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# Structure of atypical representations of the Lie superalgebras $\mathrm{sl}(m / n)$ 

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#### Abstract

In contrast to Lie algebras, Lie superalgebras contain typical and atypical irreducible representations. Character formulae for typical representations of classical Lie superalgebras have been known for a long time, but the structure of atypical representations remained unsolved. In this paper we construct the character formula for atypical representations of the special linear Lie superalgebras $\operatorname{sl}(m / n)$.


## 1. Introduction

Since the work of Corwin et al (1975), Lie superalgebras have become increasingly important in theoretical physics. We only mention their evidence in supersymmetry (Fayet and Ferrara 1977), supergravity (van Nieuwenhuizen 1981) and nuclear physics (Iachello 1980) here. Simple Lie superalgebras were classified completely (Kac 1977, Scheunert 1979) and it was shown that classical Lie superalgebras can be described by a Cartan matrix or, equivalently, by a Kac-Dynkin diagram.

Finite-dimensional irreducible representations (irreps) of Lie superalgebras were studied by Kac (1978) and categorised as being either typical or atypical. Typical representations have properties analogous to those of finite-dimensional irreps of Lie algebras and are relatively easy to handle. Kac (1978) constructed a character formula for typical irreps of classical Lie superalgebras. Atypical representations are much harder to deal with. Various techniques have been used in order to get a deeper insight into the structure of atypical representations, but still no completely general character formulae exist. Among the techniques used we mention methods based on superfields (Farmer and Jarvis 1983, 1984), tensors and supertableaux (Balantekin and Bars 1981, 1982, King 1983), shift operators (Hughes 1981, Van der Jeugt 1984, 1985a), weight spaces (Hurni and Morel 1982, 1983) and generating functions (Sharp et al 1985). Atypical representations are important for several reasons. Firstly, adjoint representations of Lie superalgebras are usually atypical. Secondly, it was shown recently (Van der Jeugt 1985b) that the phenomenon of multiplet shortening in supergravity models (Freedman and Nicolai 1984) is explained by the structure of atypical representations of the underlying superalgebra.

In the present paper, we study finite-dimensional atypical representations of the Lie superalgebras $\operatorname{sl}(m / n)$, sometimes denoted by $\mathrm{SU}(m / n), \operatorname{spl}(m, n)$ or $\mathrm{A}(m-1$, $n-1$ ). The main result is the construction of a character formula for atypical representations, given by (4.31) when precisely one of the atypicality conditions is satisfied.

[^0]The outline of the paper is as follows: in § 2, the algebra is defined and the Chevalley basis is constructed. It is also shown how the atypicality conditions are derived. In $\S 3$, a very useful theorem is proven, which shows the existence of a special highest weight vector in the Verma module $\tilde{V}(\Lambda)$, if $\Lambda$ corresponds to the highest weight of an atypical representation. This property is used in § 4 in order to obtain the character formula for atypical representations of $\operatorname{sl}(m / n)$. Section 5 gives an example for $\mathrm{L}=\operatorname{sl}(3 / 2)$.

## 2. Conventions and notation for $\operatorname{sl}(m / n)$

In matrix notation, the standard form of the Lie superalgebra $\operatorname{sl}(m / n)$ is defined as follows (Kac 1977):

$$
\operatorname{si}(m / n)=\left\{\left.x=\left(\begin{array}{ll}
a & b  \tag{2.1}\\
c & d
\end{array}\right) \right\rvert\, \operatorname{str} x=\operatorname{Tr} a-\operatorname{Tr} d=0\right\}
$$

where $a$ is a complex ( $m \times m$ ) matrix, $b$ a ( $m \times n$ ) matrix, $c$ a $(n \times m)$ matrix and $d$ a ( $n \times n$ ) matrix. The Lie superalgebra $\operatorname{sl}(m / n)$ is simple if $m \neq n$. If $m=n$, it contains the non-trivial ideal $\mathbb{C} \cdot I_{2 m}$, and then $\operatorname{sl}(n / n) / \mathbb{C} \cdot I_{2 m}$ is simple. In this paper, we shall prove the results only for $m \neq n$, but they can easily be deduced for $m=n$. The even subspace of $\operatorname{sl}(m / n)$ is the space for which $b$ and $c$ are zero and the odd subspace consists of matrices with $a=0$ and $d=0$.

Let $e_{i j}$ be the matrix with entry 1 in the $i$ th row and $j$ th column, and 0 elsewhere. We set

$$
\begin{align*}
& I_{m}=\sum_{i=1}^{m} e_{i i}  \tag{2.2}\\
& I_{n}=\sum_{i=m+1}^{m+n} e_{i i} . \tag{2.3}
\end{align*}
$$

The Cartan subalgebra H of $\mathrm{sl}(m / n)$ is spanned by

$$
\begin{align*}
& h_{i}=e_{i i}+(1 / n) I_{n} \quad(i=1, \ldots, m) \\
& h_{m+i}=(1 / m) I_{m}+e_{m+i, m+i} \quad(i=1, \ldots, n)  \tag{2.4}\\
& h_{0}=(1 / m) I_{m}+(1 / n) I_{n} .
\end{align*}
$$

Only ( $m+n-1$ ) of the elements (2.4) are independent since

$$
\begin{align*}
& \sum_{i=1}^{m} h_{i}-m h_{0}=0 \\
& \sum_{i=m+1}^{m+n} h_{i}-n h_{0}=0 \tag{2.5}
\end{align*}
$$

H is contained in the set of diagonal $(m \times n) \times(m+n)$ matrices $\mathrm{D}=$ $\left\{d=\operatorname{diag}\left(d_{11}, \ldots, d_{m+n, m+n}\right)\right\}$ and the linear forms $\varepsilon_{i}(i=1, \ldots, m)$ and $\delta_{j}(j=$ $1, \ldots, n$ ) are defined as follows:

$$
\begin{align*}
& \varepsilon_{i}(d)=d_{i i}  \tag{2.6}\\
& \delta_{j}(d)=d_{m+j, m+j} .
\end{align*}
$$

The same notation is used for the restrictions of $\varepsilon_{i}$ and $\delta_{j}$ to H . Then, the following relations can be checked for $i=0,1, \ldots, m+n$ :

$$
\left[h_{i}, e_{j k}\right]=\left(\varepsilon_{j}-\varepsilon_{k}\right)\left(h_{i}\right) e_{j k} \quad 1 \leqslant j, k \leqslant m
$$

$$
\begin{array}{ll}
{\left[h_{i}, e_{m+j, m+k}\right]=\left(\delta_{i}-\delta_{k}\right)\left(h_{i}\right) e_{m+j, m+k}} & 1 \leqslant j, k \leqslant n  \tag{2.7}\\
{\left[h_{i}, e_{j, m+k}\right]=\left(\varepsilon_{j}-\delta_{k}\right)\left(h_{i}\right) e_{j, m+k}} & 1 \leqslant j \leqslant m, 1 \leqslant k \leqslant n \\
{\left[h_{i}, e_{m+j, k}\right]=\left(-\varepsilon_{k}+\delta_{j}\right)\left(h_{i}\right) e_{m+j, k}} & 1 \leqslant j \leqslant n, 1 \leqslant k \leqslant m .
\end{array}
$$

This shows that the even (resp odd) roots are given by

$$
\begin{align*}
& \Delta_{0}=\left\{\varepsilon_{j}-\varepsilon_{k}(1 \leqslant j, k \leqslant m), \delta_{p}-\delta_{q}(1 \leqslant p, q \leqslant n)\right\}  \tag{2.8}\\
& \left(\operatorname{resp} \Delta_{1}=\left\{ \pm\left(\varepsilon_{j}-\delta_{k}\right)(1 \leqslant j \leqslant m, 1 \leqslant k \leqslant n)\right\}\right)
\end{align*}
$$

A set of simple roots is given by
$\left(\alpha_{1}, \ldots, \alpha_{m+n-1}\right)=\left(\varepsilon_{1}-\varepsilon_{2}, \ldots, \varepsilon_{m-1}-\varepsilon_{m}, \varepsilon_{m}-\delta_{1}, \delta_{1}-\delta_{2}, \ldots, \delta_{n-1}-\delta_{n}\right)$
where $\alpha_{m}=\varepsilon_{m}-\delta_{1}$ is the only odd simple root. The root vectors $E_{i}$ and $F_{1}$ corresponding with $\alpha_{i}$ and $-\alpha_{i}$, respectively, are equal to

$$
\begin{array}{ll}
E_{i}=e_{i, i+1} & F_{i}=e_{i+1, i} \quad i=1, \ldots, m-1 \\
E_{m}=e_{m, m+1} & F_{m}=e_{m+1, m}  \tag{2.10}\\
E_{m+i}=e_{m+i, m+i+1} & F_{m+i}=e_{m+i+1, m+i} \quad i=1, \ldots, n-1 .
\end{array}
$$

