# On the image of a generalized $d$-plane transform on $\mathbb{R}^{n}$ 

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#### Abstract

The generalized $d$-plane transform of a function $f$ on $\mathbb{R}^{n}$ is defined on a set $\mathcal{E}$ of $d$-dimensional affine subspaces (" $d$-planes") of $\mathbb{R}^{n}$ by integration of $f$ over each subspace in $\mathcal{E}$. In general, it renders less information about the unknown function $f$ than in the special case of the well known Radon $d$-plane transform, where $\mathcal{E}$ contains every $d$-plane in $\mathbb{R}^{n}$. We study the case where $\mathcal{E}$ appears as an orbit of a matrix group and characterize the range of the spaces of Schwartz functions and of smooth ones with compact support.


## Introduction

The Radon $d$-plane transform of a function $f$ on $\mathbb{R}^{n}$ is a function on the set of $d$-dimensional affine subspaces (" $d$-planes") of $\mathbb{R}^{n}$, obtained by integration of $f$ over each such subspace. For $d=n-1$, the case of hyperplanes, it was first introduced and studied by Radon in his famous work of 1917 [12]. For $d<n-1$, the idea goes back to John, who considered line integrals of point functions in $\mathbb{R}^{3}$ in connection with the ultrahyperbolic differential equation [10]. However, the first one who systematically studied this transform was Helgason. Beside inversion formulas, he investigated the range of certain function spaces, such as $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, the smooth functions with compact support, or $\mathcal{S}\left(\mathbb{R}^{n}\right)$, the Schwartz functions (see [9]). The range of $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ under the Radon $d$-plane transform has been described by Helgason in terms of the so called "moment conditions". For $d=n-1$ (the Radon transform), these characterize the range of $\mathcal{S}\left(\mathbb{R}^{n}\right)$, too [9]. However, this is no longer true for $d<n-1$, as Gonzalez later showed [6].

In the already mentioned article of 1938 John showed that the line functions in $\mathbb{R}^{3}$, which are obtained by integration of point functions over straight lines, satisfy certain differential equations, which characterize this class of line functions [10]. The idea was further developed by Gel'fand, Gindikin and Graev in 1980, who gave a characterization of the range of $\mathcal{S}\left(\mathbb{R}^{n}\right)$ in terms of a system of second order partial differential equations but without complete proof [3]. In 1984 Grinberg
used both moment conditions and the differential operators of Gel'fand et al. to determine the range [7]. The proof of the assertion of Gel'fand et al. was given in detail by Richter in 1986 [13]. It was his great achievement to observe that the differential operators arise from infinitesimal operators of the group of euclidean motions in $\mathbb{R}^{n}$ [14]. Later, Gonzalez and Kurusa gave other proofs, the former in the group theoretic frame of Richter, the latter following John's method of parametrizing the $d$-planes [11], [6].

In this article we study the generalized situation, where the $d$-plane transform is only defined on certain subsets $\mathcal{E}$ of the set $G(d, n)$ of all $d$-planes in $\mathbb{R}^{n}$. In practice, particularly in computer tomography, where the theory applies decisively, measurements can only be made for a finite set of $d$-planes, so it is of natural interest to consider the whole problem in the restriction to subclasses of $d$-planes. It is of great practical importance to see how the main facts such as inversion formulae and range theorems depend on the geometry of $\mathcal{E}$. Gel'fand, Graev and others have studied the complex case, where $\mathcal{E}$ is an $n$-dimensional analytic submanifold of the set of all (complex) $d$-planes in $\mathbb{C}^{n}$, and have termed $\mathcal{E}$ "admissible", if it is possible to recapture a function from its integrals over the planes in $\mathcal{E}$ (see for instance [4]). Felix has treated some more general settings [2]. In this work we show that the group theoretic approach, as it has been introduced by Helgason and applied by Richter and Gonzalez with great success to the range problem, bears its fruits in a greater generality of cases, provided that $\mathcal{E}$ is an orbit of a matrix group. So we restrict ourselves to such cases. Moreover, the group theoretic treatment of the generalized $d$-plane transform reveals the meaning of the otherwise mysterious construction of the differential operators which characterize the range in [14] (here: remark 2 to lemma 5.2).

In the sequel, by "classical case" we shall mean the case $\mathcal{E}=G(d, n)$. The Radon $d$-plane transform, which refers to it, has been extensively studied through the past decades by Helgason, Richter and Gonzalez, as already mentioned.

The organization of the paper is as follows. The set $\mathcal{E}$ is introduced in section 1 . In section 2 the notion of rapidly decreasing functions on $\mathcal{E}$ and the Fourier transform are defined. Section 3 deals with the group theoretic nature of the differential operators which serve to characterize the range. Here we follow Richter [14] and correct the way in which he defines these operators. This part can be skipped during the first reading. The $d$-plane transform, the moment conditions and the differential operators in connection to it are introduced in section 4. In 5 , the range of $\mathcal{S}(V)$, the set of Schwartz functions on the vector space $V$, is characterized by both moment conditions and differential operators. The main work is done here. Section 6 covers a lot of special cases (the classical case included), where the differential operators alone characterize the range of $\mathcal{S}(V)$. In section 7 it turns out that for $d=n-1$ the only example that fits into our framework is the classical case itself. The range of $C_{c}^{\infty}(V)$ is determined with the help of a "polar coordinate" version of the classical Paley-Wiener theorem, due to Helgason [8]. This is done in section 8. Finally, the theory is applied to some concrete situations (section 9). The precise characterization of those cases, where we can dispense with the moment conditions, is still an open problem.

The author would like to stress that the main facts concerning the range of $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ and $\mathcal{S}\left(\mathbb{R}^{n}\right)$ in the classical case are here recovered as special cases. In
this sense, the classical results are truly generalized.

## 1. Fundamental concepts

Let $H$ be a Lie group, $\rho$ a $\left(C^{\infty}-\right)$ representation of $H$ on a finite-dimensional vector space $V$. The representation $\rho$ gives rise to a semi-direct product $H \times{ }_{s} V$, where multiplication is given by

$$
\left(h_{1}, v_{1}\right) \cdot\left(h_{2}, v_{2}\right):=\left(h_{1} h_{2}, v_{1}+\rho\left(h_{1}\right) v_{2}\right)
$$

(from now on we write $h \cdot v$ instead of $\rho(h) v)$.
The Lie group $H \times{ }_{s} V$ operates transitively on $V$ :

$$
(h, v) \cdot v_{0}:=v+h \cdot v_{0} .
$$

It also operates transitively on $G(d, V)$, the set of $d$-dimensional affine subspaces of $V(1 \leq d \leq n-1, n=\operatorname{dim} V)$ :

$$
(h, v) \cdot E:=v+h \cdot E \quad(h \cdot E=\{h \cdot z \mid z \in E\}) .
$$

We fix a d-dimensional linear subspace $E$ of $V$ and denote by $H_{E}$ the isotropy subgroup of $E$ in $H$. Then $H_{E} \times_{s} E$ is the isotropy subgroup of $E$ in $H \times{ }_{s} V$. The orbit of $E$ will be identified with $H \times{ }_{s} V /{ }_{H_{E} \times{ }_{s} E}$ and will be given the differentiable structure of this homogeneous space. ${ }^{H \times \times_{s} V} / H_{E} \times_{s} E$ is in fact a vector bundle over ${ }^{H} / H_{E}$ : If $\sigma: W \rightarrow H$ is a local cross section of an open set $W \subseteq{ }^{H} / H_{E}$ into $H$ and $U$ a complement of $E$ in $V$, we obtain a local trivialization by the mapping $(\dot{h}, u) \mapsto(\sigma(\dot{h}), \sigma(\dot{h}) \cdot u)\left(H_{E} \times_{s} E\right)\left(\dot{h}=h H_{E}\right)$ of $W \times U$ into ${ }^{H \times{ }_{s} V} / H_{E \times}{ }_{s} E$. We put $\mathcal{E}:={ }^{H \times{ }_{s} V} / H_{E} \times_{s} E$ and $(\dot{h}, x)$ for $(h, x)\left(H_{E} \times{ }_{s} E\right)$.

The contragredient representation $\rho^{*}$,

$$
\left[\rho^{*}(h) \xi\right](v):=\xi\left[\rho(h)^{-1} v\right] \quad\left(v \in V, \xi \in V^{*}\right),
$$

leads analogously to a semi-direct product $H \times{ }_{s} V^{*}$. We write $h^{*} \cdot \xi$ instead of $\rho^{*}(h) \xi$. If $E^{\perp} \subseteq V^{*}$ denotes the orthogonal subspace to $E$, then the set

$$
\mathcal{E}^{*}:=\left\{(\dot{h}, \xi) \mid \xi \in h^{*} \cdot E^{\perp}\right\}
$$

can be given the structure of a vector bundle over ${ }^{H} / H_{E}$ via the mappings $(\dot{h}, w) \mapsto$ $\left(\dot{h}, \sigma(\dot{h})^{*} \cdot w\right)$ of $W \times E^{\perp}$ into $\mathcal{E}^{*}$ (notation as above). (Note that $h^{*} \cdot E^{\perp}$ does not depend on $h \in h H_{E}$.) The set $\mathcal{E}^{*}$ will be called "dual bundle" (its fibres can be naturally identified with the dual spaces of the fibres of $\mathcal{E}$ ).

