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Autor(en): Reimann, H.M.<br>Objekttyp: Article<br>Zeitschrift: Commentarii Mathematici Helvetici

Band (Jahr): 57 (1982)

PDF erstellt am: 05.07.2024
Persistenter Link: https://doi.org/10.5169/seals-43894

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## Invariant differential operators in hyperbolic space

H. M. Reimann

## 1. Introduction

The conformal mappings in real higher dimensional space $\mathbf{R}^{n}, n \geq 3$, are the proper Möbiustransformations. The group GM(n) of Möbiustransformations acts on $\hat{\mathbf{R}}^{n}=\mathbf{R}^{n} \cup\{\infty\}$ and there is a subgroup isomorphic to $G M(n-1)$ which stabilizes the unit ball $B$. It is the action of this group $G M(n-1)$ and the induced action on functions on the hyperbolic space $B$ that will be studied.

The differentiation process leads from functions to vectorfields and tensorfields of higher order. There is a natural setting which reduces the analysis of at least the symmetric tensors with vanishing traces to the study of functions on a bigger space $X$. Whereas the hyperbolic space $B$ is isomorphic to $O_{ \pm}(1, n) / O(n)$ this space $X$ is isomorphic to the quotient space $O_{ \pm}(1, n) / O(n-1)$. Geometrically it can be described as the cosphere bundle of the hyperbolic space $B$. The action of the Möbius group $G M(n-1)$ on $X$ essentially is the action of $G M(n-1)$ on the cotangent space of $B$.

The approach described here, whereby certain tensorfields on $B$ are interpreted as functions on $X$, is inspired by a similar construction for the sphere (see Levine [4]). The purpose of that construction was the characterization of invariant systems of singular differential operators on the sphere. In both cases the conformal structure seems to be essential.

The space $C(X)$ of functions on $X$ can be split into a direct sum of subspaces

$$
C(X)=\oplus_{k=0}^{\infty} E^{k}
$$

The functions in $E^{k}$ have an interpretation as tensorfields of symmetric tensors with vanishing traces. Their analysis is in a certain sense complementary to the analysis of differential forms, which in the tensor language is a theory of antisymmetric tensors. Certain striking analogies are apparent. The invariant operators $S_{k}$ and $S_{k}^{*}$ defined in Section 5, Theorem 7, are generalizations of the operators grad and div. They play a role similar to the operators $d$ and $d^{*}$ for differential forms (see Theorems 7 and 9). In particular, $S_{k}$ maps $E^{k}$ into $E^{k+1}$
and $S_{k}^{*}$ maps $E^{k}$ into $E^{k-1}$. It is shown that the operators $S_{1}$ and $S_{2}^{*}$ coincide with certain operators studied by Ahlfors [1] (see Theorem 8).

There exists an invariant differential operator $D_{Z}$ on $X$ which is of first order. As a consequence, the space of solutions of $D_{\mathrm{z}} f=0$ is an algebra. The functions $f \in E^{1}$ which satisfy $D_{\mathrm{z}} f=0$ are exactly those which correspond to vectorfields $v$ in the Lie algebra of the Möbius group (Theorem 6).

The algebra of invariant differential operators on $X$ is not commutative. It is generated by 1 , the first order differential operator $D_{Z}$ and a further differential operator $D_{|\mathbf{Y}|^{2}}$ of second order (Theorems 1 and 2). The operator $D_{|Y|^{2}}$ is basically the Laplace-operator on the sphere $O(n) / O(n-1)$. The spaces $E^{k}$ appear as eigenspaces of $D_{|Y|}$. The Laplace-Casimir operator $\Delta_{X}$ on $X$ preserves the eigenspaces (Theorem 9).

## 2. The Möbius group and its Lie algebra

The Möbius group $G M(n)$ is the transformation group of $\hat{\mathbf{R}}^{n}=\mathbf{R}^{n} \cup\{\infty\}$ which is generated by reflections in the spheres and hyperplanes of $\mathbf{R}^{n}$. The group is isomorphic to $O_{ \pm}(1, n+1)$, the subgroup of $O(1, n+1)$ which preserves the positive cone:

$$
\left\{y \in \mathbf{R}^{n+2}:\langle y, y\rangle=y_{0}^{2}-\sum_{i=1}^{n+1} y_{i}^{2}>0, y_{0}>0\right\}
$$

(see Mostow [5]). The isomorphism is constructed in the following way: The group $O(1, n+1)$ leaves invariant the quadratic form $\langle y, y\rangle=y_{0}^{2}-\sum_{i=1}^{n+1} y_{i}^{2}$ and in particular the cone $\left\{y \in \mathbf{R}^{n+2}:\langle y, y\rangle=0\right\}$. If inhomogeneous coordinates $\eta_{i}=y_{i} / y_{0}$ are introduced, the group becomes a transformation group of the sphere $\Sigma=$ $\left\{\eta \in \mathbf{R}^{n+1}:|\eta|=1\right\}$ and the elements $g$ and $-g$ give rise to the same transformation. Stereographic projection from the point $\varepsilon_{n}=(0, \ldots, 0,1)$ onto the plane $\eta_{n+1}=0$ then leads to the realization of $O_{ \pm}(1, n+1)$ as a transformation group of $\hat{\mathbf{R}}^{n}$. The subgroup of the Möbius group $G M(n)$, which stabilizes the unit ball $B \subset \mathbf{R}^{n}$ is isomorphic to the Möbius-group $G M(n-1)$ of one lower dimension. This group which acts on $B$ will again be denoted by $G M(n-1)$. Observe that under the above isomorphism this is exactly the subgroup $O_{ \pm}(1, n)$ of $O_{ \pm}(1, n+1)$ which stabilizes the lower half space in $\mathbf{R}^{n+2}$. The elements in matrix notation have the special form

$$
g=\left(\begin{array}{cccc} 
& & & 0 \\
& g_{i j} & & \vdots \\
& & & 0 \\
0 & \ldots & 0 & 1
\end{array}\right) \quad i, j=0,1, \ldots, n \quad g_{00}>0
$$

Our main concern is with this group $G=G M(n-1), n \geq 3$, which is the group of conformal and anti-conformal mappings of the unit ball $B \subset \mathbf{R}^{n}$ onto itself. Referring to the isomorphism $G M(n-1) \cong O_{ \pm}(1, n)$ we will speak about the geometric realization of the group, if we consider it as a transformation group of $B$. The algebraic realization then refers to the group as a matrix group.

The unit ball $B$ has the structure of a symmetric space (the hyperbolic space) $B=G / K$ with the invariant metric $d s^{2}=\rho^{2}|d x|^{2}, \rho(x)=\left(1-|x|^{2}\right)^{-1}$. The stabilizer $K$ of the origin is the orthogonal group. We start with an explicit description of the action of $G=G M(n-1)$ on $B \subset \mathbf{R}^{n}$.

The stereographic projection of the sphere $\Sigma=\left\{\eta \in R^{n+1}:|\eta|=1\right\}$ onto the plane $\eta_{n+1}=0$ is given by the formula

$$
x_{i}=\frac{\eta_{i}}{1-\eta_{n+1}} \quad i=1, \ldots, n
$$

and the inverse mapping is

$$
\begin{aligned}
\eta_{i} & =\frac{2 x_{i}}{1+|x|^{2}} \quad i=1, \ldots, n \\
\eta_{n+1} & =\frac{|x|^{2}-1}{|x|^{2}+1}
\end{aligned}
$$

Let $\mathrm{g}=\left(\mathrm{g}_{\mathrm{ij}}\right)$ be an element in $O_{ \pm}(1, n)$ and consider $O_{ \pm}(1, n)$ as the subgroup of $O_{ \pm}(1, n+1)$ which stabilizes the unit vector $e_{n+1}=(0, \ldots, 0,1) \in \mathbf{R}^{n+2}$. The image of the half line $y=t\left(e_{0}-e_{n+1}\right) t>0$ is the half line

$$
t\left(g e_{0}-g e_{n+1}\right)=t\left(g_{00}, \ldots, g_{n 0},-1\right)
$$

which in turn is mapped onto the point

$$
\eta=\frac{1}{g_{00}}\left(g_{10}, \ldots, g_{n 0},-1\right)
$$

Under stereographic projection this point projects onto

$$
\begin{equation*}
x=\frac{1}{1+g_{00}}\left(g_{10}, \ldots, g_{n 0}\right) \in B \tag{2.1}
\end{equation*}
$$

If $g$ is in the subgroup $O(n)$ of $O_{ \pm}(1, n)$, then $g_{00}=1$ and the corresponding point
on the ball $B$ is the center $x=0$. This establishes the isomorphism

$$
B \cong O_{ \pm}(1, n) / O(n)
$$

The group $O_{ \pm}(1, n)$ acts on the quotient space by left translation. The Möbiustransformation corresponding to the element $g \in O_{ \pm}(1, n)$ will be denoted by $\tau_{\mathrm{g}}$. It is a conformal mapping if $\mathrm{g} \in \operatorname{SO}_{ \pm}(1, n)$

$$
S O_{ \pm}(1, n)=\left\{g \in O_{ \pm}(1, n): \operatorname{det} g>0\right\},
$$

otherwise it is an anti-conformal mapping.
Consider the one parameter subgroup

$$
a_{t}=\exp t\left(\begin{array}{lll} 
& & 1  \tag{2.2}\\
& 0 & \\
1 & &
\end{array}\right)=\left(\begin{array}{lllll}
\mathrm{Ch} t & & & & \mathrm{Sh} t \\
& 1 & & & \\
& & \cdot & & \\
& & & 1 & \\
\mathrm{Sh} t & & & & \mathrm{Ch} t
\end{array}\right)
$$

in $O_{ \pm}(1, n)$. The curve $x_{t}=\tau_{a_{1}}(0)$ in the ball $B$ is given by

$$
x_{\mathrm{t}}=\frac{\operatorname{Sh} t}{1+\operatorname{Ch} t} e_{n} \quad e_{n}=(0, \ldots, 0,1) \in \mathbf{R}^{n}
$$

The tangent vector to the curve at the origin is the vector

$$
\left.\frac{d x_{t}}{d t}\right|_{t=0}=-e_{n} / 2
$$

The element $\tau_{\mathrm{g}}, \mathrm{g} \in \mathrm{O}_{ \pm}(1, n)$, maps this curve onto the curve $z_{t}=\tau_{\mathrm{g}} \tau_{a_{t}}(0)$

$$
\left(z_{t}\right)_{i}=\frac{g_{i 0} \operatorname{Ch} t+g_{i n} \operatorname{Sh} t}{1+g_{00} \operatorname{Ch} t+g_{0 n} \operatorname{Sh} t} \quad i=1, \ldots, n
$$

whose tangent vector at $\tau_{\mathrm{g}}(0)$ is given by

$$
\left.\frac{d z_{t}}{d t}\right|_{t=0}=\frac{-g_{0 n}}{\left(1+g_{00}\right)^{2}}\left(g_{10}, \ldots, g_{n 0}\right)+\frac{1}{1+g_{00}}\left(g_{i n}, \ldots, g_{n n}\right)
$$

The tangent vector $\varepsilon_{n}=(0, \ldots, 0,1)$ at the origin is therefore mapped onto the
tangent vector $\xi$ at $x=\left(1 /\left(1+g_{00}\right)\right)\left(g_{10}, \ldots, g_{n 0}\right)$ with coordinates

$$
\begin{equation*}
\xi_{i}=\frac{2 g_{0 n} g_{i 0}}{\left(1+g_{00}\right)^{2}}-\frac{2 g_{i n}}{1+g_{00}} \quad i=1, \ldots, n \tag{2.3}
\end{equation*}
$$