Then we have

$$
\begin{equation*}
\left[E_{i}, F_{j}\right]=\delta_{i j} H_{j} \tag{2.11}
\end{equation*}
$$

with

$$
\begin{align*}
& H_{i}=h_{i}-h_{i+1} \quad i=1, \ldots, m-1, m+1, \ldots, m+n-1 \\
& H_{m}=h_{m}+h_{m+1}-h_{0} . \tag{2.12}
\end{align*}
$$

The Lie superalgebra sl $(m / n)$ is generated by the set $\left\{H_{i}, E_{i}, F_{i}\right\}(1 \leqslant i \leqslant m+n-1)$, and the generators satisfy

$$
\begin{equation*}
\left[H_{i}, E_{j}\right]=c_{i j} E_{j} \quad\left[H_{i}, F_{j}\right]=-c_{i j} F_{j} \tag{2.13}
\end{equation*}
$$

with


The non-degenerate supersymmetric bilinear form $(x, y)=\operatorname{str}(x y)$ on $\operatorname{sl}(m / n)$ is still non-degenerate on H , and hence induces a non-degenerate bilinear form (, ) on the dual space $\mathrm{H}^{*}$ :

$$
\begin{align*}
& \left(\varepsilon_{i}, \varepsilon_{j}\right)=\delta_{i j}-1 /(m-n) \\
& \left(\varepsilon_{i}, \delta_{k}\right)=-1 /(m-n)  \tag{2.15}\\
& \left(\delta_{k}, \delta_{i}\right)=-\delta_{k l}-1 /(m-n) \quad 1 \leqslant i, j \leqslant m, 1 \leqslant k, l \leqslant n
\end{align*}
$$

where $\delta_{i j}$ or $\delta_{k l}$ is the Kronecker symbol.
Note that (2.9) determines the positive roots, namely

$$
\begin{align*}
& \Delta_{0}^{+}=\left\{\varepsilon_{j}-\varepsilon_{k}(1 \leqslant j<k \leqslant m), \delta_{p}-\delta_{q}(1 \leqslant p<q \leqslant n)\right\}  \tag{2.16}\\
& \Delta_{1}^{+}=\left\{\varepsilon_{j}-\delta_{k}(1 \leqslant j \leqslant m, 1 \leqslant k \leqslant n)\right\} .
\end{align*}
$$

Hence, in the notation of Kac (1978)

$$
\begin{equation*}
\bar{\Delta}_{1}^{+}=\Delta_{1}^{+} . \tag{2.17}
\end{equation*}
$$

As usual, $\rho$ is defined by

$$
\begin{equation*}
\rho=\rho_{0}-\rho_{1}=\frac{1}{2} \sum_{\alpha \in \Delta_{0}^{+}} \alpha-\frac{1}{2} \sum_{\beta \in \Delta_{1}^{+}} \beta \tag{2.18}
\end{equation*}
$$

which gives explicitly

$$
\begin{equation*}
\rho=\frac{1}{2} \sum_{i=1}^{m}(m-n-2 i+1) \varepsilon_{i}+\frac{1}{2} \sum_{j=1}^{n}(m+n-2 j+1) \delta_{j} . \tag{2.19}
\end{equation*}
$$

Furthermore, if $\Lambda \in \mathrm{H}^{*}$, we set

$$
\begin{equation*}
a_{i}=\Lambda\left(H_{i}\right) \quad i=1, \ldots, m+n-1 . \tag{2.20}
\end{equation*}
$$

With every $\Lambda \in \mathrm{H}^{*}$, an irreducible module $V(\Lambda)$ with highest weight $\Lambda$ is associated as follows. First, the elements $\left\{H_{i}, E_{i}\right\}(i=1, \ldots, m+n-1)$ generate a Borel subalgebra B of $\mathrm{sl}(m / n)$. Then, $\left\langle v_{\mathrm{A}}\right\rangle$ is a one-dimensional B module by

$$
\begin{align*}
& H_{i} v_{\mathrm{A}}=a_{i} v_{\mathrm{A}} \\
& E_{i} v_{\mathrm{A}}=0 . \tag{2.21}
\end{align*}
$$

Following Kac (1978), we set $\tilde{V}(\Lambda)=\operatorname{Ind}_{\mathrm{B}}^{\mathrm{L}}\left\langle v_{\Lambda}\right\rangle$ (with $\mathrm{L}=\operatorname{sl}(m / n)$ ) and $V(\Lambda)=$ $\tilde{V}(\Lambda) / I(\Lambda)$, where $I(\Lambda)$ is the unique maximal submodule of $\tilde{V}(\Lambda)$. Kac proved that all finite-dimensional irreps of $\operatorname{sl}(m / n)$ are of type $V(\Lambda)$ and that $V(\Lambda)$ is finitedimensional if and only if

$$
\begin{equation*}
a_{i} \in \mathbb{Z}_{+} \quad \text { for } i \neq m \tag{2.22}
\end{equation*}
$$

In general, an irreducible highest weight representation $V(\Lambda)$ is typical if

$$
\begin{equation*}
(\Lambda+\rho, \beta) \neq 0 \quad \forall \beta \in \bar{\Delta}_{1}^{+} . \tag{2.23}
\end{equation*}
$$

Making use of (2.15), (2.18) and (2.20), we deduce that

$$
\begin{gather*}
\left(\Lambda+\rho, \varepsilon_{l}-\delta_{j}\right)=a_{l}+a_{l+1}+\ldots+a_{m}-a_{m+1}-a_{m+2}-\ldots-a_{m+j-1}+m-l-j+1 \\
(1 \leqslant l \leqslant m, 1 \leqslant j \leqslant n) . \tag{2.24}
\end{gather*}
$$

These expressions give rise to the atypicality conditions for $\mathrm{sl}(\mathrm{m} / n)=\mathrm{A}(m-1, n-1)$. For finite-dimensional typical representations, a character formula is well known (Kac 1978). In the following sections, we shall obtain a character formula for atypical irreps.

## 3. Existence of an eigenvector of B

In the present section we shall prove that if $(\Lambda+\rho, \beta)=0$ for a certain $\beta$ of $\bar{\Delta}_{1}^{+}$, then there exists a weight vector $v_{\lambda-\beta}$ of $\tilde{V}(\Lambda)$ such that $v_{\lambda-\beta}$ is an eigenvector of B . The proof is very simple for $\beta=\alpha_{m}$, the unique odd simple root.

Lemma 1. There exists a weight vector $w_{m, 1}$ of $\tilde{V}(\Lambda)$, with corresponding weight $\Lambda-\alpha_{m}=\Lambda-\left(\varepsilon_{m}-\delta_{1}\right)$, such that

$$
\begin{align*}
& E_{i} w_{m, 1}=0 \quad \text { for } i \neq m \\
& E_{m} w_{m, 1}=a_{m} v_{1} . \tag{3.1}
\end{align*}
$$

Proof. Let $w_{m, 1}=F_{m} v_{1}$. This is a weight vector with weight $\Lambda-\alpha_{m}$, and for $i \neq m$ we have

$$
E_{i} w_{m, 1}=F_{m} E_{i} v_{\Lambda}=0
$$

whereas for $i=m$

$$
E_{m} w_{m, 1}=\left(-F_{m} E_{m}+H_{m}\right) v_{1}=a_{m} v_{1} .
$$

A similar property for the roots $\varepsilon_{m}-\delta_{j}(j=2, \ldots, n)$ of $\bar{\Delta}_{1}^{+}$is given in lemma 2. The proof of lemma 2 is more technical, although we only make use of the generating relations (2.11) and (2.13).

