_{\alpha} \sum_{n \geq 0} p^{-n} \Theta_{-n \alpha} & \text { if } \alpha>0 \\ \sum_{n \geq 0} p^{n} \Theta_{n \alpha} & \text { if } \alpha<0\end{cases}
$$

Then

$$
\left(1-p \Theta_{\alpha}\right)^{-1}=\left\{\begin{array}{ll}
\left(1-p \Theta_{\alpha}\right)_{+}^{-1} & \text { if } \alpha>0, \alpha x^{-1}<0 \\
\left(1-p \Theta_{\alpha}\right)_{-}^{-1} & \text { otherwise }
\end{array} \text { or if } \alpha<0, \alpha x^{-1}>0\right.
$$

The map $\tilde{\Theta}_{s}, s \in S_{0}$, can be regarded as an element of $\operatorname{Hom}_{\mathcal{H}}\left(\hat{\mathcal{H}}_{x}, \hat{\mathcal{H}}_{x s}\right)$. For $x \in W_{0}$, $s=s_{\alpha}, \alpha \in \Pi$, we can therefore define an operator

$$
\Theta_{s}^{+}=\mathbf{e}(\alpha)^{-1} \tilde{\Theta}_{s}
$$

belonging to $\operatorname{Hom}_{\mathcal{H}}\left(\hat{\mathcal{H}}_{x}, \hat{\mathcal{H}}_{x s}\right)$. By definition, we have $\left.\Theta_{s}^{+}\right|_{\mathcal{H}^{0}}=\Theta_{s}$. The corresponding operator in $\operatorname{Hom}_{\mathcal{H}}\left(\hat{\mathcal{M}}_{x}, \hat{\mathcal{M}}_{x s}\right)$ will be denoted by $\theta_{s}^{+}$.

Let $w_{0}=s_{1} \cdots s_{n}, s_{i} \in S_{0}, 1 \leq i \leq n$, be a reduced expression, and define

$$
\Theta_{w_{0}}^{+}=\Theta_{s_{n}}^{+} \cdots \Theta_{s_{1}}^{+} \in \operatorname{Hom}_{\mathcal{H}}\left(\hat{\mathcal{H}}, \hat{\mathcal{H}}_{w_{0}}\right)
$$

(which can be shown to be independent of the chosen reduced expression). Each $\Theta_{s_{i}}$, $1 \leq i \leq n$, is a map in $\operatorname{Hom}_{\mathcal{H}}\left(\hat{\mathcal{H}}_{x}, \hat{\mathcal{H}}_{x s_{i}}\right)$ for some $x \in W_{0}$ such that $x s_{i}>x$. Hence $\left(-\alpha_{i}\right)\left(x s_{i}\right)^{-1}>0$ (where $\alpha_{i} \in \Pi$ is such that $\left.s_{i}=s_{\alpha_{i}}\right)$, and

$$
\Theta_{s_{i}}^{+}=\left(1-v^{c_{\alpha_{i}}+\tilde{c}_{\alpha_{i}}} \Theta_{-\alpha_{i}}\right)_{+}^{-1}\left(1+v^{c_{\alpha_{i}}-\tilde{c}_{\alpha_{i}}} \Theta_{-\alpha_{i}}\right)_{+}^{-1} \tilde{\Theta}_{s_{i}}
$$

Since we also have $\left(1-p \Theta_{-\alpha}\right)^{-1}=\left(1-p \Theta_{-\alpha}\right)_{+}^{-1}$ on $\hat{\mathcal{H}}_{w_{0}}$ for any $\alpha>0$, we deduce, using Proposition 2.1.1, that

$$
\Theta_{w_{0}}^{+}=\prod_{\alpha>0} \mathrm{e}(\alpha)^{-1} \tilde{\Theta}_{w_{0}}
$$

where $\tilde{\Theta}_{w_{0}}=\tilde{\Theta}_{s_{n}} \cdots \tilde{\Theta}_{s_{1}}$. We now set

$$
\Psi=v^{-2 m\left(w_{0}\right)} \psi \Theta_{w_{0}}^{+}
$$

an $\mathcal{H}$-antilinear map on $\hat{\mathcal{H}}$.
Let $\preceq$ be the partial order on $W_{a}$ given by $y \preceq w, y, w \in W_{a}$, if and only if $y\left(A^{-}\right) \leq w\left(A^{-}\right)$. We point out that, since $x\left(A^{-}\right)=A^{-} x$ for $x \in W_{a}$, this definition agrees with the definition of $\preceq$ on $Q \subset W_{a}$.

Lemma 3.1.3. For $w \in W_{a}$, we have

$$
\Psi\left(\check{T}_{w}\right)=v^{-2 h(w)} \sum_{y \preceq w}(-1)^{g(w)-g(y)} \mathcal{R}_{y\left(A^{-}\right), w\left(A^{-}\right)} \check{T}_{y}
$$

(This follows from the correspondence between $\Psi$ and the $\mathcal{H}$-antilinear map $\hat{\Phi}_{\gamma}$ on $\hat{\mathcal{M}}$.)
3.2. Alcove polynomials. We introduce elements $E_{C} \in \hat{\mathcal{M}}, C \in X$, whose coefficients generically equal the Kazhdan-Lusztig polynomials.

Definition 3.2.1. The alcove polynomials $\hat{P}_{y, w} \in \mathcal{A}^{+}, y, w \in W_{a}$, are defined by the conditions
(i) $v^{-2 h(w)} \hat{P}_{y, w}=\sum_{y \preceq z \preceq w} v^{-2 h(y)} \overline{\mathcal{R}_{z\left(A^{+}\right), y\left(A^{+}\right)}} \overline{\hat{P}_{z, w}}$ and
(ii) the degree of $\hat{P}_{y, w}$ is less than $h(w)-h(y)$ if $y \prec w$ and $\hat{P}_{w, w}=1$.

For $A=y\left(A^{-}\right)$and $C=w\left(A^{-}\right)$, we set $\hat{P}_{A, C}=\hat{P}_{y, w}$.
Let

$$
D=\left\{w \in W_{a} \mid w\left(A^{-}\right) \subset \mathcal{C}^{+}\right\}
$$

Elements in $N D$ are called dominant. If $y, w \in D$ and $y\left(A^{-}\right), w\left(A^{-}\right)$are sufficiently far from the walls of $\mathcal{C}^{+}$, it turns out that $\hat{P}_{y, w}=P_{y, w}$.

For $\alpha \in \Phi$, put

$$
\mathbf{f}(\alpha)= \begin{cases}1-v^{-2 c_{\alpha}} \theta_{-\alpha} & \text { if } c_{\alpha}=\tilde{c}_{\alpha} \\ \left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \theta_{-\alpha}\right) & \text { otherwise }\end{cases}
$$

(The condition $c_{\alpha}=\tilde{c}_{\alpha}$ is equivalent to the condition $\check{\alpha} \notin 2 U$ where $U$ is as in the proof of Proposition 2.1.1.)

Since all $\mathbf{f}(\alpha), \alpha \in \Phi^{+}$, are invertible on $\hat{\mathcal{M}}$, we can define elements $E_{C} \in \hat{\mathcal{M}}$, $C \in X$, by

$$
E_{C}=\prod_{\alpha>0} \mathbf{f}(\alpha)^{-1} D_{C}
$$

In the remainder of this section, we prove the following theorem.
Theorem 3.2.2. (comp. [9], Theorem 3.5) For any $C \in X$, we have

$$
E_{C}=\sum_{A \leq C} \hat{P}_{A, C} A
$$

Let $C \in X$. We set

$$
E_{C}^{\prime}=\sum_{A \leq C} \hat{P}_{A, C} A \in \hat{\mathcal{M}}
$$

and we define an $\mathcal{H}$-antilinear map $\tilde{\Phi}$ on $\hat{\mathcal{M}}$ given by

$$
\tilde{\Phi}(A)=\sum_{B \in X} v^{-2 \gamma(B)} \overline{\mathcal{R}_{A w_{0}, B w_{0}}} B
$$

for $A \in X$.
The element $E_{C}^{\prime}$ can then be characterized by the conditions
(i) $\tilde{\Phi}\left(E_{C}^{\prime}\right)=v^{-2 \gamma(C)} E_{C}^{\prime}$ and
(ii) $E_{C}^{\prime}=\sum_{A \leq C} P_{A, C}^{\prime} A$ where $P_{A, C}^{\prime} \in \mathcal{A}^{+}$has degree less than $c(A, C)$ if $A<C$ and $P_{C, C}^{\prime}=1$.
We first show that $E_{C}$ satisfies (ii).
Let $\alpha>0$. Remember that the inverse of $1-p \theta_{-\alpha}, p \in \mathcal{A}$, on $\hat{\mathcal{M}}$ is given by

$$
\left(1-p \theta_{-\alpha}\right)^{-1}=\sum_{n \geq 0} p^{n} \theta_{-n \alpha} .
$$

First, suppose $c_{\alpha}=\tilde{c}_{\alpha}$. The degree of the coefficient of $A p_{-n \alpha}, A \in X, n \geq 0$, in $v^{-2 n c_{\alpha}} \theta_{-n \alpha} A$ is $c\left(A p_{-n \alpha}, A\right)-2 n c_{\alpha}$.

Next, suppose $c_{\alpha} \neq \tilde{c}_{\alpha}$. The degree of the coefficient of $A p_{-(m+n) \alpha}, A \in X, m, n \geq$ 0 , in $v^{-m\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \theta_{-m \alpha}(-1)^{n} v^{-n\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \theta_{-n \alpha} A$ is $c\left(A p_{-(m+n) \alpha}, A\right)-m\left(c_{\alpha}+\tilde{c}_{\alpha}\right)-n\left(c_{\alpha}-\tilde{c}_{\alpha}\right)$.

So if $\beta=\sum_{\alpha>0} n_{\alpha} \alpha \neq 0, n_{\alpha} \geq 0$, then the degree of the coefficient of $A p_{\beta}$ in $\prod_{\alpha>0} \mathbf{f}(\alpha)^{-1} A, A \in X$, is less than $c\left(A p_{\beta}, A\right)$. (Also, notice that in both cases the degree is non-negative.) This, together with the facts that $Q_{A, C}$ has degree less than $c(A, C)$ if $A<C$ and $Q_{C, C}=1$ shows that $E_{C}$ satisfies condition (ii).

Condition (i) is implied by the following theorem.

Theorem 3.2.3. (comp. [9], Theorem 3.7) We can write $\tilde{\Phi}$ as

$$
\tilde{\Phi}=\prod_{\alpha>0} \mathbf{f}(\alpha)^{-1} \hat{\Phi}_{\gamma} \prod_{\alpha>0} \mathbf{f}(\alpha) .
$$

The rest of this section will be taken by the proof of this theorem.
We have an $\mathcal{A}$-valued symmetric bilinear form $(\cdot, \cdot)$ on $\mathcal{H}$ satisfying

$$
\left(\check{T}_{w}, \check{T}_{y}\right)=v^{2 h(w)} \delta_{y, w}
$$

where $y, w \in W_{a}$. We also write $(\cdot, \cdot)$ for the corresponding bilinear form on $\mathcal{M}$.
Let $\mathcal{H}^{*}$ be the $\mathcal{A}$-module consisting of all formal sums $\sum_{w \in W_{a}} a_{w} T_{w}, a_{w} \in \mathcal{A}$ for $w \in W_{a}$. The domain of $(\cdot, \cdot)$ can be extended to $\mathcal{H}^{*} \times \mathcal{H}$ (and $\mathcal{H} \times \mathcal{H}^{*}$ ), and we can consider $\mathcal{H}^{*}$ as the dual space of $\mathcal{H}$ via $f(g)=(f, g)$ for $f \in \mathcal{H}^{*}, g \in \mathcal{H}$. For an arbitrary $\mathcal{A}$-linear map $\pi: \mathcal{H} \rightarrow \mathcal{H}^{*}$, the formal adjoint operator $\pi^{*}: \mathcal{H} \rightarrow \mathcal{H}^{*}$ is then defined by

$$
\left(\pi\left(\check{T}_{w}\right), \check{T}_{y}\right)=\left(\check{T}_{w}, \pi^{*}\left(\check{T}_{y}\right)\right)
$$

for $y, w \in W_{a}$. We will also need to extend the domain of $\pi^{*}$ to $\hat{\mathcal{H}}_{x}$ for some $x \in W_{0}$. If, for example, $\pi$ is of the form

$$
\pi\left(\check{T}_{w}\right)=\sum_{y \succeq w} a_{y, w} \check{T}_{y}
$$

with $a_{y, w} \in \mathcal{A}, y, w \in W_{a}$, we have

$$
\pi^{*}\left(\check{T}_{y}\right)=v^{2 h(y)} \sum_{w \preceq y} a_{y, w} v^{-2 h(w)} \check{T}_{w}
$$

for $y \in W_{a}$. Hence $\pi^{*}$ can be regarded as an element of $\operatorname{End}_{\mathcal{A}}(\hat{\mathcal{H}})$.
Formal adjoint operators on $\mathcal{M}$ (or $\hat{\mathcal{M}}_{x}, x \in W_{0}$ ) are defined accordingly.
We define an involutive ring automorphism $j$ on $\mathcal{H}$ by

$$
j\left(\sum_{w \in W_{a}} a_{w} T_{w}\right)=\sum_{w \in W_{a}} \overline{a_{w}}(-1)^{l(w)} v^{2 m(w)}
$$

$a_{w} \in \mathcal{A}$ for $w \in W_{a}$. The corresponding map in $\operatorname{Aut}_{\mathbb{Z}}(\mathcal{M})$ will also be denoted by $j$.
Let $w=x p_{\lambda} \in W_{a}, x \in W_{0}, \lambda \in Q$, and let $\mu \in Q \cap T^{+}$be such that $\lambda+\mu \in T^{+}$. We compute

$$
\begin{aligned}
j\left(\check{T}_{w}\right) & =j\left(T_{x} \check{T}_{\lambda}\right)=j\left(T_{x} T_{\lambda+\mu} T_{\mu}^{-1}\right) \\
& =(-1)^{l(x)+l\left(p_{\lambda+\mu}\right)-l\left(p_{\mu}\right)} v^{-2 m(x)-2 m\left(p_{\lambda+\mu}\right)+2 m\left(p_{\mu}\right)} T_{x} T_{\lambda+\mu} T_{\mu}^{-1} \\
& =(-1)^{g(w)} v^{-2 h(w)} \check{T}_{w} .
\end{aligned}
$$

It follows that

$$
j(A)=(-1)^{\delta(A)} v^{-2 \gamma(A)} A
$$

for $A \in X$.