The invariance of the quadratic form $\langle y, y\rangle$ implies

$$
\begin{align*}
1 & =g_{00}^{2}-\sum_{i=1}^{n} g_{i 0}^{2} \\
-1 & =g_{0 k}^{2}-\sum_{i=1}^{n} g_{i k}^{2} \quad k=1, \ldots, n \\
0 & =g_{00} g_{0 k}-\sum_{i=1}^{n} g_{i k} g_{i 0} \tag{2.4}
\end{align*}
$$

and it follows that

$$
\begin{align*}
& |x|^{2}=\left(1+g_{00}\right)^{-2} \sum_{i=1}^{n} g_{i 0}^{2}=\frac{g_{00}-1}{g_{00}+1} \\
& \frac{2}{1+g_{00}}=1-|x|^{2} \tag{2.5}
\end{align*}
$$

if $x=\tau_{g}(0)$. The length $|\xi|=\left(\sum_{i=1}^{n} \xi_{i}^{2}\right)^{1 / 2}$ of the tangent vector $\xi$ can now easily be calculated to be $1-|x|^{2}$

$$
\begin{align*}
\frac{1}{\left(1-|x|^{2}\right)^{2}}|\xi|^{2} & =4^{-1}\left(1+g_{00}\right)^{2}|\xi|^{2} \\
& =\left(1+g_{00}\right)^{-2} g_{0 n}^{2}\left(g_{00}^{2}-1\right)-\left(1+g_{00}\right)^{-1} 2 g_{0 n}^{2} g_{00}+g_{0 n}^{2}+1 \\
& =\left(1+g_{00}\right)^{-1} g_{0 n}^{2}\left(g_{00}-1-2 g_{00}\right)+g_{0 n}^{2}+1=1 \\
|\xi| & =1-|x|^{2} \tag{2.6}
\end{align*}
$$

This proves the invariance of the metric

$$
d s^{2}=\rho^{2}|d x|^{2} \quad \rho=\left(1-|x|^{2}\right)^{-1}
$$

and the conformality (or anti-conformality) of the transformations $\tau_{\mathrm{g}}$.
Next we define the subgroup $M$ of the Möbius group $G=G M(n-1)$ as the stabilizer of both the origin and the tangent vector $\varepsilon_{n}$ at the origin in $B . M$ is a
subgroup of $K$. In the algebraic picture this is the orthogonal group

$$
\left.O(n-1)=\left\{g \in O_{ \pm}(1, n): g=\left(\begin{array}{l}
1  \tag{2.7}\\
\\
\\
\\
\\
\end{array}\right)_{1}\right)\right\} \cong M
$$

The cosets are parametrized by the geometric parameters $x=\tau_{\mathrm{g}}(0)$ and $\xi=$ $d \tau_{\mathrm{g}}(0) \varepsilon_{n}$. We call the pair $(x, \xi)$ the coordinates for the coset $g O(n-1)$. The equations (2.1) and (2.3) express these coordinates by the matrix elements $\mathrm{g}_{\mathrm{ij}}$ of g . Geometrically, the quotient space $G / M$ can be realized as the cosphere bundle $X$ of $B$. Since $|\xi|=1-|x|^{2}$, the group $G M(n-1)$ acts on

$$
\begin{equation*}
X=\left\{(x, \xi) \in B \times \mathbf{R}^{n}:|\xi|=1-|x|^{2}\right\} \tag{2.8}
\end{equation*}
$$

and the action is seen to be transitive. It can be described by the formula

$$
\begin{equation*}
(x, \xi) \rightarrow\left(\tau_{\mathrm{g}} x, d \tau_{\mathrm{g}}(x) \xi\right) \tag{2.9}
\end{equation*}
$$

where $d \tau_{\mathrm{g}}(x)$ is the cotangent mapping which maps the cotangent space at $x$ onto the cotangent space at $z=\tau_{\mathrm{g}} x$.

We now turn to a description of the Lie algebra $\mathfrak{g}$ of $O_{ \pm}(1, n)$. Let $E_{i j} \in$ $G L(n+1)$ denote the matrix with element 1 at the place $i, j$ and zero otherwise. A basis for the Lie algebra of $O_{ \pm}(1, n)$ is given by the matrices

$$
\begin{equation*}
X_{0 i}=E_{0 j}+E_{j 0} \quad j=1, \ldots, n \tag{2.10}
\end{equation*}
$$

and

$$
X_{i j}=E_{i j}-E_{j i} \quad 1 \leq i<j \leq n
$$

We set

$$
\begin{align*}
& X_{i}=X_{0 i} \\
& Z=X_{0 n}  \tag{2.11}\\
& Y_{i}=X_{i n} \quad i=1, \ldots, n-1 \\
&
\end{align*}
$$

The stabilizer $O(n)$ of $e_{0} \in \mathbf{R}^{n+1}$ is a maximal compact subgroup in $O_{ \pm}(1, n)$ and $O_{ \pm}(1, n) / O(n) \cong B$ is a symmetric space of rank one. In the Cartan decomposition $\mathfrak{g}=\mathfrak{f}+\mathfrak{p}$ the subalgebra $\mathfrak{f}$ has the vectorspace basis $\left\{X_{i j}: 1 \leq i<j \leq n\right\}$ and $\mathfrak{p}$ is the
linear subspace with basis $\left\{X_{0 j}: j=1, \ldots, n\right\}$. The commutator relations

$$
\begin{align*}
& {[\mathfrak{p}, \mathfrak{p}] \subset \mathfrak{f}}  \tag{2.12}\\
& {[\mathfrak{p}, \mathfrak{f}] \subset \mathfrak{p}} \tag{2.13}
\end{align*}
$$

hold. A maximal abelian subalgebra in $\mathfrak{p}$ is given by $\mathfrak{a}=\mathbf{R} Z$, it is one dimensional. If the corresponding subgroup is denoted by $A$, then the subgroup $O(n-1) \cong M$ defined above (2.7) is the centralizer of $A$ in $O(n) \cong K$. Its Lie algebra $m$ has the basis $\left\{X_{i j}: 1 \leq i<j \leq n-1\right\}$

The commutator relations are as follows

$$
\begin{array}{lll}
{[\mathrm{m}, Z]=0} & \\
{\left[X_{i}, Z\right]=Y_{i}} & {\left[X_{i}, X_{i j}\right]=X_{j}} & 1 \leq i<j \leq n-1 \\
{\left[Y_{i}, Z\right]=X_{i}} & {\left[Y_{i}, X_{i j}\right]=Y_{j}} & 1 \leq i<j \leq n-1  \tag{2.14}\\
{\left[X_{i}, X_{j}\right]=X_{i j}} & {\left[Y_{i}, Y_{j}\right]=-X_{i j}} & 1 \leq i<j \leq n-1 \\
{\left[X_{i}, Y_{j}\right]=\delta_{i j} Z} & & i, j=1, \ldots, n-1
\end{array}
$$

In particular it should be noted that if $\mathfrak{q}$ is the linear subspace with basis $\left\{X_{1}, \ldots, X_{n-1}, Z, Y_{1}, \ldots, Y_{n-1}\right\}$ then

$$
\begin{equation*}
[\mathfrak{q}, \mathrm{m}] \subset \mathfrak{q} \tag{2.15}
\end{equation*}
$$

which shows that $G / M$ is a reductive coset space (see Section 3 ). $\left\{X_{i}-Y_{i}: i=\right.$ $1, \ldots, n-1\}$ is a basis of the $\alpha$-root space $n$ of the pair $(g, a)$ :

$$
\left[t Z, X_{i}-Y_{i}\right]=t\left(X_{i}-Y_{i}\right), \quad \alpha(t Z)=t
$$

whereas $\overline{\mathfrak{n}}$ is given by $\left\{X_{i}+Y_{i}: i=1, \ldots, n-1\right\}$.
The Weyl group $W=O^{\prime}(n-1) / O(n-1)$ where $O^{\prime}(n-1)$ and $O(n-1)$ are the normalizer and centralizer of $A$ in $O(n)=K$ consists of two elements only. They are represented by the identity and the matrix

$$
w=\left(\begin{array}{llll}
1 & & &  \tag{2.16}\\
& \cdot & & \\
& & \cdot & \\
& & & -1
\end{array}\right)
$$

The mapping $\omega$, which maps the cosets $g O(n-1)$ onto the cosets $g w O(n-1)$ can geometrically be described by the formula

$$
\begin{equation*}
\omega(x, \xi)=(x,-\xi) \tag{2.17}
\end{equation*}
$$

This mapping is not a Möbiustransformation on $X$.
Geometrically, the Lie algebra of $G=G M(n-1)$ is given by the vectorfields on $B$ which generate the one parameter subgroups $\tau_{\mathrm{g}_{1}}$ of $G$. The vectorfields are determined by the equation

$$
v(x)=\left.\frac{d}{d t} \tau_{8_{t}}(x)\right|_{t=0}
$$

Conversely, the one parameter subgroup $\tau_{\mathrm{g}_{\mathrm{t}}}$ is obtained from the vectorfield $v$ by solving the differential equation

$$
\frac{d z}{d t}=v(z)
$$

with initial condition $z(0)=x$. The one parameter subgroup is then given by $\tau_{\mathrm{g}_{\mathrm{t}}}(x)=z(t)$.

In a first step the vectorfields on $\mathbf{R}^{n}$ are determined, which are the infinitesimal generators of the one parameter subgroups of the group $G M(n)$ acting on $\hat{\mathbf{R}}^{n}$. The vectorfields in the Lie algebra of $G M(n-1)$ are then singled out by the condition

$$
\begin{equation*}
(v(x), x)=0 \quad \text { for } \quad|x|=1 \tag{2.18}
\end{equation*}
$$

The vectorfield $v$ has to be tangent to the boundary of $B \subset \mathbf{R}^{n}$. The vectorfields in the Lie algebra of $G M(n)$ are

$$
\begin{equation*}
v(x)=a+B x+\lambda x+c|x|^{2}-2 x(c, x) \tag{2.19}
\end{equation*}
$$

with $a, c$ constant vectors in $\mathbf{R}^{n}, B$ a constant matrix with $B^{\prime}=-B$ and $\lambda \in \mathbf{R}$. The vectorfields $B x$ account for the rotations in $\mathbf{R}^{n}$ (the subgroup $M$ with respect to $G M(n)$ ), the constant vectors $a$ for the translations (the subgroup $N$ ) and $\lambda x$ for the dilations (the subgroup $A$ ). The remaining vectorfields $c|x|^{2}-2 x(c, x)$ generate the one parameter subgroups $\tau_{\mathrm{g}}$ conjugate to the translations (the subgroup $\bar{N}$ ):

$$
s \circ \tau_{\mathrm{g}_{\mathrm{t}} \circ} s(x)=x+c t
$$

where $s$ is the reflection in the unit sphere. The vectorfields in the Lie algebra of $G M(n-1)$ can easily be singled out by condition (2.18). The restrictions are $\lambda=0$ and

$$
(a, x)-(c, x)=0 \quad \text { for } \quad|x|=1
$$

The Lie algebra of $G M(n-1)$ is therefore described by the vectorfields

$$
\begin{equation*}
v(x)=B x+c\left(1+|x|^{2}\right)-2 x(c, x) \tag{2.20}
\end{equation*}
$$

The vectorfields $B x$ now correspond to the subalgebra $\mathfrak{f} \subset \mathfrak{g}$ and the remaining vectorfields to the complementary subspace $\mathfrak{p} \subset \mathfrak{g}$.

## 3. Invariant differential operators

The group $O_{ \pm}(1, n)$ is not connected. The connected component of the identity is the subgroup $S O_{ \pm}(1, n)$. The spaces $O_{ \pm}(1, n) / O(n-1)$ and $S O_{ \pm}(1, n) / S O(n-1)$ are isomorphic coset spaces with in the first instance the group $O_{ \pm}(1, n)$, in the second the group $S O_{ \pm}(1, n)$ acting by left translations.

DEFINITION (Nomizu [6]). Let $G$ be a connected Lie group with Lie algebra $g$ and denote the adjoint representation of $G$ on $g$ by Ad (g). Assume that $M$ is a closed subgroup with Lie algebra $m$. The coset space $G / M$ is reductive, if there exists a subspace $q$ of $g$, complementary to $m$, such that $\operatorname{Ad}(m) q \subset q$ for all $m \in M$.