## 2. Rapidly decreasing functions on $\mathcal{E}$ and the Fourier transform

Gonzalez introduced a class of functions on $G\left(d, \mathbb{R}^{n}\right)$, the so called rapidly decreasing functions [6]. They can be defined on $\mathcal{E}$ in a similar way. By a local cross section of a compact set $M \subseteq{ }^{H} / H_{E}$ into $H$ we shall always mean a local $C^{\infty}$ cross section of an open neighborhood of $M$ into $H$. Let $U$ be a complement of $E$ in $V$ and introduce linear coordinates $\left(u_{1}, \ldots, u_{n-d}\right)$ on $U$. The vector bundle $\mathcal{E}$ is of fibre type $U$, as already shown.

Definition 2.1. A $C^{\infty}$-function $\varphi$ on $\mathcal{E}$ is called rapidly decreasing if for all multiindices $\alpha, \beta \in \mathbb{Z}_{+}^{n-d}$, all differential operators $T$ on ${ }^{H} / H_{E}$ and all compact sets $M \subseteq{ }^{H} / H_{E}$ admitting a local cross section $\sigma$ into $H$

$$
\sup _{\dot{h} \in M, u \in U}\left|u^{\alpha} T_{\dot{h}} \frac{\partial^{\beta}}{\partial u^{\beta}} \varphi(\dot{h}, \sigma(\dot{h}) \cdot u)\right|<\infty .
$$

(We have put $u^{\alpha}:=\prod_{i=1}^{n-d} u_{i}^{\alpha_{i}}, \frac{\partial^{\beta}}{\partial u^{\beta}}:=\prod_{i=1}^{n-d} \frac{\partial^{\beta_{i}}}{\partial u_{i}^{\beta_{i}}}$.)
The above definition does not depend on the choice of $U$ or of its basis. Rapidly decreasing functions on $\mathcal{E}^{*}$ can be defined similarly ( $E^{\perp}$ replaces $U$ ). We write $\mathcal{S}(\mathcal{E}), \mathcal{S}\left(\mathcal{E}^{*}\right)$ for these classes of functions.

Let $f \in C^{\infty}\left(V^{*}\right)$ be a function of the Schwartz class. We define a function $\varphi_{f}$ on $\mathcal{E}^{*}$ by

$$
\varphi_{f}(\dot{h}, \xi):=f(\xi)
$$

Lemma 2.2. $\varphi_{f} \in \mathcal{S}\left(\mathcal{E}^{*}\right)$.

Proof. $\quad \varphi_{f} \in C^{\infty}\left(\mathcal{E}^{*}\right)$ by the local trivializations of the vector bundle $\mathcal{E}^{*}$. We keep the notation of definition 2.1 and consider a covering of $M$ by open, relatively compact subsets $O_{1}, \ldots, O_{k}$ having the additional property that each $\overline{O_{i}}$ is contained in a chart of ${ }^{H} / H_{E}$ and that there exists a local cross section on the union of all $\overline{O_{i}}$. Under these assumptions it suffices to verify the property of definition 2.1 for each $O_{i}$ separately. In what follows, $\overline{O_{i}}$ will be regarded as a compact set in an $\mathbb{R}^{m}$.

Let $\|\cdot\|_{V^{*}}$ be a norm on $V^{*},\|\cdot\|$ the corresponding norm of linear operators on the space of endomorphisms of $V^{*}$. If $T$ is a differential operator of order $r$ on ${ }^{H} / H_{E}$, let

$$
C:=\max \left\{\sup _{0 \leq|s| \leq r} \sup _{\dot{h} \in O_{i}}\left\|\partial^{s} \rho^{*}(\sigma(\dot{h}))\right\|, \sup _{\dot{h} \in O_{i}}\left\|\rho^{*}(\sigma(\dot{h}))^{-1}\right\|\right\}<\infty
$$

$\left(s=\left(s_{1}, \ldots, s_{m}\right) \in \mathbb{Z}_{+}^{m}\right.$ is a multiindex, $|s|:=\sum_{i=1}^{m} s_{i}$, and $\partial^{s}$ is defined as in definition $2.1\left(O_{i} \subseteq \mathbb{R}^{m}\right)$ ). Let $\left(w_{1}, \ldots, w_{n-d}\right)$ denote the vector of coordinates with respect to a basis in $E^{\perp}$. The expression

$$
\begin{gathered}
\left|w^{\alpha} T_{\dot{h}}\left(\frac{\partial}{\partial w_{1}}, \ldots, \frac{\partial}{\partial w_{n-d}}\right)^{\beta} \varphi_{f}\left(\dot{h}, \sigma(\dot{h})^{*} \cdot w\right)\right| \\
=\left|\left(\left(\sigma(\dot{h})^{-1}\right)^{*} \cdot\left(\sigma(\dot{h})^{*} \cdot w\right)\right)^{\alpha}\left(\frac{\partial}{\partial w_{1}}, \ldots, \frac{\partial}{\partial w_{n-d}}\right)^{\beta} T_{\dot{h}} f\left(\sigma(\dot{h})^{*} \cdot w\right)\right|
\end{gathered}
$$

can then be estimated by a linear combination of terms of the form

$$
C^{l}\left\|\sigma(\dot{h})^{*} \cdot w\right\|_{V^{*}}^{l^{\prime}} \cdot\left|D f\left(\sigma(\dot{h})^{*} \cdot w\right)\right|
$$

( $D$ a constant coefficient differential operator, $l, l^{\prime} \in \mathbb{Z}_{+}$). However, since $f \in$ $\mathcal{S}\left(V^{*}\right)$, the latter are bounded for $\dot{h} \in O_{i}, w \in E^{\perp}$.

Let $p_{U}$ denote the projection of $V$ onto $U$, according to the decomposition $V=U \oplus E$. In order to define a Fourier transform on $\mathcal{E}$, we have to assume that a translation-invariant integral $d u$ on $U$ is left invariant by the mappings

$$
u \longmapsto p_{U}(s \cdot u), s \in H_{E} .
$$

Definition 2.3. The Fourier transform $\tilde{\varphi}$ of a rapidly decreasing function $\varphi$ on $\mathcal{E}$ is a function on the dual bundle $\mathcal{E}^{*}$, defined by

$$
\tilde{\varphi}(\dot{h}, \xi):=\int_{U} \varphi(\dot{h}, h \cdot u) e^{-i\langle h \cdot u, \xi\rangle} d u, \quad \xi \in h^{*} \cdot E^{\perp},\langle h \cdot u, \xi\rangle:=\xi(h \cdot u) .
$$

(The right hand side does not depend on $h \in h H_{E}$ because of the assumption above.)

In fact, the Fourier transform on $\mathcal{E}$ is the classical Fourier transform on its fibres. It therefore can be extended to functions which are integrable on the fibres of $\mathcal{E}$.

Proposition 2.4. $\quad \varphi \in \mathcal{S}(\mathcal{E}) \Rightarrow \tilde{\varphi} \in \mathcal{S}\left(\mathcal{E}^{*}\right)$.
The proof is straightforward and is left to the reader.
3. The left regular representation of $H \times{ }_{s} V$ on $C^{\infty}(V)$ and the kernel of its differential

As we have seen in the first section, the Lie group $H \times{ }_{s} V$ operates transitively on $V$. The subgroup $H$ (identified with $H \times_{s}\{0\}$ ) is the isotropy subgroup of $0 \in V$, so $V$ can be identified (set-theoretically and as a manifold) with ${ }^{H \times{ }_{s} V} /{ }_{H}$. The left regular representation $\lambda$ of $H \times{ }_{s} V$ on $C^{\infty}(V)$ is defined by

$$
[\lambda(h, v) f]\left(v^{\prime}\right)=f\left[(h, v)^{-1} \cdot v^{\prime}\right]=f\left(h^{-1} \cdot\left(v^{\prime}-v\right)\right) \quad\left(f \in C^{\infty}(V)\right) .
$$

It follows that the subspace of Schwartz functions $\mathcal{S}(V)$ is $\lambda$-invariant. The representation $\lambda$ can be differentiated on $C^{\infty}(V)$. For $Y \in \mathfrak{h} \times_{s} V$, the Lie algebra of $H \times_{s} V$, and $f \in C^{\infty}(V)$ we have

$$
d \lambda(Y) f(v)=\left.\frac{d}{d t} f(\exp (-t Y) \cdot v)\right|_{t=0} .
$$

If $f \in \mathcal{S}(V)$, the right hand side remains in $\mathcal{S}(V)$. This shows that $\mathcal{S}(V)$ is $d \lambda$-invariant, so $\lambda$ can be differentiated as a representation on $\mathcal{S}(V)$. Since $d \lambda$ is a representation of $\mathfrak{h} \times{ }_{s} V$, it can be extended to a representation of the universal enveloping algebra $U\left(\mathfrak{h} \times{ }_{s} V\right)$.