We will also consider $j$ as a map on $\hat{\mathcal{H}}_{x}$ and $\hat{\mathcal{M}}_{x}, x \in W_{0}$. Note that $j^{*}=j$.
The following three lemmas contain several relations needed for the proof of Theorem 3.2.3. The calculations are straightforward, and we only show the first lemma.

Lemma 3.2.4. (comp. [9], Lemma 3.10) On $\mathcal{H}$ resp. $\mathcal{M}$, we have
(i) $j \Theta_{\lambda} j=\Theta_{\lambda}$ for $\lambda \in Q$,
(ii) $j \tilde{\Theta}_{s} j=-v^{-2 c_{s}} \tilde{\Theta}_{s}$ for $s \in S_{0}$,
(iii) $j \psi j=(-1)^{l\left(w_{0}\right)} v^{-2 m\left(w_{0}\right)} \psi$ and
(iv) $j \varphi j=(-1)^{l\left(w_{0}\right)} v^{-2 m\left(w_{0}\right)} \varphi$.

Moreover, (i) holds on $\hat{\mathcal{H}}_{x}, x \in W_{0}$, (ii) is an equality of operators in $\operatorname{Hom}_{\mathcal{H}}\left(\hat{\mathcal{H}}_{x}, \hat{\mathcal{H}}_{x s}\right)$, $x \in W_{0}$, and (iii) resp. (iv) is an equality of $\mathcal{H}$-antilinear operators in $\operatorname{Hom}_{\mathbb{Z}}\left(\hat{\mathcal{H}}, \hat{\mathcal{H}}_{w_{0}}\right)$. resp. $\operatorname{Hom}_{\mathbb{Z}}\left(\hat{\mathcal{M}}, \hat{\mathcal{M}}_{w_{0}}\right)$.

Proof. (i) Since

$$
j\left(\hat{T}_{\lambda}\right)=v^{h(\lambda)}(-1)^{2 \mathrm{ht}(\lambda)} v^{-2 h(\lambda)} \check{T}_{\lambda}=\hat{T}_{\lambda}
$$

we obtain

$$
j \Theta_{\lambda} j(f)=j\left(j(f) \hat{T}_{\lambda}\right)=f j\left(\hat{T}_{\lambda}\right)=\Theta_{\lambda}(f)
$$

for $f \in \hat{\mathcal{H}}_{x}, x \in W_{0}$.
(ii) If $s=s_{\alpha}, \alpha \in \Pi$, we have

$$
\begin{aligned}
j\left(\tilde{I}_{s}\right) & =-v^{-2 c_{s}} T_{s}\left(1-\hat{T}_{-2 \alpha}\right)-\hat{T}_{-\alpha}\left(v^{-\left(c_{s}+\tilde{c}_{s}\right)}-v^{-\left(c_{s}-\tilde{c}_{s}\right)}\right)+1-v^{2 c_{s}} \\
& =-v^{-2 c_{s}}\left(T_{s}\left(1-\hat{T}_{-2 \alpha}\right)-\hat{T}_{-\alpha}\left(v^{c_{s}+\tilde{c}_{s}}-v^{c_{s}-\tilde{c}_{s}}\right)+1-v^{-2 c_{s}}\right) \\
& =-v^{-2 c_{s}} \tilde{I}_{s} .
\end{aligned}
$$

For $f \in \hat{\mathcal{H}}_{x}, x \in W_{0}$, we conclude

$$
j \tilde{\Theta}_{s} j(f)=j\left(j(f) \tilde{I}_{s}\right)=-v^{-2 c_{s}} \tilde{\Theta}_{s}(f)
$$

(iii) We calculate

$$
\begin{aligned}
j \psi j\left(\check{T}_{w}\right) & =j \psi\left((-1)^{g(w)} v^{-2 h(w)} \check{T}_{w}\right) \\
& =j\left((-1)^{g(w)} v^{2 h(w)} \check{T}_{w} T_{w_{0}}\right) \\
& =(-1)^{l\left(w_{0}\right)} v^{-2 m\left(w_{0}\right)} \check{T}_{w} T_{w_{0}} \\
& =(-1)^{l\left(w_{0}\right)} v^{-2 m\left(w_{0}\right)} \psi\left(\check{T}_{w}\right)
\end{aligned}
$$

for $w \in W_{a}$.
(iv) This is equivalent to (iii).

Lemma 3.2.5. (comp. [9], Lemma 3.11) The map $\tilde{\Phi}$ can be expressed as

$$
\tilde{\Phi}=v^{-2 m\left(w_{0}\right)} j\left(j \varphi \hat{\Phi}_{\gamma} \varphi\right)^{*}
$$

Notice that $\left(\varphi \hat{\Phi}_{\gamma} \varphi\right)^{*}$ and therefore $j\left(j \varphi \hat{\Phi}_{\gamma} \varphi\right)^{*}$ can be considered as $\mathcal{H}$-antilinear elements of $\operatorname{End}_{\mathbb{Z}}(\hat{\mathcal{M}})$.

Lemma 3.2.6. (comp. [9], Lemma 3.13) On $\mathcal{H}$ resp. $\mathcal{M}$ and moreover as elements of $\operatorname{Hom}_{\mathbb{Z}}\left(\hat{\mathcal{H}}, \hat{\mathcal{H}}_{w_{0}}\right)$ resp. $\operatorname{Hom}_{\mathbb{Z}}\left(\hat{\mathcal{M}}, \hat{\mathcal{M}}_{w_{0}}\right)$, we have
(i) $(j \psi)^{*}=(-1)^{l\left(w_{0}\right)} j \psi$ and
(ii) $(j \varphi)^{*}=(-1)^{l\left(w_{0}\right)} j \varphi$.

Let

$$
\tilde{\Psi}=v^{-2 m\left(w_{0}\right)} j(j \psi \Psi \psi)^{*}
$$

be the operator on $\hat{\mathcal{H}}$ corresponding to $\tilde{\Phi}$. By the previous lemma, we get

$$
\begin{aligned}
\tilde{\Psi} & =j\left(j \psi \psi \Theta_{w_{0}}^{+} \psi\right)^{*} \\
& =j(j \psi)^{*}\left(j \Theta_{w_{0}}^{+} j\right)^{*} \\
& =(-1)^{l\left(w_{0}\right)} \psi\left(j \Theta_{w_{0}}^{+} j\right)^{*}
\end{aligned}
$$

For $s \in S_{0}$, we define $\mathcal{A}$-linear maps $\beta_{s}$ and $\rho_{s}$ on $\mathcal{H}$ by

$$
\beta_{s}\left(T_{x} \check{T}_{\lambda}\right)=T_{x} T_{s} \check{T}_{\lambda}
$$

and

$$
\rho_{s}\left(T_{x} \check{T}_{\lambda}\right)=v^{h(\lambda-\lambda s)} T_{x} \check{T}_{\lambda s}
$$

where $x \in W_{0}$ and $\lambda \in Q$. The map $\beta_{s}$ can be considered as an element of $\operatorname{End}_{\mathcal{A}}\left(\hat{\mathcal{H}}_{x}\right)$ and $\rho_{s}$ as an element of $\operatorname{Hom}_{\mathcal{A}}\left(\hat{\mathcal{H}}_{x}, \hat{\mathcal{H}}_{x s}\right), x \in W_{0}$. Let $s=s_{\alpha}, \alpha \in \Pi$. For $x \in W_{0}$, $\lambda \in Q$, we have

$$
\begin{aligned}
& \left\{\left(\beta_{s}\left(1-\Theta_{-2 \alpha}\right)-\Theta_{-\alpha}\left(v^{c_{s}+c_{\tilde{s}}}-v^{c_{s}-c_{\tilde{s}}}\right)+1-v^{2 c_{s}}\right) \rho_{s}\right\}\left(T_{x} \check{T}_{\lambda}\right) \\
& \quad=v^{h(\lambda-\lambda s)}\left(\beta_{s}\left(1-\Theta_{-2 \alpha}\right)-\Theta_{-\alpha}\left(v^{c_{s}+c_{\tilde{s}}}-v^{c_{s}-c_{\tilde{s}}}\right)+1-v^{2 c_{s}}\right)\left(T_{x} \check{T}_{\lambda s}\right) \\
& \quad=v^{h(\lambda-\lambda s)} T_{x}\left(T_{s} \check{T}_{\lambda s}\left(1-\hat{T}_{-2 \alpha}\right)-\check{T}_{\lambda s} \hat{T}_{-\alpha}\left(v^{c_{s}+c_{\bar{s}}}-v^{c_{s}-c_{\tilde{s}}}\right)+\check{T}_{\lambda s}\left(1-v^{2 c_{s}}\right)\right) \\
& \quad=v^{h(\lambda-\lambda s)} T_{x} \tilde{I}_{s} \check{T}_{\lambda s} \\
& \quad=T_{x} \check{T}_{\lambda} \tilde{I}_{s}
\end{aligned}
$$

where the last equality follows from Proposition 2.1.1. So

$$
\tilde{\Theta}_{s}=\left(\beta_{s}\left(1-\Theta_{-2 \alpha}\right)-\Theta_{-\alpha}\left(v^{c_{s}+c_{\tilde{s}}}-v^{c_{s}-c_{\tilde{s}}}\right)+1-v^{2 c_{s}}\right) \rho_{s}
$$

Arguments analogous to the ones in loc. cit. then demonstrate the following lemma.
Lemma 3.2.7. (comp. [9], Lemma 3.15) On $\mathcal{H}$ and as elements of $\operatorname{Hom}_{\mathcal{A}}\left(\hat{\mathcal{H}}_{x}, \hat{\mathcal{H}}_{x s}\right)$, $x \in W_{0}, s \in S_{0}$, we have

$$
\left(\tilde{\Theta}_{s}\right)^{*}=\tilde{\Theta}_{s}
$$

The next lemma contains a few more equalities which are all easy consequences of the definitions.

Lemma 3.2.8. On $\mathcal{H}$, we have
(i) $\Theta_{\lambda}^{*}=\Theta_{-\lambda}$ for $\lambda \in Q$,
(ii) $\tilde{\Theta}_{s} \Theta_{\alpha}=\Theta_{\alpha s} \tilde{\Theta}_{s}$ for $s \in S_{0}, \alpha \in \Phi$, and
(iii) $\psi \Theta_{\alpha}=\Theta_{-\alpha} \psi$ for $\alpha \in \Phi$.

Furthermore, (i) and (ii) hold on $\hat{\mathcal{H}}_{x}, x \in W_{0}$, and (iii) on $\hat{\mathcal{H}}$.

We remark that (ii) implies

$$
\begin{equation*}
\Theta_{w_{0}}^{+} \Theta_{\alpha}=\Theta_{-\alpha} \Theta_{w_{0}}^{+} \tag{*}
\end{equation*}
$$

for $\alpha \in \Phi$.
We now prove Theorem 3.2.3 or rather the equivalent on $\hat{\mathcal{H}}$.
Recall that $\Theta_{w_{0}}^{+}=\prod_{\alpha>0} \mathbf{e}(\alpha)^{-1} \tilde{\Theta}_{w_{0}}$ and that for $p \in \mathcal{A}$ and $\alpha>0$, we have

$$
\left(1-p \Theta_{-\alpha}\right)^{-1}=-p \Theta_{-\alpha} \sum_{n \geq 0} p^{-n} \Theta_{n \alpha}
$$

on $\hat{\mathcal{H}}_{w_{0}}$ and

$$
\left(1-p \Theta_{\alpha}\right)^{-1}=-p \Theta_{\alpha} \sum_{n \geq 0} p^{-n} \Theta_{-n \alpha}
$$

on $\hat{\mathcal{H}}$. Therefore, using Lemma 3.2.4 (ii) and (i), Lemma 3.2.7 and Lemma 3.2.8 (i), we can write

$$
\left(j \Theta_{w_{0}}^{+} j\right)^{*}=(-1)^{l\left(w_{0}\right)} v^{-2 m\left(w_{0}\right)} \tilde{\Theta}_{w_{0}} \prod_{\alpha>0}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1}
$$

Thus

$$
\begin{aligned}
\tilde{\Psi}= & (-1)^{l\left(w_{0}\right)} \psi\left(j \Theta_{w_{0}}^{+} j\right)^{*} \\
= & v^{2 m\left(w_{0}\right)} \psi \tilde{\Theta}_{w_{0}} \prod_{\alpha>0}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1} \\
= & v^{2 m\left(w_{0}\right)} \psi \prod_{\alpha>0}\left(1-v^{c_{\alpha}+\tilde{c}_{\alpha}} \Theta_{-\alpha}\right)\left(1+v^{c_{\alpha}-\tilde{c}_{\alpha}} \Theta_{-\alpha}\right) \Theta_{w_{0}}^{+} \\
& \quad \prod_{\alpha>0}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1}
\end{aligned}
$$

Set $\Phi_{1}^{+}=\left\{\alpha \in \Phi^{+} \mid c_{\alpha}=\tilde{c}_{\alpha}\right\}$ and $\Phi_{2}^{+}=\left\{\alpha \in \Phi^{+} \mid c_{\alpha} \neq \tilde{c}_{\alpha}\right\}$. Then

$$
\begin{aligned}
& \tilde{\Psi}= v^{2 m\left(w_{0}\right)} \psi \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{2 c_{\alpha}} \Theta_{-\alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{c_{\alpha}+\tilde{c}_{\alpha}} \Theta_{-\alpha}\right)\left(1+v^{c_{\alpha}-\tilde{c}_{\alpha}} \Theta_{-\alpha}\right) \\
& \cdot \Theta_{w_{0}}^{+} \prod_{\alpha \in \Phi_{1}^{+}}\left(1+\Theta_{\alpha}\right) \prod_{\alpha \in \Phi_{1}^{+}}\left(1+\Theta_{\alpha}\right)^{-1} \\
&=\prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-2 c_{\alpha}} \Theta_{\alpha}\right)^{-1} \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{\alpha}\right)^{-1} \\
&=v^{2 m\left(w_{0}\right)} \psi \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-2 c_{\alpha}} \Theta_{-\alpha}\right)^{-1} \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1} \\
& \cdot \Theta_{w_{0}}^{+} \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{2 c_{\alpha}} \Theta_{\alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{c_{\alpha}+\tilde{c}_{\alpha}} \Theta_{\alpha}\right)\left(1+v^{c_{\alpha}-\tilde{c}_{\alpha}} \Theta_{\alpha}\right)
\end{aligned}
$$