Lemma 2. For $j \in\{2, \ldots, n\}$, there exists a weight vector $w_{m, j}$ of $\tilde{V}(\Lambda)$ with weight $\Lambda-\varepsilon_{m}+\delta_{j}$, such that

$$
\begin{align*}
& E_{1} w_{m, j}=0 \quad \forall i \neq m  \tag{3.2}\\
& E_{m} w_{m, j}=\left(a_{m}-a_{m+1}-\ldots-a_{m+j-1}-j+1\right) w_{m, j-1}^{+}
\end{align*}
$$

where $w_{m, j-1}^{+}$is a weight vector with weight $\Lambda-\delta_{1}+\delta_{j}$.
Proof. First, we check the result for $j=2$. We put

$$
\begin{equation*}
x=F_{m} F_{m+1} v_{1} \quad y=F_{m+1} F_{m} v_{1} . \tag{3.3}
\end{equation*}
$$

Obviously, $E_{i} x=E_{i} y=0$ for all $i \neq m, m+1$. For $i=m+1$, we find, making use of (2.11), (2.13) and (2.21),

$$
\begin{aligned}
& E_{m+1} x=a_{m+1} w_{m, 1} \\
& E_{m+1} y=\left(a_{m+1}+1\right) w_{m, 1} .
\end{aligned}
$$

Hence we let

$$
\begin{equation*}
w_{m, 2}=\left(a_{m+1}+1\right) x-a_{m+1} y \tag{3.4}
\end{equation*}
$$

for which $E_{m+1} w_{m, 2}=0$. Then, we see

$$
\begin{align*}
E_{m} w_{m, 2} & =\left[\left(a_{m+1}+1\right)\left(a_{m}-1\right)-a_{m+1} a_{m}\right] F_{m+1} v_{1} \\
& =\left(a_{m}-a_{m+1}-1\right) F_{m+1} v_{\Lambda} \tag{3.5}
\end{align*}
$$

and $F_{m+1} v_{1}$ is a weight vector with weight $\Lambda-\delta_{1}+\delta_{2}$. Moreover, we introduce the following notation:
if $\quad w=\sum c_{i_{1} \ldots i_{k}}\left(a_{1}, \ldots, a_{m+n-2}\right) F_{i_{1}} F_{i_{2}} \ldots F_{i_{k}} v_{A} \quad\left(\right.$ all $\left.i_{k} \leqslant m+n-2\right)$
then $w^{+}=\sum c_{i_{1}, \ldots i_{h}}\left(a_{2}, \ldots, a_{m+n-1}\right) F_{i_{1}+1} \ldots F_{i_{k}+1} v_{\Lambda}$
where $c_{i_{1} \ldots i_{h}}$ are polynomial functions in $a_{1}, \ldots, a_{m+n-2}$. Hence, $w^{+}$is deduced from $w$ by increasing all the $a_{i}$ and $F_{i}$ indices by 1 . With this convention, we see that

$$
\begin{equation*}
E_{m} w_{m, 2}=\left(a_{m}-a_{m+1}-1\right) w_{m, 1}^{+} . \tag{3.7}
\end{equation*}
$$

Lemma 2 is proven in the case $j=2$. For $j>2$, we shall prove the property by means of induction on $j$. Hence, we may start from the following induction hypothesis: ( $2 \leqslant j \leqslant n-1$ ).

Let

$$
\begin{equation*}
w_{m, j}=\sum_{(k)} \lambda_{k}\left(a_{m+1}, \ldots, a_{m+j-1}\right) P_{k}\left(F_{m}, F_{m+1}, \ldots, F_{m+j-1}\right) v_{\Lambda} \tag{3.8}
\end{equation*}
$$

Here, $k$ belongs to a subset $S$ of the set $P=\{\sigma \mid \sigma$ is a permutation of $(m, m+1, \ldots, m+$ $j-1)\}$. In order to describe $S$, let $\cong$ be the equivalence relation on $P$ defined by
$\sigma_{1} \simeq \sigma_{2} \Leftrightarrow F_{\sigma_{1}(m)} \ldots F_{\sigma_{1}(m+j-1)}=F_{\sigma_{2}(m)} \ldots F_{\sigma_{2}(m+j-1)} \quad$ in $\mathrm{U}(\mathrm{sl}(m / n))$.
Then $S$ is the set of equivalence classes in $P$ for the relation $\cong$; when an equivalence class is represented by a particular representative, $S$ can be seen as a subset of $P$. Since $\left[F_{p}, F_{q}\right]=0$ if $q \neq p \pm 1$ in $\operatorname{sl}(m / n)$, we see that the cardinality of $S$ equals $2^{j-1}$ (whereas the cardinality of $P$ is, of course, $j!$ ). For $k \in S, F_{k(m)} F_{k(m+1)} \ldots F_{k(m+j-1)}$ is denoted by $P_{k}\left(F_{m}, F_{m+1}, \ldots, F_{m+j-1}\right)$ and $\lambda_{k}\left(a_{m+1}, \ldots, a_{m+j-1}\right)$ is the corresponding coefficient which is a polynomial in $a_{m+1}, \ldots, a_{m+j-1}$. The vector (3.8) satisfies the following properties:

$$
\begin{align*}
& E_{i} w_{m, j}=0 \quad \forall i \neq m \\
& E_{m} w_{m, j}=\left(a_{m}-a_{m+1}-\ldots-a_{m+j-1}-j+1\right) w_{m, j-1}^{+} \tag{3.10}
\end{align*}
$$

This describes the induction hypothesis completely and it is easy to verify it for $j=2$.
Now we will prove the property for $(j+1)$. For this purpose, define, in a similar way as in (3.3),

$$
\begin{align*}
X & =F_{m} w_{m j}^{+} \\
& =F_{m} \sum_{k} \lambda_{k}\left(a_{m+2}, \ldots, a_{m+j}\right) P_{k}\left(F_{m+1}, \ldots, F_{m+j}\right) v_{\Lambda}  \tag{3.11}\\
Y & =\sum_{k} \lambda_{k}\left(a_{m+2}, \ldots, a_{m+j}\right) P_{k}\left(F_{m+1}, \ldots, F_{m+j}\right) F_{m} v_{\Lambda}
\end{align*}
$$

Then, for $i \neq m, m+1$,

$$
E_{i} X=F_{m} E_{i} w_{m j}^{+}
$$

However since $i>m+1, E_{i} w_{m j}^{+}$is equal to $\left(E_{i-1} w_{m j}\right)^{+}$, which vanishes according to (3.10). Hence

$$
\begin{equation*}
E_{i} X=0 \quad \text { for } i \neq m, m+1 \tag{3.12}
\end{equation*}
$$

Making use of the explicit form of $X$, the action of $E_{i}$ implies that all $F_{i}$ are replaced by $H_{i}$ :

$$
\begin{align*}
0=E_{i} X=F_{m} & \sum_{k} \lambda_{k}\left(a_{m+2}, \ldots, a_{m+j}\right) P_{k}\left(F_{m+1}, \ldots, H_{i}, \ldots, F_{m+j}\right) v_{A} \\
& i \neq m, m+1 . \tag{3.13}
\end{align*}
$$

In (3.13), $H_{i}$ is substituted by its eigenvalue $\left(a_{1}+c\right)$, where the number $c$ is 0,1 or 2 , if no, one or both of the operators $F_{i-1}, F_{i+1}$ are placed behind $H_{i}$ in the explicit form of $P_{k}$, respectively. But then ( $i \neq m, m+1$ )
$E_{i} Y=\sum_{k} \lambda_{k}\left(a_{m+2}, \ldots, a_{m+j}\right) P_{k}\left(F_{m+1}, \ldots, H_{i}, \ldots, F_{m+j}\right) F_{m} v_{A}=0$
because $H_{i}$ is substituted by the same eigenvalue, since $\left[H_{i}, F_{m}\right.$ ] $=0$. Now, consider the case of $E_{i}=E_{m+1}$. Obviously

$$
E_{m+1} X=F_{m} E_{m+1} w_{m j}^{+} .
$$

The vector $E_{m+1} w_{m, j}^{+}$is not equal right away to $\left(E_{m} w_{m j}\right)^{+}$, since the $m$ th row of the Cartan matrix $C$ (2.14) differs from the next rows. But that $E_{m+1} w_{m, j}^{+}$is still related to $\left(E_{m} w_{m j}\right)^{+}$can be seen as follows.
(i) In order to compute $E_{m+1} w_{m j}^{+}$, all $F_{m+1}$ in $w_{m j}^{+}$are replaced by $H_{m+1}$ and then $H_{m+1}$ is substituted by $a_{m+1}-c_{m+1, m+2}=a_{m+1}+1$ if $F_{m+2}$ appears behind $H_{m+1}$, and by $a_{m+1}$ otherwise.
(ii) In order to compute $E_{m} w_{m j}$, all $F_{m}$ in $w_{m j}$ are replaced by $H_{m}$ and then $H_{m}$ is substituted by $a_{m}-c_{m, m+1}=a_{m}-1$ if $F_{m+1}$ appears behind $H_{m}$, and by $a_{m}$ otherwise. Hence, if we take $E_{m} w_{m j}$, replace all $a_{m}$ by $-a_{m}$, then change the signs of all coefficients, and finally increase all indices by 1 , we must obtain $E_{m+1} w_{m j}^{+}$. But $E_{m} w_{m j}$ is given in (3.10), hence

$$
\begin{equation*}
E_{m+1} w_{m j}^{+}=\left(a_{m+1}+a_{m+2}+\ldots+a_{m+j}+j-1\right) w_{m, j-1}^{++} \tag{3.15}
\end{equation*}
$$

or

$$
\begin{equation*}
E_{m+1} X=\left(a_{m+1}+a_{m+2}+\ldots+a_{m+j}+j-1\right) F_{m} w_{m, j-1}^{++} \tag{3.16}
\end{equation*}
$$