Let $X_{1}, \ldots, X_{m}$ be a basis of $\mathfrak{h}, Z_{1}, \ldots, Z_{n}$ one of $V$ (being identified with its Lie algebra). By the Poincaré-Birkhoff-Witt theorem the elements $X_{1}^{i_{1}} \ldots X_{m}^{i_{m}} Z_{1}^{j_{1}} \ldots Z_{n}^{j_{n}}, i_{k}, j_{l} \in \mathbb{Z}_{+}$, form a basis of $U\left(\mathfrak{h} \times_{s} V\right)$.

Definition 3.1. $\quad P \in U\left(\mathfrak{h} \times_{s} V\right)$ is said to have $X$-degree $M \geq 0$ if $P$ can be written as a linear combination of basis vectors $X_{1}^{i_{1}} \ldots X_{m}^{i_{m}} Z_{1}^{j_{1}} \ldots Z_{n}^{j_{n}}$ with $i_{1}+\ldots+i_{m} \leq M$, where at least one vector appears with $i_{1}+\ldots+i_{m}=M$.

This definition does not depend on the (ordered) basis $X_{1}, \ldots, X_{m}$ of $\mathfrak{h}$. It comes from the canonical filtration of $U(\mathfrak{h})$. An $X$-degree is assigned to every nonzero element $P \in U\left(\mathfrak{h} \times{ }_{s} V\right)$.

The adjoint representation Ad of $H \times{ }_{s} V$ as well as its differential ad extend to representations on $U\left(\mathfrak{h} \times{ }_{s} V\right)$.

Lemma 3.2. If $P \in U\left(\mathfrak{h} \times_{s} V\right)$ has $X$-degree $M>0$, then for every $Z \in V$ $a d Z(P)$ has $X$-degree $<M$ or is equal to 0 .

The proof follows immediately from the relations $\operatorname{ad} Z\left(X_{k}\right) \in V(1 \leq k \leq$ $m$ ) for the derivations ad $Z$.

Let $\pi_{V}: U\left(\mathfrak{h} \times_{s} V\right) \rightarrow U(V)$ be the canonical projection, defined by $\pi_{V}\left(X_{1}^{i_{1}} \ldots X_{m}^{i_{m}} Z_{1}^{j_{1}} \ldots Z_{n}^{j_{n}}\right)=0$ for $i_{1}+\ldots+i_{m}>0, \pi_{V}\left(Z_{1}^{j_{1}} \ldots Z_{n}^{j_{n}}\right)=Z_{1}^{j_{1}} \ldots Z_{n}^{j_{n}}$. We define a subspace $\mathcal{K}\left(\mathfrak{h} \times{ }_{s} V\right)$ of $U\left(\mathfrak{h} \times{ }_{s} V\right)$ as follows: ${ }^{1}$

An element $P \in U\left(\mathfrak{h} \times{ }_{s} V\right)$ of $X$-degree $M \geq 0$ belongs to $\mathcal{K}\left(\mathfrak{h} \times{ }_{s} V\right)$ iff for all $j_{1}, \ldots, j_{n} \in \mathbb{Z}_{+}, j_{1}+\ldots+j_{n} \leq M, \pi_{V}\left(\operatorname{ad} Z_{1}\right)^{j_{1}} \ldots\left(\operatorname{ad} Z_{n}\right)^{j_{n}}(P)=$ 0 . Moreover, $0 \in \mathcal{K}\left(\mathfrak{h} \times{ }_{s} V\right)$.

For later purposes we write $\mathcal{K}^{1}\left(\mathfrak{h} \times{ }_{s} V\right)$ for $\left\{P \in \mathcal{K}\left(\mathfrak{h} \times{ }_{s} V\right) \mid P\right.$ has $X$-degree 1$\} \cup$ $\{0\}$.

Now we are ready to state and prove the main result of this section.
Theorem 3.3. $\quad \operatorname{ker}(d \lambda)=\mathcal{K}\left(\mathfrak{h} \times{ }_{s} V\right)$.
Proof. We first prove the following statement:
If an element $P$ with $\pi_{V}(P)=0$ satisfies

$$
d \lambda(P)=d \lambda(\operatorname{Ad}(\exp Z) P)
$$

for every $Z \in V$, then $P \in \operatorname{ker}(d \lambda)$.
Let $P \neq 0$ be such an element,

$$
P=\sum_{J} c_{J} P_{J}, \quad J=\left(i_{1}, \ldots, j_{n}\right), \quad P_{J}=X_{1}^{i_{1}} \ldots X_{m}^{i_{m}} Z_{1}^{j_{1}} \ldots Z_{n}^{j_{n}}, \quad c_{J} \in \mathbb{C} \backslash\{0\}
$$

For $f \in C^{\infty}(V)$ we have

$$
\begin{gathered}
d \lambda\left(P_{J}\right)(f)(v)=\left.\frac{\partial^{|J|}}{\partial t^{J}} f(v(t)+h(t) \cdot v)\right|_{t=0}, \quad \text { where } \\
t=\left(t_{1}, \ldots, t_{m}, t_{m+1}, \ldots, t_{m+n}\right), \quad h(t)=\exp \left(-t_{m} X_{m}\right) \ldots \exp \left(-t_{1} X_{1}\right) \in H \\
v(t)=\exp \left(-t_{m+n} Z_{n}\right) \ldots \exp \left(-t_{m+1} Z_{1}\right)=-\sum_{i=1}^{n} t_{m+i} Z_{i} \in V
\end{gathered}
$$

[^0]$$
\frac{\partial^{|J|}}{\partial t^{J}}=\frac{\partial^{i_{1}+\ldots+i_{m}+j_{1}+\ldots+j_{n}}}{\partial t_{1}^{i_{1}} \ldots \partial t_{m}^{i_{m}} \partial t_{m+1}^{j_{1}} \ldots \partial t_{m+n}^{j_{n}}}
$$

The relation $\exp (\operatorname{Ad}(g) Y)=g \exp (Y) g^{-1}, Y \in \mathfrak{h} \times{ }_{s} V, g \in H \times{ }_{s} V$, implies

$$
\begin{equation*}
d \lambda(\operatorname{Ad}(g) P)(f)(v)=d \lambda(P)\left(f \circ L_{g}\right)\left(g^{-1} \cdot v\right) \quad \text { for } P \in U\left(\mathfrak{h} \times_{s} V\right) \tag{1}
\end{equation*}
$$

$L_{g}$ denoting left translation by $g$. We obtain

$$
d \lambda\left(\operatorname{Ad}(\exp v) P_{J}\right)(f)(v)=d \lambda\left(P_{J}\right)\left(f \circ L_{v}\right)(0)=\left.\frac{\partial^{|J|}}{\partial t^{J}} f(v+v(t)+h(t) \cdot 0)\right|_{t=0}=0
$$

since $\frac{\partial^{|J|}}{\partial t^{J}}$ contains at least one derivative with respect to a $t_{i}, 1 \leq i \leq m$. The assumption on $P$ leads to $d \lambda(P)(f)(v)=0$, which holds for all $v \in V$, $f \in C^{\infty}(V)$. This proves the statement.

Let now $P \in \mathcal{K}\left(\mathfrak{h} \times_{s} V\right)$ having $X$-degree $M$. For all $Z \in V$ and all $j_{1}, \ldots, j_{n} \in \mathbb{Z}_{+}, j_{1}+\ldots+j_{n}=M-1$, we obtain by the use of $\operatorname{Ad}(\exp Z)=e^{a d Z}$ :

$$
d \lambda\left(\operatorname{Ad}(\exp Z)\left(\operatorname{ad} Z_{1}\right)^{j_{1}} \ldots\left(\operatorname{ad} Z_{n}\right)^{j_{n}} P\right)=d \lambda\left(\left(\operatorname{ad} Z_{1}\right)^{j_{1}} \ldots\left(\operatorname{ad} Z_{n}\right)^{j_{n}} P\right)
$$

From the above statement it follows that the right hand side is equal to 0 . We repeat this argument with $j_{1}+\ldots+j_{n}=M-2, M-3$ etc. and obtain inductively $d \lambda(P)=0$.