For the above calculations, we used consequence (*) of Lemma 3.2.8 (ii). Lemma 3.2.8 (iii) now yields

$$
\begin{gathered}
\tilde{\Psi}=v^{2 m\left(w_{0}\right)} \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{2 c_{\alpha}} \Theta_{\alpha}\right)^{-1} \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{c_{\alpha}+\tilde{c}_{\alpha}} \Theta_{\alpha}\right)^{-1}\left(1+v^{c_{\alpha}-\tilde{c}_{\alpha}} \Theta_{\alpha}\right)^{-1} \\
\cdot \psi \Theta_{w_{0}}^{+} \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{2 c_{\alpha}} \Theta_{\alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{c_{\alpha}+\tilde{c}_{\alpha}} \Theta_{\alpha}\right)\left(1+v^{c_{\alpha}-\tilde{c}_{\alpha}} \Theta_{\alpha}\right) .
\end{gathered}
$$

Using

$$
\sum_{\alpha \in \Phi_{1}^{+}} 2 c_{\alpha}+\sum_{\alpha \in \Phi_{2}^{+}}\left(c_{\alpha}+\tilde{c}_{\alpha}+c_{\alpha}-\tilde{c}_{\alpha}\right)=\sum_{\alpha>0} 2 c_{\alpha}=2 m\left(w_{0}\right)
$$

we can rewrite the last expression for $\tilde{\Psi}$ as

$$
\begin{aligned}
& \tilde{\Psi}= v^{2 m\left(w_{0}\right)} v^{-2 m\left(w_{0}\right)}(-1)^{l\left(w_{0}\right)} \prod_{\alpha \in \Phi_{1}^{+}} \Theta_{\alpha}^{-1} \prod_{\alpha \in \Phi_{2}^{+}} \Theta_{2 \alpha}^{-1} \\
& \cdot \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-2 c_{\alpha}} \Theta_{-\alpha}\right)^{-1} \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1} \\
& \cdot \psi \Theta_{w_{0}}^{+} v^{2 m\left(w_{0}\right)}(-1)^{l\left(w_{0}\right)} \prod_{\alpha \in \Phi_{1}^{+}} \Theta_{\alpha} \prod_{\alpha \in \Phi_{2}^{+}} \Theta_{2 \alpha} \\
&=\prod^{-2 m\left(w_{0}\right)} \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-2 c_{\alpha}} \Theta_{-\alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right) \\
&\left.\quad \cdot \psi \Theta_{-\alpha}^{+}\right)^{-1} \prod_{w_{0}} \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1} \\
&=\left.\Theta_{-\alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right) .
\end{aligned}
$$

We conclude

$$
\begin{aligned}
\tilde{\Psi}= & \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-2 c_{\alpha}} \Theta_{-\alpha}\right)^{-1} \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1}\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)^{-1} \Psi \\
& \cdot \prod_{\alpha \in \Phi_{1}^{+}}\left(1-v^{-2 c_{\alpha}} \Theta_{-\alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(1-v^{-\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \Theta_{-\alpha}\right)
\end{aligned}
$$

which is equivalent to Theorem 3.2.3.
3.3. The Kazhdan-Lusztig polynomials. We now generalize the formula in [9] for Kazhdan-Lusztig polynomials. They appear as coefficients of elements $F_{C} \in \hat{\mathcal{M}}$, $C \in X$. At the end of the section, we describe an analogue of Kostant's partition function, generalizing part of [8].

We start with introducing a third right action of $W_{0}$ on $X$. Let $A=y\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$, $y \in W_{a}, \lambda \in T$. For $s \in S_{0}$ with $c_{s}=c_{\tilde{s}}$, we set

$$
A \circledast s=y\left(A_{(\lambda+\rho) s-\rho}^{+}\right) \subseteq \Pi_{(\lambda+\rho) s-\rho} .
$$

For $s \in S_{0}$ with $c_{s} \neq c_{\tilde{s}}$, we set

$$
A \circledast s=y\left(A_{(\lambda+2 \rho) s-2 \rho}^{+}\right) \subseteq \Pi_{(\lambda+2 \rho) s-2 \rho} .
$$

Hence for $w \in W_{0}$, we have

$$
A \circledast w=y\left(A_{(\lambda+\tilde{\rho}) w-\tilde{\rho}}^{+}\right) \subseteq \Pi_{(\lambda+\tilde{\rho}) w-\tilde{\rho}}
$$

where

$$
\tilde{\rho}=\frac{1}{2} \sum_{\alpha \in \Phi_{1}^{+}} \alpha+\sum_{\alpha \in \Phi_{2}^{+}} \alpha .
$$

If $A=A^{-} y, C=A^{-} w, y, w \in W_{a}$, we let $P_{A, C}=P_{y, w}$. The next theorem and corollary contain the main statements of this section.

Theorem 3.3.1. (comp. [9], Theorem 4.2) Let $C$ be an alcove in $\mathcal{C}^{+}$. The coefficient of any alcove $A \subset \mathcal{C}^{+}$in

$$
F_{C}=\sum_{x \in W_{0}}(-1)^{l(x)} v^{c(C \circledast x, C)} E_{C \circledast x}
$$

equals $P_{A, C}$.
Remark 3.3.2. As in loc. cit., the above theorem implies that $P_{A, C}, A, C \subset \mathcal{C}^{+}$, can be calculated from $P_{A^{\prime}, C^{\prime}}$ for $A^{\prime}, C^{\prime} \subset \Pi$. Furthermore, for equal parameters, the above formula generalizes the $q$-analogue of Kostant's weight multiplicity formula found in [8].

Corollary 3.3.3. (comp. [9], Corollary 4.3) For any alcoves $A, C \subset \mathcal{C}^{+}$, we have

$$
P_{A, C}=\sum_{x \in W_{0}}(-1)^{l(w)} v^{c(C \circledast x, C)} \hat{P}_{A, C \circledast x} .
$$

The corollary is an immediate consequence of the preceding theorem. In order to prove the theorem, we first have to generalize a result of Andersen in [1]. Let $\mathcal{N}$ be the $\mathcal{H}$-submodule of $\mathcal{M}$ generated by $A^{+}+A^{+} s$ where $s$ runs over $S_{0}$, and let

$$
\hat{\mathcal{N}}^{\perp}=\left\{f \in \hat{\mathcal{M}}_{w_{0}} \mid(n, f)=0 \text { for all } n \in \mathcal{N}\right\}
$$

Instead of considering $\mathcal{M}$, Andersen introduces a module $\mathcal{M}_{-}$which is a free $\mathcal{A}$ module on $\mathcal{C}^{-}$and has an $\mathcal{H}$-module structure given by

$$
T_{s} A= \begin{cases}-A & \text { if } s A \notin \mathcal{C}^{-} \\ s A & \text { if } s A \in \mathcal{C}^{-} \\ v^{2 c_{s}} s A+\left(v^{2 c_{s}}-1\right) A & \text { if } s A<A\end{cases}
$$

for $s \in S_{a}$ and $A \subset \mathcal{C}^{-}$. It turns out that $\mathcal{M}_{-}$is isomorphic to $\mathcal{M} / \mathcal{N}$. By modifying [14], Andersen obtains elements $D_{C}^{\prime}$ for all alcoves $C \subset \mathcal{C}^{-}$whose coefficients $Q_{A, C}^{\prime}$ are exactly inverse Kazhdan-Lusztig polynomials. The polynomials $Q_{A, C}^{\prime}$ satisfy an induction formula which coincides with a formula for the composition factors of Weyl modules. Andersen obtains the latter formula by "inverting" the Lusztig conjecture in [13].

The result we need is the following one, which for equal parameters reformulates [1], Theorem 7.2.

Theorem 3.3.4. For alcoves $A, B \subset \mathcal{C}^{-}$, set

$$
\mathcal{R}_{B, A}^{\prime}=\sum_{x \in W_{0}} R_{y, w x}
$$

where $y, w \in D$ are such that $y\left(A^{+}\right)=A$ and $w\left(A^{+}\right)=B$. Then

$$
\left(\hat{\Phi}_{\gamma}(A), f\right)=v^{-2 \gamma(A)} \sum_{B \subset \mathcal{C}^{-}, B \leq A}(-1)^{d(A, B)} \mathcal{R}_{B, A}^{\prime}(B, f)
$$

for any $f \in \hat{\mathcal{N}}^{\perp}$.
Note that $\mathcal{R}_{B, A}^{\prime}=0$ unless $B \leq A$.
We comment on those parts of the proof of the above theorem which require real changes.

Let $\alpha \in \Pi$. We denote by $\mathcal{N}_{\alpha}$ the $\mathcal{H}$-submodule of $\mathcal{M}$ generated by $A^{+}+A^{+} s_{\alpha}$, and for $n \in \mathbb{Z}$, we set

$$
U_{n}=\{v \in V \mid n<\langle v, \check{\alpha}\rangle<n+1\} .
$$

(Note loc. cit., Remark 7.4.) For any alcove $A$ in $U_{0}$, we define alcoves $A_{n}, n \in \mathbb{Z}$, inductively by setting $A=A_{0}$ and $A_{n}=A_{n-1} \sigma_{\alpha, n}$. We write

$$
f_{A_{n}}=A_{n}+A_{n+1}
$$

where $n \in \mathbb{Z}$ and

$$
\begin{aligned}
g_{n}(A) & =v^{2 c\left(A, A_{n-1}\right)} A_{-n}+\sum_{i=1}^{n-1}\left(v^{2 c\left(A, A_{n-i}\right)}-v^{2 c\left(A, A_{n-i-1}\right)}\right) A_{2 i-n-1}+A_{n-1} \\
& =v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}-\sum_{i=1}^{n-1} v^{2 c\left(A, A_{n-i-1}\right)}\left(f_{A_{2 i-n-1}}-f_{A_{2 i-n}}\right)
\end{aligned}
$$

where $n \in \mathbb{N}$.

Lemma 3.3.5. (comp. [1], Lemma 4.3)
(i) The elements $g_{n}(A), A$ an alcove in $U_{0}$ and $n \in \mathbb{N}$, constitute a basis for $\mathcal{N}_{\alpha}$ as an $\mathcal{A}$-module.
(ii) Let $A$ be an alcove in $U_{0}$. We have

$$
v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}+f_{A_{n-2}} \in \mathcal{N}_{\alpha}
$$

if either $c_{\alpha}=\tilde{c}_{\alpha}$ and $n \in \mathbb{N}$ arbitrary, or if $c_{\alpha} \neq \tilde{c}_{\alpha}$ and $n \in \mathbb{N}$ odd or equal to 2. If $c_{\alpha} \neq \tilde{c}_{\alpha}$ and $n \in \mathbb{N}$ even, we have

$$
v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}+v^{2 c\left(A, A_{1}\right)}\left(v^{2 c\left(A_{1}, A_{n-2}\right)} f_{A_{-n+2}}+f_{A_{n-4}}\right)+f_{A_{n-2}} \in \mathcal{N}_{\alpha} .
$$

Proof. (i) for any $\alpha \in \Pi$ and (ii) for $\alpha \in \Pi \cap \Phi_{1}^{+}$follow as in loc. cit.
Suppose $\alpha \in \Pi \cap \Phi_{2}^{+}$. For $n \in \mathbb{N}$, we denote the respective expression in (ii) by $F_{n}(A)$. (We take the first expression for $n=2$. The second expression then equals $2 F_{2}(A)$.) It is easily checked that

$$
\begin{aligned}
F_{1}(A) & =v^{2 c\left(A, A_{0}\right)} f_{A_{-1}}+f_{A_{-1}}=2 g_{1}(A) \in \mathcal{N}_{\alpha}, \\
F_{2}(A) & =v^{2 c\left(A, A_{1}\right)} f_{A_{-2}}+f_{A}=g_{2}(A)+g_{1}(A) \in \mathcal{N}_{\alpha}, \\
F_{3}(A) & =v^{2 c\left(A, A_{2}\right)} f_{A_{-3}}+f_{A_{1}}=g_{3}(A)+F_{2}(A)-v^{2 c\left(A, A_{1}\right)} g_{1}(A) \in \mathcal{N}_{\alpha}, \text { and } \\
F_{4}(A) & =v^{2 c\left(A, A_{3}\right)} f_{A_{-4}}+v^{2 c\left(A, A_{1}\right)}\left(v^{2 c\left(A_{1}, A_{2}\right)} f_{A_{-2}}+f_{A}\right)+f_{A_{2}} \\
& =g_{4}(A)+F_{3}(A)+v^{2 c\left(A, A_{1}\right)} g_{1}(A) \in \mathcal{N}_{\alpha} .
\end{aligned}
$$

For $n \geq 5$ we claim

$$
F_{n}(A)= \begin{cases}g_{n}(A)+\left(F_{n-1}(A)-v^{2 c\left(A, A_{1}\right)} F_{n-2}(A)\right)-v^{2 c\left(A, A_{2}\right)} g_{n-4}(A) & \text { if } n \text { odd } \\ g_{n}(A)+\left(F_{n-1}(A)+v^{2 c\left(A, A_{1}\right)} F_{n-3}(A)\right)-v^{2 c\left(A, A_{2}\right)} g_{n-4}(A) & \text { if } n \text { even. }\end{cases}
$$