Then we find

$$
\begin{equation*}
E_{m+1} Y=\sum_{k} \lambda_{k}\left(a_{m+2}, \ldots, a_{m+j}\right) P_{k}\left(H_{m+1}, F_{m+2}, \ldots, F_{m+j}\right) F_{m} v_{\lambda} \tag{3.17}
\end{equation*}
$$

The expression on the right-hand side of (3.17) is the same as in the one for $E_{m+1} X$, but with $F_{m}$ placed at the end. This implies that, compared to $E_{m+1} X$, all $a_{m+1}$ are replaced by $a_{m+1}-c_{m+1, m}=a_{m+1}+1$. This leads to

$$
\begin{equation*}
E_{m+1} Y=\left(a_{m+1}+a_{m+2}+\ldots+a_{m+j}+j\right) F_{m} w_{m, j-1}^{++} \tag{3.18}
\end{equation*}
$$

Hence, if we put

$$
\begin{equation*}
w_{m, j+1}=\left(a_{m+1}+a_{m+2}+\ldots+a_{m+j}+j\right) X-\left(a_{m+1}+a_{m+2}+\ldots+a_{m+j}+j-1\right) Y \tag{3.19}
\end{equation*}
$$

$w_{m, j+1}$ is precisely of the form (3.8) with $j \rightarrow j+1$, and from (3.12), (3.14), (3.16) and (3.18) we see that

$$
\begin{equation*}
E_{i} w_{m, j+1}=0 \quad \forall i \neq m \tag{3.20}
\end{equation*}
$$

Finally, since

$$
E_{m} X=H_{m} w_{m j}^{+}=\left(a_{m}-1\right) w_{m j}^{+}
$$

and

$$
E_{m} Y=\sum_{k} \lambda_{k}\left(a_{m+2}, \ldots, a_{m+j}\right) P_{k}\left(F_{m+1}, \ldots, F_{m+j}\right) H_{m} v_{1}=a_{m} w_{m j}^{+}
$$

we find

$$
\begin{align*}
E_{m} w_{m, j+1}= & {\left[\left(a_{m+1}+a_{m+2}+\ldots+a_{m+i}+j\right)\left(a_{m}-1\right)-\left(a_{m+1}+\ldots+a_{m+j}+j-1\right) a_{m}\right] w_{m, j}^{+} } \\
& =\left(a_{m}-a_{m+1}-\ldots-a_{m+j}-j\right) w_{m j}^{+} \tag{3.21}
\end{align*}
$$

and obviously $w_{m j}^{+}$is a weight vector with weight $\Lambda-\delta_{1}+\delta_{j+1}$.
In the next lemma, the property (3.2) will be extended for the case of weight vectors with weight $\Lambda-\left(\varepsilon_{l}-\delta_{l}\right)$ with $1 \leqslant l \leqslant m-1$.
Lemma 3. For $l \in\{1, \ldots, m\}$ and $j \in\{1, \ldots, n\}$, there exists a weight vector $w_{l j}$ of $\tilde{V}(\Lambda)$ with weight $\Lambda-\varepsilon_{l}+\delta j$, such that

$$
\begin{align*}
& E_{i} w_{l j}=0 \quad \forall i \neq l  \tag{3.22}\\
& E_{l} w_{l j}=\left(a_{l}+a_{l+1}+\ldots+a_{m}-a_{m+1}-\ldots-a_{m+j-1}+m-l-j+1\right) w_{l+1, j}
\end{align*}
$$

where $w_{l+1, j}$ is a weight vector with weight $\Lambda-\varepsilon_{l+1}+\delta_{j}$ for $l<m$.
Proof. From lemma 2 we see that (3.22) is true for $l=m$. Hence we will prove lemma 3 by induction on $l$. Thus, we may suppose that there exists a vector $w_{l j}$ of the following form:

$$
\begin{equation*}
w_{l j}=\sum_{(k)} \lambda_{k} P_{k}\left(F_{l}, F_{l+1}, \ldots, F_{m+j-1}\right) v_{1} \tag{3.23}
\end{equation*}
$$

such that (3.22) holds. Herein, $\lambda_{k}$ is a coefficient dependent on $a_{l+1}, a_{l+2}, \ldots, a_{m+j-1}$. We put

$$
\begin{align*}
& x=F_{l-1} w_{l j}=F_{l-1} \sum_{k} \lambda_{k} P_{k}\left(F_{l}, \ldots, F_{m+j-1}\right) v_{\lambda} \\
& y=\sum_{k} \lambda_{k} P_{k}\left(F_{l}, \ldots, F_{m+j-1}\right) F_{i-1} v_{\lambda} . \tag{3.24}
\end{align*}
$$

Then, for $i \neq l, l-1$, we find that

$$
\begin{equation*}
E_{i} x=F_{l-1} E_{i} w_{l j}=0 \tag{3.25}
\end{equation*}
$$

and in a similar way as in the proof of lemma 2, this implies

$$
\begin{equation*}
E_{i} y=0 \quad \text { for } i \neq l, l-1 . \tag{3.26}
\end{equation*}
$$

Next, we take $E_{i}=E_{i}$. Then the induction hypothesis leads to

$$
\begin{equation*}
E_{l} x=\left(a_{1}+a_{l+1}+\ldots+a_{m}-a_{m+1}-\ldots-a_{m+j-1}+m-l-j+1\right) F_{l-1} w_{l+1, j} . \tag{3.27}
\end{equation*}
$$

Similarly as in (3.17)-(3.18), we see that $E_{l} y$ can be deduced from $E_{l} x$ by replacing $a_{l}$ by $a_{1}-c_{l, 1-1}=a_{1}+1$. Hence

$$
\begin{equation*}
E_{l} y=\left(a_{l}+a_{l+1}+\ldots+a_{m}-a_{m+1}-\ldots-a_{m+j-1}+m-l-j+2\right) F_{l-1} w_{l+1, j} . \tag{3.28}
\end{equation*}
$$

So, we define

$$
\begin{align*}
& w_{l-1, j}=\left(a_{l}+\ldots+a_{m}-a_{m+1}-\ldots-a_{m+j-1}+m-l-j+2\right) x \\
&-\left(a_{l}+\ldots+a_{m}-a_{m+1}-\ldots-a_{m+j-1}+m-l-j+1\right) y \tag{3.29}
\end{align*}
$$

Then, also $E_{l} w_{l-1, j}=0$, and since $E_{l-1} x=H_{l-1} w_{l j}=\left(a_{j-1}+1\right) w_{l j}$ and $E_{l-1} y=a_{l-1} w_{l j}$, one obtains
$E_{l-1} w_{l-1, j}=\left(a_{l-1}+a_{l}+\ldots+a_{m}-a_{m+1}-\ldots-a_{m+j-1}+m-l-j+2\right) w_{l j}$.
This proves lemma 3.

Theorem 1. Let $\beta \in \bar{\Delta}_{1}^{+}=\Delta_{1}^{+}$. Then $(\Lambda+\rho, \beta)=0$ if and only if there exists a weight vector $v_{1-\beta}$ of $\dot{V}(\Lambda)$ with weight $\Lambda-\beta$ such that $v_{\lambda-\beta}$ is an eigenvector of the Borel subalgebra B. Moreover, such a vector $v_{1-\beta}$ is unique (up to a non-zero factor).