For the converse we observe that $\operatorname{ker}(d \lambda)$ is $\operatorname{Ad}\left(H \times_{s} V\right)$-invariant, as equation (1) shows. By differentiation we obtain $(\operatorname{ad} Z)(\operatorname{ker}(d \lambda)) \subseteq \operatorname{ker}(d \lambda)$ for every $Z \in V$ (in fact, $\operatorname{Ad}(g)$ and $\operatorname{ad}(Z)$ operate on the finite dimensional subspaces of elements of $X$-degrees $\leq M)$. Therefore it suffices to show that if $P \in \operatorname{ker}(d \lambda)$ it holds $\pi_{V}(P)=0$. Let $P_{0}=\pi_{V}(P)$. For $f \in C^{\infty}(V), v \in V$ we obtain:

$$
\begin{gathered}
0=d \lambda(P)(f)(v)=d \lambda[\operatorname{Ad}(\exp v) P](f)(v)=d \lambda(P)\left(f \circ L_{\exp v}\right)(0) \text { by }(1) \\
=d \lambda\left(P_{0}\right)\left(f \circ L_{\exp v}\right)(0)+d \lambda\left(P-P_{0}\right)\left(f \circ L_{\exp v}\right)(0)=d \lambda\left(P_{0}\right)\left(f \circ L_{\exp v}\right)(0) \\
=d \lambda\left[\operatorname{Ad}(\exp v) P_{0}\right](f)(v)=d \lambda\left(P_{0}\right)(f)(v),
\end{gathered}
$$

so $d \lambda\left(P_{0}\right)=0$ and $P_{0}=0$, since $d \lambda$ is injective on $U(V)$.

Remark . The definition of the subspace $\mathcal{K}\left(\mathfrak{h} \times{ }_{s} V\right)$ is suggested by Corollary 3.6 in [1].

## 4. The $d$-plane transform

We return to the setting of the first two sections. $E$ is a $d$-dimensional subspace of $V(1 \leq d \leq n-1), U$ a complement of $E$ in $V$. The following considerations contain the Radon transform $(d=n-1)$ as a special case.

Let us introduce translation invariant integrals $d z$ and $d u$ on $E$ and $U$ respectively. The integral $d y=d z d u$ is then translation invariant on $V$. From now on we assume that $H_{E}$ acts by unimodular automorphisms on $E$ (that is of determinant $\pm 1$ ). This, combined with the condition needed for the Fourier transform on $\mathcal{E}$, leads to the assumption that $H_{E}$ acts unimodularly on both $E$ and $V$. Moreover, we can assume that $H$ operates by unimodular automorphisms on $V$; if this is not the case, we only have to replace $\rho(h)$ by $\delta(h)^{-1 / n} \rho(h)$, where $\delta(h)$ denotes the modulus ( $=$ absolute value of the determinant) of $\rho(h)$.

Definition 4.1. Let $f$ be an integrable function on $V$. The $d$-plane transform of $f$ is a function $\hat{f}$ on $\mathcal{E}$, defined by

$$
\hat{f}(\dot{h}, x)=\int_{E} f(x+h \cdot z) d z
$$

The new function is well defined almost everywhere, due to the unimodularity assumption above. The following lemma, often called "Projection-Slice theorem", is of great importance.

Lemma 4.2. Let $f \in \mathcal{S}(V), \tilde{f} \in \mathcal{S}\left(V^{*}\right)$ being the (classical) Fourier transform of $f$. The d-plane transform $\hat{f}$ is then integrable on every fibre of $\mathcal{E}$, and we have:

$$
\tilde{\hat{f}}(\dot{h}, \xi)=\tilde{f}(\xi), \quad \xi \in h^{*} \cdot E^{\perp}
$$

Proof. Clearly, $\hat{f}$ is everywhere defined and integrable on the fibres by the Fubini theorem. We have:

$$
\begin{gathered}
\tilde{\hat{f}}(\dot{h}, \xi)=\int_{U} \hat{f}(\dot{h}, h \cdot u) e^{-i\langle h \cdot u, \xi\rangle} d u=\int_{U} \int_{E} f(h \cdot(u+z)) e^{-i\langle h \cdot(u+z), \xi\rangle} d z d u \\
(\langle h \cdot z, \xi\rangle=0) \\
=\int_{V} f(h \cdot y) e^{-i\langle h \cdot y, \xi\rangle} d y=\int_{V} f(y) e^{-i\langle y, \xi\rangle} d y=\tilde{f}(\xi)
\end{gathered}
$$

Lemma 4.3. For $f \in \mathcal{S}(V)$ we have $\hat{f} \in \mathcal{S}(\mathcal{E})$.
Proof. Clearly, $\hat{f} \in C^{\infty}(\mathcal{E})$. From lemmas 4.2 and 2.2 it follows that $\tilde{\hat{f}} \in$ $\mathcal{S}\left(\mathcal{E}^{*}\right)$. The desired relation is now obtained from the analogue of proposition 2.4 for the inverse Fourier transform.

An important property of the $d$-plane transforms $\hat{f}$ is that they satisfy the so called "moment conditions".

Proposition 4.4. Let $\varphi \in \mathcal{S}(V)^{\wedge}$. For every $m \in \mathbb{Z}_{+}$there exists a homogeneous polynomial $P_{m}$ of degree $m$ on $V^{*}$, such that for every $h \in{ }^{H} / H_{E}$, $\xi \in h^{*} \cdot E^{\perp}$ :

$$
\begin{equation*}
\int_{U} \varphi(\dot{h}, h \cdot u) \cdot\langle h \cdot u, \xi\rangle^{m} d u=P_{m}(\xi) . \tag{2}
\end{equation*}
$$

Proof. Let $\varphi=\hat{f}, f \in \mathcal{S}(V)$. We have:

$$
\begin{gathered}
\int_{U} \varphi(\dot{h}, h \cdot u)\langle h \cdot u, \xi\rangle^{m} d u=\int_{U} \int_{E} f(h \cdot(u+z))\langle h \cdot(u+z), \xi\rangle^{m} d z d u \\
=\int_{V} f(h \cdot y)\langle h \cdot y, \xi\rangle^{m} d y=\int_{V} f(y)\langle y, \xi\rangle^{m} d y
\end{gathered}
$$

Let $\mathcal{S}_{M}(\mathcal{E}):=\{\varphi \in \mathcal{S}(\mathcal{E}) \mid \varphi$ satisfies $(2)\}$. We have shown: $\mathcal{S}(V)^{\wedge} \subseteq$ $\mathcal{S}_{M}(\mathcal{E})$.

The functions in $\mathcal{S}(V)^{\wedge}$ have a further common property. They are annihilated by certain differential operators.

Let $\nu$ be the left regular representation of $H \times{ }_{s} V$ on $C^{\infty}(\mathcal{E})$. For $f \in \mathcal{S}(V)$ we compute:

$$
\begin{gathered}
{[\nu(h, v) \hat{f}]\left(\dot{h}^{\prime}, x^{\prime}\right)=\hat{f}\left[(h, v)^{-1} \cdot\left(\dot{h}^{\prime}, x^{\prime}\right)\right]=\hat{f}\left(h^{-1} \dot{h}^{\prime}, h^{-1} \cdot x^{\prime}-h^{-1} \cdot v\right)} \\
=\int_{E} f\left(h^{-1} \cdot x^{\prime}-h^{-1} \cdot v+h^{-1} h^{\prime} \cdot z\right) d z=\int_{E} f\left((h, v)^{-1} \cdot\left(x^{\prime}+h^{\prime} \cdot z\right)\right) d z \\
=\int_{E}[\lambda(h, v) f]\left(x^{\prime}+h^{\prime} \cdot z\right) d z=[\lambda(h, v) f]^{\wedge}\left(\dot{h}^{\prime}, x^{\prime}\right)
\end{gathered}
$$

( $\lambda$ denoting the left regular representation of $H \times{ }_{s} V$ on $C^{\infty}(V)$ ), so

$$
\begin{equation*}
[\lambda(g) f]^{\wedge}=\nu(g) \hat{f}, \quad g \in H \times_{s} V, \tag{3}
\end{equation*}
$$

which means that the operator $f \mapsto \hat{f}$ intertwines $\lambda$ and $\nu$. Differentiation of (3) leads to

$$
[d \lambda(Y) f]^{\wedge}=d \nu(Y) \hat{f}
$$

for $Y$ in $\mathfrak{h} \times{ }_{s} V$ and then in $U\left(\mathfrak{h} \times{ }_{s} V\right)$. Thus, if $Y \in \operatorname{ker}(d \lambda), d \nu(Y) \hat{f}=0$ for all $f \in \mathcal{S}(V)^{\wedge}$.

Let $\mathcal{S}_{D}(\mathcal{E}):=\{\varphi \in \mathcal{S}(\mathcal{E}) \mid d \nu(Y) \varphi=0$ for $Y \in \operatorname{ker}(d \lambda)\}$. We have shown: $\mathcal{S}(V)^{\wedge} \subseteq \mathcal{S}_{D}(\mathcal{E})$. Thus, $\mathcal{S}(V)^{\wedge} \subseteq \mathcal{S}_{D}(\mathcal{E}) \cap \mathcal{S}_{M}(\mathcal{E})$.