Indeed, if $n$ is odd we obtain

$$
\begin{aligned}
\mathrm{RHS}= & v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}-\sum_{i=1}^{n-1} v^{2 c\left(A, A_{n-i-1}\right)}\left(f_{A_{2 i-n-1}}-f_{A_{2 i-n}}\right) \\
& +v^{2 c\left(A, A_{n-2}\right)} f_{A_{-n+1}}+v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+3}}+v^{2 c\left(A, A_{1}\right)} f_{A_{n-5}}+f_{A_{n-3}} \\
& \quad-v^{2 c\left(A, A_{1}\right)} v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+2}}-v^{2 c\left(A, A_{1}\right)} f_{A_{n-4}} \\
& \quad-v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+4}}+\sum_{i=3}^{n-3} v^{2 c\left(A, A_{n-i-1}\right)}\left(f_{A_{2 i-n-1}}-f_{A_{2 i-n}}\right) \\
= & v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}-v^{2 c\left(A, A_{n-2}\right)}\left(f_{A_{-n+1}}-f_{A_{-n+2}}\right)-v^{2 c\left(A, A_{n-3}\right)}\left(f_{A_{-n+3}}-f_{A_{-n+4}}\right) \\
& \quad-v^{2 c\left(A, A_{1}\right)}\left(f_{A_{n-5}}-f_{A_{n-4}}\right)-\left(f_{A_{n-3}}-f_{A_{n-2}}\right) \\
& \quad+v^{2 c\left(A, A_{n-2}\right)} f_{A_{-n+1}}+v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+3}}+v^{2 c\left(A, A_{1}\right)} f_{A_{n-5}}+f_{A_{n-3}} \\
= & \quad v^{2 c\left(A, A_{n-2}\right)} f_{A_{-n+2}}-v^{2 c\left(A, A_{1}\right)} f_{A_{n-4}}-v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+4}} \\
= & v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}+f_{A_{n-2}} \\
= & F_{n}(A) .
\end{aligned}
$$

For $n$ even, we obtain

$$
\begin{aligned}
\mathrm{RHS}= & v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}-\sum_{i=1}^{n-1} v^{2 c\left(A, A_{n-i-1}\right)}\left(f_{A_{2 i-n-1}}-f_{A_{2 i-n}}\right) \\
& +v^{2 c\left(A, A_{n-2}\right)} f_{A_{-n+1}}+f_{A_{n-3}}+v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+3}}+v^{2 c\left(A, A_{1}\right)} f_{A_{n-5}} \\
& \quad-v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+4}}+\sum_{i=3}^{n-3} v^{2 c\left(A, A_{n-i-1}\right)}\left(f_{A_{2 i-n-1}}-f_{A_{2 i-n}}\right) \\
= & v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}-v^{2 c\left(A, A_{n-2}\right)}\left(f_{A_{-n+1}}-f_{A_{-n+2}}\right)-v^{2 c\left(A, A_{n-3}\right)}\left(f_{A_{-n+3}}-f_{A_{-n+4}}\right) \\
& \quad-v^{2 c\left(A, A_{1}\right)}\left(f_{A_{n-5}}-f_{A_{n-4}}\right)-\left(f_{A_{n-3}}-f_{A_{n-2}}\right) \\
& +v^{2 c\left(A, A_{n-2}\right)} f_{A_{-n+1}}+f_{A_{n-3}}+v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+3}}+v^{2 c\left(A, A_{1}\right)} f_{A_{n-5}} \\
& \quad-v^{2 c\left(A, A_{n-3}\right)} f_{A_{-n+4}} \\
= & v^{2 c\left(A, A_{n-1}\right)} f_{A_{-n}}+v^{2 c\left(A, A_{n-2}\right)} f_{A_{-n+2}}+v^{2 c\left(A, A_{1}\right)} f_{A_{n-4}}+f_{A_{n-2}} \\
= & F_{n}(A)
\end{aligned}
$$

So $F_{n}(A) \in \mathcal{N}_{\alpha}$ for any $n \in \mathbb{N}$.

We need this lemma for the following proposition.
Proposition 3.3.6. (comp. [1], Proposition 5.1) Let $\lambda \in T$ and $w \in W_{0}$. We have

$$
e_{\lambda}-(-1)^{l(w)} v^{c\left(A_{\lambda w}^{+}, A_{\lambda}^{+}\right)} e_{\lambda w} \in \mathcal{N}
$$

Proof. As in loc.cit., it is enough to prove that

$$
e_{\lambda}+v^{c\left(A_{\lambda \sigma_{\alpha}}^{+}, A_{\lambda}^{+}\right)} e_{\lambda \sigma_{\alpha}} \in \mathcal{N}_{\alpha} \subseteq \mathcal{N}
$$

for all $\alpha \in \Pi$. Let $\alpha \in \Pi$ and $n \in \mathbb{Z}$ be such that $\lambda \in H_{\alpha, n}$. We write

$$
e_{\lambda}=\sum_{B \in X, \bar{B} \ni \lambda} B=\sum_{B \in X, \bar{B} \ni \lambda, B<B \sigma_{\alpha, n}}\left(B+B \sigma_{\alpha, n}\right)=\sum_{A} f_{A_{n-1}}
$$

where the last sum runs over all alcoves $A$ in $U_{0}$ such that $\lambda \in \overline{A_{n-1}}$. Similarly, we have

$$
e_{\lambda \sigma_{\alpha}}=\sum_{A} f_{A_{-n-1}}
$$

where the last sum runs over the same set of alcoves $A$ as before.
According to the definition of $T$, if $c_{\alpha} \neq \tilde{c}_{\alpha}$ the integer $n$ is necessarily even. So, since $c\left(A_{\lambda \sigma_{\alpha}}^{+}, A_{\lambda}^{+}\right)=2 c\left(A, A_{\lambda}^{+}\right)$where $A$ is the alcove in $U_{0}$ such that $A_{n}=A_{\lambda}^{+}$(hence $\left.A_{\lambda \sigma_{\alpha}}^{+}=A_{-n}\right)$, the claim follows from the previous lemma.

Notice that this proposition shows that the operator $\theta_{w} \in \operatorname{Aut} \mathcal{H}^{( }\left(\mathcal{M}^{0}\right), w \in W_{0}$, maps $\mathcal{M}^{0} \cap \mathcal{N}$ into itself and therefore induces an $\mathcal{H}$-linear automorphism of $\mathcal{M}^{0} / \mathcal{M}^{0} \cap$ $\mathcal{N}$. Theorem 3.3.4 can now be proven as in loc.cit.

Back to the proof of Theorem 3.3.1, it remains to check that

$$
j \varphi\left(F_{C}\right) \in \hat{\mathcal{N}}^{\perp}
$$

(The other arguments in [9] have a straightforward generalization.)
The left ideal $\mathcal{I}$ of $\mathcal{H}$ generated by $T_{w_{0}}+T_{w_{0} s}$ for any $s \in S_{0}$ is the submodule of $\mathcal{H}$ corresponding to $\mathcal{N}$. Hence the condition

$$
j \varphi(f) \in \hat{\mathcal{N}}^{\perp}
$$

for $f \in \hat{\mathcal{M}}$ is equivalent to the condition

$$
(\mathcal{I}, j \psi(g))=0
$$

for $g=\phi^{-1}(f) \in \hat{\mathcal{H}}$. And since

$$
j \psi\left(T_{w_{0}}+T_{w_{0} s}\right)=j\left(1+T_{s}\right)=1-v^{-2 c_{s}} T_{s},
$$

the latter condition is equivalent to

$$
\left(\mathcal{H}\left(T_{s}-v^{2 c_{s}}\right), g\right)=0
$$

If $s=s_{\alpha}, \alpha \in \Pi$, the operator $\tilde{\Theta}_{s}-\mathbf{e}(\alpha)$ on $\mathcal{H}$ is right multiplication with

$$
\begin{aligned}
& T_{s}\left(1-\hat{T}_{-2 \alpha}\right)-\hat{T}_{-\alpha}\left(v^{c_{s}+c_{\bar{s}}}-v^{c_{s}-c_{\bar{s}}}\right)+1-v^{2 c_{s}}-\left(1-\hat{T}_{-\alpha} v^{c_{s}+c_{\bar{s}}}\right)\left(1+\hat{T}_{-\alpha} v^{c_{s}-c_{\bar{s}}}\right) \\
& \quad=T_{s}\left(1-\hat{T}_{-2 \alpha}\right)-v^{2 c_{s}}+\hat{T}_{-2 \alpha} v^{2 c_{s}},
\end{aligned}
$$

and therefore $\left(1-\Theta_{-2 \alpha}\right)^{-1}\left(\tilde{\Theta}_{s}-\mathbf{e}(\alpha)\right)$ is equal to the map $h \mapsto h\left(T_{s}-v^{2 c_{s}}\right), h \in \mathcal{H}$. Hence the formal adjoint operator of this map is

$$
\left(\left(1-\Theta_{-2 \alpha}\right)^{-1}\left(\tilde{\Theta}_{s}-\mathbf{e}(\alpha)\right)^{*}=\left(\tilde{\Theta}_{s}-\mathbf{e}(-\alpha)\right)\left(1-\Theta_{2 \alpha}\right)^{-1}\right.
$$

which we consider as a map on $\hat{\mathcal{H}}$. The condition on $g$ becomes

$$
\left(\left(\tilde{\Theta}_{s}-\mathbf{e}(\alpha)\right)\left(1-\Theta_{2 \alpha}\right)^{-1} g, \mathcal{H}\right)=0
$$

i.e.

$$
\tilde{\Theta}_{s}\left(1-\Theta_{2 \alpha}\right)^{-1} g=\mathbf{e}(-\alpha)\left(1-\Theta_{2 \alpha}\right)^{-1} g
$$

We can apply $1-\Theta_{-2 \alpha}=-\Theta_{-2 \alpha}\left(1-\Theta_{2 \alpha}\right)$ on both sides and obtain the equivalent condition

$$
\tilde{\Theta}_{s} g=-\mathbf{e}(-\alpha) \Theta_{-2 \alpha} g
$$

Now, $\tilde{\Theta}_{s}=\mathbf{e}(\alpha) \Theta_{s}^{+}$and consequently

$$
\begin{aligned}
\tilde{\Theta}_{s} g & =\left(1-\Theta_{-\alpha} v^{c_{s}+c_{\tilde{s}}}\right)\left(1+\Theta_{-\alpha} v^{c_{s}-c_{\tilde{s}}}\right) \Theta_{s}^{+} g \\
& =-\Theta_{-2 \alpha} v^{2 c_{s}}\left(1-v^{-\left(c_{s}+c_{\tilde{z}}\right)} \Theta_{\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\tilde{s}}\right)} \Theta_{\alpha}\right) \Theta_{s}^{+} g \\
& =-\Theta_{-2 \alpha} v^{2 c_{s}} \Theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\tilde{s}}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\tilde{z}}\right)} \Theta_{-\alpha}\right) g
\end{aligned}
$$

On the other hand,

$$
-\mathbf{e}(-\alpha) \Theta_{-2 \alpha} g=v^{2 c_{s}}\left(1-v^{-\left(c_{s}+c_{\mathfrak{3}}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\mathfrak{z}}\right)} \Theta_{-\alpha}\right) g
$$

So we can write the condition on $g \in \hat{\mathcal{H}}$ in the form

$$
\begin{aligned}
& \Theta_{-2 \alpha} \Theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\bar{s}}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\overline{3}}\right)} \Theta_{-\alpha}\right) g \\
& =-\left(1-v^{-\left(c_{s}+c_{\tilde{s}}\right)} \Theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\tilde{s}}\right)} \Theta_{-\alpha}\right) g
\end{aligned}
$$

and the initial condition $j \varphi(f) \in \hat{\mathcal{N}}^{\perp}$ is equivalent to

$$
\begin{aligned}
& \theta_{-2 \alpha} \theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\bar{z}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\tilde{s}}\right)} \theta_{-\alpha}\right) f \\
& \quad=-\left(1-v^{-\left(c_{s}+c_{\bar{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\tilde{z}}\right)} \theta_{-\alpha}\right) f .
\end{aligned}
$$

We defined $E_{C}=\prod_{\beta>0} \mathbf{f}(\beta)^{-1} D_{C}$ for $C \in X$, and according to Proposition 2.3.12 we have $\theta_{s} D_{C}=v^{c(C * s, C)} D_{C * s}$.

Suppose $C=y\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}, y \in W_{a}, \lambda \in T$. We first consider the case that $c_{s}=c_{\tilde{s}}$. We have

$$
C \circledast s=y\left(A_{(\lambda+\tilde{\rho}) s-\tilde{\rho}}^{+}\right)=y\left(A_{\lambda s-\alpha}^{+}\right)=(C * s) p_{-\alpha}
$$

and

$$
c(C \circledast s, C)=\gamma\left(A^{+} p_{\alpha}\right)+c(C * s, C)=h\left(p_{\alpha}\right)+c(C * s, C) .
$$

Hence

$$
\begin{aligned}
\theta_{-2 \alpha} & \theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\mathfrak{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\mathfrak{s}}\right)} \theta_{-\alpha}\right) E_{C} \\
& =\theta_{-\alpha}\left(1+\theta_{\alpha}\right) \theta_{-\alpha} \theta_{s}^{+} \prod_{\beta>0, \beta \neq \alpha} \mathbf{f}(\beta)^{-1} D_{C} \\
& =\left(1+\theta_{-\alpha}\right) \prod_{\beta>0, \beta \neq \alpha} \mathbf{f}(\beta)^{-1} v^{c(C * s, C)-h\left(p_{-\alpha}\right)} D_{(C * s) p_{-\alpha}} \\
& =\left(1+\theta_{-\alpha}\right) \prod_{\beta>0, \beta \neq \alpha} \mathbf{f}(\beta)^{-1} v^{c(C \circledast s, C)} D_{C \circledast s} \\
& =\left(1-v^{-\left(c_{s}+c_{s}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\bar{s}}\right)} \theta_{-\alpha}\right) v^{c(C \circledast s, C)} E_{C \circledast s} .
\end{aligned}
$$

Secondly, if $c_{s} \neq c_{\tilde{s}}$ we have

$$
C \circledast s=y\left(A_{\lambda s-2 \alpha}^{+}\right)=(C * s) p_{-2 \alpha}
$$

and

$$
c(C \circledast s, C)=h\left(p_{2 \alpha}\right)+c(C * s, C) .
$$

Hence

$$
\begin{aligned}
& \theta_{-2 \alpha} \theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\mathfrak{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{s}\right)} \theta_{-\alpha}\right) E_{C} \\
& \quad=\prod_{\beta>0, \beta \neq \alpha} \mathbf{f}(\beta)^{-1} v^{c(C * s, C)-h\left(p_{-2 \alpha}\right)} D_{(C * s) p_{-2 \alpha}} \\
& \quad=\left(1-v^{-\left(c_{s}+c_{\mathfrak{z}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\mathfrak{s}}\right)} \theta_{-\alpha}\right) v^{c(C \circledast s, C)} E_{C \circledast s} .
\end{aligned}
$$

In both cases, we thereby obtain

$$
\begin{aligned}
& \theta_{-2 \alpha} \theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\bar{z}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\bar{s}}\right)} \theta_{-\alpha}\right) F_{C} \\
& \quad=\theta_{-2 \alpha} \theta_{s}^{+}\left(1-v^{-\left(c_{s}+c_{\bar{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\mathfrak{s}}\right)} \theta_{-\alpha}\right) \sum_{x \in W_{0}}(-1)^{l(x)} v^{c(C \oplus x, C)} E_{C \circledast x} \\
& \quad=\left(1-v^{-\left(c_{s}+c_{\bar{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\bar{s}}\right)} \theta_{-\alpha}\right) \sum_{x \in W_{0}}(-1)^{l(x)} v^{c(C \circledast x, C)+c(C \circledast x \circledast s, C \circledast x)} E_{C \circledast x \circledast s} \\
& \quad=-\left(1-v^{-\left(c_{s}+c_{\bar{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\bar{s}}\right)} \theta_{-\alpha}\right) \sum_{x \in W_{0}}(-1)^{l(x s)} v^{c(C \oplus x s, C)} E_{C \circledast x s} \\
& \quad=-\left(1-v^{-\left(c_{s}+c_{\bar{s}}\right)} \theta_{-\alpha}\right)\left(1+v^{-\left(c_{s}-c_{\bar{s}}\right)} \theta_{-\alpha}\right) F_{C} .
\end{aligned}
$$

Thus $j \varphi\left(F_{C}\right) \in \hat{\mathcal{N}}^{\perp}$, which finishes the proof of Theorem 3.3.1.