Proof. If $(\Lambda+\rho, \beta)=0$ then (2.24) and the previous lemmas show the existence of a vector $v_{1-\beta}$ with the required properties. On the other hand, let $v_{1-\beta}$ be a weight vector, which is an eigenvector of $B$. Since $\beta \in \Delta_{1}^{+}$, we can put $\beta=\varepsilon_{l}-\delta_{j}(1 \leqslant l \leqslant m, 1 \leqslant$ $j \leqslant n)$. Then, the most general form of a weight vector in $\tilde{V}(\Lambda)$ with weight $\Lambda-\beta$ is equal to

$$
\begin{equation*}
v_{1-\beta}=\sum_{k} \mu_{k} P_{k}\left(F_{l}, F_{l+1}, \ldots, F_{m+j-1}\right) v_{A} \tag{3.31}
\end{equation*}
$$

with $\mu_{k}$ arbitrary coefficients and ( $k$ ) belonging to the subset $S$ of permutations of $(l, l+1, \ldots, m+j-1)$, as described in the proof of lemma 2. Since $v_{1-\beta}$ is an eigenvector of B , we have to require that $E_{i} v_{\lambda-\beta}=0$, for all $i$. But

$$
\begin{equation*}
E_{i} v_{A-\beta}=0 \quad \text { for } i=l+1, l+2, \ldots, m+j-1 \tag{3.32}
\end{equation*}
$$

are precisely the conditions that determine the coefficients $\mu_{k}$ (up to a factor), as can be seen from the proof of lemmas 2 and 3. Hence, $\mu_{k}=c \lambda_{k}$ for all $k$ and $v_{1-\beta}=c w_{l j}$. This shows that $v_{A-\beta}$ is unique and

$$
\begin{equation*}
E_{l} v_{1-\beta}=(\Lambda+\rho, \beta) v^{\prime} \tag{3.33}
\end{equation*}
$$

implies that $(\Lambda+\rho, \beta)=0$, since $v_{1-\beta}$ is an eigenvector of B .
Note that theorem 1 can partly be deduced from theorem 3 of Kac (1978). The main difference, however, is that we have given an explicit construction of the vector $v_{1-\beta}$.

## 4. The character formula for atypical representations

Let $w \in W$ be the elements of the Weyl group of $\operatorname{sl}(m / n)_{\overline{0}}=\operatorname{sl}(m) \oplus \operatorname{sl}(n) \oplus \mathbb{C}$. Then the character of a typical finite-dimensional irrep is given by (Kac 1978)
$\operatorname{ch} V(\Lambda)=\frac{\Pi_{\beta \in \Delta_{i}^{+}}[\exp (\beta / 2)+\exp (-\beta / 2)]}{\Pi_{\alpha \in \Delta_{0}^{+}}[\exp (\alpha / 2)-\exp (-\alpha / 2)]} \sum_{w \in W} \varepsilon(w) \exp [w(\Lambda+\rho)]$.
Here, $\exp (\alpha)$ is the notation for the formal exponential. Since $\Delta_{1}^{+}=$ $\left\{\varepsilon_{i}-\delta_{j}(i=1, \ldots, m ; j=1, \ldots, n)\right\}$, we may denote this set by $\Delta_{1}^{+}=\left\{\beta_{1}, \beta_{2}, \ldots, \beta_{m n}\right\}$. Also, since

$$
\begin{equation*}
\rho_{1}=\frac{1}{2} \sum_{i=1}^{m n} \beta_{i} \tag{4.2}
\end{equation*}
$$

the following equality holds:

$$
\begin{align*}
\prod_{i=1}^{m n}\left[\exp \left(\beta_{i} / 2\right)\right. & \left.+\exp \left(-\beta_{i} / 2\right)\right]=\exp \left(\rho_{1}\right)+\sum_{i=1}^{m n} \exp \left(\rho_{1}-\beta_{i}\right) \\
& +\sum_{1 \leqslant i<j \leqslant m n} \sum_{i \leqslant p} \exp \left(\rho_{1}-\beta_{i}-\beta_{j}\right)+\ldots+\exp \left(-\rho_{1}\right) . \tag{4.3}
\end{align*}
$$

On the right-hand side of (4.3) there appear $(m n+1)$ parts with, respectively, $1, m n$, $m n(m n-1) / 2, m n(m n-1)(m n-2) / 3!, \ldots, 1$ term(s). The $(k+1)$ part in (4.3) is of the form

$$
\begin{equation*}
\sum_{1 \leqslant i_{1}<i_{2}<\ldots<i_{k} \leqslant m n} \exp \left(\rho_{1}-\beta_{i_{1}}-\beta_{i_{2}}-\ldots-\beta_{i_{h}}\right) . \tag{4.4}
\end{equation*}
$$

Since $w\left(\rho_{1}\right)=\rho_{1}$ (see Kac 1978, proposition 1.7 c ), $\Delta_{1}^{+}$is invariant under $W$ and $W$ acts transitively on $\Delta_{1}^{+}$, we find that

$$
\begin{equation*}
\sum \exp \left(\rho_{1}-\beta_{i_{1}}-\ldots-\beta_{i_{k}}\right)=\sum \exp \left(\rho_{1}-w\left(\beta_{i_{1}}\right)-\ldots-w\left(\beta_{i_{k}}\right)\right) \tag{4.5}
\end{equation*}
$$

for any element $w \in W$, where the summation in (4.5) is the same as in (4.4). Hence, (4.3) is $W$ invariant. Then, combining (4.1)-(4.5), one obtains
ch $V(\Lambda) \prod_{\alpha \in \Delta_{0}^{+}}[\exp (\alpha / 2)-\exp (-\alpha / 2)]$

$$
\begin{align*}
= & \sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}\right)\right]+\sum_{i=1}^{m n}\left(\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-\beta_{i}\right)\right]\right)+\ldots \\
& +\sum_{1 \leqslant i_{1}<\ldots<i_{1} \leqslant m n}\left(\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-\beta_{i_{1}}-\ldots-\beta_{i_{k}}\right)\right]\right)+\ldots \\
& +\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-2 \rho_{1}\right)\right] . \tag{4.6}
\end{align*}
$$

Expression (4.6) contains $2^{m n}$ non-zero $\mathrm{sl}(m) \oplus \mathrm{sl}(n)$ characteristics if and only if all $\Lambda-\beta_{i_{1}}-\ldots-\beta_{i_{h}}$ correspond to highest weights of $\operatorname{sl}(m) \oplus \operatorname{sl}(n)$, i.e., if and only if

$$
\begin{array}{ll}
\text { all } a_{i} \geqslant n & \text { for } 1 \leqslant i \leqslant m-1 \\
\text { all } a_{i} \geqslant m & \text { for } m+1 \leqslant i \leqslant m+n-1 . \tag{4.7}
\end{array}
$$

In that case, it follows from (4.6) that the irrep $V(\Lambda)$ splits in $2^{m n}$ irreps of the Lie subalgebra $\operatorname{sl}(m) \oplus \operatorname{sl}(n)$.

Since $\operatorname{sl}(m / n)$ is a Lie superalgebra of type I , we have the following $\mathbb{Z}$ gradation:

$$
\begin{equation*}
\operatorname{sl}(m / n)=\mathrm{L}_{-1} \oplus \mathrm{~L}_{0} \oplus \mathrm{~L}_{+1} \tag{4.8}
\end{equation*}
$$

where $\mathrm{L}_{0}=\mathrm{L}_{\overline{0}}$ and $\mathrm{L}_{1}$ (resp $\mathrm{L}_{-1}$ ) is the subalgebra of $\operatorname{sl}(m / n)$ spanned by the matrices $\left(\begin{array}{ll}0 & b \\ 0 & 0\end{array}\right)\left(\operatorname{resp}\left(\begin{array}{ll}0 & 0 \\ c & 0\end{array}\right)\right.$ ) in (2.1). Put

$$
\begin{equation*}
\mathrm{P}=\mathrm{L}_{0} \oplus \mathrm{~L}_{+1} . \tag{4.9}
\end{equation*}
$$

Let $\Lambda$ be a weight such that it is dominant integral for the Cartan subalgebra of $\operatorname{sl}(m) \oplus \operatorname{sl}(n)$. Then $\Lambda$ determines unambiguously an $\mathrm{L}_{\overline{0}}$ module $V^{0}(\Lambda)$ which is irreducible and finite dimensional. We extend $V^{0}(\Lambda)$ to a $P$ module by putting $L_{+1}\left(V^{0}(\Lambda)\right)=0$. Then we define, following Kac (1978), the induced L module

$$
\begin{equation*}
\bar{V}(\Lambda)=\operatorname{Ind}_{\mathrm{P}}^{\mathrm{L}} V^{0}(\Lambda) . \tag{4.10}
\end{equation*}
$$

Under these conditions, $\bar{V}(\Lambda)$ is finite dimensional and it contains a unique maximal submodule $\bar{I}(\Lambda)$ such that

$$
\begin{equation*}
V(\Lambda)=\bar{V}(\Lambda) / \bar{I}(\Lambda) . \tag{4.11}
\end{equation*}
$$