Upon taking $G=S O_{ \pm}(1, n)$ and $M=S O(n-1)$ one finds that the subspace $q$ with basis $\left\{X_{1}, \ldots, X_{n-1}, Z, Y_{1}, \ldots, Y_{n-1}\right\}$ is complementary to the Lie algebra $m$ of $M$ and that furthermore $[m, q] \subset q$ (see (2.11) and (2.15)). Since $M$ is connected, this implies $\operatorname{Ad}(m) q \subset q$ for all $m \in M$. The coset space $X=$ $S O_{ \pm}(1, n) / S O(n-1)$ (with $S O_{ \pm}(1, n)$ acting on it by left translation) is therefore reductive.

By definition, the differential operator $D$ on $G / M$ is invariant (with respect to left translations $\tau^{8} f(x)=f\left(\tau_{g-1} x\right)$ ) if $D \tau^{8} f=\tau^{8} D f$ for all $f \in C_{c}(G / M)$ and for all $g \in G$. The algebra of invariant differential operators is denoted by $\underset{\sim}{D}(G / M)$. It can be determined on the base of a theorem of Helgason [3]. For this purpose let $I(\mathfrak{q})$ denote the polynomials in the symmetric algebra $S(\mathbf{q})$ over $\mathfrak{q}$, which are invariant under Ad $(m)$ for all $m \in M$. The polynomials in $S(q)$ are polynomials in the variables $Z_{1}, \ldots, Z_{k}$ where $\left\{Z_{1}, \ldots, Z_{k}\right\}$ is a basis in $M$.

The symmetrization mapping $\lambda$ associates with every polynomial $Q \in S(q)$ a differential operator on the group $G$. Symmetrization is a linear mapping, which maps the elements $Y_{1} Y_{2} \cdots Y_{p} \in S(q)$ (where the $Y_{j}$ are elements in the subspace $q$ of $\mathfrak{g}, j=1, \ldots, p$ ) onto the differential operator

$$
\lambda\left(Y_{1} Y_{2} \cdots Y_{p}\right)=\frac{1}{p!} \sum_{\sigma} Y_{\sigma(1)} \cdot Y_{\sigma(2)} \cdots \cdot Y_{\sigma(p)}
$$

In this sum $\sigma$ runs over the symmetric group on $p$ letters. In particular, $\lambda(Y)$ is the differential operator defined by the Lie algebra element $Y \in \mathbf{g}$

THEOREM (Helgason). Let $G / M$ be a reductive coset space, $\mathfrak{g}=\mathfrak{m}+\mathfrak{q}$, $\operatorname{Ad}(m) \mathfrak{q} \subset \mathfrak{q}$ for all $m \in M$. Then there exists a linear bijection of $I(\mathfrak{q})$ onto $D(G / M)$. It associates to the polynomial $Q\left(Z_{1}, \ldots, Z_{k}\right) \in I(q)$ the differential operator $D_{Q}$ which can be determined by one of the equivalent methods:
(1) $D_{Q} f(x)=\left.Q\left(\frac{\partial}{\partial t_{1}}, \ldots, \frac{\partial}{\partial t_{k}}\right) f \circ \pi\left(g \exp \sum_{i=1}^{k} t_{i} Z_{i}\right)\right|_{t=0}$
where $\pi$ is the canonical projection of $G$ onto $G / M, \pi(\mathrm{~g})=x$.
(2) $\lambda(Q)(f \circ \pi)=D_{\mathrm{Q}} f \circ \pi$

This formula defines $D_{\mathrm{Q}} f$, since $\lambda(Q)(f \circ \pi)$ is constant on each coset $g M$ if $f \in C_{c}^{\infty}(G / M)$.

THEOREM 1. Let $G=S O_{ \pm}(1, n), M=S O(n-1)$ and $\mathfrak{g}=\mathfrak{m}+\mathfrak{q}$ with the specified basis $\left\{X_{1}, \ldots, X_{n-1}, Z, Y_{1}, \ldots, Y_{n-1}\right\}$ for $q$ (see Section 2 ). Then the algebra $I(\mathfrak{q})$ of $\operatorname{Ad}(M)$ invariant polynomials is generated by the polynomials

$$
1, \quad Z, \quad|X|^{2}=\sum_{i=1}^{n-1} X_{i}^{2}, \quad(X, Y)=\sum_{i=1}^{n-1} X_{i} Y_{i}, \quad|Y|^{2}=\sum_{i=1}^{n-1} Y_{i}^{2} .
$$

We calculate the action of $\operatorname{Ad}(m)$. If $X \in \mathfrak{m}, Y \in \mathfrak{q}$ then
$\operatorname{Ad}(\exp t X) Y=e^{t a d} X=\sum_{n=0}^{\infty} \frac{t^{n}}{n!}(\operatorname{ad} X)^{n} Y$

Set $X=X_{i j} \in \mathfrak{m}$ and $Y=Z_{i}$, which stands for $X_{i}$ or $Y_{i} \in \mathfrak{q}$. Then

$$
\begin{aligned}
& \left(\operatorname{ad} X_{i j}\right) Z_{i}=\left[X_{i j}, Z_{i}\right]=-Z_{i} \\
& \begin{aligned}
\left(\operatorname{ad} X_{i j}\right) Z_{i}=Z_{i}, \quad\left(\operatorname{ad} X_{i j}\right) Z_{k}=0 \quad k \neq i, j
\end{aligned} \\
& \begin{aligned}
& \operatorname{Ad}\left(\exp t X_{i j}\right) Z_{i}= \sum_{n=0}^{\infty} \frac{(-1)^{n} t^{2 n}}{(2 n)!} Z_{i}-\sum_{n=0}^{\infty} \frac{(-1)^{n} t^{2 n+1}}{(2 n+1)!} Z_{i} \\
& \quad=Z_{i} \cos t-Z_{j} \sin t
\end{aligned}
\end{aligned}
$$

$\operatorname{Ad}\left(\exp t X_{i j}\right) Z_{j}=Z_{i} \sin t+Z_{i} \cos t$

It follows that
$\operatorname{Ad}(m) X_{k}=\sum_{h=1}^{n-1} m_{h k} X$
for $m=\exp t X_{i j} \in S O(n-1) \subset G L(n-1)$ with $m=\left(m_{n k}\right)$. This equation therefore holds for all $m \in S O(n-1)$. Furthermore, if $X=\sum_{k=1}^{n-1} x_{k} X_{k}$, then $\operatorname{Ad}(m) X=$ $\sum_{h=1}^{n-1} x_{h}^{\prime} X_{h}$ with $x^{\prime}=m x$. Similarly, if $Y=\sum_{k=1}^{n-1} y_{k} Y_{k}$ then $\operatorname{Ad}(m) Y=\sum_{h=1}^{n-1} y_{h}^{\prime} Y_{h}$ with $y^{\prime}=m y$. Finally, since $\operatorname{Ad}(m) z Z=z Z(z \in \mathbf{R})$, the action of $\operatorname{Ad}(m)$ on the polynomials $P(x, y, z)$ in the variables $x, y \in \mathbf{R}^{n-1}, z \in \mathbf{R}$ is given by

$$
\operatorname{Ad}(m) P(x, y, z)=P(m x, m y, z)
$$

Assume now that the polynomial $Q$ is invariant under the action of $\operatorname{Ad}(M)$. It can then be written as a finite sum

$$
Q(x, y, z)=\sum_{k} z^{k} Q_{k}(x, y)
$$

with invariant polynomials $Q_{k}(x, y)$. It is well known (see e.g. Weyl [7] p. 31 ff .) that the invariant polynomials in the variables $x, y$ under the action

$$
(x, y) \rightarrow(m x, m y) \quad m \in S O(n-1)
$$

are generated by the polynomials $1,|x|^{2}=\sum_{i=1}^{n-1} x_{i}^{2},(x, y)=\sum_{i=1}^{n-1} x_{i} y_{i}$ and $|y|^{2}=$ $\sum_{i=1}^{n-1} y_{i}^{2}$. This proves the theorem.

The invariant operators $1, D_{\mathbf{Z}}, D_{|X|^{2}}, D_{(X, Y)}$ and $D_{|Y|^{2}}$ generate the whole
algebra $\underset{\sim}{D}(G / M)$. This follows from the fact that

$$
D_{P_{1} P_{2}}=D_{P_{1}} \cdot D_{P_{2}}+D
$$

where the order of the invariant differential operator $D$ is less than the sum of the degrees of the polynomials $P_{1}$ and $P_{2}$ (see Helgason [3] p. 269). In the present situation there is however more that can be said:

THEOREM 2. The differential operators satisfy the following commutator relations:

$$
\begin{gather*}
{\left[D_{\mathrm{Z}}, D_{|X|^{2}}\right]=-2 D_{(X, Y)}}  \tag{3.3}\\
{\left[D_{Z}, D_{|Y|^{2}}\right]=-2 D_{(X, Y)}}  \tag{3.4}\\
{\left[D_{Z}, D_{(X, Y)}\right]=-D_{|X|^{2}}-D_{|Y|^{2}}} \tag{3.5}
\end{gather*}
$$

Consequently, $\underline{D}(G / M)$ is generated by $1, D_{Z}$ and $D_{|Y|^{2}}$ (or by $1, D_{Z}$ and $D_{|X|^{2}}$ ).
The proof relies on the symmetrization mapping $\lambda$. The differential operator $D_{Z|Y|^{2}}$ is obtained from the differential operator on $G$ which is given by

$$
\lambda\left(Z|Y|^{2}\right)=\frac{1}{3!} \sum_{i=1}^{n-1} 2\left(Y_{i} \cdot Y_{i} \cdot Z+Y_{i} \cdot Z \cdot Y_{i}+Z \cdot Y_{i} \cdot Y_{i}\right)
$$

The commutator relations for the Lie algebra (2.14) then imply

$$
\begin{aligned}
& \lambda\left(Z|Y|^{2}\right)=\sum_{i=1}^{n-1} Y_{i} \cdot Y_{i} \cdot Z-\frac{1}{2} \sum_{i=1}^{n-1}\left(X_{i} \cdot Y_{i}+Y_{i} \cdot X_{i}\right)+\frac{n-1}{6} Z \\
& \lambda\left(Z|Y|^{2}\right)=\sum_{i=1}^{n-1} Z \cdot Y_{i} \cdot Y_{i}+\frac{1}{2} \sum_{i=1}^{n-1}\left(X_{i} \cdot Y_{i}+Y_{i} \cdot X_{i}\right)+\frac{n-1}{6} Z
\end{aligned}
$$

It follows that

$$
D_{|Y|^{2}} D_{Z}-D_{(X, Y)}+\frac{n-1}{6} D_{Z}=D_{Z} D_{|Y|^{2}}+D_{(X, Y)}+\frac{n-1}{6} D_{Z}
$$

which proves the first equality. The second is proved in the same way and the
third is a consequence of the following equations:

$$
\begin{aligned}
& \begin{aligned}
\lambda\left(\sum_{i=1}^{n-1} X_{i} Y_{i} Z\right)= & \frac{1}{3} \sum_{i=1}^{n-1}\left(X_{i} \cdot Y_{i} \cdot Z+Y_{i} \cdot X_{i} \cdot Z+X_{i} \cdot Z \cdot Y_{i}+Y_{i} \cdot Z \cdot X_{i}\right. \\
& \left.+Z \cdot X_{i} \cdot Y_{i}+Z \cdot Y_{i} \cdot X_{i}\right) \\
= & \frac{1}{2} \sum_{i=1}^{n-1}\left(X_{i} \cdot Y_{i}+Y_{i} \cdot X_{i}\right) \cdot Z-\frac{1}{2} \sum_{i=1}^{n-1}\left(X_{i} \cdot X_{i}+Y_{i} \cdot Y_{i}\right) \\
= & \frac{1}{2} \sum_{i=1}^{n-1} Z \cdot\left(X_{i} \cdot Y_{i}+Y_{i} \cdot X_{i}\right)+\frac{1}{2} \sum_{i=1}^{n-1}\left(X_{i} \cdot X_{i}+Y_{i} \cdot Y_{i}\right)
\end{aligned} \\
& D_{(X, Y)} D_{Z}-\frac{1}{2} D_{|X|^{2}}-\frac{1}{2} D_{|Y|^{2}}=D_{Z} D_{(X, Y)}+\frac{1}{2} D_{|X|^{2}}+\frac{1}{2} D_{|Y|^{2}}
\end{aligned}
$$