We conclude this section with the generalization of Theorem 1.8 in [8].
If $K=\left(k_{\alpha}\right)_{\alpha} \in \mathbb{Z}_{+}^{\left|\Phi_{1}^{+}\right|}, M=\left(m_{\alpha}\right)_{\alpha} \in \mathbb{Z}_{+}^{\left|\Phi_{2}^{+}\right|}$and $N=\left(n_{\alpha}\right)_{\alpha} \in \mathbb{Z}_{+}^{\left|\Phi_{2}^{+}\right|}$, we set

$$
\sigma(N)=\sum_{\alpha \in \Phi_{2}^{+}} n_{\alpha}
$$

and

$$
\sigma^{\prime}(K, M, N)=2 \sum_{\alpha \in \Phi_{1}^{+}} k_{\alpha} c_{\alpha}+\sum_{\alpha \in \Phi_{2}^{+}}\left(m_{\alpha}\left(c_{\alpha}+\tilde{c}_{\alpha}\right)+n_{\alpha}\left(c_{\alpha}-\tilde{c}_{\alpha}\right)\right)
$$

For $\kappa \in T, \kappa \geq 0$, we then define

$$
\mathcal{P}^{*}(\kappa ; v)=\sum_{K, M, N}(-1)^{\sigma(N)} v^{\sigma^{\prime}(K, M, N)}
$$

where the sum runs over all $K=\left(k_{\alpha}\right) \in \mathbb{Z}_{+}^{\left|\Phi_{1}^{+}\right|}, M=\left(m_{\alpha}\right) \in \mathbb{Z}_{+}^{\left|\Phi_{2}^{+}\right|}, N=\left(n_{\alpha}\right) \in \mathbb{Z}_{+}^{\left|\Phi_{2}^{+}\right|}$ such that

$$
\sum_{\alpha \in \Phi_{1}^{+}} k_{\alpha} \alpha+\sum_{\alpha \in \Phi_{2}^{+}}\left(m_{\alpha}+n_{\alpha}\right) \alpha=\kappa .
$$

Let $n_{\lambda} \in \widehat{W}_{a}, \lambda \in T^{+}$, be the element of maximal length in $W_{0} p_{\lambda} W_{0}$.
Theorem 3.3.7. Let $\lambda, \mu \in T^{+}$with $\lambda \geq \mu$. We have

$$
v^{h(\lambda-\mu)} P_{n_{\mu}, n_{\lambda}}\left(v^{-1}\right)=\sum_{x \in W_{0}}(-1)^{l(x)} \mathcal{P}^{*}((\lambda+\tilde{\rho}) x-(\mu+\tilde{\rho}) ; v)
$$

Proof. Since

$$
P_{n_{\mu}, n_{\lambda}}=P_{A^{-} n_{\mu}, A^{-} n_{\lambda}}=P_{A_{\mu}^{+}, A_{\lambda}^{+}}
$$

we calculate the coefficient of $A_{\mu}^{+}$in $F_{A_{\lambda}^{+}}$.
Remember that

$$
F_{A_{\lambda}^{+}}=\sum_{x \in W_{0}}(-1)^{l(x)} v^{c\left(A_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}, A_{\lambda}^{+}\right)} E_{A_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}^{+}}
$$

and

$$
\begin{aligned}
E_{A_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}^{+}}= & \prod_{\alpha>0} \mathbf{f}(\alpha)^{-1} e_{(\lambda+\tilde{\rho}) x-\tilde{\rho}} \\
= & \prod_{\alpha \in \Phi_{1}^{+}}\left(\sum_{k \geq 0} v^{-2 k c_{\alpha}} \theta_{-k \alpha}\right) \prod_{\alpha \in \Phi_{2}^{+}}\left(\sum_{m \geq 0} v^{-m\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \theta_{-m \alpha}\right) \\
& \cdot \prod_{\alpha \in \Phi_{2}^{+}}\left(\sum_{n \geq 0}(-1)^{n} v^{-n\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \theta_{-n \alpha}\right) e_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}
\end{aligned}
$$

We have

$$
\prod_{\alpha \in \Phi_{1}^{+}} \theta_{-k_{\alpha} \alpha} \prod_{\alpha \in \Phi_{2}^{+}} \theta_{-m_{\alpha} \alpha} \prod_{\alpha \in \Phi_{2}^{+}} \theta_{-n_{\alpha} \alpha}\left(A_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}^{+}\right)=A_{\mu}^{+}
$$

for some $K=\left(k_{\alpha}\right)_{\alpha} \in \mathbb{Z}_{+}^{\left|\Phi_{1}^{+}\right|}, M=\left(m_{\alpha}\right)_{\alpha} \in \mathbb{Z}_{+}^{\left|\Phi_{2}^{+}\right|}, N=\left(n_{\alpha}\right)_{\alpha} \in \mathbb{Z}_{+}^{\left|\Phi_{2}^{+}\right|}$if and only if

$$
\sum_{\alpha \in \Phi_{1}^{+}} k_{\alpha} \alpha+\sum_{\alpha \in \Phi_{2}^{+}}\left(m_{\alpha}+n_{\alpha}\right) \alpha=(\lambda+\tilde{\rho}) x-(\mu+\tilde{\rho})
$$

and the contribution to the coefficient of $A_{\mu}^{+}$in $E_{A_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}^{+}}$is

$$
\begin{gathered}
\prod_{\alpha \in \Phi_{2}^{+}}(-1)^{n_{\alpha}} \prod_{\alpha \in \Phi_{1}^{+}} v^{k_{\alpha} h(\alpha)-2 k_{\alpha} c_{\alpha}} \prod_{\alpha \in \Phi_{2}^{+}} v^{m_{\alpha} h(\alpha)-m_{\alpha}\left(c_{\alpha}+\tilde{c}_{\alpha}\right)} \prod_{\alpha \in \Phi_{2}^{+}} v^{n_{\alpha} h(\alpha)-n_{\alpha}\left(c_{\alpha}-\tilde{c}_{\alpha}\right)} \\
=(-1)^{\sigma(N)} v^{h((\lambda+\tilde{\rho}) x-(\mu+\tilde{\rho}))-\sigma^{\prime}(K, M, N)} .
\end{gathered}
$$

Since

$$
c\left(A_{(\lambda+\tilde{\rho}) x-\tilde{\rho}}^{+}, A_{\lambda}^{+}\right)+h((\lambda+\tilde{\rho}) x-(\mu+\tilde{\rho}))=c\left(A_{\mu}^{+}, A_{\lambda}^{+}\right),
$$

the claim follows.

## 4. Some results on generalized cells

In the first section, we briefly introduce generalized cells and the $\boldsymbol{a}$-function. In the second section, we examine a particular generalized two-sided cell of $W_{a}$. Here, Lemma 2.2.7 will play an important role. The third section contains a description of the generalized left cells in a Coxeter group for one class of parameters. We conclude this chapter with an example.
4.1. Definitions and basic properties. For equal parameters, the following definitions and facts can be found in [18]. In [16], Lusztig generalizes the concept of cells to unequal parameters.

Let $W$ be a Coxeter group with generating set $S$ and parameters $c_{s}, s \in S$, subject to the condition that $c_{s}=c_{t}$ whenever $s$ and $t$ are conjugate in $S$. The definition of the corresponding Hecke algebra $\mathcal{H}$ is analogous to the definition in Section 1.2. Here, it will be more convenient to work with the $\mathcal{A}$-basis $\left\{\tilde{T}_{w} \mid w \in W\right\}$ (instead of $\left\{T_{w} \mid w \in W\right\}$ ) where

$$
\tilde{T}_{w}=v^{-m(w)} T_{w}
$$

Multiplication in terms of these elements is given by

$$
\left(\tilde{T}_{s}-v^{c_{s}}\right)\left(\tilde{T}_{s}+v^{-c_{s}}\right)=0
$$

for $s \in S$ and

$$
\tilde{T}_{w} \tilde{T}_{w^{\prime}}=\tilde{T}_{w w^{\prime}}
$$

for $w, w^{\prime} \in W, l(w)+l\left(w^{\prime}\right)=l\left(w w^{\prime}\right)$.
The $\mathcal{A}$-basis of $\mathcal{H}$ consisting of the elements

$$
C_{w}=\sum_{y \leq w}(-1)^{l(w)-l(y)} v^{m(w)-m(y)} P_{y, w}\left(v^{-1}\right) \tilde{T}_{y}
$$

for $w \in W$ will be denoted by $\mathcal{C}$ (see loc. cit.).

For $y, w \in W$ and $s \in S$ such that $s y<y<w<s w$, we inductively define polynomials $M_{y, w}^{s} \in \mathcal{A}$ by
(i) $\sum_{y \leq z<w, s z<z} v^{m(z)-m(y)} P_{y, z} M_{z, w}^{s}-v^{c_{s}+m(w)-m(y)} P_{y, w} \in \mathbb{Z}\left[v^{-1}\right]$ and
(ii) $\frac{y \leq z<w, s z<z}{M_{y, w}^{s}}=M_{y, w}^{s}$.

It turns out that

$$
C_{s} C_{w}= \begin{cases}-\left(v^{c_{s}}+v^{-c_{s}}\right) C_{w} & \text { if } s w<w  \tag{*}\\ C_{s w}-\sum_{z<w, s z<z}(-1)^{l(w)-l(z)} M_{z, w}^{s} C_{z} & \text { if } s w>w\end{cases}
$$

We remark that

$$
C_{w} \in \tilde{T}_{w}+v \sum_{y<w} \mathcal{A}^{+} \tilde{T}_{y}
$$

and

$$
\tilde{T}_{w} \in C_{w}+v \sum_{y<w} \mathcal{A}^{+} C_{y} .
$$

For $y, w \in W, y \leq w$, define polynomials $Q_{y, w} \in \mathcal{A}^{+}$by

$$
\sum_{y \leq z \leq w}(-1)^{l(z)-l(y)} P_{y, z} Q_{z, w}=\delta_{y, w}
$$

and set

$$
D_{y}=\sum_{y \leq w} v^{m(w)-m(y)} Q_{y, w}\left(v^{-1}\right) \tilde{T}_{w},
$$

which is an element in the set $\mathcal{H}^{*}$ of formal $\mathcal{A}$-linear combinations of the elements $\tilde{T}_{w}, w \in W$.

We have an $\mathcal{A}$-linear map $\tau: \mathcal{H}^{*} \rightarrow \mathcal{A}$, given by

$$
\tau\left(\sum_{w \in W} a_{w} \tilde{T}_{w}\right)=a_{e},
$$

$a_{w} \in \mathcal{A}$ for $w \in W$. It is easy to check that

$$
\tau\left(\tilde{T}_{x} \tilde{T}_{y}\right)=\delta_{x, y^{-1}}
$$

and

$$
\tau\left(C_{x} D_{y}\right)=\tau\left(D_{y} C_{x}\right)=\delta_{x, y^{-1}}
$$

for $x, y \in W$.
Let $\leq_{L}$ be the preorder on $W$ which is generated by $x \leq_{L} y$ for $x, y \in W_{a}$, if there exists some $s \in S$ such that $\tau\left(C_{s} C_{y} D_{x^{-1}}\right) \neq 0$. The associated equivalence relation is denoted by $\sim_{L}$, and the equivalence classes with respect to $\sim_{L}$ are called generalized
left cells. Similarly, we define $\leq_{R}, \sim_{R}$, and generalized right cells. We say $x \leq_{L R} y$ for $x, y \in W$, if and only if there exists a sequence

$$
x=x_{0}, x_{1}, \ldots, x_{n}=y
$$

such that for all $0 \leq i \leq n-1$ we have $x_{i} \leq_{L} x_{i+1}$ or $x_{i} \leq_{R} x_{i+1}$. We write $\sim_{L R}$ for the associated equivalence relation, and the equivalence classes are called generalized two-sided cells. (We use the attribute generalized whenever unequal parameters are involved.)

For $y \in W$, define

$$
\mathcal{L}(y)=\{s \in S \mid s y<y\}
$$

and

$$
\mathcal{R}(y)=\{s \in S \mid y s<y\} .
$$

Remark 4.1.1. Let $x, y \in W$.
(i) ([31], Corollary 1.20) If $x \leq_{L} y$, then $\mathcal{R}(x) \supseteq \mathcal{R}(y)$. Therefore, $x \sim_{L} y$ implies $\mathcal{R}(x)=\mathcal{R}(y)$.
(ii) If $C_{x} D_{y} \neq 0$, then $y^{-1} \leq_{L} x$. (Use [31], Corollary 1.15 (a).)