Now we make use of the following property (Kac 1978, proposition 2.1.c).
Proposition 1. Let $\mathrm{H} \subset \mathrm{L}$ be a subalgebra of L containing $\mathrm{L}_{\overline{0}}$ and $g_{1}, \ldots, g_{i}$ be odd elements of L whose projections onto $\mathrm{L} / \mathrm{H}$ form a basis. Let Z be an H module. Then, one has the vector space decomposition

$$
\begin{equation*}
\operatorname{Ind}_{H}^{L} Z=\bigoplus_{i \leqslant i_{1}<\ldots<i_{r} \leqslant i} g_{i_{1}} \ldots, g_{i,} Z . \tag{4.12}
\end{equation*}
$$

We can apply this for $\mathrm{H}=\mathrm{P}$ and $\mathrm{Z}=V^{0}(\Lambda)$ :
$\bar{V}(\Lambda)=\bigoplus_{1 \leqslant i_{1}<\ldots<i_{s} \leqslant m n} g_{i_{1}} \ldots g_{i_{s}} V^{0}(\Lambda)$
where $g_{i}$ is a root vector of the root space $\mathrm{L}^{-\beta_{i}}$. The character of $V^{0}(\Lambda)$ is given by Weyl's character formula. Hence (4.13) implies
ch $\bar{V}(\Lambda)=\prod_{i=1}^{m n}\left[1+\exp \left(-\beta_{i}\right)\right] \frac{\Sigma_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}\right)\right]}{\Pi_{\alpha \in \Delta_{0}^{+}}[\exp (\alpha / 2)-\exp (-\alpha / 2)]}$.
Making use of arguments as in (4.4)-(4.5), this leads to

$$
\text { ch } \begin{align*}
\bar{V}(\Lambda) \prod_{\alpha \in \Delta_{0}^{+}} & {[\exp (\alpha / 2)-\exp (-\alpha / 2)] } \\
& =\sum_{1 \leqslant i_{1}<\ldots<i_{k} \leqslant m m}\left(\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-\beta_{i_{1}}-\ldots-\beta_{i_{k}}\right)\right]\right) . \tag{4.15}
\end{align*}
$$

Note that the right-hand sides of (4.6) and (4.15) are equal, but (4.6) is valid only for typical weights $\Lambda$, whereas (4.15) can also be used for atypical weights (in both cases $\Lambda$ is the dominant integral for $\operatorname{sl}(m) \oplus \mathrm{sl}(n))$. In fact, this shows that $\mathrm{ch} V(\Lambda)=\mathrm{ch} \bar{V}(\Lambda)$ for $\Lambda$ typical, which is obvious, since Kac (1978) proved that in this case $V(\Lambda)=\bar{V}(\Lambda)$. When $\Lambda$ is an atypical highest weight, $\bar{V}(\Lambda)$ will not be irreducible, and then the question is: what is the maximal invariant subspace of $\vec{V}(\Lambda)$ ?

Obviously, $\bar{V}(\Lambda)$ is also an $\mathrm{L}_{\overline{0}}$ module and, according to (4.15), it splits in $2^{m n}$ parts with highest weights given by $\Lambda-\beta_{i_{1}}-\ldots-\beta_{i_{k}}\left(1 \leqslant i_{1}<\ldots<i_{k} \leqslant m n\right)$. Let $S$ be the set of sequences $s=\left(i_{1}, \ldots, i_{k}\right)$ :

$$
\begin{equation*}
S=\left\{s=\left(i_{1}, \ldots, i_{k}\right) \mid 1 \leqslant i_{1}<\ldots<i_{k} \leqslant m n\right\} . \tag{4.16}
\end{equation*}
$$

With every element $s$ of $S$ there corresponds a weight $\Lambda_{s}$ under the mapping $\sigma: S \rightarrow$ $H^{*}: s \rightarrow \Lambda_{s}=\Lambda-\beta_{i_{1}}-\ldots-\beta_{i k}$. The elements $\Lambda_{s}(s \in S)$ are called the weights of $S$. Note that $S$ contains $2^{m n}$ elements, but to different elements of $S$ there may correspond equal weights (for instance, if $\beta_{i}+\beta_{j}=\beta_{k}+\beta_{l}$ ). The weights of $S$ are dominant (and hence correspond to highest weights of irreducible finite $L_{\overline{0}}$ modules) if (4.7) is satisfied. From now on, we shall suppose that this is the case. Let $v_{\Lambda_{1}}^{0}(s \in S)$ be the highest weight vectors of the $\mathrm{L}_{\overline{0}}$ modules $V^{0}\left(\Lambda_{5}\right)$. Since these highest weight vectors are elements of $\bar{V}(\Lambda)$, they can be written in the following form:

$$
\begin{equation*}
v_{\Lambda_{1}}^{0}=\sum g_{j_{1}} \ldots g_{j_{k}} Q_{(j)} \otimes v_{\Lambda} \quad s=\left(i_{1}, \ldots, i_{k}\right) . \tag{4.17}
\end{equation*}
$$

Here, the summation is over certain root vectors $g_{j_{1}}, \ldots, g_{j_{k}}$ (corresponding to negative odd roots $\left.-\beta_{j_{1}}, \ldots,-\beta_{j_{k}}\right)$ and elements $Q_{(j)} \in \mathrm{U}\left(\mathrm{L}_{\overline{0}}\right)$ with weight $-\Sigma n_{\alpha} \alpha\left(\alpha \in \Delta_{0}^{+}, n_{\alpha} \geqslant\right.$ 0 ) such that

$$
\begin{equation*}
\beta_{i_{1}}+\ldots+\beta_{i_{h}}=\beta_{j_{1}}+\ldots+\beta_{j_{k}}+\sum n_{\alpha} \alpha \tag{4.18}
\end{equation*}
$$

is satisfied. Note that due to (2.8) the number of odd roots on both sides of (4.18) is always the same. Also, (4.17) always contains a non-vanishing term of the form $g_{i_{1}} \ldots g_{i_{h}} \otimes v_{1}$.

Let us now consider the situation in which precisely one atypicality condition is fulfilled:

$$
\begin{align*}
& \left(\Lambda+\rho, \beta_{p}\right)=0 \\
& \left(\Lambda+\rho, \boldsymbol{\beta}_{i}\right) \neq 0 \quad \text { for } i \neq p \tag{4.19}
\end{align*}
$$

Then theorem 1 shows the existence of a $B$ eigenvector $v_{1-\beta_{p}}^{0}$ of $\bar{V}(\Lambda)$ :

$$
\begin{equation*}
v_{1-\beta_{r}}^{0}=R \otimes v_{1} \quad R \in \mathrm{U}(\mathrm{~L}) . \tag{4.20}
\end{equation*}
$$

In fact, an explicit construction of $R$ has been given in $\S 3$. Note that $R$ can be written in the following form:

$$
\begin{equation*}
R=g_{p}+g_{p^{\prime}} T_{-\alpha_{1}}+g_{p^{\prime}} T_{-\alpha_{1}^{\prime}} T_{-\alpha_{2}^{\prime}}+\ldots \tag{4.21}
\end{equation*}
$$

where $\alpha_{1}, \alpha_{1}^{\prime}, \alpha_{2}^{\prime}, \ldots \in \Delta_{0}^{+}, T_{\alpha}$ are the corresponding root vectors and

$$
\begin{equation*}
\beta_{p}=\beta_{p}+\alpha_{1}=\beta_{p^{\prime \prime}}+\alpha_{1}^{\prime}+\alpha_{2}^{\prime}=\ldots \tag{4.22}
\end{equation*}
$$

It is obvious that $v_{1-\beta_{n}}^{0}$ can be identified with an element of $\bar{V}(\Lambda)$, also given by (4.20):

$$
\begin{equation*}
v_{1-\beta_{p}}^{0} \in \bar{V}(\Lambda) \tag{4.23}
\end{equation*}
$$

The weight $\Lambda-\beta_{p}$ determines an $L_{\overline{0}}$ module $V^{0}\left(\Lambda-\beta_{p}\right)$ with highest weight vector $u_{\Lambda-\beta_{p}}$. Then, just as in (4.10), $\bar{V}\left(\Lambda-\beta_{p}\right)$ is well defined. The L module $\bar{V}\left(\Lambda-\beta_{p}\right)$ splits into a number of irreducible $L_{\overline{0}}$ modules:

$$
\bar{V}\left(\Lambda-\beta_{p}\right)=X_{1} \oplus \ldots \oplus X_{N} \quad N=2^{m n}
$$

Let $x_{k}$ be the $\mathrm{L}_{\overline{0}}$ highest weight vector of $X_{k}$. The weight of $x_{k}$ is of the form $\sigma\left(s_{k}\right)-\beta_{p}$ for some $s_{k} \in S$. Suppose $s_{k}=\left(i_{1}, \ldots, i_{1}\right)$, then the analogue of (4.17) is

$$
\begin{equation*}
x_{k}=\sum g_{j_{1}} \ldots g_{j l} Q_{(j)}^{\prime} \otimes u_{\lambda-\beta_{p}} . \tag{4.24}
\end{equation*}
$$

We define a mapping $A$ from the elements $x_{k}$ in (4.24) to $\bar{V}(\Lambda)$ by replacing $u_{A-\beta_{r}}$ by $R \otimes \mathcal{V}$, on the rhs of (4.24). Then $A$ is extended to elements of the form $q x_{k}, q \in \mathrm{U}\left(\mathrm{L}_{\overline{0}}\right)$ by putting $A\left(q x_{k}\right)=q\left(A x_{k}\right)$. In this way, $A$ is defined for all elements of $\bar{V}\left(\Lambda-\beta_{p}\right)$. Also, every element of $\bar{V}\left(\Lambda-\beta_{p}\right)$ can be uniquely written as a sum of elements of the form $q x_{k}$.