The Killing form on the Lie algebra of $S O(1, n)$ is given by

$$
\begin{aligned}
& B(X, X)=2(n-1)\left\{\sum_{i=1}^{n} x_{i}^{2}-\sum_{1 \leq i<j \leq n} x_{i j}^{2}\right\} \\
& X=\sum_{i=1}^{n} x_{i} X_{i}+\sum_{1 \leq i<j \leq n} x_{i j} X_{i j}
\end{aligned}
$$

(see the definitions (2.10) and (2.11) in section 2). The Killing form is invariant under $\operatorname{Ad}(g)$ for all $g \in S O(1, n)$ and in particular for $g \in S O(n)$ or $S O(n-1)$. The Casimir operator restricted to $B \cong S O_{ \pm}(1, n) / S O(n)$ is

$$
\begin{equation*}
\Delta_{K}=D_{|X|^{2}}+D_{Z^{2}} \tag{3.6}
\end{equation*}
$$

and restricted to $X \cong S O_{ \pm}(1, n) / S O(n-1)$ it is

$$
\begin{equation*}
\Delta_{M}=D_{|X|^{2}}+D_{Z^{2}}-D_{|Y|^{2}} \tag{3.7}
\end{equation*}
$$

It follows that the operators $\Delta_{K}$ and $\Delta_{M}$, considered as operators in $\underset{\sim}{D}(G / M)$ commute. In fact, $\Delta_{M}$ commutes with every differential operator in $\underset{\sim}{D}(G / M)$.

In the next section it will be shown that the operators in $\underset{(G / M)}{ }$ are invariant under the whole group $O_{ \pm}(1, n)$ and not only under the subgroup $S O_{ \pm}(1, n)$.

## 4. The calculations for some operators

In this section the geometric versions of the operators $D_{Z}, D_{|Y|^{2}}$ and $D_{(X, Y)}$ will be calculated. This means that the operators will be expressed as differential
operators in the variables $(x, \xi)$. Recall that

$$
\begin{equation*}
x_{i}=\left(1+g_{00}\right)^{-1} g_{i 0} \tag{2.1}
\end{equation*}
$$

and

$$
\begin{align*}
\xi_{i} & =2 g_{0 n} g_{i 0}\left(1+g_{00}\right)^{-2}-2 g_{i n}\left(1+g_{00}\right)^{-1} \\
& =2\left(g_{0 n} x_{i}-g_{i n}\right)\left(1+g_{00}\right)^{-1} \tag{2.3}
\end{align*}
$$

$i=1, \ldots, n$ are the coordinates for the coset $g O(n-1)$. The matrices $\left(g_{i j}\right)$ representing $g$ satisfy the relations (2.4) and in particular

$$
\begin{equation*}
2\left(1+g_{00}\right)^{-1}=1-|x|^{2}=|\xi|^{2} \tag{2.5}
\end{equation*}
$$

and

$$
\begin{align*}
(x \mid \xi) & =\sum_{i=1}^{n} x_{i} \xi_{i}=\frac{2}{1+g_{00}}\left(g_{0 n}|x|^{2}-\left(1+g_{00}\right)^{-1} \sum_{i=1}^{n} g_{i 0} g_{i n}\right) \\
& =2 g_{0 n}\left(1+g_{00}\right)^{-1}\left(|x|^{2}-g_{00}\left(1+g_{00}\right)^{-1}\right) \\
& =-\frac{1}{2} g_{0 n}\left(1-|x|^{2}\right)^{2}=-2 g_{0 n}\left(1+g_{00}\right)^{-2} \tag{4.1}
\end{align*}
$$

Let $a_{t}=\exp t Z$ denote the one parameter subgroup of $O_{ \pm}(1, n)$ defined by $Z$. In order to calculate $D_{z} f$ at the point $(x, \xi)$ (coordinates of the coset $g O(n-1)$ ), the definition of Lie derivatives is used:

$$
\begin{equation*}
D_{\mathrm{Z}} f(x, \xi)=\left.\frac{d}{d t} f\left(x_{t}, \xi_{\mathrm{t}}\right)\right|_{\mathrm{t}=0} \tag{4.2}
\end{equation*}
$$

where $\left(x_{t}, \xi_{t}\right)$ are the coordinates of the coset $g a_{t} O(n-1)$ :

$$
\begin{align*}
\left(x_{t}\right)_{i}= & \left(g_{i 0} \mathrm{Ch} t+\mathrm{g}_{i n} \operatorname{Sh} t\right)\left(1+\mathrm{g}_{00} \mathrm{Ch} t+\mathrm{g}_{0 n} \operatorname{Sh} t\right)^{-1}  \tag{4.3}\\
\left(\xi_{t}\right)_{i}= & 2\left(g_{00} \operatorname{Sh} t+\mathrm{g}_{0 n} \mathrm{Ch} t\right)\left(g_{i 0} \mathrm{Ch} t+\mathrm{g}_{\text {in }} \operatorname{Sh} t\right)\left(1+\mathrm{g}_{00} \mathrm{Ch} t+\mathrm{g}_{0 n} \operatorname{Sh} t\right)^{-2} \\
& -2\left(\mathrm{~g}_{i 0} \operatorname{Sh} t+\mathrm{g}_{\mathrm{in}} \mathrm{Ch} t\right)\left(1+\mathrm{g}_{00} \mathrm{Ch} t+\mathrm{g}_{0 n} \operatorname{Sh} t\right)^{-1} \tag{4.4}
\end{align*}
$$

It follows that

$$
\begin{equation*}
\left.\frac{d\left(x_{t}\right)_{i}}{d t}\right|_{t=0}=g_{i n}\left(1+g_{00}\right)^{-1}-g_{i 0} g_{0 n}\left(1+g_{00}\right)^{-2}=-\frac{1}{2} \xi_{i} \tag{4.5}
\end{equation*}
$$

and after some calculations

$$
\begin{align*}
\left.\frac{d\left(\xi_{t}\right)_{i}}{d t}\right|_{t=0} & =-2 g_{0 n}\left(1+g_{00}\right)^{-1} \xi_{i}-2\left(1+g_{00}\right)^{-1} x_{i} \\
& =\left(1-|x|^{2}\right)^{-1}\left(2(x \mid \xi) \xi_{i}-|\xi|^{2} x_{i}\right) \tag{4.6}
\end{align*}
$$

The operator $D_{\mathcal{Z}}$ can be expressed by the formula

$$
\begin{equation*}
D_{\mathrm{Z}} f(x, \xi)=\left.\sum_{i=1}^{n} f_{x_{i}} \frac{\left(d x_{t}\right)_{i}}{d t}\right|_{t=0}+\left.\sum_{i=1}^{n} f_{\xi_{i}} \frac{\left(d \xi_{t}\right)_{i}}{d t}\right|_{t=0} \tag{4.7}
\end{equation*}
$$

## THEOREM 3.

$$
\begin{equation*}
D_{Z} f(x, \xi)=-\frac{1}{2} \sum_{i=1}^{n} f_{x_{i}} \xi_{i}+\left(1-|x|^{2}\right)^{-1} \sum_{i=1}^{n} f_{\xi_{i}}\left(2(x \mid \xi) \xi_{i}-|\xi|^{2} x_{i}\right) \tag{4.8}
\end{equation*}
$$

This operator is invariant under the group GM(n-1) of Möbiustransformations on $X$. Under the mapping $\omega(x, \xi)=(x,-\xi)$ it transforms into the operator $-D_{Z}$.

The group $G M(n-1)$ has two components. By construction, the operator $D_{z}$ is invariant under proper Möbius transformations. It suffices to prove its invariance for a single transformation $\tau_{g}, g \notin S O_{ \pm}(1, n)$. Such a transformation is

$$
\begin{array}{ll}
y_{1}=-x_{1} & \eta_{1}=-\xi_{1} \\
y_{k}=x_{k} & \eta_{k}=\xi_{k} \quad k=2, \ldots, n \tag{4.9}
\end{array}
$$

The transformed operator is

$$
\begin{aligned}
D \frac{\mathrm{z}}{\mathrm{Z}} f(y, \eta)= & \frac{1}{2} \sum_{i, j=1}^{n}\left(f_{y_{i}} \frac{\partial y_{j}}{\partial x_{i}}+f_{\eta_{i}} \frac{\partial \eta_{j}}{\partial x_{i}}\right) \xi_{i} \\
& +\left(1-|x|^{2}\right)^{-1} \sum_{i, j=1}^{n}\left(f_{y_{1}} \frac{\partial y_{j}}{\partial \xi_{i}}+f_{\eta_{i}} \frac{\partial \eta_{j}}{\partial \xi_{i}}\right)\left(2(x \mid \xi) \xi_{i}-|\xi|^{2} x_{i}\right) \\
= & -\frac{1}{2} \sum_{i=1}^{n} f_{y_{i}} \eta_{i}+\left(1-|y|^{2}\right)^{-1} \sum_{i=1}^{n} f_{\eta_{i}}\left(2(y \mid \eta) \eta_{i}-|\eta|^{2} y_{i}\right)
\end{aligned}
$$

It coincides with $D_{z}$. The same calculation shows that the mapping $\omega$ (see (2.17)) transforms $D_{Z}$ into the operator $-D_{\mathbf{Z}}$.

A remark about the derivatives $f_{x_{i}}, f_{\xi_{1}} i=1, \ldots, n$ is appropriate. The function $f$ is defined on

$$
X=\left\{(x, \xi) \in \mathbf{R}^{2 n}:|\xi|^{2}=1-|x|^{2}\right\}
$$

In order that the derivatives with respect to $x$ and $\xi$ have some meaning, the domain of definition for $f$ first has to be extended into a neighbourhood of $X$ in $\mathbf{R}^{2 n}$. The resulting operator $D_{Z}$ is however known to depend only on the values of $f$ on $X$. It is independent of the particular extension of $f$.

The calculation of the remaining operators $D_{|Y|^{2}}$ and $D_{(X, Y)}$ is based on the theorem of Helgason (section 3). For fixed $g$ with coordinates $(x, \xi)$ and for a given function $f \in C_{c}(G / M)$ consider the function

$$
\begin{equation*}
\tilde{f}(s, t)=f \circ \pi\left(g \exp \sum_{i=1}^{n-1}\left(s_{i} X_{i}+t_{i} Y_{i}\right)\right) \tag{4.10}
\end{equation*}
$$

$\pi$ is the canonical projection and $(x(s, t), \xi(s, t))$ are the coordinates of $\pi\left(g \exp \sum_{i=1}^{n-1}\left(s_{i} X_{i}+t_{i} Y_{i}\right)\right)$. Take as an example the operator $D_{(X, Y)}$. We then have

$$
\begin{equation*}
D_{(X, Y)} f(x, \xi)=\left.\sum_{i=1}^{n-1} \frac{\partial^{2}}{\partial s_{i} \partial t_{i}} \tilde{f}(s, t)\right|_{s=t=0} \tag{4.11}
\end{equation*}
$$

The chain rule for the second derivative $\tilde{f}_{s_{i}, t}$ gives

$$
\begin{align*}
\tilde{f}_{s_{i},}= & \sum_{m, l=1}^{n} f_{x_{l} x_{m}} \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial x_{m}}{\partial t_{j}}+\sum_{m, l=1}^{n} f_{x_{l} \xi_{m}} \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial \xi_{m}}{\partial t_{j}} \\
& +\sum_{m, l=1}^{n} f_{\xi_{l}} x_{m} \frac{\partial \xi_{l}}{\partial s_{j}} \frac{\partial x_{m}}{\partial t_{j}}+\sum_{m, l=1}^{n} f_{x_{l} \xi_{m}} \frac{\partial \xi_{l}}{\partial s_{j}} \frac{\partial \xi_{m}}{\partial t_{j}}+\sum_{l=1}^{n} f_{x_{1}} \frac{\partial^{2} x_{l}}{\partial s_{j} \partial t_{j}}+\sum_{l=1}^{n} f_{\xi_{l}} \frac{\partial^{2} \xi_{l}}{\partial s_{j} \partial t_{j}} \tag{4.12}
\end{align*}
$$

The partial derivatives of $f$ with respect to $x$ and $\xi$ have the same interpretation as above. In addition, the calculations will show that the derivatives of the coordinate functions at $s=t=0$ are functions on the group. However the resulting operator maps functions on $X$ into functions on $X$. It can be expressed in the variables $x$ and $\xi$.