## 5. The image of the $d$-plane transform

In this section we prove the relation

$$
\begin{equation*}
\mathcal{S}(V)^{\wedge}=\mathcal{S}_{D}(\mathcal{E}) \cap \mathcal{S}_{M}(\mathcal{E}) \tag{4}
\end{equation*}
$$

For this we have to impose the following two conditions on the action of the Lie group $H$ :
(i) The mapping from $H \times E^{\perp}$ into $V^{*},(h, w) \mapsto h^{*} \cdot w$, is surjective and its restriction to $H \times\left(E^{\perp} \backslash\{0\}\right)$ a submersion.
(ii) If $h^{*} \cdot w_{1}=w_{2}, w_{1}, w_{2} \in E^{\perp}, h \in H$, then there exists $s \in H_{E}$ such that $s^{*} \cdot w_{1}=w_{2}$.

Both conditions are satisfied in the classical case $\left(H=S O(n), V=\mathbb{R}^{n}\right) .{ }^{2}$

[^1]Lemma 5.1. The $d$-plane transform is an injective mapping on $\mathcal{S}(V)$.
The proof follows from lemma 4.2 and the surjectivity in assumption (i) above.

For (4) it remains to be proven that for $\varphi \in \mathcal{S}_{D}(\mathcal{E}) \cap \mathcal{S}_{M}(\mathcal{E})$ there exists $f \in \mathcal{S}(V)$ such that $\varphi=\hat{f}$. We proceed in three steps:

1. There exists a function $F \in C^{\infty}\left(V^{*} \backslash\{0\}\right)$ satisfying $\tilde{\varphi}(\dot{h}, \xi)=F(\xi)$.
2. $F \in C^{\infty}\left(V^{*}\right)$.
3. $F \in \mathcal{S}\left(V^{*}\right)$.

It then holds for $f \in \mathcal{S}(V)$ with $\tilde{f}=F$ :

$$
\tilde{\varphi}(\dot{h}, \xi)=F(\xi)=\tilde{f}(\xi)=\tilde{\hat{f}}(\dot{h}, \xi)
$$

by lemma 4.2 , so $\tilde{\varphi}=\tilde{\hat{f}}$, and by inverse Fourier transform, $\varphi=\hat{f}$.

## First step

Let $X_{1}, \ldots, X_{m}$ be a basis of $\mathfrak{h}$. For every $w \in V^{*}$ we denote by $H_{w}$ the isotropy subgroup of $w$ in $H$ and by $\mathfrak{h}_{w}$ its Lie algebra.

Lemma 5.2. Let $l=\inf \left\{\operatorname{dim} H_{w} \mid w \in V^{*} \backslash\{0\}\right\}, O=\left\{w \in V^{*} \backslash\{0\} \mid \operatorname{dim} H_{w}=\right.$ $l\}$.
a) $O$ is a dense open subset of $V^{*} \backslash\{0\}$.
b) Let $w_{0} \in O$. For every $w$ in a dense open neighborhood of $w_{0}$ in $O$ there exists a basis

$$
Y_{j}(w)=\sum_{k=1}^{m} a_{j k}(w) X_{k}, \quad 1 \leq j \leq l,
$$

of $\mathfrak{h}_{w}$, such that the coefficients $a_{j k}$ are homogeneous polynomials in $w \in V^{*}$.
Proof. If we differentiate with respect to $t$ the equation $\rho^{*}\left[\exp t\left(a_{1} X_{1}+\ldots+\right.\right.$ $\left.\left.a_{m} X_{m}\right)\right] w-w=0$, we obtain:

$$
\begin{equation*}
Y=\sum_{k=1}^{m} a_{k} X_{k} \in \mathfrak{h}_{w} \Longleftrightarrow d \rho^{*}(Y) w=0 \Longleftrightarrow \sum_{k=1}^{m} a_{k} d \rho^{*}\left(X_{k}\right) w=0 \tag{5}
\end{equation*}
$$

a) $O \subseteq V^{*}$ is the set of vectors $w \in V^{*}$, where the system $\left(d \rho^{*}\left(X_{1}\right) w, \ldots\right.$, $\left.d \rho^{*}\left(X_{m}\right) w\right)$ has maximal rank, that is $m-l$.
b) Let us assume that $\left(d \rho^{*}\left(X_{k}\right) w\right)_{1 \leq k \leq m-l}$ is a basis of $\left\langle d \rho^{*}\left(X_{k}\right) w\right| 1 \leq k \leq$ $m\rangle$ for $w$ in a dense open neighborhood of $w_{0}$ in $O$. If we bring the terms with $m-l<k \leq m$ in (5) to the right hand side (and choose a basis on $V^{*}$ ), we obtain for every choice of $a_{m-l+1}, \ldots, a_{m}$ an $n \times(m-l)$-linear system, whose equations are linear in $w$. The statement of the lemma follows by applying Cramer's rule to an $(m-l) \times(m-l)$-part of the system and multiplying the vector $\left(a_{1}, \ldots, a_{m}\right)$ by the common denominator.

Remark . 1. According to the construction given in the proof, the polynomials $a_{j k}$ have the same degree of homogeneity.
2. If we consider the polynomials $a_{j k}$ as elements of the symmetric ( $=$ universal enveloping) algebra of $V$, then the $Y_{j}$ become the operators in $\mathcal{K}^{1}\left(\mathfrak{h} \times{ }_{s} V\right)$ (see section 3):

$$
P_{j}:=\sum_{k=1}^{m} X_{k} a_{j k} \in \mathcal{K}^{1}\left(\mathfrak{h} \times_{s} V\right) .
$$

This can be seen as follows $\left(Z_{1}, \ldots, Z_{n}\right.$ a basis of $\left.V\right)$ :

$$
0=\sum_{k=1}^{m} a_{j k}(\xi) d \rho^{*}\left(X_{k}\right) \xi
$$

implies

$$
\begin{gathered}
0=\left\langle\sum_{k=1}^{m} a_{j k}(\xi) d \rho^{*}\left(X_{k}\right) \xi, Z_{l}\right\rangle=\sum_{k=1}^{m} a_{j k}(\xi)\left\langle d \rho^{*}\left(X_{k}\right) \xi, Z_{l}\right\rangle \\
=-\sum_{k=1}^{m} a_{j k}(\xi)\left\langle\xi, d \rho\left(X_{k}\right) Z_{l}\right\rangle=\sum_{k=1}^{m} a_{j k}(\xi)\left\langle\xi,\left[Z_{l}, X_{k}\right]\right\rangle=\left\langle\xi,\left(\mathrm{ad} Z_{l}\right) P_{j}\right\rangle .
\end{gathered}
$$

This gives a geometrical meaning to $\operatorname{dim} \mathcal{K}^{1}\left(\mathfrak{h} \times{ }_{s} V\right)=l$.
We observe that for $\varphi \in \mathcal{S}(\mathcal{E}), T \in \mathfrak{h} \times_{s} V, d \nu(T) \varphi$ remains in $\mathcal{S}(\mathcal{E})$. Therefore we can define the Fourier transform of every $d \nu(T), T \in \mathfrak{h} \times_{s} V$, and then for $T \in U\left(\mathfrak{h} \times{ }_{s} V\right)$, by iteration:

$$
d \widetilde{\nu(T)} \tilde{\varphi}:=[d \nu(T) \varphi]^{\sim}
$$

We calculate $d \widetilde{\nu(X)}, \widetilde{\nu \nu(Z)}$ for $X \in \mathfrak{h}, Z \in V$ :

$$
\begin{gather*}
d \widetilde{\nu(X)} \tilde{\varphi}(\dot{h}, \xi)=\int_{U} d \nu(X) \varphi(\dot{h}, h \cdot u) e^{-i\langle h \cdot u, \xi\rangle} d u \\
=\left.\int_{U} \frac{d}{d t} \varphi(\exp (-t X) \dot{h}, \exp (-t X) h \cdot u)\right|_{t=0} \cdot e^{-i\langle h \cdot u, \xi\rangle} d u \\
=\left.\frac{d}{d t} \int_{U} \varphi(\exp (-t X) \dot{h}, \exp (-t X) h \cdot u) \cdot e^{-i\left\langle\exp (-t X) h \cdot u, \exp (-t X)^{*} \cdot \xi\right\rangle} d u\right|_{t=0} \\
=\left.\frac{d}{d t} \tilde{\varphi}\left(\exp (-t X) \dot{h}, \exp (-t X)^{*} \cdot \xi\right)\right|_{t=0} ;  \tag{6}\\
d \widetilde{\nu(Z)} \tilde{\varphi}(\dot{h}, \xi)=\int_{U} d \nu(Z) \varphi(\dot{h}, h \cdot u) e^{-i\langle h \cdot u, \xi\rangle} d u \\
=\left.\int_{U} \frac{d}{d t} \varphi(\dot{h},-t Z+h \cdot u)\right|_{t=0} \cdot e^{-i\langle h \cdot u, \xi\rangle} d u \\
=\left.\int_{U} \frac{d}{d t} \varphi\left(\dot{h},-t h \cdot p_{U}\left(h^{-1} \cdot Z\right)+h \cdot u\right)\right|_{t=0} \cdot e^{-i\langle h \cdot u, \xi\rangle} d u
\end{gather*}
$$