## A function

$$
\boldsymbol{a}: W \rightarrow \mathbb{N}_{0} \cup\{\infty\}
$$

is defined as follows (cf. [18] for equal parameters and [23] for unequal parameters). Let $w \in W$. For $x, y \in W$, express $\tilde{T}_{x} \tilde{T}_{y}$ with respect to the basis $\mathcal{C}$, and consider the coefficient of $C_{w^{-1}}$. If the order of the pole at 0 of these coefficients is bounded as $x$ and $y$ vary, we set $\boldsymbol{a}(w)$ equal to the largest such order. Otherwise, $\boldsymbol{a}(w)=\infty$. For equal parameters, this function is an important tool in the study of representations of Hecke algebras in [18].

Remark 4.1.2. Let $M$ be an abelian group acting on $W$ in a way such that $m(S)=S$ for all $m \in M$. Consider the semidirect product $W^{\prime}=M \ltimes W$. The definitions in this section can be extended to $W^{\prime}$ (comp. [21]). Generalized cells in $W^{\prime}$ are of the form $\{(m, w) \mid m \in M, w \in \Gamma\}$ where $\Gamma$ is a generalized cell in $W$.
4.2. The lowest generalized two-sided cell in $W_{a}$. Suppose $W=W_{a}$. (We use the notation of Chapter 1.)

Let $\nu=l\left(w_{0}\right), \tilde{\nu}=m\left(w_{0}\right)$, and set $\xi_{s}=v^{c_{s}}-v^{-c_{s}}$ for $s \in S$. For $x, y \in W_{a}$, we write

$$
\tilde{T}_{x} \tilde{T}_{y}=\sum_{z \in W_{a}} m_{x, y, z} \tilde{T}_{z}{ }^{-1}
$$

Note that $\tilde{T}_{s}^{2}=\xi_{s} \tilde{T}_{s}+1$ for $s \in S$, so any $m_{x, y, z}$ is a polynomial in $\xi_{s}, s \in S$.

Theorem 4.2.1. Let $x, y, z \in W_{a}$.
(i) As a polynomial in $\xi_{s}, s \in S_{a}$, the degree of $m_{x, y, z}$ is at most $\nu$, and the coefficients are non-negative integers.
(ii) The degree of $m_{x, y, z}$ in $v$ is at most $\tilde{\nu}$.

The proof of part (i) is analogous to the proof of Theorem 7.2 in [18]. Part (ii) is approached in a similar way, using Lemma 2.2.7 instead of [14], Lemma 4.3.

We will need the following two corollaries, whose proofs are again analogous to the ones in [18].

Corollary 4.2.2. (comp. [18], Corollary 7.3) For $w \in W_{a}$, we have $\boldsymbol{a}(w) \leq \tilde{\nu}$.
Corollary 4.2.3. (comp. [18], Corollary 7.10) For any $x, y, z \in W_{a}$, the elements $v^{\bar{\nu}} \tau\left(\tilde{T}_{x} \tilde{T}_{y} \tilde{T}_{z}\right)$ and $v^{\tilde{\nu}} \tau\left(\tilde{T}_{x} \tilde{T}_{y} D_{z}\right)$ are in $\mathcal{A}^{+}$and have the same constant term.

Now, for any $w, x, y \in W_{a}$, the notation $w=x \cdot y$ means that $w=x y$ and $l(w)=l(x)+l(y)$ (and similarly for $w=x \cdot y \cdot z, w, x, y, z \in W_{a}$.) For $\lambda \in T$, let

$$
M_{\lambda}=\left\{z \in W_{a} \mid w_{\lambda} z=w_{\lambda} \cdot z\right\}
$$

and define

$$
N_{\lambda, z}=\left\{w \in W_{a} \mid w=z^{\prime} \cdot w_{\lambda} \cdot z, z^{\prime} \in W_{a}\right\}
$$

where $z \in M_{\lambda}$. Note that according to [28], Lemma 3.2, the condition $w_{\lambda} z=w_{\lambda} \cdot z$ implies $z^{-1} w_{\lambda} z=z^{-1} \cdot w_{\lambda} \cdot z$.

Theorem 4.2.4. Let $\lambda \in T$ and $z \in M_{\lambda}$. The set $N_{\lambda, z}$ is contained in a generalized left cell.

Proof. Let $x \in W_{a}$, and suppose $x=s_{k} \cdots s_{1}, s_{i} \in S_{a}$ for $1 \leq i \leq k$, is a reduced expression. Let $y \in W_{a}$. We denote by $\mathcal{I}_{y}$ the collection of all $I=\left\{i_{1}, \ldots, i_{p_{I}}\right\}$ such that $1 \leq i_{1}<\cdots<i_{p_{I}} \leq k$ and

$$
s_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{1}} \cdots s_{1} y<\hat{s}_{i_{t}} \cdots \hat{s}_{i_{t-1}} \cdots \hat{s}_{i_{1}} \cdots s_{1} y
$$

for all $1 \leq t \leq p_{I}$. For $I \in \mathcal{I}_{y}$, we write

$$
\tilde{T}_{I}=\tilde{T}_{s_{k} \cdots \hat{s}_{i_{p_{I}}} \cdots \hat{s}_{i_{1}} \cdots s_{1} y}
$$

Induction on $k$ shows that

$$
\tilde{T}_{x} \tilde{T}_{y}=\sum_{I \in \mathcal{I}_{y}}\left(\prod_{j=1}^{p_{I}} \xi_{s_{i}}\right) \tilde{T}_{I}
$$

Now take $y \in N_{\lambda, z}, y=z^{\prime} \cdot w_{\lambda} \cdot z$, for some $z^{\prime} \in W_{a}$, and $x=z^{-1} w_{\lambda} z$. (So $x \in N_{\lambda, z}$.) Say $x=s_{k} \cdots s_{1}$ is a reduced expression such that $w_{\lambda}=s_{m} \cdots s_{n}$ and $s_{n-1} \cdots s_{1}=z$
for some $k \geq m \geq n \geq 1$. If we set $I=\{n, \ldots, m\}$, we have $I \in \mathcal{I}_{y^{-1}}$ and $\tilde{p}_{I}=\tilde{\nu}$ (see Lemma 2.2.7 for the definition of $\left.\tilde{p}_{I}\right)$. Hence

$$
v^{\tilde{\nu}} \tau\left(\tilde{T}_{x} \tilde{T}_{y^{-1}} \tilde{T}_{\left(z^{-1} z y^{-1}\right)^{-1}}\right)=v^{\tilde{\nu}} \tau\left(\tilde{T}_{x} \tilde{T}_{y^{-1}} \tilde{T}_{y}\right)
$$

has non-zero constant term. Using

$$
\tau\left(\tilde{T}_{x} \tilde{T}_{y^{-1}} \tilde{T}_{y}\right)=\tau\left(\tilde{T}_{y^{-1}} \tilde{T}_{y} \tilde{T}_{x}\right)=\tau\left(\tilde{T}_{y} \tilde{T}_{x} \tilde{T}_{y^{-1}}\right)
$$

and Corollary 4.2.3, we see that the polynomials $v^{\tilde{\nu}} \tau\left(C_{y^{-1}} C_{y} D_{x}\right)$ and $v^{\tilde{\nu}} \tau\left(C_{y} C_{x} D_{y^{-1}}\right)$ have non-zero constant terms, in particular, $C_{y} D_{x}$ and $C_{x} D_{y^{-1}}$ are non-zero. Thus, by Remark 4.1.1 (ii), we have $x=x^{-1} \leq_{L} y$ and $y \leq_{L} x$, i.e. $y \sim_{L} x$ for all $y \in N_{\lambda, z}$.

The next corollary is an immediate consequence of the above theorem for $z=e$ and Remark 4.1.1 (i). For equal parameters, this is Corollary 8.5 in [18].

For $\lambda \in T$, let $S_{\lambda}=W_{\lambda} \cap S_{a}$ and write $\mathcal{R}\left(S_{\lambda}\right)=\left\{w \in W_{a} \mid w s<w\right.$ for all $\left.s \in S_{\lambda}\right\}$.
Corollary 4.2.5. The set $\mathcal{R}\left(S_{\lambda}\right), \lambda \in T$, is a generalized left cell in $W_{a}$.
We also get the following result, which for Case 2 is Theorem 3.22 in [31]. Let

$$
W_{T}=\left\{w \in W_{a} \mid w=z^{\prime} \cdot w_{\lambda} \cdot z, z, z^{\prime} \in W_{a}, \lambda \in T\right\}
$$

Corollary 4.2.6. The set $W_{T}$ is a generalized two-sided cell.
Proof. We only need to look at Case 1 , in which case

$$
W_{T}=\bigcup_{z \in M_{\lambda}} N_{\lambda, z}
$$

for some fixed $\lambda \in T$. Let $z^{\prime} w_{\lambda} z$ and $y^{\prime} w_{\lambda} y$ be elements in $W_{T}, z^{\prime} w_{\lambda} z=z^{\prime} \cdot w_{\lambda} \cdot z$ and $y^{\prime} w_{\lambda} y=y^{\prime} \cdot w_{\lambda} \cdot y$. Using Theorem 4.2.4, together with its version for generalized right cells, we obtain

$$
z^{\prime} w_{\lambda} z \sim_{L} w_{\lambda} z \sim_{R} w_{\lambda} y \sim_{L} y^{\prime} w_{\lambda} y
$$

Thus $W_{T}$ is contained in a generalized two-sided cell. The other inclusion is proven in the same way as in loc.cit.

We remark that $W_{T}$ is the lowest generalized two-sided cell in $W_{a}$ with respect to $\leq_{L R}$.

Theorem 4.2.7. The set $W_{T}$ contains at most $\left|W_{0}\right|$ generalized left cells.
Proof. For $\lambda \in T$, let

$$
M_{\lambda}^{\prime}=\left\{z \in W_{a} \mid w_{\lambda} z=w_{\lambda} \cdot z, s w_{\lambda} z \notin W_{T} \text { for all } s \in S_{\lambda}\right\}
$$

Following [28], we choose a set of representatives for the $\Omega$-orbits on $T$ and denote it by $R$. Then

$$
W_{T}=\bigcup_{\lambda \in R, z \in M_{\lambda}^{\prime}} N_{\lambda, z}
$$

so the number of generalized left cells in $W_{T}$ is at most the number of pairs $(\lambda, z)$, $\lambda \in R, z \in M_{\lambda}^{\prime}$.

As in loc. cit., we see that $z \in M_{\lambda}^{\prime}$ for some $\lambda \in T$ implies $z^{-1}\left(A_{\lambda}^{+}\right) \subseteq \Pi_{\lambda}$. Since all $z^{-1}\left(A_{\lambda}^{+}\right), \lambda \in R, z \in M_{\lambda}^{\prime}$, are different, the number of pairs $(\lambda, z)$ is thereby bounded by the cardinality of $\left\{A \in X \mid A \subseteq \Pi_{\lambda}\right.$ for some $\left.\lambda \in R\right\}$. The latter set is easily seen to be a fundamental domain for the action of the translation subgroup of $\Omega$ on $X$, so its cardinality is $\left|W_{0}\right|$. The assertion follows.

There is another way to describe the set $W_{T}$, which is similar to the description in [3]. Let $i \in \overline{\mathcal{F}}$, and denote by $\mathcal{F}_{i}$ the set of all hyperplanes $H \in \mathcal{F}$ of direction $i$ such that $c_{H}=c_{i}$. The connected components of

$$
V-\bigcup_{H \in \mathcal{F}_{i}} H
$$

are called strips. We write

$$
\mathcal{U}(A)=\bigcup_{U \text { strip }, U \supseteq A} U
$$

for $A \in X$.
Lemma 4.2.8. We have

$$
W_{T}=\left\{w \in W_{a} \mid w\left(A^{+}\right) \nsubseteq \mathcal{U}\left(A^{+}\right)\right\}
$$

Notice that instead of $A^{+}$, we could have chosen any other alcove $A \in X$ since $\mathcal{U}\left(A^{+} v\right)=\mathcal{U}\left(A^{+}\right) v$ for any $v \in W_{a}$.

Proof. First, let $w \in W_{a}$ be such that $w\left(A^{+}\right) \nsubseteq \mathcal{U}\left(A^{+}\right)$. The alcove $w\left(A^{+}\right)$lies in some connected component $\mathcal{C}$ of $V-\cup_{\alpha \in \Phi}{ }^{+} H_{\alpha, 0}$. This quarter $\mathcal{C}$ with vertex 0 can be described as

$$
\mathcal{C}=\{x \in V \mid\langle x, \check{\alpha}\rangle>0 \text { for } \alpha \in y . \Pi\}
$$

where $y \in W_{0}$ maps $\mathcal{C}^{+}$to $\mathcal{C}$. So there are $r$ linearly independent positive roots $\alpha_{1}, \ldots, \alpha_{r}$ and some $k \in\{1, \ldots, r\}$ such that

$$
\mathcal{C}=\left\{x \in V \mid\left\langle x, \check{\alpha}_{i}\right\rangle<0 \text { for } 1 \leq i \leq k,\left\langle x, \check{\alpha}_{i}\right\rangle>0 \text { for } k+1 \leq i \leq r\right\}
$$

We remove from $\mathcal{C}$ all alcoves which lie in $\mathcal{U}\left(A^{+}\right)$and obtain the quarter

$$
\mathcal{C}^{\prime}=\left\{x \in V \mid\left\langle x, \check{\alpha}_{i}\right\rangle<0 \text { for } 1 \leq i \leq k,\left\langle x, \check{\alpha}_{i}\right\rangle>b_{i} \text { for } k+1 \leq i \leq r\right\}
$$

where

$$
b_{i}= \begin{cases}1 & \text { if } c_{\alpha_{i}}=\tilde{c}_{\alpha_{i}} \\ 2 & \text { otherwise }\end{cases}
$$

for $k+1 \leq i \leq r$. (The set $\mathcal{C}^{\prime}$ is a translate of $\mathcal{C}$.) So $w\left(A^{+}\right) \subset \mathcal{C}^{\prime}$.
Let $\lambda$ be the vertex of $\mathcal{C}^{\prime}$, and let $z \in W_{a}$ be such that $z\left(A^{+}\right)$is the unique alcove in $\mathcal{C}^{\prime}$ containing $\lambda$ in its closure. Since $\mathcal{C}^{\prime} \subset \mathcal{C}$, any other element $v \in W_{a}$ with $v\left(A^{+}\right) \subset \mathcal{C}^{\prime}$ satisfies $v=v^{\prime} \cdot z$ for some $v^{\prime} \in W_{a}$, in particular $w=w^{\prime} \cdot z$ for some $w^{\prime} \in W_{a}$.