The first derivatives of the coordinate functions
Let $e_{1}, \ldots, e_{n-1}$ be the canonical basis in the parameter spaces $\mathbf{R}^{n-1}$ for the $s$
and $t$ variables. If $h \in \mathbf{R}$ then

$$
\begin{aligned}
& x_{m}\left(h e_{j}, 0\right)=\frac{g_{m 0} \operatorname{Ch} h+g_{m j} \operatorname{Sh} h}{1+g_{00} \operatorname{Ch} h+g_{0 j} \operatorname{Sh} h} \\
& x_{m}\left(0, h e_{j}\right)=\frac{g_{m 0}}{1+g_{00}} \\
& \xi_{m}\left(h e_{j}, 0\right)=\frac{2\left(g_{m 0} C h h+g_{m j} \operatorname{Sh} h\right) g_{0 n}}{\left(1+g_{00} C h h+g_{0 j} \operatorname{Sh} h\right)^{2}}-\frac{2 g_{m n}}{1+g_{00} \operatorname{Ch} h+g_{0 j} \operatorname{Sh} h} \\
& \xi_{m}\left(0, h e_{j}\right)=\frac{2 g_{m 0}\left(g_{0 j} \sin h+g_{0 n} \cos h\right)}{\left(1+g_{00}\right)^{2}}-\frac{2\left(g_{m j} \sin h+g_{m n} \cos h\right)}{1+g_{00}}
\end{aligned}
$$

The partial derivatives at $(s, t)=(0,0)$ are

$$
\begin{aligned}
& \frac{\partial x_{m}}{\partial s_{j}}=\left.\frac{d}{d h} x_{m}\left(h e_{j}, 0\right)\right|_{h=0}=\frac{g_{m j}}{1+g_{00}}-\frac{g_{m 0} g_{0 j}}{\left(1+g_{00}\right)^{2}} \\
& \frac{\partial x_{m}}{\partial t_{j}}=0 \\
& \frac{\partial \xi_{m}}{\partial s_{j}}=-2 \frac{2 g_{m 0} g_{0 n} g_{0 j}}{\left(1+g_{00}\right)^{3}}+\frac{2 g_{m j} g_{0 n}}{\left(1+g_{00}\right)^{2}}+\frac{2 g_{m n} g_{0 j}}{\left(1+g_{00}\right)^{2}} \\
& \frac{\partial \xi_{m}}{\partial t_{j}}=\frac{2 g_{m 0} g_{0 j}}{\left(1+g_{00}\right)^{2}}-\frac{2 g_{m j}}{1+g_{00}}=-2 \frac{\partial x_{m}}{\partial s_{j}}
\end{aligned}
$$

The following expressions are needed for the differential operators:

$$
\begin{align*}
\sum_{i=1}^{n-1} \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial x_{m}}{\partial s_{j}} & =-\frac{1}{2} \sum_{i=1}^{n-1} \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial \xi_{m}}{\partial t_{j}}=\frac{1}{4} \sum_{j=1}^{n-1} \frac{\partial \xi_{l}}{\partial t_{j}} \frac{\partial \xi_{m}}{\partial t_{j}} \\
& =-\frac{1}{4} \xi_{l} \xi_{m}+\frac{1}{4} \delta_{l m}|\xi|^{2} \tag{4.13}
\end{align*}
$$

$$
\begin{equation*}
\sum_{j=1}^{n-1} \frac{\partial \xi_{l}}{\partial s_{j}} \frac{\partial \xi_{m}}{\partial t_{j}}=-\frac{(x \mid \xi)}{1-|x|^{2}}\left(2 \xi_{l} \xi_{m}-\delta_{l m}|\xi|^{2}-\frac{x_{m} \xi_{l}}{(x \mid \xi)}|\xi|^{2}\right) \tag{4.14}
\end{equation*}
$$

As an example, the calculation of formula (4.13) is given:

$$
\begin{aligned}
& \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial x_{m}}{\partial s_{j}}=\left(\frac{g_{m j}}{1+g_{00}}-\frac{g_{m 0} g_{0 j}}{\left(1+g_{00}\right)^{2}}\right)\left(\frac{g_{l j}}{1+g_{00}}-\frac{g_{l 0} g_{0 j}}{\left(1+g_{00}\right)^{2}}\right) \\
&=\left(1+g_{00}\right)^{-2} g_{m j} g_{l j}-\left(1+g_{00}\right)^{-3}\left(g_{l 0} g_{m j} g_{0 j}+g_{m 0} g_{l j} g_{0 j}\right) \\
&+\left(1+g_{00}\right)^{-4} g_{l 0} g_{m 0} g_{0 j}^{2} \\
& \begin{aligned}
\sum_{j=1}^{n-1} \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial x_{m}}{\partial s_{j}}= & \left(1+g_{00}\right)^{-2}\left(\delta_{l m}+g_{l 0} g_{m 0}-g_{l n} g_{m n}\right) \\
& \quad-\left(1+g_{00}\right)^{-3}\left(g_{l 0}\left(g_{00} g_{m 0}-g_{0 n} g_{m n}\right)+g_{m 0}\left(g_{00} g_{l 0}-g_{0 n} g_{l n}\right)\right) \\
& \quad+\left(1+g_{00}\right)^{-4} g_{l 0} g_{m 0}\left(g_{00}^{2}-1-g_{0 n}^{2}\right)
\end{aligned}
\end{aligned}
$$

The expression $\frac{1}{4} \xi_{l} \xi_{m}$ has the value

$$
\left(1+g_{00}\right)^{-4} g_{0 n}^{2} g_{i 0} g_{m 0}-g_{0 n}\left(1+g_{00}\right)^{-3}\left(g_{l 0} g_{m n}+g_{m 0} g_{i n}\right)+\left(1+g_{00}\right)^{-2} g_{l n} g_{m n}
$$

Therefore

$$
\sum_{j=1}^{n-1} \frac{\partial x_{l}}{\partial s_{j}} \frac{\partial x_{m}}{\partial s_{j}}=-\frac{1}{4} \xi_{l} \xi_{m}+\delta_{l m}\left(1+g_{00}\right)^{-2}=\frac{1}{4}\left(-\xi_{l} \xi_{m}+\delta_{l m}|\xi|^{2}\right)
$$

(All partial derivatives are taken at $s=t=0$.)
The second derivatives of the coordinate functions
The second derivatives are calculated according to the formulas

$$
\begin{aligned}
\left.\frac{\partial^{2} x}{\partial s_{j} \partial t_{j}}\right|_{s=t=0} & =\lim _{h \rightarrow 0} h^{-2}\left(x\left(h e_{j}, h e_{j}\right)-x\left(h e_{j}, 0\right)-x\left(0, h e_{j}\right)+x(0,0)\right) \\
\left.\frac{\partial^{2} \xi}{\partial t_{j}^{2}}\right|_{s=t=0} & =\lim _{h \rightarrow 0} h^{-2}\left(\xi\left(0, h e_{j}\right)+\xi\left(0,-h e_{j}\right)-2 \xi(0,0)\right)
\end{aligned}
$$

Up to third order terms

$$
\begin{aligned}
x_{m}\left(h e_{j}, h e_{j}\right) \simeq & \frac{1}{N}\left(g_{m 0}\left(1+h^{2} / 2\right)+g_{m j} h-g_{m n} h^{2} / 2\right) \\
\xi_{m}\left(h e_{j}, h e_{j}\right) \simeq & \frac{2}{N}\left(g_{00} h^{2} / 2+g_{0 j} h+g_{0 n}\left(1-h^{2} / 2\right)\right) x_{m}\left(h e_{j}, h e_{j}\right) \\
& -\frac{2}{N}\left(g_{m 0} h^{2} / 2+g_{m j} h+g_{m n}\left(1-h^{2} / 2\right)\right)
\end{aligned}
$$

with

$$
N=1+g_{00}\left(1+h^{2} / 2\right)+g_{0 j} h-g_{0 n} h^{2} / 2
$$

The resulting expressions (at $s=t=0$ ) are

$$
\begin{aligned}
& \frac{\partial^{2} x_{m}}{\partial s_{j} \partial t_{j}}=\frac{1}{4} \xi_{m} \\
& \frac{\partial^{2} x_{m}}{\partial t_{j}^{2}}=0 \\
& \frac{\partial^{2} \xi_{m}}{\partial s_{j} \partial t_{j}}=2 g_{m 0}\left(g_{0 n}^{2}-2 g_{0 j}^{2}\right)\left(1+g_{00}\right)^{-3}+\left(-2 g_{m n} g_{0 n}-g_{m 0}+4 g_{0 j} g_{m j}\right)\left(1+g_{00}\right)^{-2} \\
& \frac{\partial^{2} \xi_{m}}{\partial t_{j}^{2}}=-\xi_{m}
\end{aligned}
$$

As above this leads to the required equations

$$
\begin{align*}
& \sum_{j=1}^{n-1} \frac{\partial^{2} x_{m}}{\partial s_{j} \partial t_{j}}=\frac{n-1}{4} \xi_{m}  \tag{4.15}\\
& \sum_{j=1}^{n-1} \frac{\partial^{2} x_{m}}{\partial t_{j}^{2}}=0  \tag{4.16}\\
& \sum_{j=1}^{n-1} \frac{\partial^{2} \xi_{m}}{\partial s_{j} \partial t_{j}}=-(n+1) \frac{(x \mid \xi)}{1-|x|^{2}} \xi_{m}-\frac{n-5}{2}\left(1-|x|^{2}\right) x_{m}  \tag{4.17}\\
& \sum_{j=1}^{n-1} \frac{\partial^{2} \xi_{m}}{\partial t_{j}^{2}}=-(n-1) \xi_{m} \tag{4.18}
\end{align*}
$$

THEOREM 4. The operator $D_{|Y|^{2}}$ on $X \cong O_{ \pm}(1, n) / O(n-1)$ is given by

$$
\begin{equation*}
D_{|Y|^{2}} f=-\sum_{\mathrm{l}, m=1}^{n} f_{\xi_{\mathrm{l}} \xi_{m}}\left(\xi_{l} \xi_{m}-\delta_{l m}|\xi|^{2}\right)-(n-1) \sum_{m=1}^{n} f_{\xi_{m}} \xi_{m} \tag{4.19}
\end{equation*}
$$

It is invariant under the Möbius group GM(n-1) and under the mapping $\omega(x, \xi)=$ $(x,-\xi)$. At the same time, $D_{|Y|^{2}}$ is the Laplace operator on the sphere $\left\{\xi \in \mathbf{R}^{n}\right.$ : $|\xi|=1\}$.