( $p_{U}$ denoting projection on $U$ parallel to $E$ )

$$
=\left.\frac{d}{d t} \int_{U} \varphi(\dot{h}, h \cdot u) e^{-i\langle h \cdot u, \xi\rangle} e^{-i\langle t Z, \xi\rangle} d u\right|_{t=0} \quad(\varphi \in \mathcal{S}(\mathcal{E}))
$$

$$
\begin{equation*}
=-i\langle Z, \xi\rangle \tilde{\varphi}(\dot{h}, \xi) \tag{7}
\end{equation*}
$$

Following the notation in lemma 5.2 we define operators $Y_{j}, 1 \leq j \leq l$, on $C^{\infty}\left(\mathcal{E}^{*}\right)$ :

$$
\left.Y_{j} \psi(\dot{h}, \xi):=\sum_{k=1}^{m} a_{j k}(\xi) d \nu \widetilde{(X}_{k}\right) \psi(\dot{h}, \xi)
$$

For $f \in \mathcal{S}(V)$ we compute:

$$
\begin{gather*}
\left.Y_{j}(\tilde{\tilde{f}})(\dot{h}, \xi)=\sum_{k=1}^{m} a_{j k}(\xi) d \nu \widetilde{\left(X_{k}\right)} \tilde{\hat{f}}\right)(\dot{h}, \xi) \\
=\left.\sum_{k=1}^{m} a_{j k}(\xi) \frac{d}{d t}(\tilde{\hat{f}})\left(\exp \left(-t X_{k}\right) \dot{h}, \exp \left(-t X_{k}\right)^{*} \cdot \xi\right)\right|_{t=0} \quad \text { by } \quad(6) \\
=\left.\sum_{k=1}^{m} a_{j k}(\xi) \frac{d}{d t} \tilde{f}\left(\exp \left(-t X_{k}\right)^{*} \cdot \xi\right)\right|_{t=0} \quad \text { by lemma } 4.2 \\
=\left.\frac{d}{d t} \tilde{f}\left(\exp \left(-t Y_{j}(\xi)\right)^{*} \cdot \xi\right)\right|_{t=0}=0, \quad 1 \leq j \leq l . \tag{8}
\end{gather*}
$$

We now proceed to a closer study of the operators $Y_{j}$. A homogeneous polynomial in $\xi \in V^{*}$ is obtained by multiplication and addition of terms of the form $\langle Z, \xi\rangle, \underline{Z} \in V$. Thus, because of (7) there exist elements $P_{j k} \in U(V)$ satisfying $\left.d \nu \widetilde{\left(P_{j k}\right)}\right) \psi(\dot{h}, \xi)=a_{j k}(\xi) \psi(\dot{h}, \xi)$. Using this fact, we obtain:

$$
\begin{gathered}
\left.Y_{j} \tilde{\varphi}(\dot{h}, \xi)=\sum_{k=1}^{m} a_{j k}(\xi) d \nu \widetilde{\left(X_{k}\right)} \tilde{\varphi}(\dot{h}, \xi)=\sum_{k=1}^{m} d \nu \widetilde{\nu P}_{j k}\right) d \nu \widetilde{\left(X_{k}\right)} \tilde{\varphi}(\dot{h}, \xi) \\
=\sum_{k=1}^{m}\left[d \nu\left(P_{j k}\right) d \nu\left(X_{k}\right) \varphi\right]^{\sim}(\dot{h}, \xi)
\end{gathered}
$$

so

$$
Y_{j} \tilde{\varphi}=\sum_{k=1}^{m}\left[d \nu\left(P_{j k}\right) d \nu\left(X_{k}\right) \varphi\right]^{\sim}=\left[d \nu\left(P_{j}\right) \varphi\right]^{\sim}=d \widetilde{\nu\left(P_{j}\right)} \tilde{\varphi}
$$

where we have put $P_{j}:=\sum_{k=1}^{m} P_{j k} X_{k} \in U\left(\mathfrak{h} \times_{s} V\right)$. Equation (8) now implies:

$$
0=\left[d \nu\left(P_{j}\right) \hat{f}\right]^{\sim} \Longrightarrow 0=d \nu\left(P_{j}\right) \hat{f}=\left[d \lambda\left(P_{j}\right) f\right]^{\wedge} \Longrightarrow d \lambda\left(P_{j}\right) f=0
$$

by (3) and lemma 5.1.
We have thus proved that there exist elements $P_{j} \in \operatorname{ker}(d \lambda)$ such that

$$
\begin{equation*}
Y_{j}=d \widetilde{\nu\left(P_{j}\right)}, \quad 1 \leq j \leq l \tag{9}
\end{equation*}
$$

Remark . It holds: $P_{j}=\sum_{k=1}^{m} X_{k} P_{j k}$. In fact, for $Q_{j}:=\sum_{k=1}^{m} X_{k} P_{j k}$ we have $Q_{j} \in \mathcal{K}^{1}\left(\mathfrak{h} \times_{s} V\right)$ (see section 3), since $\operatorname{ad}(Z) Q_{j}=\operatorname{ad}(Z) P_{j}=0$ for every $Z \in V$ $\left(P_{j} \in \mathcal{K}^{1}\left(\mathfrak{h} \times_{s} V\right)\right.$ ). It follows: $Q_{j} \in \operatorname{ker}(d \lambda) \Rightarrow Q_{j}-P_{j} \in \operatorname{ker}(d \lambda) \Rightarrow Q_{j}-P_{j}=0$, because $Q_{j}-P_{j} \in U(V)$ and $d \lambda$ is injective on $U(V)$.

We come to the main result of this subsection:

Proposition 5.3. For every $\varphi \in \mathcal{S}_{D}(\mathcal{E})$ there exists a function $F \in C^{\infty}\left(V^{*} \backslash\right.$ $\{0\})$ such that $\tilde{\varphi}(\dot{h}, \xi)=F(\xi)$ for $\xi \in h^{*} \cdot E^{\perp}, \xi \neq 0$.

Proof. The restriction of $\pi: \mathcal{E}^{*} \rightarrow V^{*},(\dot{h}, \xi) \mapsto \xi$, to $\pi^{-1}\left(V^{*} \backslash\{0\}\right)$ is by assumption a submersion. Thus it is enough to show that $\tilde{\varphi}$ is constant on every subset $\pi^{-1}(\xi), \xi \neq 0$. For this we recall assumption (ii) in the beginning of this section.

Let $\xi=h_{1}^{*} \cdot w_{1}=h_{2}^{*} \cdot w_{2}, w_{1}, w_{2} \in E^{\perp} \backslash\{0\}$. Then there exists $s \in H_{E}$ such that $s^{*} \cdot w_{1}=w_{2}$, and we have:

$$
h_{1}^{*} \cdot w_{1}=\left(h_{2} s\right)^{*} \cdot w_{1} \Leftrightarrow h_{1}^{-1} h_{2} s \in H_{w_{1}} \Leftrightarrow h_{2} s \in h_{1} H_{w_{1}}=H_{\xi} h_{1}
$$

and $h_{2} s H_{E}=\dot{h}_{2}$, so there is $r \in H_{\xi}$ such that $r \cdot\left(\dot{h}_{1}, \xi\right):=\left(r \dot{h}_{1}, r^{*} \cdot \xi\right)=\left(\dot{h}_{2}, \xi\right)$. It therefore remains to show that the function

$$
\begin{equation*}
r \longmapsto \tilde{\varphi}\left(r \dot{h}, r^{*} \cdot \xi\right) \tag{10}
\end{equation*}
$$

is constant on $H_{\xi}$.
Relation (9) implies $Y_{j} \tilde{\varphi}=\left[d \nu\left(P_{j}\right) \varphi\right]^{\sim}=0,1 \leq j \leq l$. On the other hand, with the help of $\Phi: H \times E^{\perp} \rightarrow \mathbb{C}, \Phi(h, w):=\tilde{\varphi}\left(\dot{h}, h^{*} \cdot w\right)$, we compute:

$$
\begin{gathered}
Y_{j} \tilde{\varphi}(\dot{h}, \xi)=\sum_{k=1}^{m} a_{j k}(\xi) d \widetilde{\nu\left(X_{k}\right) \tilde{\varphi}(\dot{h}, \xi)} \\
=\left.\sum_{k=1}^{m} a_{j k}(\xi) \frac{d}{d t} \tilde{\varphi}\left(\exp \left(-t X_{k}\right) \dot{h}, \exp \left(-t X_{k}\right)^{*} \cdot \xi\right)\right|_{t=0} \quad \text { by } \quad(6) \\
=\left.\sum_{k=1}^{m} a_{j k}(\xi) \frac{d}{d t} \Phi\left(\exp \left(-t X_{k}\right) h, w\right)\right|_{t=0} \quad \text { if } \xi=h^{*} \cdot w \\
=\left.\frac{d}{d t} \Phi\left(\exp \left(-t Y_{j}(\xi)\right) h, w\right)\right|_{t=0}=\left.\frac{d}{d t} \tilde{\varphi}\left(\exp \left(-t Y_{j}(\xi)\right) \dot{h}, \xi\right)\right|_{t=0}
\end{gathered}
$$