Let $s \in S_{\lambda}$, so there is a wall of $\mathcal{C}$ containing a face of type $s$. For each wall of $\mathcal{C}^{\prime}$, the alcoves $A^{+}$and $z\left(A^{+}\right)$lie on different sides. Since $l(v)$ for $v \in W_{a}$ counts the number of hyperplanes such that $A^{+}$and $v\left(A^{+}\right)$lie on different sides, we obtain $s z<z$. We conclude $z=w_{\lambda} \cdot z^{\prime}$ for some $z^{\prime} \in W_{a}$ and $w=w^{\prime} \cdot w_{\lambda} \cdot z^{\prime}$.

Conversely, let $w \in W_{T}, w=z^{\prime} \cdot w_{\lambda} \cdot z$ for some $z^{\prime}, z \in W_{a}$ and $\lambda \in T$. If we are in Case 1, we set $A=z^{-1}\left(A^{+}\right)$. Using [28], Lemma 4.2, we get

$$
A=z^{-1}\left(A^{+}\right) \subset \mathcal{C}^{+}
$$

and

$$
w(A)=z^{\prime} w_{\lambda}\left(A^{+}\right)=z^{\prime}\left(A^{-}\right) \subset \mathcal{C}^{-}
$$

Hence $\langle x, \check{\alpha}\rangle>0$ and $\langle y, \check{\alpha}\rangle<0$ for all $x \in A, y \in w(A), \alpha \in \Phi^{+}$, which implies $w(A) \nsubseteq \mathcal{U}(A)$.

If we are in Case 2, we set $A=z^{-1}\left(A_{\lambda}^{+}\right)$. We obtain $\langle x, \check{\alpha}\rangle \neq\left\langle x^{\prime}, \check{\alpha}\right\rangle$ for all $x \in A$, $x^{\prime} \in w(A), \alpha \in \Phi^{+}$, and this again implies $w(A) \nsubseteq \mathcal{U}(A)$.

Note that the connected components of $V-\mathcal{U}\left(A^{+}\right)$turn out to be precisely the quarters of the form $\mathcal{C}^{\prime}$.

Remark 4.2.9. Define

$$
W_{(\tilde{\nu})}=\left\{w \in W_{a} \mid \boldsymbol{a}(w)=\tilde{\nu}\right\}
$$

For $w \in W_{T}$, we obtain from the proof of Theorem 4.2.4 that $\boldsymbol{a}(w) \geq \tilde{\nu}$ and from Corollary 4.2.2 that $\boldsymbol{a}(w) \leq \tilde{\nu}$. Thus $W_{T} \subseteq W_{(\tilde{\nu})}$. At the end of the next section, we will see that this is actually an equality.

We also point out that we can use the procedure in [19] to attach a based ring $J_{W_{T}}$ to the lowest generalized two-sided cell $W_{T}$. For fixed $T$, the ring $J_{W_{T}}$ does not depend on the parameters.
4.3. Parameters coming from graph automorphisms. We now place ourselves in the setting of [15], Section 8 . Let $\widetilde{W}$ be a finite or affine Weyl group with a set $\widetilde{S}$ of simple reflections and $\alpha: \widetilde{W} \rightarrow \widetilde{W}$ a non-trivial automorphism such that $\alpha(\widetilde{S})=\widetilde{S}$. Let $W$ be the fixed point set of $\widetilde{W}$ under $\alpha$. This is again a finite or affine Weyl group and has a set $S$ of simple reflections corresponding to the orbits of $\alpha$ on $\widetilde{S}$. For $s \in S$, let $c_{s}$ be the length of the longest element in the subgroup of $\widetilde{W}$ generated by the orbit corresponding to $s$. We thereby get parameters $c_{s} \in \mathbb{N}$ such that $c_{s}=c_{t}$ whenever $s$ and $t$ are conjugate in $W$.

We will write $P_{y, w}, C_{w}, \boldsymbol{a}$ etc. for $W$ and $\widetilde{P}_{y, w}, \widetilde{C}_{w}, \widetilde{\boldsymbol{a}}$ etc. for $\widetilde{W}$.
For $x, y, z \in W$, set

$$
h_{x, y, z}=\tau\left(C_{x} C_{y} D_{z^{-1}}\right),
$$

and correspondingly

$$
\widetilde{h}_{x, y, z}=\widetilde{\tau}\left(\widetilde{C}_{x} \widetilde{C}_{y} \widetilde{D}_{z^{-1}}\right)
$$

for $x, y, z \in \widetilde{W}$.
Let $x, y, z \in W$ be such that $x \leq y$. According to [11], [16], and [24], the coefficients of $\widetilde{P}_{x, y}$ and $\widetilde{h}_{x, y, z}$ can (up to a sign) be interpreted as dimensions of certain vector spaces on which $\alpha$ acts, and the corresponding coefficients of $P_{x, y}$ and $h_{x, y, z}$ are the traces of $\alpha$ on these vector spaces. We need the following facts which are consequences of these interpretations.
4.3.1. We have $\operatorname{deg} P_{x, y} \leq \operatorname{deg} \widetilde{P}_{x, y}$.
4.3.2. If the coefficient of $v^{i}, i \in \mathbb{Z}$, in $h_{x, y, z}$ is non-zero, then the coefficient of $v^{i}$ in $\widetilde{h}_{x, y, z}$ is non-zero as well. In particular, $\boldsymbol{a}(z) \leq \widetilde{\boldsymbol{a}}(z)$.
4.3.3. If the coefficient of $v^{i}, i \in \mathbb{Z}$, in $\widetilde{h}_{x, y, z}$ is $\pm 1$, then the coefficient of $v^{i}$ in $h_{x, y, z}$ is non-zero.

For $x, y, z \in W$, let $c_{x, y, z}$, be the integer given by

$$
v^{a(z)} h_{x, y, z^{-1}}-c_{x, y, z} \in v \mathcal{A}^{+},
$$

and write $\delta(z)$ for the degree of the polynomial $P_{e, z}$ in $v$. Similarly, we define $\widetilde{c}_{x, y, z}$ and $\widetilde{\delta}(z)$ for $x, y, z \in \widetilde{W}$, and we set $\widetilde{\gamma}_{x, y, z}=(-1)^{\widetilde{a}(z)} \widetilde{c}_{x, y, z}$.

Lemma 4.3.4. We have $\boldsymbol{a}(z) \leq m(z)-\delta(z)$ for any $z \in W$.
Proof. According to [19], the corresponding inequality holds for $\widetilde{W}$. In view of 4.3.2 and 4.3.1, we can therefore conclude

$$
\begin{equation*}
\boldsymbol{a}(z) \leq \widetilde{\boldsymbol{a}}(z) \leq \widetilde{l}(z)-\widetilde{\delta}(z) \leq m(z)-\delta(z) \tag{*}
\end{equation*}
$$

Remark 4.3.5. Let

$$
\widetilde{\mathcal{D}}=\{z \in \widetilde{W} \mid \widetilde{\boldsymbol{a}}(z)=\widetilde{l}(z)-\widetilde{\delta}(z)\}
$$

be the set of distinguished involutions in $\widetilde{W}$. We set

$$
\mathcal{D}=\{z \in W \mid \boldsymbol{a}(z)=m(z)-\delta(z)\} .
$$

From (*), we derive that $\boldsymbol{a}(d)=\widetilde{\boldsymbol{a}}(d)$ for all $d \in \mathcal{D}$ and that $\mathcal{D} \subset \widetilde{\mathcal{D}}$.
The following two results can be found in [19].
4.3.6. Any left cell $\widetilde{\Gamma}$ of $\widetilde{W}$ contains a unique $d \in \widetilde{\mathcal{D}}$. For $z \in \widetilde{\Gamma}$, we have $\widetilde{\gamma}_{z^{-1}, z, d}=1$, and $d$ is the unique element in $\widetilde{W}$ such that $\widetilde{\gamma}_{z^{-1}, z, d} \neq 0$.
4.3.7. We have $\widetilde{\gamma}_{x, y, z}=\widetilde{\gamma}_{y, z, x}$ for all $x, y, z \in \widetilde{W}$.

Theorem 4.3.8. The generalized left cells of $W$ are the fixed point sets of the left cells of $\widetilde{W}$ under $\alpha$.

Proof. Let $z \in W$. It follows from the definition of $C_{w}, w \in W$, that

$$
\tau\left(C_{z^{-1}} C_{z}\right) \in 1+v \mathcal{A}^{+}
$$

On the other hand,

$$
\begin{aligned}
\tau\left(C_{z^{-1}} C_{z}\right) & =\sum_{w \in W} h_{z^{-1}, z, w} \tau\left(C_{w}\right) \\
& =\sum_{w \in W} h_{z^{-1}, z, w}(-1)^{l(w)} v^{m(w)} P_{e, w}\left(v^{-1}\right)
\end{aligned}
$$

Lemma 4.3.4 shows that there exists some $d \in \mathcal{D}$ such that $c_{z^{-1}, z, d} \neq 0$. Since $\boldsymbol{a}(d)=\widetilde{\boldsymbol{a}}(d)$, we conclude $\widetilde{\boldsymbol{c}}_{z^{-1}, z, d} \neq 0$ and hence $\widetilde{\gamma}_{z^{-1}, z, d} \neq 0$. We have $d \in \widetilde{\mathcal{D}}$, and using 4.3.6 and 4.3.7, we arrive at

$$
\widetilde{\gamma}_{z^{-1}, z, d}=\widetilde{\gamma}_{z, d, z^{-1}}=1
$$

Thus, according to 4.3 .3 , both elements $C_{z} D_{d}$ and $C_{d} D_{z^{-1}}$ are non-zero, which in turn implies $d \leq_{L} z \leq_{L} d$ (cf. Remark 4.1.1 (ii)), i.e. $d \sim_{L} z$. This means that any generalized left cell in $W$ contains a distinguished involution.

As a consequence of 4.3.2, any generalized left cell in $W$ is contained in the fixed point set of some left cell in $\widetilde{W}$ (comp. [15]). Since there is only one distinguished involution in each left cell in $\widetilde{W}$, the assertion follows.

In [19], Lusztig proves that if $\widetilde{W}$ is an affine Weyl group, $\widetilde{W}$ consists of only finitely many left cells. So we obtain the following statement.

Corollary 4.3.9. If $W$ is an affine Weyl group (with parameters as above) the number of generalized left cells in $W$ is finite.

Note that the proof of Theorem 4.3 .8 shows that $\mathcal{D}$ is the fixed point set of $\widetilde{\mathcal{D}}$ and that $\boldsymbol{a}(z)=\widetilde{\boldsymbol{a}}(z)$ for any $z \in W$. Several statements in [19] then carry over to $W$. We will need the following one for the proof of Theorem 4.3.13.

Corollary 4.3.10. (comp. [19], Corollary 1.9 (b)) For any $z^{\prime}, z \in W$, if $z^{\prime} \leq_{L} z$ and $\boldsymbol{a}\left(z^{\prime}\right)=\boldsymbol{a}(z)$ then $z^{\prime} \sim_{L} z$.

Remark 4.3.11. As we will see in the next section, a generalized two-sided cell in $W$ does not have to be the fixed point set of some two-sided cell in $\widetilde{W}$.

We now again use the notations of Chapter 1 (taking $W=W_{a}$ ).
Theorem 4.3.12. If we have parameters coming from a graph automorphism, $W_{T}$ consists of exactly $\left|W_{0}\right|$ generalized left cells.

Proof. Let

$$
N=\left\{z^{-1} w_{\lambda} z \mid \lambda \in R, z \in M_{\lambda}^{\prime}\right\} .
$$

As in [28], we see that $N \subseteq \mathcal{D}$. In view of Theorem 4.2.7, it therefore suffices to show that $|N|=\left|W_{0}\right|$.

It follows from the proof of Lemma 4.2 .8 that $V-\mathcal{U}\left(A^{+}\right)$has $\left|W_{0}\right|$ connected components. According to the same lemma, we have

$$
V-\mathcal{U}\left(A^{+}\right)=W_{T} \cdot A^{+}
$$

which in turn equals

$$
\bigcup_{\lambda \in R, z \in M_{\lambda}^{\prime}} N_{\lambda, z} \cdot A^{+} .
$$

In the proof of Theorem 4.2.7, we saw that $\left|\left\{(\lambda, z) \mid \lambda \in R, z \in M_{\lambda}^{\prime}\right\}\right| \leq\left|W_{0}\right|$, and since $N_{\lambda, z} \cdot A^{+}$is connected, we obtain $\left|\left\{(\lambda, z) \mid \lambda \in R, z \in M_{\lambda}^{\prime}\right\}\right|=\left|W_{0}\right|$. We have $z^{-1} w_{\lambda} z \cdot A^{+} \subseteq N_{\lambda, z} \cdot A^{+}$for $\lambda \in R, z \in M_{\lambda}^{\prime}$, so all $z^{-1} w_{\lambda} z . A^{+}$and thus all $z^{-1} w_{\lambda} z$ are different. We conclude $|N|=\left|W_{0}\right|$.

Finally, we prove the equality indicated in Remark 4.2.9. (This holds for any parameters.)

Theorem 4.3.13. We have

$$
W_{(\tilde{\nu})}=W_{T}
$$

Proof. It remains to show that $W_{(\tilde{\nu})} \subseteq W_{T}$.
We first assume that we have parameters coming from a graph automorphism. Let $z \in W_{a}$ be such that $\boldsymbol{a}(z)=\tilde{\nu}$. Choose some $\lambda \in T$, and let $y \in W_{a}$ be the element of minimal length in $W_{\lambda} z$. Then $z=x \cdot y$ for some $x \in W_{\lambda}$, and by induction on
$l\left(w_{\lambda}\right)-l(x)$ we see that we can find $s_{1}, \ldots, s_{n} \in S_{\lambda}$ such that $s_{n} \cdots s_{1} z \in W_{T}$ and $s_{i} \cdots s_{1} z>s_{i-1} \cdots s_{1} z$ for all $1 \leq i \leq n$, i.e. there exists an element $z^{\prime} \in W_{T}$ such that $z^{\prime} \leq_{L} z$. By Corollary 4.3.10, we get $z^{\prime} \sim_{L} z$, thus $z \in W_{T}$.