Consider the stabilizer $K \cong O(n)$ of the sphere

$$
\sum=\{(x, \xi) \in X: x=0,|\xi|=1\}
$$

The Lie algebra elements $Y_{1}, \ldots, Y_{n-1}$ (see (2.11)) are in the Lie algebra 1 $O(n)$. The invariant differential operator $D_{|Y|^{2}}$ is therefore a differential operator on the subgroup $O(n)$. Furthermore, it is the restriction of the Casimir operator $\sum_{1 \leq i<j \leq n} X_{i j}^{2}$ of $\mathfrak{f}$ onto the quotient space $\Sigma \cong O(n) / O(n-1)$. This operator is the Laplace operator on the sphere.

According to the preceeding formulas (4.13)-(4.18), the operator $D_{|Y|^{2}}$ on $X$ has the explicit form given in the theorem. In particular it is seen to be independent of the $x$ coordinate (apart from the restriction $|\xi|=1-|x|^{2}$ ).

The invariance of the operator $D_{|Y|^{2}}$ under the whole Möbius group $G M(n-1)$ and under the mapping $\omega$ can be established with the same method which was used in connection with the operator $D_{Z}$.

COROLLARY. All differential operators on $X$ which are invariant under the group of special Möbius transformations $S M(n-1) \cong S O_{ \pm}(1, n)$ are invariant under the whole group $G M(n-1)$. The operators $D_{|Y|^{2}}$ and $D_{|X|^{2}}$ are also invariant under the mapping $\omega$, yet $\omega$ transforms $D_{Z}$ and $D_{(X, Y)}$ into $-D_{Z}$ and $-D_{(X, Y)}$ respectively.

THEOREM 5. The operator $D_{(X, Y)}$ is given by

$$
\begin{align*}
D_{(X, Y)} f=\frac{1}{2} & \sum_{l, m=1}^{n} f_{x_{l} \xi_{m}}\left(\xi_{l} \xi_{m}-\delta_{l m}|\xi|^{2}\right)+\frac{n-1}{4} \sum_{m=1}^{n} f_{x_{m}} \xi_{m} \\
& -\sum_{l, m=1}^{n} f_{\xi_{l} \xi_{m}}\left[\frac{(x \mid \xi)}{1-|x|^{2}}\left(2 \xi_{l} \xi_{m}-\delta_{l m}|\xi|^{2}\right)-\left(1-|x|^{2}\right) x_{l} \xi_{m}\right] \\
& -\sum_{m=1}^{n} f_{\xi_{m}}\left[(n+1) \frac{(x \mid \xi)}{1-|x|^{2}} \xi_{m}+\frac{n-5}{2}\left(1-|x|^{2}\right) x_{m}\right] \tag{4.20}
\end{align*}
$$

## 5. Spherical harmonics and the operators $S_{k}$ and $S_{k}^{*}$

A spherical harmonic of degree $k$ on the sphere $\Sigma=\left\{\xi \in \mathbf{R}^{n}:|\xi|=1\right\}$ is the restriction of a harmonic polynomial in $\mathbf{R}^{n}$ which is homogeneous of degree $k$. The space of spherical harmonics of degree $k$ will be denoted by $H^{k}$. Alternatively, it can be described as the eigenspace with eigenvalue $-k(k+n-2)$ of the Laplace operator $\Delta_{\Sigma}$ on the sphere. The system of spherical harmonics is
complete in $L^{2}(\Sigma)$. It gives a decomposition of this space as a direct orthogonal Hilbert sum

$$
L^{2}(\Sigma)=\bigoplus_{k=0}^{\infty} H^{k}
$$

DEFINITION. A spherical harmonic of degree $k$ on $X \cong O_{ \pm}(1, n) / O(n-1)$ is an eigenfunction of the operator $D_{|Y|^{2}}$ with eigenvalue $-k(k+n-2)$.

$$
\begin{equation*}
E^{k}(X)=\left\{f \in C^{\infty}(X): D_{|Y|} f=-k(k+n-2) f\right\} \tag{5.1}
\end{equation*}
$$

If a function $f \in C^{\infty}(X)$ is an eigenfunction of the operator $D_{|Y|}$, then for every fixed $x$

$$
-\sum_{l, m=1}^{n} f_{\xi \xi_{i \xi m}}\left(\xi_{l} \xi_{m}-\delta_{l m}|\xi|^{2}\right)-(n-1) \sum_{m=1}^{n} f_{\xi m} \xi_{m}=\lambda f
$$

But the left hand side is the spherical Laplace operator $\Delta_{\Sigma}$ applied to $f(x, \xi)$ with $x$ fixed. Therefore the eigenvalue $\lambda$ is of the form $-k(k+n-2)$ for some non negative integer $k$. If $\left\{h_{k 1}, \ldots, h_{k d}\right\}, d=d(k)$, is an orthogonal basis in $H^{k}$, then

$$
f(x, \xi)=\sum_{j=1}^{d} c_{k j}(x) h_{k j}(\xi)
$$

with coefficients $c_{k j} j=1, \ldots, d$ which will depend (smoothly) on $x$. Conversely, any such function is in $E^{k}$.

From the completeness property of the system of spherical harmonics on $\Sigma$ we conclude that any function $f \in C^{\infty}(X)$ has an expansion of the form

$$
\begin{equation*}
f(x, \xi)=\sum_{k=0}^{\infty} \sum_{j=1}^{d} c_{k j}(x) h_{k j}(\xi) \tag{5.2}
\end{equation*}
$$

which converges for every fixed $x$ in $L^{2}(\Sigma)$.
A harmonic polynomial $p$ of degree $k$ in $\left(\mathbf{R}^{n}\right)$ defines a symmetric tensor $t$ of order $k$ with vanishing traces

$$
\begin{align*}
& p(\xi)=\sum_{i_{1}, \ldots, i_{k}=1}^{n} t_{i_{1} \cdots i_{k}} \xi_{i_{1} \ldots \xi_{i_{k}}} \\
& t_{i_{1} \cdots i_{k}}=t_{i_{\sigma(1)} \cdots i_{\sigma(k)}} \text { for any permutation } \sigma \text { on the indices }  \tag{5.3}\\
& \sum_{i=1}^{n} t_{i i i_{3} \cdots i_{k}=0}
\end{align*}
$$

Conversely, to any such tensor the formula associates a harmonic polynomial $p$ which is homogeneous of degree $k$. The functions $f \in E^{k}$ therefore can be viewed as tensorfields of order $k$ on the hyperbolic space $B=O_{ \pm}(1, n) / O(n)$ :

$$
\begin{equation*}
E^{k}=\left\{f \in C^{\infty}(X): f(x, \xi)=\left(1-|x|^{2}\right)^{-2 k} \sum_{i_{1}, \ldots, i_{k}} t_{i_{1} \cdots i_{k}}(x) \xi_{i_{1}} \cdots \xi_{i_{k}}\right\} \tag{5.4}
\end{equation*}
$$

In this representation $t(x)=t_{i_{1}} \cdots i_{k}(x)$ is a tensorfield of symmetric tensors with vanishing trace. The factor $\left(1-|x|^{2}\right)^{-2 k}$ is a normalizing factor.

The type of the tensorfield $t$ is given by its transformation behaviour under Möbius transformations. Recall that the action of $G M(n-1)$ on $X$ is defined by

$$
\begin{equation*}
(x, \xi) \rightarrow\left(\tau_{8} x, d \tau_{8} \xi\right) \tag{2.9}
\end{equation*}
$$

The action on $C(X)$ then becomes

$$
\begin{equation*}
f^{\mathrm{g}^{-1}}(x, \xi)=f\left(\tau_{\mathrm{g}} x, d \tau_{\mathrm{g}} \xi\right) \tag{5.5}
\end{equation*}
$$

First consider the special case of a vectorfield

$$
\begin{aligned}
& f(x, \xi)=\left(1-|x|^{2}\right)^{-2} \sum_{i=1}^{n} v_{i}(x) \xi_{i} \in E^{1} \\
& f^{\mathrm{g}-1}(x, \xi)=\left(1-\left|\tau_{\mathrm{g}} x\right|^{2}\right)^{-2} \sum_{i=1}^{n} v_{i}\left(\tau_{\mathrm{g}} x\right)\left(d \tau_{\mathrm{g}} \xi\right)_{i}
\end{aligned}
$$

We set $y=\tau_{8} x$. Since $d s^{2}=\left(1-|x|^{2}\right)^{-2}|d x|^{2}$ is an invariant metric, the Jacobian determinant of the matrix

$$
G(x)=\left(\frac{\partial y_{i}}{\partial x_{k}}(x)\right)
$$

representing the tangent mapping $d \tau_{\mathrm{g}}(x)$ is given by

$$
\operatorname{det} G(x)= \pm\left(1-|y|^{2}\right)^{n}\left(1-|x|^{2}\right)^{-n}
$$

The conformality (or anti-conformality) implies that $\left(\left(1-|x|^{2}\right) /\left(1-|y|^{2}\right)\right) G(x)$ is an orthogonal matrix. In particular

$$
\begin{equation*}
G^{-1}(x)=\left(1-|x|^{2}\right)^{2}\left(1-|y|^{2}\right)^{-2} G^{t}(x) \tag{5.6}
\end{equation*}
$$

( $G^{t}$ is the transposed matrix). It then follows that

$$
\begin{aligned}
f^{g^{-1}}(x, \xi) & =\left(1-|x|^{2}\right)^{-2} \sum_{k=1}^{n} \sum_{i=1}^{n} v_{i}(y) \frac{\partial y_{i}}{\partial x_{k}}(x) \xi_{k}\left(1-|x|^{2}\right)^{2}\left(1-|y|^{2}\right)^{-2} \\
& =\left(1-|x|^{2}\right)^{-2} \sum_{k=1}^{n} v_{k}^{\mathrm{g}^{-1}}(x) \xi_{k}
\end{aligned}
$$

with

$$
\begin{equation*}
v^{\mathrm{g}-1}(x)=G^{-1}(x) v\left(\tau_{\mathrm{g}} x\right) \tag{5.7}
\end{equation*}
$$

Next assume that $f \in E^{k}$,

$$
f(x, \xi)=\left(1-|x|^{2}\right)^{-2 k} \sum_{i_{1}, \cdots, i_{k}} t_{i_{1} \cdots i_{k}}(x) \xi_{i_{1}} \cdots \xi_{i_{k}}
$$

Then the same calculations show that

$$
\begin{equation*}
f^{g^{-1}}(x, \xi)=\left(1-|x|^{2}\right)^{-2 k} \sum_{i_{1}, \cdots, i_{k}} t_{i_{1}}^{\mathrm{g}^{-1} \cdots i_{k}}(x) \xi_{i_{1}} \cdots \xi_{i_{k}} \tag{5.8}
\end{equation*}
$$

with

$$
t_{i_{1} \cdots i_{k}}^{\mathrm{g}-1}(x)=\sum_{j_{1}, \ldots, j_{k}} a_{i_{1} j_{1}} \cdots a_{i_{k} j_{k}} t_{j_{1}} \cdots j_{k}\left(\tau_{\mathrm{g}} x\right)
$$

where the $a_{k j}$ are the components of the matrix $G^{-1}(x)$. The transformation behaviour of the tensors is influenced by the choice of the normalizing factor $\left(1-|x|^{2}\right)^{-2 k}$. To illustrate this set

$$
\begin{equation*}
f(x, \xi)=\left(1-|x|^{2}\right)^{-2} \sum_{i, k} \varphi_{i k}(x) \xi_{i} \xi_{k} \tag{5.9}
\end{equation*}
$$

Here, $\Phi(x)=\left(\varphi_{i k}(x)\right)$ is a symmetric matrix with vanishing trace. The same calculations as above then show that

$$
f^{\mathrm{g}^{-1}}(x, \xi)=\left(1-|x|^{2}\right)^{-2} \sum_{i, k} \varphi_{i k}^{\mathrm{g}^{-1}}(x) \xi_{i} \xi_{k}
$$

where the transformed matrix is given by

$$
\begin{equation*}
\Phi^{g^{-1}}(x)=G^{-1}(x) \Phi\left(\tau_{8} x\right) G(x) \tag{5.10}
\end{equation*}
$$

This transformation behaviour differs from the preceding by a factor $(\operatorname{det} G(x))^{2 / n}$.