Since $Y_{j}(\xi), 1 \leq j \leq l$, form a basis of $\mathfrak{h}_{\xi}$ for $\xi$ in a dense open set $O^{\prime} \subseteq V^{*} \backslash\{0\}$, (10) is constant on $H_{\xi}, \xi \in O^{\prime}$. Thus there exists a function $F \in C^{\infty}\left(O^{\prime}\right)$ satisfying $\tilde{\varphi}(\dot{h}, \xi)=F(\xi), \xi \in O^{\prime}$. For $0 \neq \xi_{0}=h_{0}^{*} \cdot w_{0} \notin O^{\prime}$ let $W_{\left(h_{0}, w_{0}\right)}$ be a submanifold of $H \times E^{\perp}$ through $\left(h_{0}, w_{0}\right)$ which is mapped by $(h, w) \mapsto$ $h^{*} \cdot w$ diffeomorphically onto an open neighborhood $W_{\xi_{0}}$ of $\xi_{0}$ in $V^{*}$, and put $F\left(h^{*} \cdot w\right):=\tilde{\varphi}\left(\dot{h}, h^{*} \cdot w\right)$ for $(h, w) \in W_{\left(h_{0}, w_{0}\right)}$. The argument shows that $F$ can be smoothly extended to an open neighborhood of every $\xi \notin O^{\prime}, \xi \neq 0$. If we consider a covering of $V^{*} \backslash\{0\}$ by such neighborhoods and by $O^{\prime}$, there exists a locally finite subcovering and a partition of unity associated with it, whereby $F$ can be extended to a $C^{\infty}$-function on $V^{*} \backslash\{0\}$. The equation $\tilde{\varphi}=F \circ \pi$ then holds on the dense subset $\pi^{-1}\left(O^{\prime}\right)$ and hence everywhere on $\pi^{-1}\left(V^{*} \backslash\{0\}\right)$.

Remark . For the proof we have only used the fact that $\varphi$ is annihilated by $d \nu(Y), Y \in \mathcal{K}^{1}\left(\mathfrak{h} \times{ }_{s} V\right)$.

## Second Step

The proof of $F \in C^{\infty}\left(V^{*}\right)$ is based on an application of a theorem of Glaeser [5], which we shall now recall.

Definition (Glaeser). Let $\omega \subseteq \mathbb{R}^{n}$ be an open set, $o \subseteq \omega, \mathcal{V} \subseteq C^{\infty}(\omega)$. A function $\Phi \in C^{\infty}(\omega)$ is said to be attached to $\mathcal{V}$ on $o$ if there is $\Psi \in \mathcal{V}$ such that $\Phi-\Psi$ and all its partial derivatives vanish on $o$. (In that case $\Psi$ will be called "interpolating function" for $\Phi$ ). The function $\Phi$ will be called bipunctually attached to $\mathcal{V}$ if $\Phi$ is attached to $\mathcal{V}$ on every two-element subset of $\omega$.

Furthermore, we consider an open set $\Omega \in \mathbb{R}^{p}, p \leq n$, and a $C^{\infty}$-map $\theta: \omega \rightarrow \Omega$. Let $A_{\theta}$ denote the set of $C^{\infty}$-functions on $\omega$ of the form $F \circ \theta$, $F \in C^{\infty}(\Omega)$.

Theorem (Glaeser). If $\theta: \omega \rightarrow \Omega$ satisfies the following 4 conditions and $\Phi \in C^{\infty}(\omega)$ is bipunctually attached to $A_{\theta}$, then $\Phi$ belongs to $A_{\theta}$.
$\Theta_{1}: \theta: \omega \rightarrow \Omega$ is real analytic.
$\Theta_{2}$ : $\theta$ has maximal rank $(=p)$ on a dense open subset of $\omega$.
$\Theta_{3}: \theta(\omega)$ is closed in $\Omega$.
$\Theta_{4}$ : For every compact set $K \subseteq \theta(\omega)$ there is a compact set $k \subseteq \omega$ such that $\theta(k)=K$.

We prepare the application of this theorem.
Lemma 5.4. $\quad$ There is a compact set $K \subseteq H$ such that the mapping $(h, w) \mapsto$ $h^{*} \cdot w$ from $K \times E^{\perp}$ into $V^{*}$ remains surjective.

Proof. Let $S_{V^{*}}$ denote the sphere of radius 1 in a norm $\|\cdot\|$ in $V^{*}$. The map $\sigma: H \times\left(E^{\perp} \backslash\{0\}\right) \rightarrow S_{V^{*}},(h, w) \mapsto \frac{h^{*} \cdot w}{\left\|h^{*} \cdot w\right\|}$ is continuous and open, since it is the composition of the submersions $(h, w) \mapsto h^{*} \cdot w$ and $v \mapsto \frac{v}{\|v\|}$. For every $\xi=\sigma(h, w) \in S_{V^{*}}$ let $U_{h}$ be an open, relatively compact neighborhood of $h$ in $H$, $W_{h}:=U_{h} \times\left(E^{\perp} \backslash\{0\}\right)$. Since $S_{V^{*}}$ is compact, it is covered by finitely many open sets $\sigma\left(W_{h_{1}}\right), \ldots, \sigma\left(W_{h_{k}}\right)$. The subset $K:=\bar{A}, A:=\cup_{i=1}^{k} U_{h_{i}}$, has the required property.

Let $\dot{K}$ be the projection of $K$ on ${ }^{H} / H_{E}$. There is a finite open covering of $\dot{K}$ by charts $\left(U_{i}, \varphi_{i}\right), 1 \leq i \leq q$, with the property that there exists on every $U_{i}$ a cross section $\sigma_{i}$ into $H$. The sets $U_{i}$ can be considered as pairwise disjoint open sets in $\mathbb{R}^{k}\left(k=\operatorname{dim}{ }^{H} / H_{E}\right)$ (if we identify $U_{i}$ with $\varphi_{i}\left(U_{i}\right)$ and take $\varphi_{i}\left(U_{i}\right)$ bounded, the $\varphi_{i}\left(U_{i}\right)$ can be disjointly distributed in $\mathbb{R}^{p}$ by the use of translations). ${ }^{3}$

Let now $\omega:=\left(\cup_{i=1}^{q} U_{i}\right) \times E^{\perp}, \Omega:=V^{*}, \theta: \omega \rightarrow V^{*}, \theta(\dot{h}, w):=\sigma(\dot{h})^{*} \cdot w$, $\left.\sigma\right|_{U_{i}}:=\sigma_{i} .{ }^{4}$ The map $\theta$ is surjective and satisfies $\Theta_{2}$ and $\Theta_{3}$ by assumption. Our manifolds has a real analytic atlas, and the cross sections $\sigma_{i}$ as well as the representation $\rho$ can be taken analytic. This establishes $\Theta_{1}$. The next lemma establishes $\Theta_{4}$.

[^2]Lemma 5.5. For every compact set $K^{\prime} \subseteq V^{*}$ there exists a compact set $k \subseteq \omega$ such that $\theta(k)=K^{\prime}$.

Proof. Similar considerations as those in the proof of lemma 5.4 lead to a compact set $\tilde{K} \subseteq \cup_{i=1}^{q} U_{i}$ with the property that the restriction of $\theta$ on $\tilde{K} \times E^{\perp}$ is surjective. For $\dot{h} \in \tilde{K}$ we put

$$
K_{\dot{h}}:=\left\{w \in E^{\perp} \mid \theta(\dot{h}, w) \in K^{\prime}\right\} .
$$

Obviously, $K_{\dot{h}}$ is compact. We finish the proof by verifying the compactness of $k:=\cup_{\dot{h} \in \tilde{K}}\{\dot{h}\} \times K_{\dot{h}}$.