For arbitrary parameters $c_{s}, s \in S_{a}$, one can check that there always are parameters $c_{s}^{\prime}, s \in S_{a}$, such that the corresponding set $W_{\left(\tilde{\nu}^{\prime}\right)}$ resp. $W_{T^{\prime}}$ equals $W_{(\tilde{\nu})}$ resp. $W_{T}$. (Use Corollary 4.2.3 for $W_{\left(\tilde{\nu}^{\prime}\right)}$.) The claim follows.
4.4. The case $\tilde{C}_{2}$. In this section, we explicitly determine the generalized left and two-sided cells for the case $W$ of type $\tilde{C}_{2}\left(=\tilde{B}_{2}\right)$ and $\widetilde{W}$ of type $\tilde{A}_{3}$. More cases can be found in the appendix.

Let $\widetilde{W}$ be a Coxeter group of type $\tilde{A}_{3}$ with Dynkin diagram


Let $\alpha$ be the automorphism on $\widetilde{W}$ which fixes $\tilde{s}_{0}$ and $\tilde{s}_{2}$ and interchanges $\tilde{s}_{1}$ and $\tilde{s}_{3}$. Then $W=\widetilde{W}^{\alpha}$ is of type $\tilde{C}_{2}$ with parameters


The cells for Weyl groups $\mathcal{A}_{r}, r \geq 2$, of type $\tilde{A}_{r-1}$ are given in [27]. In particular, Lusztig and Shi established a bijective correspondence between the two-sided cells of $\mathcal{A}_{r}$ and the partitions of $r$, which we describe now.

We realize $\mathcal{A}_{r}$ as the group of all permutations $\sigma$ on $\mathbb{Z}$ such that $\sigma(i+r)=\sigma(i)+r$ for all $i \in \mathbb{Z}$ and $\sum_{i=1}^{r} \sigma(i)=\sum_{i=1}^{r} i$ by letting the simple reflection $s_{i}, 0 \leq i \leq r-1$ act as

$$
s_{i}(j)= \begin{cases}j+1 & \text { if } j \equiv i(\bmod r) \\ j-1 & \text { if } j \equiv i+1(\bmod r) \\ j & \text { if } j \not \equiv i, i+1(\bmod r)\end{cases}
$$

A map $\pi$ from $\mathcal{A}_{r}$ to the set of partitions of $r$ is defined as follows. Let $w \in \mathcal{A}_{r}$. We denote by $d_{k}, k \in \mathbb{N}$, the maximal cardinality of a subset of $\mathbb{Z}$ whose elements are incongruent to each other mod $r$ and which is a disjoint union of $k$ subsets each of which has its natural order reversed by $w$. (We set $d_{0}=0$.) Let $n \in \mathbb{N}$ be minimal such that $d_{n}=r$. For $1 \leq k \leq n$, we define $\lambda_{k}=d_{k}-d_{k-1}$. Then $\sum_{k=1}^{n} \lambda_{k}=r$ and $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{n}$ (cf. [27]). We set $\pi(w)=\left(\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{n}\right)$. The two-sided cells in $\mathcal{A}_{r}$ then coincide with the fibers of $\pi$.

Now take $\widetilde{W}$ and $W$ as at the beginning of this section. We use the above description of the two-sided cells in order to determine the generalized left cells in $W$. The fixed point sets of the two-sided cells in $\widetilde{W}$ are unions of generalized left and two-sided cells in $W$.

The lowest generalized two-sided cell $\underline{c}_{0}$ consists of the eight generalized left cells

$$
N_{0,1}, N_{0, s_{0}}, N_{0, s_{0} s_{1}}, N_{0, s_{0} s_{1} s_{2}}, N_{\frac{1}{2} \alpha_{0}, 1}, N_{\frac{1}{2} \alpha_{0}, s_{2}}, N_{\frac{1}{2} \alpha_{0}, s_{2} s_{1}}, N_{\frac{1}{2} \alpha_{0}, s_{2} s_{1} s_{0}} .
$$

The two-sided cell corresponding to the partition $(1 \geq 1 \geq 1 \geq 1)$ contains only the identity element. The fixed point set $\underline{c}_{1}=\{1\}$ is a generalized left and two-sided cell.

As noted in [12], the two-sided cell corresponding to $(2 \geq 1 \geq 1)$ consists of all elements in $\widetilde{W}$ with a unique reduced expression. Its set of fixed points is therefore $\left\{s_{0}, s_{2}\right\}$. Remark 4.1.1 (i) implies that $\underline{c}_{2}=\left\{s_{0}\right\}$ and $\underline{c}_{3}=\left\{s_{2}\right\}$ are generalized left and right cells and hence generalized two-sided cells. (This is an example where the generalized two-sided cells in $W$ are not the fixed point sets of the two-sided cells in $\widetilde{W}$.)

We are left with the elements in $\underline{c}_{4}=\underline{c}_{4}^{1} \cup \underline{c}_{4}^{2} \cup \underline{c}_{4}^{3} \cup \underline{c}_{4}^{4}$ where

$$
\begin{aligned}
& \underline{c}_{4}^{1}=\left\{s_{1}\left(s_{0} s_{2} s_{1}\right)^{n}, s_{2} s_{1}\left(s_{0} s_{2} s_{1}\right)^{n}, s_{0} s_{1}\left(s_{0} s_{2} s_{1}\right)^{n},\left(s_{0} s_{2} s_{1}\right)^{n+1} ; n \geq 0\right\} \\
& \underline{c}_{4}^{2}=\left\{s_{2} s_{0}\left(s_{1} s_{2} s_{0}\right)^{n},\left(s_{1} s_{2} s_{0}\right)^{n+1}, s_{2}\left(s_{1} s_{2} s_{0}\right)^{n+1}, s_{0}\left(s_{1} s_{2} s_{0}\right)^{n+1} ; n \geq 0\right\} \\
& \underline{c}_{4}^{3}=\left\{s_{1} s_{2}\left(s_{0} s_{1} s_{2}\right)^{n}, s_{2} s_{1} s_{2}\left(s_{0} s_{1} s_{2}\right)^{n},\left(s_{0} s_{1} s_{2}\right)^{n+1}, s_{2}\left(s_{0} s_{1} s_{2}\right)^{n+1} ; n \geq 0\right\} \\
& \underline{c}_{4}^{4}=\left\{s_{1} s_{0}\left(s_{2} s_{1} s_{0}\right)^{n}, s_{0} s_{1} s_{0}\left(s_{2} s_{1} s_{0}\right)^{n},\left(s_{2} s_{1} s_{0}\right)^{n+1}, s_{0}\left(s_{2} s_{1} s_{0}\right)^{n+1} ; n \geq 0\right\}
\end{aligned}
$$

and $\underline{c}_{5}=\underline{c}_{5}^{1} \cup \underline{c}_{5}^{2} \cup \underline{c}_{5}^{3} \cup \underline{c}_{5}^{4}$ where

$$
\begin{aligned}
& \underline{c}_{5}^{1}=\left\{s_{1} s_{2} s_{1}\left(s_{0} s_{1} s_{2} s_{1}\right)^{n},\left(s_{0} s_{1} s_{2} s_{1}\right)^{n+1}, s_{1}\left(s_{0} s_{1} s_{2} s_{1}\right)^{n+1}, s_{2} s_{1}\left(s_{0} s_{1} s_{2} s_{1}\right)^{n+1} ; n \geq 0\right\}, \\
& \underline{c}_{5}^{2}=\left\{s_{1} s_{0} s_{1}\left(s_{2} s_{1} s_{0} s_{1}\right)^{n},\left(s_{2} s_{1} s_{0} s_{1}\right)^{n+1}, s_{1}\left(s_{2} s_{1} s_{0} s_{1}\right)^{n+1}, s_{0} s_{1}\left(s_{2} s_{1} s_{0} s_{1}\right)^{n+1} ; n \geq 0\right\} \\
& \underline{c}_{5}^{3}=\left\{\left(s_{1} s_{0} s_{1} s_{2}\right)^{n}, s_{2}\left(s_{1} s_{0} s_{1} s_{2}\right)^{n}, s_{1} s_{2}\left(s_{1} s_{0} s_{1} s_{2}\right)^{n}, s_{0} s_{1} s_{2}\left(s_{1} s_{0} s_{1} s_{2}\right)^{n} ; n \geq 1\right\} . \\
& \underline{c}_{5}^{4}=\left\{\left(s_{1} s_{2} s_{1} s_{0}\right)^{n}, s_{0}\left(s_{1} s_{2} s_{1} s_{0}\right)^{n}, s_{1} s_{0}\left(s_{1} s_{2} s_{1} s_{0}\right)^{n}, s_{2} s_{1} s_{0}\left(s_{1} s_{2} s_{1} s_{0}\right)^{n} ; n \geq 1\right\} .
\end{aligned}
$$

It can be shown by induction on $n$ that $\pi$ sends the elements in $\underline{c}_{4}$ to $(2 \geq 2)$ and the elements in $\underline{c}_{5}$ to $(3 \geq 1)$. (For example, the element $\left(s_{0} s_{2} s_{1}\right)^{n}, n \in \mathbb{N}$, as permutation on $\mathbb{Z}$ maps

$$
\begin{aligned}
& 1 \text { to } 1+2 n, \\
& 2 \text { to } 2-2 n, \\
& 3 \text { to } 3+2 n, \\
& 4 \text { to } 4-2 n
\end{aligned}
$$

So sets of integers which are incongruent mod 4 and whose order is reversed by $\left(s_{0} s_{2} s_{1}\right)^{n}$ are of the form

$$
\{1+4 k, 2+4 l\},\{1+4 k, 4 l\},\{2+4 k, 3+4 l\}, \text { or }\{3+4 k, 4 l\}
$$

for certain $k, l \in \mathbb{Z}$. Thus $\pi\left(\left(s_{0} s_{2} s_{1}\right)^{n}\right)=(2 \geq 2)$.)
We conclude that $\tilde{\boldsymbol{a}}$ is constant on $\underline{c}_{4}$ and $\underline{c}_{5}$, respectively, and that $\underline{c}_{4}$ and $\underline{c}_{5}$ are unions of generalized left and two-sided cells. (It is shown in [18] that $\widetilde{\boldsymbol{a}}$ is constant on two-sided cells.)

It follows from the multiplication formula $(*)$ in Section 4.1 that

$$
s_{1}\left(s_{0} s_{2} s_{1}\right)^{n} \geq_{L}\left\{\begin{array}{l}
s_{2} s_{1}\left(s_{0} s_{2} s_{1}\right)^{n} \\
s_{0} s_{1}\left(s_{0} s_{2} s_{1}\right)^{n}
\end{array}\right\} \geq_{L}\left(s_{0} s_{2} s_{1}\right)^{n+1} \geq_{L} s_{1}\left(s_{0} s_{2} s_{1}\right)^{n+1}
$$

for $n \geq 0$. Similar relationships hold for the elements in $\underline{c}_{4}^{i}, i \geq 2$, and $\underline{c}_{5}^{j}, j \geq 1$. We derive from Corollary 4.3.10 that elements $x, y \in W$ with $x \leq_{L} y$ and $\pi(x)=\pi(y)$ lie in the same generalized left cell. Therefore, any $\underline{c}_{4}^{i}$ and any $\underline{c}_{5}^{j}$ is contained in a generalized left cell. Remark 4.1.1 (i) then implies that $\underline{c}_{4}^{i}$ for $i \geq 1$ and $\underline{c}_{5}^{j}$ for $j \geq 3$ are generalized left cells.

According to [27], Theorem 16.1.2, elements in $\mathcal{A}_{r}, r \geq 2$, lie in the same left cells if and only if they have the same generalized right $\tau$-invariant (for the definition see [27], p.18). Since the elements in $\mathcal{A}_{4}$ corresponding to $s_{1} s_{2} s_{1} \in \underline{c}_{5}^{1}$ and $s_{1} s_{0} s_{1} \in \underline{c}_{5}^{2}$ do not have the same generalized right $\tau$-invariant (take the right star operation with respect to $\left\{\tilde{s}_{2}, \tilde{s}_{3}\right\}$ ), they do not lie in the same left cell. It follows that $\underline{c}_{5}^{1}$ and $\underline{c}_{5}^{2}$ are generalized left cells.

Since generalized two-sided cells are unions of generalized left and right cells and since $\left(\underline{c}_{n}^{1}\right)^{-1} \cap \underline{c}_{n}^{i} \neq \emptyset$ for $n \in\{4,5\}$ and all $i \geq 1$, the sets $\underline{c}_{4}$ and $\underline{c}_{5}$ are generalized two-sided cells.

We remark that in this case, as well as in the cases listed in the appendix, each generalized two-sided cell has a non-trivial intersection with a maximal parabolic subgroup $W_{\lambda}$ of $W$ for some special point $\lambda$.

## Appendix

For $W$ of type $\tilde{C}_{2}$, we describe generalized left and two-sided cells using the realization of $W$ in terms of alcoves as in Section 2.3. The generalized left cells are obtained as the connected components after removing the thick lines. The generalized two-sided cells are the unions of all generalized left cells with the same label $\underline{c}_{i}$, $i \geq 0$. In all cases, the elements in the lowest generalized two-sided cell are labeled $\underline{c}_{0}$. The alcove corresponding to the identity element is labeled $\underline{c}_{1}$.

We start with the equal parameter case, which is taken from [18]. In the unequal parameter cases, we indicate $\widetilde{W}$ as well as the non-trivial actions of $\alpha$. (The notation is that of Section 4.3.) The calculations are similar to those in Section 4.4. (For $\widetilde{W}$ of type $\tilde{B}_{3}$ resp. $\tilde{D}_{4}$, we use [6] resp. [5], [30].)
$W$ :

$$
\stackrel{1}{1}
$$

The 4 two-sided cells consist of $8+1+3+4=16$ left cells (in this order).

$\widetilde{W}$ :

$W$ :


The 6 generalized two-sided cells consist of $8+1+1+1+4+4=19$ generalized left cells.


$W$ :


The 6 generalized two-sided cells consist of $8+1+2+4+4+4=23$ generalized left cells.

$\widetilde{W}:$

$\widetilde{W}:$

$W$ :

$W$ :


The 7 generalized two-sided cells consist of $8+1+1+1+1+4+4=20$ generalized left cells.

$\widetilde{W}$ :

$$
\begin{aligned}
& W: \\
& \\
& \\
& \\
& \\
& \hline
\end{aligned}
$$

The 7 generalized two-sided cells consist of $8+1+1+2+4+4+4=24$ generalized left cells.


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