THEOREM 6. A function $f(x, \xi)=\left(1-|x|^{2}\right)^{-2} \sum_{i=1}^{n} v_{i}(x) \xi_{i} \in E^{1}$ satisfies $D_{z} f=$ 0 if and only if $v$ is a vectorfield in the Lie algebra of $\operatorname{GM}(n-1)$.

The vectorfields $v$ in the Lie algebra of $G M(n-1)$ are of the form

$$
\begin{equation*}
v(x)=B x+c\left(1+|x|^{2}\right)-2 x(c, x) \tag{2.20}
\end{equation*}
$$

with $B^{t}=-B$ and $c \in \mathbf{R}^{n}$. Direct verification shows that the functions $f \in E^{1}$ which are associated to these vectorfields satisfy $D_{\mathrm{z}} f=0$. Conversely, assume that $f \in E^{1}$ satisfies

$$
\begin{aligned}
D_{\mathrm{Z}} f & =-\left(1-|x|^{2}\right)^{-2} \frac{1}{2} \sum_{i, j} v_{i, j} \xi_{i} \xi_{j}-\left(1-|x|^{2}\right)^{-3}(v, x) \sum_{i, j} \delta_{i j} \xi_{i} \xi_{j} \\
& =0
\end{aligned}
$$

for all $(x, \xi) \in X\left(v_{i, j}\right.$ is the notation for the partial derivative $\left.\left(\partial v_{i} / \partial x_{\mathrm{j}}\right)(x)\right)$. It follows that

$$
\begin{array}{ll}
v_{i, j}=-v_{i, i} & i \neq j \\
v_{i, i}=-2\left(1-|x|^{2}\right)^{-1}(v(x), x) & i=1, \ldots, n
\end{array}
$$

and in particular

$$
v_{i, i}=v_{j, j}
$$

Assume now that $i, j$ and $k$ are different indices. Then the differentiated equations

$$
\begin{aligned}
& v_{i, j k}+v_{j, i k}=0 \\
& v_{k, i j}+v_{i, k j}=0 \\
& v_{j, k i}+v_{k, j i}=0
\end{aligned}
$$

show that $v_{i, j k}=0$. Similarly

$$
v_{i, i i j}=v_{k, k i j}=0
$$

and therefore

$$
v_{k, i i j}=0 \quad v_{i, k k k}=0
$$

This shows that all third order derivatives vanish. The vectorfield is therefore given by a second order polynomial

$$
v_{i}(x)=\frac{1}{2} \sum_{k, l} a_{i k l} x_{k} x_{l}+\sum_{k} b_{i k} x_{k}+c_{i} \quad i=1, \ldots, n
$$

and it can be assumed that

$$
a_{i k l}=a_{i l k}=v_{i, k l}
$$

A comparison of the coefficients in the equations

$$
\left(1-|x|^{2}\right) v_{i, i}=-2(v, x) \quad i=1, \ldots, n
$$

with

$$
\begin{gathered}
v_{i, i}=\frac{1}{2} \sum_{k}\left(a_{i k i}+a_{i i k}\right) x_{k}+b_{i i} \\
(v, x)=\frac{1}{2} \sum_{i, k} a_{i k l} x_{i} x_{k} x_{l}+\sum_{i, k} b_{i k} x_{i} x_{k}+\sum_{i} c_{i} x_{i}
\end{gathered}
$$

now results in the equations

$$
\begin{aligned}
b_{i i} & =0 \\
a_{i i k} & =-2 c_{k} \quad i, k=1, \ldots, n \\
b_{i k} & =-b_{k i}
\end{aligned}
$$

Since it is already known that

$$
a_{i j k}=0 \quad \text { if } \quad i \neq j \neq k \neq i
$$

and

$$
a_{k i i}=-a_{i k i}=2 c_{k} \quad \text { if } \quad k \neq i
$$

it can be concluded that

$$
\begin{aligned}
v_{i}(x) & =\sum_{k=1}^{n} a_{i k i} x_{k} x_{i}-\frac{1}{2} a_{i i i} x_{i}^{2}+\frac{1}{2} \sum_{k \neq i} a_{i k k} x_{k}^{2}+\sum_{k=1}^{n} b_{i k} x_{k}+c_{i} \\
& =-2 x_{i} \sum_{k=1}^{n} c_{k} x_{k}+c_{i} x_{i}^{2}+\sum_{k \neq i} c_{i} x_{k}^{2}+\sum_{k=1}^{n} b_{i k} x_{k}+c_{i}
\end{aligned}
$$

This shows that

$$
v(x)=c\left(1+|x|^{2}\right)-2 x(c, x)+B x \quad B^{t}=-B
$$

It should be noted that the theorem is still true for the dimension $n=2$, yet for this case the proof has to be modified slightly.

The theorem shows that the operator $D_{Z}$ applied to vectorfields (i.e. to the spherical harmonics of degree 1 on $X$ ) singles out exactly the Lie algebra of the Möbius group $G M(n-1)$.

The space of functions $f \in C^{\infty}(X)$ satisfying $D_{Z} f=0$ is an algebra, since $D_{Z}$ is a first order differential operator. If $\left\{v^{(1)}, \ldots, v^{(d)}\right\}, d=\frac{1}{2} n(n+1)$, is a basis of the Lie algebra of $G M(n-1)$ and if

$$
f_{j}(x, \xi)=\left(1-|x|^{2}\right)^{-2} \sum_{i} v_{i}^{(i)}(x) \xi_{i} \in E^{1} \quad j=1, \ldots, d
$$

then any convergent power series in $f_{1}, \ldots, f_{d}$ will be a solution of $D_{\mathrm{z}} f=0$.

## THEOREM 7. The operator

$$
\begin{equation*}
S_{k}=D_{(X, Y)}+\left(\frac{1}{2}-\left(\frac{n}{2}+k-1\right)\right) D_{Z} \tag{5.11}
\end{equation*}
$$

maps $E^{k}$ into $E^{k+1}$, and the operator

$$
\begin{equation*}
S_{k}^{*}=D_{(X, Y)}+\left(\frac{1}{2}+\left(\frac{n}{2}+k-1\right)\right) D_{Z} \tag{5.12}
\end{equation*}
$$

maps $E^{k}$ into $E^{k-1}, k=1,2,3 \ldots$

COROLLARY. The operators $D_{Z}$ and $D_{(X, Y)}$ on $E^{k}$ take the form

$$
\begin{align*}
D_{Z} & =-(n+2 k-2)^{-1} S_{k}+(n+2 k-2)^{-1} S_{k}^{*}  \tag{5.13}\\
2 D_{(X, Y)} & =\left(1+(n+2 k-2)^{-1}\right) S_{k}+\left(1-(n+2 k-2)^{-1}\right) S_{k}^{*} \tag{5.14}
\end{align*}
$$

For the proof of the theorem the operator $D_{(X, Y)}+c D_{Z}, c \in \mathbf{R}$, is applied to the function

$$
f(x, \xi)=\rho^{r} \sum_{i_{1}, \ldots, i_{k}} t_{i_{1} \cdots i_{k}}(x) \xi_{i_{1}} \cdots \xi_{i_{k}}
$$

where $t$ is a symmetric tensor with vanishing traces, $r \in \mathbf{R}$ and $\rho(x)=\left(1-|x|^{2}\right)^{-1}$. The summation convention will be applied (summation over indices which appear twice). The derivatives of the components of $t$ are denoted by

$$
\frac{\partial}{\partial x_{m}} t_{i_{1} \cdots i_{k}}=t_{i_{1} \cdots i_{k}, m}
$$

and these are no longer the components of a symmetric tensor. The result is as follows:

$$
\begin{align*}
D_{(X, Y)} f+c D_{Z} f= & \frac{k}{2} \rho^{r-2} t_{i_{1} \cdots i_{k-1} m, m} \xi_{i_{1}} \cdots \xi_{i_{k-1}}+A \rho^{r} t_{i_{1} \cdots i_{k-1} m, l} \xi_{i_{1}} \cdots \xi_{i_{k-1}} \xi_{m} \xi_{l} \\
& +B \rho^{r-1} t_{i_{1} \cdots i_{k-1} m} \xi_{i_{1}} \cdots \xi_{i_{k-1} x_{m}}+C \rho^{r+1} t_{i_{1} \cdots i_{k}} \xi_{i_{1}} \cdots \xi_{i_{k}}(x \mid \xi) \tag{5.15}
\end{align*}
$$

with

$$
\begin{aligned}
& A=-\frac{k}{2}-\frac{n-1}{4}+\frac{c}{2} \\
& B=k r-k(k-1)+k \frac{n-5}{2}+k c \\
& C=-k r+2 k(k-1)-\frac{n-1}{2} r+k(n+1)+c(r-2 k)
\end{aligned}
$$

The first observation is that $C=0$ if $r=2 k$. This motivates the normalizing factor
$\left(1-|x|^{2}\right)^{-2 k}$ occurring in the description (5.4). Having fixed $r=2 k$, the operators $S_{k}^{*}$ and $S_{k}$ are now defined by the equations $A=0$ and $B=0$ respectively. The constant $c$ for the operator $S_{k}$ is determined by the equations $r=2 k, B=$ 0 . It follows that

$$
\begin{align*}
c & =\frac{1}{2}-\left(\frac{n}{2}+k-1\right) \\
A & =-\left(\frac{n}{2}+k-1\right)  \tag{5.16}\\
S_{k} f & =\frac{k}{2} \rho^{2 k-2} t_{i_{1} \cdots i_{k} m, m} \xi_{i_{1}} \cdots \xi_{i_{k-1}}-\left(\frac{n}{2}+k-1\right) \rho^{2 k} t_{i_{1} \cdots i_{k-1} m, l} \xi_{i_{1}} \cdots \xi_{i_{k-1}} \xi_{m} \xi_{l}
\end{align*}
$$

It remains to be shown that $S_{k} f \in E^{k+1}$. For this purpose set

$$
\begin{equation*}
q_{i_{1} \cdots i_{k+1}}=\frac{1}{k+1} \sum_{i=1}^{k+1} t_{i_{1} \cdots i_{1} \cdots i_{k+1}, i_{i}} \tag{5.17}
\end{equation*}
$$

(the symbol $\hat{i}_{j}$ indicates that the index $i_{j}$ is omitted). $q$ is a symmetric tensor and

$$
\begin{equation*}
t_{i_{1} \cdots i_{k+1}} \xi_{i_{1}} \cdots \xi_{i_{k+1}}=q_{i_{1} \cdots i_{k+1}} \xi_{i_{1}} \cdots \xi_{i_{k+1}} \tag{5.18}
\end{equation*}
$$

However in general the traces of $q$ will not vanish:

$$
\begin{equation*}
q_{i_{1} \cdots i_{k-i j}}=\frac{2}{k+1} t_{i_{1} \cdots i_{k-1}, j} \tag{5.19}
\end{equation*}
$$

Consider the symmetric tensor $z$

$$
\begin{equation*}
z_{i_{1} \cdots i_{k+1}}=\delta_{i_{1} i_{2}} q_{j i i_{3} \cdots i_{k+1}}+\delta_{i_{1} i_{3}} q_{j i_{2} i_{4} \cdots i_{k+1}}+\cdots+\delta_{i_{k} i_{k+1}} q_{i_{1} \cdots i_{k-1}-1 j} \tag{5.20}
\end{equation*}
$$

Summation gives

$$
\begin{equation*}
\delta_{i_{1} i_{2}} q_{j i i_{3} \cdots i_{k+1}} \xi_{i_{1}} \cdots \xi_{i_{k+1}}=|\xi|^{2} q_{i_{1} \cdots i_{k-1 i j}} \xi_{i_{1}} \cdots \xi_{i_{k-1}} \tag{5.21}
\end{equation*}
$$

Since there are $\frac{k(k+1)}{2}$ terms in the definition of $z$, the equations (5.19), (5.20) and (5.21) show that

$$
\begin{equation*}
z_{i_{1} \cdots i_{k+1}} \xi_{i_{1}} \cdots \xi_{i_{k+1}}=k|\xi|^{2} t_{i_{1}} \cdots i_{k-1, i j} \xi_{i_{1}} \cdots \xi_{i_{k-1}} \tag{5.22}
\end{equation*}
$$

This implies

$$
\begin{equation*}
S_{k} f=\frac{1}{2} \rho^{2 k}\left(-(n+2 k-2) q_{i_{1} \cdots i_{k+1}}+z_{i_{1} \cdots i_{k+1}}\right) \xi_{i_{1}} \cdots \xi_{i_{k+1}} \tag{5.23}
\end{equation*}
$$

and it can now be shown that $S_{k} f$ is defined by a tensor with vanishing traces:

$$
\begin{aligned}
z_{i j i_{3} \cdots i_{k+1}}= & n q_{i j i_{3} \cdots i_{k+1}}+q_{i i_{3} i_{4} \cdots i_{k+1}}+\cdots \\
& +q_{i 3 j i i_{4} \cdots i_{k+1}}+\cdots \\
& +0 \\
= & (n+2(k-1)) q_{j i i_{3} \cdots i_{k+1}}
\end{aligned}
$$

(Observe that e.g. $q_{i j i} \cdots i_{k-1 i}=0$ if $k \geq 3$ ). This completes the proof for the fact that $S_{k} f \in E^{k+1}$ if $f \in E^{k}$.