Lemma 4. The mapping $A: \bar{V}\left(\Lambda-\beta_{p}\right) \rightarrow \bar{V}(\Lambda)$ satisfies the following properties:
(a) it is a homomorphism of L modules: if $u, u^{\prime} \in \bar{V}\left(\Lambda-\beta_{p}\right), x \in \mathrm{~L}$ and $x \cdot u=u^{\prime}$, then $x \cdot A u=A u^{\prime}$;
(b) $x_{k}$ with weight $\sigma\left(s_{k}\right)-\beta_{p}$ is mapped into the zero vector if and only if $\sigma\left(s_{k}\right)-$ $\beta_{p} \notin \sigma(S)$.

Proof. It is enough to prove (a) for elements of the form $u=q x_{k}, q \in \mathrm{U}\left(\mathrm{L}_{\overline{0}}\right)$. But for such elements (a) follows from the fact that $v_{A-\beta_{p}}$ is a highest weight vector for $L$ and the definitions of $\bar{V}\left(\Lambda-\beta_{p}\right)$ and $\bar{V}(\Lambda)$. As a consequence, $A$ preserves the weight of weight vectors. To prove (b), let $s_{k}=\left(i_{1}, \ldots, i_{i}\right)$ and $\sigma\left(s_{k}\right)-\beta_{p} \notin \sigma(S)$. According to (a), $A x_{k}$ must be a highest weight vector for $\mathrm{L}_{\overline{0}}$. But all highest weight vectors for $\mathrm{L}_{\overline{0}}$ contained in $\bar{V}(\Lambda)$ are given by $v_{\Lambda_{\dot{c}}}^{0}$. Hence $\sigma\left(s_{k}\right)-\beta_{p} \notin \sigma(S)$ implies $A x_{k}=0$. Conversely, if $\sigma\left(s_{k}\right)-\beta_{p} \in \sigma(S)$, then $p \notin\left\{i_{1}, \ldots, i_{l}\right\}$. Now $x_{k}$ contains a unique term of the form $g_{i_{1}} \ldots g_{i_{1}} \otimes u_{\lambda-\beta_{p}}$. Hence, according to (4.19) $A x_{k}$ contains a unique term of the form

$$
g_{i_{1}} \ldots g_{i,} g_{p} \otimes v_{\Lambda}
$$

which shows that $A x_{k} \neq 0$.
The image of $\bar{V}\left(\Lambda-\beta_{p}\right)$ under $A, \operatorname{Im} \bar{V}\left(\Lambda-\beta_{p}\right)=J\left(\Lambda-\beta_{p}\right)$, is a submodule of $\bar{V}(\Lambda)$ with highest weight $\Lambda-\beta_{p}$. From lemma 4 , it follows that the $L_{\bar{o}}$ components $V^{0}(\lambda)$
of $J\left(\Lambda-\beta_{p}\right)$ are those for which $\lambda \in \sigma(S)$ and for which $\lambda$ is a highest weight of an $\mathrm{L}_{\overline{0}}$ module contained in $\bar{V}\left(\Lambda-\beta_{p}\right)$. Hence

$$
\text { ch } \begin{align*}
J\left(\Lambda-\beta_{p}\right) & \prod_{\alpha \in \Delta_{0}^{+}}[\exp (\alpha / 2)-\exp (-\alpha / 2)] \\
& =\sum_{\substack{1 \leqslant i_{1}<\ldots<i_{k} \leq m n \\
i_{1} \neq p, \ldots, i_{k} \neq p}}\left(\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda-\beta_{p}+\rho_{0}-\beta_{i_{1}}-\ldots-\beta_{i_{k}}\right)\right]\right) . \tag{4.25}
\end{align*}
$$

Any other submodule of $\bar{V}(\Lambda)$ must have a highest weight from $S$. But from theorem 1 it follows that no weight vector of $\bar{V}(\Lambda)$ with weight $\Lambda-\beta_{i}(i \neq p)$ can be a highest weight vector for L , and from $\operatorname{Kac}(1978$, theorem $3 b$ ) one can conclude that no other weights of the form $\Lambda-\beta_{i_{1}}-\ldots-\beta_{i_{h}}\left(i_{1} \neq p, \ldots, i_{k} \neq p\right)$ correspond to highest weight vectors for L. Hence

$$
\begin{equation*}
\bar{I}(\Lambda)=J\left(\Lambda-\beta_{p}\right) . \tag{4.26}
\end{equation*}
$$

Then $\operatorname{ch} V(\Lambda)=\operatorname{ch} \bar{V}(\Lambda)-\operatorname{ch} J\left(\Lambda-\beta_{p}\right)$, or explicitly

$$
\begin{align*}
\operatorname{ch} V(\Lambda) \prod_{\alpha \in \Delta_{0}^{+}} & {[\exp (\alpha / 2)-\exp (-\alpha / 2)] } \\
= & \sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}\right)\right]+\sum_{\substack{i=1 \\
i \neq p}}^{m n}\left(\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-\beta_{i}\right)\right]\right)+\ldots \\
& +\sum_{\substack{1 \leqslant i_{i}<\ldots, i_{n}=m n \\
i_{1} \neq p_{, \ldots,} \neq p}}\left(\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-\beta_{i_{1}}-\ldots-\beta_{i_{h}}\right)\right]\right)+\ldots \\
& +\sum_{w \in W} \varepsilon(w) \exp \left[w\left(\Lambda+\rho_{0}-2 \rho_{1}+\beta_{p}\right)\right] . \tag{4.2}
\end{align*}
$$

We shall now obtain a more convenient form of (4.27). For this purpose, let $w_{1}$ be a fixed element of $W$ and consider in all the parts of (4.27) the term with $w=w_{1}$. This is equal to

$$
\begin{equation*}
\varepsilon\left(w_{0}\right) \exp \left[w_{1}\left(\Lambda+\rho_{0}\right)\right]\left(1+\sum_{i \neq p} \exp \left[-w_{1}\left(\beta_{i}\right)\right]+\ldots+\exp \left[-w_{1}\left(2 \rho_{1}-\beta_{p}\right)\right]\right) . \tag{4.28}
\end{equation*}
$$

The right-hand side is

$$
\begin{align*}
& \prod_{\substack{i=1 \\
i \neq p}}^{m n}\left\{1+\exp \left[-w_{1}\left(\beta_{i}\right)\right]\right\} \\
&= w_{1}\left(\prod_{i=1}^{m n}\left[1+\exp \left(-\beta_{i}\right)\right]\left[1+\exp \left(-\beta_{p}\right)\right]^{-1}\right) \\
&= w_{1}\left(\exp \left(-\rho_{1}+\beta_{p} / 2\right) \prod_{i=1}^{m n}\left[\exp \left(\beta_{i} / 2\right)+\exp \left(-\beta_{i} / 2\right)\right]\right. \\
&\left.\times\left[\exp \left(\beta_{p} / 2\right)+\exp \left(-\beta_{p} / 2\right)\right]^{-1}\right) . \tag{4.29}
\end{align*}
$$

But $w_{1}\left(\rho_{1}\right)=\rho_{1}$ and $\Pi\left[\exp \left(\beta_{i} / 2\right)+\exp \left(-\beta_{i} / 2\right)\right]$ is $W$ invariant. Hence, (4.29) equals

$$
\begin{equation*}
\exp \left[-\rho_{1}+w_{1}\left(\beta_{p} / 2\right)\right] \frac{\Pi_{i=1}^{m n}\left[\exp \left(\beta_{i} / 2\right)+\exp \left(-\beta_{i} / 2\right)\right]}{\left\{\exp \left[w_{1}\left(\beta_{p} / 2\right)\right]+\exp \left[-w_{1}\left(\beta_{p} / 2\right)\right]\right\}} \tag{4.30}
\end{equation*}
$$