The constant $c$ for the operator $S_{k}^{*}$ is determined by the equations $r=2 k, A=$ 0 . It follows that

$$
\begin{align*}
c & =\frac{1}{2}+\left(\frac{n}{2}+k-1\right)  \tag{5.24}\\
S_{k}^{*} f & =\frac{k}{2} \rho^{2 k-2} t_{i_{1} \cdots i_{k-1} m, m} \xi_{i_{1}} \cdots \xi_{i_{k-1}}+k(n+2 k-2) \rho^{2 k-1} t_{i_{1} \cdots i_{k-1} m} \xi_{i_{1}} \cdots \xi_{i_{k-1}} x_{m}
\end{align*}
$$

This clearly shows that $S_{k}^{*} f \in E^{k-1}$.
The operator $S_{k}^{*}$ can be put into a different form:

$$
\begin{equation*}
S_{k}^{*} f=\frac{k}{2} \rho^{-n} \sum_{i_{1}, \ldots, i_{k+1}} \sum_{m=1}^{n} \frac{\partial}{\partial x_{m}}\left(\rho^{n+2 k-2} t_{i_{1} \cdots i_{k-1} m}\right) \xi_{i_{1}} \cdots \xi_{i_{k-1}} \tag{5.25}
\end{equation*}
$$

The case $k=0$ is special. The functions $f \in E^{0}$ are identified with the functions on
the hyperbolic space $B$. The operators $D_{Z}$ and $D_{(X, Y)}$ map $E^{0}$ into $E^{1}$ :

$$
\begin{gathered}
D_{Z} f=-\frac{1}{2} \sum_{i=1}^{n} f_{x_{1}} \xi_{i} \\
D_{(X, Y)} f=\frac{n-1}{4} \sum_{i=1}^{n} f_{x_{i}} \xi_{i}
\end{gathered}
$$

and $S_{0}$ can be defined by the formula

$$
\begin{align*}
& S_{0} f=D_{(X, Y)} f-\frac{n+1}{2} D_{Z} f  \tag{5.26}\\
& S_{0} f=\frac{n}{2} \sum_{i=1}^{n} f_{X_{1}} \xi_{i}=\frac{n}{2}\left(1-|x|^{2}\right)^{-2} \sum_{i=1}^{n}\left(1-|x|^{2}\right)^{2} f_{x_{1}} \xi_{i}
\end{align*}
$$

The operator $S_{1}^{*} S_{0}$ then takes the form

$$
\begin{align*}
S_{1}^{*} S_{0} f & =\frac{1}{2} \frac{n}{2} \rho^{-n} \sum_{m=1}^{n} \frac{\partial}{\partial x_{m}}\left(\rho^{n}\left(1-|x|^{2}\right)^{2} f_{x_{m}}\right) \\
& =\frac{n}{4} \rho^{-n} \operatorname{div}\left(\rho^{n-2} \operatorname{grad} f\right) \tag{5.27}
\end{align*}
$$

This is (a multiple of) the Laplace operator for the hyperbolic space $\boldsymbol{B}$.
Following Ahlfors [1] the invariant operator $P$ mapping vectorfields $v$ on $B$ into tensorfields $\varphi$ is defined by the equation

$$
\begin{equation*}
\rho^{-n}(\boldsymbol{P} v)_{i j}=\frac{1}{2}\left(v_{i, j}+v_{j, i}\right)-\delta_{i j} \frac{1}{n} \sum_{k=1}^{n} v_{k, k} \tag{5.28}
\end{equation*}
$$

The tensors $P v(x)$ are symmetric and have zero trace. The operator $P^{*}$ mapping such tensorfields into vectorfields is defined by the formula

$$
\begin{equation*}
\left(P^{*} \varphi\right)_{i}=\rho^{-n-2} \sum_{j=1}^{n} \varphi_{i j, j} \tag{5.29}
\end{equation*}
$$

THEOREM 8. The operator $S_{1}$ on $E^{1}$ coincides with the operator $-(n / 2) \rho^{-n} P$
on vectorfields and $S_{2}^{*}$ on $E^{2}$ coincides with $P^{*} \rho^{n}$ provided the following identifications are made:
(1) The vectorfield $v$ on $B$ is identified with the function

$$
V(x, \xi)=\sum_{i=1}^{n}\left(1-|x|^{2}\right)^{-2} v_{i}(x) \xi_{i} \in E^{1}
$$

(2) The tensorfield $\varphi$ on $B$ is identified with the function

$$
\Phi(x, \xi)=\sum_{i, j=1}^{n}\left(1-|x|^{2}\right)^{-2} \varphi_{i j}(x) \xi_{i} \xi_{j} \in E^{2} .
$$

In particular it follows that $S_{2}^{*} S_{1}$ is the same operator as $-(n / 2) P^{*} P$.
The operator $S_{1}$ is applied to the function $V(x, \xi) \in E^{1}$ :

$$
\begin{aligned}
S_{1} V & =\frac{1}{2} \rho^{2}|\xi|^{2} \sum_{m=1}^{n} v_{m, m}-\frac{n}{2} \rho^{2} \sum_{m, l=1}^{n} v_{m, l} \xi_{m} \xi_{l} \\
& =-\frac{n}{2}\left(1-|x|^{2}\right)^{-2} \sum_{m, l=1}^{n}\left(\frac{1}{2}\left(v_{m, l}+v_{l, m}\right)-\frac{1}{n} \delta_{l m} \sum_{k=1}^{n} v_{k, k}\right) \xi_{l} \xi_{m}
\end{aligned}
$$

This shows that $S_{1}$ corresponds to $-(n / 2) \rho^{-n} P$.
Similarly, if $S_{2}^{*}$ is applied to $\Phi$, it follows from (5.25) that

$$
\begin{aligned}
S_{2}^{*} \Phi & =S_{2}^{*} \sum_{i, j=1}^{n}\left(1-|x|^{2}\right)^{-4}\left(1-|x|^{2}\right)^{2} \varphi_{i j} \xi_{i} \xi_{j} \\
& =\rho^{-n} \sum_{i, m=1}^{n} \frac{\partial}{\partial x_{m}}\left(\rho^{n+4-2}\left(1-|x|^{2}\right)^{2} \varphi_{i m}\right) \xi_{i} \\
& =\rho^{-n-2}\left(1-|x|^{2}\right)^{-2} \sum_{i, m=1}^{n} \frac{\partial}{\partial x_{m}}\left(\rho^{n} \varphi_{i m}\right) \xi_{i}
\end{aligned}
$$

This completes the proof of Theorem 8.
Equation (5.7) gives the transformation behaviour of the vectorfields under Möbiustransformations. The transformation of the tensorfields $\left(\varphi_{i j}(x)\right)$ is described by (5.10). These formulas coincide with formulas (1.5) and (1.7) in [1].

COROLLARY (Ahlfors [2], equation (2.1)). The solutions of $S_{1} f=0, f(x, \xi)=$ $\left(1-|x|^{2}\right)^{-2} \sum_{i=1}^{n} v_{i}(x) \xi_{i} \in E^{1}$ are of the form

$$
v(x)=a+B x+\lambda x+c|x|^{2}-2 x(c, x), \quad \lambda \in \mathbf{R}, \quad a, c \in \mathbf{R}^{n}, \quad B^{t}=-B
$$

The solutions of $S_{1} f=0, f \in E^{1}$ describe exactly the Lie algebra of the Möbius group $M(n)$ as a transformation group of $\mathbf{R}^{n}$ (see equation (2.19)).

THEOREM 9. For all $f \in E^{k}, k=1,2, \ldots$ there is equality

$$
\begin{equation*}
D_{|X|^{2}+Z^{2}+|Y|^{2}} f=-(n+2 k-2)^{-1}\left(S_{k+1}^{*} S_{k} f-S_{k-1} S_{k}^{*} f\right) \tag{5.30}
\end{equation*}
$$

COROLLARY. $\Delta_{K}=D_{|X|^{2}+Z^{2}}$ and $\Delta_{X}=D_{|X|^{2}+Z^{2}-|Y|^{2}}$ map $E^{k}$ into $E^{k}, k=$ $0,1,2, \ldots$

For the proof of the theorem let us calculate the commutator

$$
D_{|X|^{2}+|\mathbf{Y}|^{2}}=\left[D_{(X, Y)}, D_{Z}\right]
$$

using equations (5.13) and (5.14). Assume that $f \in E^{k}, k \in \mathbb{N}$.

$$
\begin{aligned}
(n+ & 2 k-2)\left[D_{(X, Y)}, D_{Z}\right] \\
= & S_{k+1} S_{k}(n+2 k)^{-1}\left(\frac{1}{2}-\left(\frac{n}{2}+k\right)-\frac{1}{2}+\left(\frac{n}{2}+k-1\right)\right) \\
& +S_{k+1}^{*} S_{k}(n+2 k)^{-1}\left(-\frac{1}{2}-\left(\frac{n}{2}+k\right)+\frac{1}{2}-\left(\frac{n}{2}+k-1\right)\right) \\
& +S_{k-1} S_{k}^{*}(n+2 k-4)^{-1}\left(-\frac{1}{2}+\left(\frac{n}{2}+k-2\right)+\frac{1}{2}+\left(\frac{n}{2}+k-1\right)\right) \\
& +S_{k-1}^{*} S_{k}^{*}(n+2 k-4)^{-1}\left(\frac{1}{2}+\left(\frac{n}{2}+k-2\right)-\frac{1}{2}-\left(\frac{n}{2}+k-1\right)\right)
\end{aligned}
$$

If the expression

$$
\begin{aligned}
& (n+2 k-2) D_{Z^{2}} \\
& \quad=(n+2 k)^{-1}\left(S_{k+1} S_{k}-S_{k+1}^{*} S_{k}\right)-(n+2 k-4)^{-1}\left(S_{k-1} S_{k}^{*}-S_{k-1}^{*} S_{k}^{*}\right)
\end{aligned}
$$

is added, the formula of the theorem follows:

$$
(n+2 k-2)\left(D_{|X|^{2}}+D_{|Y|^{2}}+D_{Z^{2}}\right)=-S_{k+1}^{*} S_{k}+S_{k-1} S_{k}^{*}
$$

The case $k=0$ reduces to the Laplace operator (5.27).

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Received November 4, 1981