Combining the previous results, we find the following.
Theorem 2. Let $\Lambda$ be the highest weight of an atypical representation for which the conditions (4.7) and (4.19) are fulfilled. Then
ch $V(\Lambda) \prod_{\alpha \in \Delta_{0}^{+}}[\exp (\alpha / 2)-\exp (-\alpha / 2)]$

$$
\begin{align*}
= & \sum_{i=1}^{m n}\left[\exp \left(\beta_{1} / 2\right)+\exp \left(-\beta_{1} / 2\right)\right] \ldots\left[\exp \left(\beta_{i} / 2\right)+\exp \left(-\beta_{i} / 2\right)\right] \ldots \\
& \ldots\left[\exp \left(\beta_{m n} / 2\right)+\exp \left(-\beta_{m n} / 2\right)\right]\left(\sum_{w \in W_{N_{i}}} \varepsilon(w) \exp \left[w\left(\Lambda+\rho+\frac{1}{2} \beta_{p}\right)\right]\right) \tag{4.31}
\end{align*}
$$

where $W_{p i}=\left\{w \in W \mid w\left(\beta_{p}\right)=\beta_{i}\right\}$ and is written above the factor to be deleted.
Note that (4.31) coincides with expression (2.11) of Sharp et al (1985) in the case of $\operatorname{sl}(m / n)$.

We have tried to find a similar expression for representations where more than one atypicality condition is fulfilled, but have not obtained a satisfactory result.

## 5. Example

In this section we shall give an example of an atypical representation of the Lie superalgebra $\mathrm{L}=\mathrm{sl}(3 / 2)$. In particular we shall consider a representation with highest weight $\Lambda$ which is 'low lying' (this means for which (4.7) is not fulfilled). Although the proof of (4.31) only works for highest weight representations satisfying (4.7), the following example indicates that (4.31) may be true for all single-atypical representations.

The simple roots of $\operatorname{sl}(3 / 2)$ are, in the same notation as $\S 2$,

$$
\begin{equation*}
\varepsilon_{1}-\varepsilon_{2} \quad \varepsilon_{2}-\varepsilon_{3} \quad \varepsilon_{3}-\delta_{1} \quad \delta_{1}-\delta_{2} \tag{5.1}
\end{equation*}
$$

For $\lambda \in \mathrm{H}^{*}$, the labels of $\lambda$ are $\left(a_{1}, a_{2}, a_{3}, a_{4}\right)$ with

$$
\begin{equation*}
a_{i}=\lambda\left(H_{i}\right) . \tag{5.2}
\end{equation*}
$$

Then one can express $\lambda$ as follows:

$$
\begin{equation*}
\lambda=\left(a_{1}+a_{2}+a_{3}-a_{4}\right) \varepsilon_{1}+\left(a_{2}+a_{3}-a_{4}\right) \varepsilon_{2}+\left(a_{3}-a_{4}\right) \varepsilon_{3}+a_{4} \delta_{1} . \tag{5.3}
\end{equation*}
$$

Note that the following equation holds in $H^{*}$ :

$$
\begin{equation*}
\left(\varepsilon_{1}+\varepsilon_{2}+\varepsilon_{3}\right)-\left(\delta_{1}+\delta_{2}\right)=0 \tag{5.4}
\end{equation*}
$$

The elements $\beta_{i}$ of $\Delta_{1}^{+}$are determined by

$$
\begin{equation*}
\left(\beta_{1}, \beta_{2}, \ldots, \beta_{6}\right)=\left(\varepsilon_{1}-\delta_{1}, \varepsilon_{1}-\delta_{2}, \varepsilon_{2}-\delta_{1}, \varepsilon_{2}-\delta_{2}, \varepsilon_{3}-\delta_{1}, \varepsilon_{3}-\delta_{2}\right) \tag{5.5}
\end{equation*}
$$

In this order, the six atypicality conditions are

$$
\begin{align*}
& a_{1}+a_{2}+a_{3}+2=0  \tag{1}\\
& a_{1}+a_{2}+a_{3}-a_{4}+1=0  \tag{2}\\
& a_{2}+a_{3}+1=0  \tag{3}\\
& a_{2}+a_{3}-a_{4}=0  \tag{4}\\
& a_{3}=0  \tag{5}\\
& a_{3}-a_{4}-1=0 . \tag{6}
\end{align*}
$$

The even part of $\operatorname{sl}(3 / 2)$ is $\operatorname{sl}(3) \oplus \operatorname{sl}(2) \oplus \mathrm{u}(1)$. In general, $\mathrm{sl}(m / n)$ contains a $\mathrm{u}(1)$ generator $h$ such that

$$
[h, x]=j x \quad \text { for } x \in L_{j}, j \in\{-1,0,+1\}
$$

where $\mathrm{L}_{-1} \oplus \mathrm{~L}_{0} \oplus \mathrm{~L}_{+1}$ is the previously mentioned $\mathbb{Z}$ gradation. An explicit expression for $h$ is given by

$$
\begin{equation*}
h=\frac{n}{n-m} \sum_{i=1}^{m} i H_{i}-\frac{m}{n-m} \sum_{i=1}^{n-1}(n-i) H_{m+i} . \tag{5.7}
\end{equation*}
$$

For $\operatorname{sl}(3 / 2)$ we find

$$
\begin{equation*}
h=-2\left(H_{1}+2 H_{2}+3 H_{3}\right)+3 H_{4} . \tag{5.8}
\end{equation*}
$$

Every dominant integral weight $\lambda$ with labels ( $a_{1}, a_{2}, a_{3}, a_{4}$ ) is also the highest weight of an irreducible $L_{\overline{0}}$ module $V^{0}(\lambda)$. This module $V^{\circ}(\lambda)$ is characterised by the sl(3) Cartan-Dynkin labels $\left(a_{1}, a_{2}\right)$, the $\operatorname{sl}(2)$ label $\left(a_{4}\right)$ and the eigenvalue of the $\mathrm{u}(1)$ generator $k=-2\left(a_{1}+2 a_{2}+3 a_{3}\right)+3 a_{4}$. Hence, $V^{0}(\lambda)$ is determined by the following set of labels:

$$
\begin{equation*}
\left(a_{1}, a_{2} / a_{4} / k\right) \tag{5.9}
\end{equation*}
$$

Let $\Lambda$ be the highest weight with labels ( $1,0,0,1$ ). Clearly $\Lambda$ satisfies exactly one of the atypicality conditions, namely (5.6(5)). Obviously, (4.7) is not fulfilled. However, after some lengthy calculations, we find that (4.31) gives rise to the following $\mathrm{L}_{\overline{0}}$ module structure of $V(\Lambda)$ (in the notation of (5.9)):

$$
\begin{equation*}
(1,0 / 1 / 1) \oplus(0,0 / 2 / 0) \oplus(0,0 / 0 / 0) \oplus(1,1 / 0 / 0) \oplus(0,1 / 1 /-1) \tag{5.10}
\end{equation*}
$$

Thus, (4.31) yields the correct character formula. Indeed, $(1,0,0,1)$ is the adjoint representation of $\mathrm{sl}(3 / 2):(1,0 / 1 / 1)$ is $\mathrm{sl}_{3} \otimes \mathrm{sl}_{2}^{*}$ (with $\mathrm{u}(1)$ value +1$),(0,0 / 2 / 0)$ is the three-dimensional adjoint representation of $\operatorname{sl}(2),(0,0 / 0 / 0)$ represents $u(1),(1,1 / 0 / 0)$ is the eight-dimensional $\operatorname{sl}(3)$ representation and $(0,1 / 1 /-1)$ is the six-dimensional $\mathrm{sl}(3) \oplus \mathrm{sl}^{(2)}$ irrep $\mathrm{sl}_{3}^{*} \otimes \mathrm{sl}_{2}$ (with $\mathrm{u}(1)$-value -1 ).

This example (and some other examples studied by the author) argues that, although the proof of (4.31) is not valid for 'low lying' representations, the final character formula of theorem 2 is true for all single-atypical representations.

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Note added in revision. It was pointed out by the referee that the definition of the mapping $A$ after equation (4.24) is not completely convincing. For example, if $q x_{k}=q^{\prime} x_{k}$ for $q \neq q^{\prime}$, one has to show that this implies $q A x_{k}=q^{\prime} A x_{k}$. It would be natural to seek for an alternative definition of the mapping $A$ where its properties are unambiguously clear but we have been unsuccessful in such an attempt.

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