Let $\left(\left(\dot{h}_{n}, w_{n}\right)\right)_{n \in \mathbb{N}}$ be a sequence in $k$ and suppose that $\dot{h}_{n}$ converges to $\dot{h} \in \tilde{K}$. We have:

$$
w_{n} \in K_{\dot{h}_{n}}=\left[\left(\sigma\left(\dot{h}_{n}\right)^{-1}\right)^{*} \cdot K^{\prime}\right] \cap E^{\perp} .
$$

Let $B$ be a compact neighborhood of $\dot{h}$ in $\cup_{i=1}^{q} U_{i}$. Almost all $w_{n}$ lie in the compact set $\left[\left(\sigma(B)^{-1}\right)^{*} \cdot K^{\prime}\right] \cap E^{\perp}$, so there is a subsequence of $\left(w_{n}\right)_{n \in \mathbb{N}}$ which converges to $w \in\left[\left(\sigma(B)^{-1}\right)^{*} \cdot K^{\prime}\right] \cap E^{\perp}$. Therefore there exists $\dot{h}_{B} \in B$ such that $\sigma\left(\dot{h}_{B}\right)^{*} \cdot w \in K^{\prime}$. If we let $B$ shrink to $\dot{h}$, we obtain a sequence $\left(\dot{h}_{B_{n}}\right)_{n \in \mathbb{N}}$ converging to $\dot{h}$ and having the property that $\sigma\left(\dot{h}_{B_{n}}\right)^{*} \cdot w=\theta\left(\dot{h}_{B_{n}}, w\right) \in K^{\prime}$. This implies $\theta(\dot{h}, w) \in K^{\prime}$, so $(\dot{h}, w) \in k$.

To apply Glaeser's above theorem to the function $\Phi$ on $\omega$, defined by

$$
\begin{equation*}
\Phi(\dot{h}, w):=\tilde{\varphi}(\dot{h}, \theta(\dot{h}, w)) \tag{11}
\end{equation*}
$$

we need to show that $\Phi$ is bipunctually attached to $A_{\theta}$. Taking into account the goal of the first step, we only have to show this for the subset $\left\{(\dot{h}, 0) \mid \dot{h} \in \cup_{i=1}^{q} U_{i}\right\}$ of $\omega$. We first prove an algebraic lemma.

Lemma 5.6. Let $P_{k}$ be a homogeneous polynomial of degree $k \geq 0$ on $V^{*}$. There is a unique $k$-linear symmetric form $T_{k}$ on $V^{*}$ such that $T_{k}(x, \ldots, x)=$ $P_{k}(x)$ for all $x \in V^{*}$.

Proof. Let $\left(x^{1}, \ldots, x^{n}\right)$ denote the vector of coordinates of $x \in V^{*}$ with respect to a fixed basis. The proof follows easily, if we consider the general expression of a symmetric $k$-linear form:

$$
T_{k}\left(x_{1}, \ldots, x_{k}\right)=\sum_{1 \leq i_{1} \leq \ldots \leq i_{k} \leq n} \frac{1}{k!} c_{i_{1} \ldots i_{k}} \sum_{\sigma \in S_{k}} x_{\sigma(1)}^{i_{1}} \ldots x_{\sigma(k)}^{i_{k}}, \quad c_{i_{1} \ldots i_{k}} \in \mathbb{C}
$$

( $S_{k}$ the symmetric group of $k$ elements).

Proposition 5.7. The function $\Phi$ defined by (11) on $\omega$ is attached to $A_{\theta}$ on $\left\{(\dot{h}, 0) \mid \dot{h} \in \cup_{i=1}^{q} U_{i}\right\}$.

Proof. Let $\left(u_{1}, \ldots, u_{n-d}\right)$ be the vector of coordinates with respect to a basis of $U,\left(w_{1}, \ldots, w_{n-d}\right)$ the one with respect to the dual basis of $E^{\perp}$. We have:

$$
\begin{aligned}
\frac{\partial}{\partial w_{j}^{\prime}} & \Phi\left(\dot{h}, w^{\prime}\right)=\frac{\partial}{\partial w_{j}^{\prime}} \int_{U} \varphi(\dot{h}, \sigma(\dot{h}) \cdot u) e^{-i\left\langle u, w^{\prime}\right\rangle} d u \\
& =-i \cdot \int_{U} \varphi(\dot{h}, \sigma(\dot{h}) \cdot u) u_{j} e^{-i\left\langle u, w^{\prime}\right\rangle} d u
\end{aligned}
$$

The moment conditions for $\varphi$ (proposition 4.4) imply:

$$
\begin{equation*}
\left.\left(\sum_{j=1}^{n-d} w_{j} \frac{\partial}{\partial w_{j}^{\prime}}\right)^{m} \Phi\left(\dot{h}, w^{\prime}\right)\right|_{w^{\prime}=0}=(-i)^{m} P_{m}\left(\sigma(\dot{h})^{*} \cdot w\right) \tag{12}
\end{equation*}
$$

( $P_{m}$ a homogeneous polynomial of degree $m$ on $V^{*}$ ).
For every $m \in \mathbb{Z}_{+}$let $T_{m}$ be the unique $m$-linear symmetric form on $V^{*}$ such that $T_{m}(x, \ldots, x)=(-i)^{m} P_{m}(x)$ (according to the preceding lemma). By the generalized theorem of Emile Borel every (formal) power series is the Taylor series of a $C^{\infty}$-function. ${ }^{5}$ Therefore there exists a function $F^{\prime} \in C^{\infty}\left(V^{*}\right)$ such that $D^{m} F^{\prime}(0)=T_{m}, m \in \mathbb{Z}_{+}$. We shall show that $F^{\prime} \circ \theta$ is an interpolating function for $\Phi$ on $\left\{(\dot{h}, 0) \mid \dot{h} \in \cup_{i=1}^{q} U_{i}\right\}$.

If we set

$$
\left(\dot{h}^{*} T_{m}\right)\left(w_{1}, \ldots, w_{m}\right):=T_{m}\left(\sigma(\dot{h})^{*} \cdot w_{1}, \ldots, \sigma(\dot{h})^{*} \cdot w_{m}\right)
$$

and denote by $D_{2}^{m} \Phi(\dot{h}, 0)$ the $m$-th derivative of $w \mapsto \Phi(\dot{h}, w)$ at 0 , (12) takes the form

$$
D_{2}^{m} \Phi(\dot{h}, 0)(w, \ldots, w)=\left(\dot{h}^{*} T_{m}\right)(w, \ldots, w)
$$

Since both sides are symmetric tensors, $D_{2}^{m} \Phi(\dot{h}, 0)=\dot{h}^{*} T_{m}$ by the preceding lemma. On the other hand, $D_{2}^{m}\left(F^{\prime} \circ \theta\right)(\dot{h}, 0)=\dot{h}^{*} T_{m}$, the left hand side denoting the $m$-th derivative of $w \mapsto F^{\prime} \circ \theta(\dot{h}, w)=F^{\prime}\left(\sigma(\dot{h})^{*} \cdot w\right)$. Hence,

$$
D_{2}^{m} \Phi(\dot{h}, 0)=D_{2}^{m}\left(F^{\prime} \circ \theta\right)(\dot{h}, 0), \quad m \in \mathbb{Z}_{+} .
$$

Since this last equation holds for all $\dot{h} \in \cup_{i=1}^{q} U_{i}$, we can apply on both sides a differential operator of constant coefficients in $\dot{h}\left(\cup_{i=1}^{q} U_{i} \subseteq \mathbb{R}^{p}\right)$ without doing any harm to the equality. This shows that $F^{\prime} \circ \theta$ is an interpolating function for $\Phi$.

By Glaeser's theorem there now exists $F^{\prime} \in C^{\infty}\left(V^{*}\right)$ such that $\Phi$ and $F^{\prime} \circ \theta$ coincide on $\omega$. Therefore, the function $F$ of the first step can be extended to the $C^{\infty}$-function $F^{\prime}$ on $V^{*}$. We shall continue to write $F$ instead of $F^{\prime}$. By continuity we have $\tilde{\varphi}(\dot{h}, \xi)=F(\xi)$ for all $\xi \in h^{*} \cdot E^{\perp}, \dot{h} \in{ }^{H} / H_{E}$.

## Third step

Proposition 5.8. The function $F$ defined above belongs to the Schwartz class $\mathcal{S}\left(V^{*}\right)$.

[^3]Proof. It remains to show that for every $p \in \mathbb{Z}_{+}$and every differential operator $D$ with constant coefficients on $V^{*}$

$$
\sup _{\xi \in V^{*}}\|\xi\|^{p}|D F(\xi)|<\infty
$$

Let $K \subseteq H$ be as in lemma 5.4, $\dot{K}$ the projection of $K$ on ${ }^{H} / H_{E}$ and $\left(U_{i}\right)_{1 \leq i \leq q}$ a covering of $\dot{K}$ by open, relatively compact sets, which admit local cross sections $\sigma_{i}$ on their topological closures. Thus we have to show that

$$
\begin{equation*}
\sup _{\dot{h} \in \dot{K}, w \in E^{\perp}}\left\|\sigma_{i}(\dot{h})^{*} \cdot w\right\|^{p} \cdot\left|D F\left(\sigma_{i}(\dot{h})^{*} \cdot w\right)\right|<\infty \tag{13}
\end{equation*}
$$

(if $\dot{h}$ is contained in $U_{i}$ ). For this we shall express $D$ in coordinates $(\dot{h}, w)$.
Let us take a basis of $E^{\perp}$, extend it to one of $V^{*}$ and denote by $\left(w_{1}, \ldots, w_{n-d}\right)$ and $\left(\xi_{1}, \ldots, \xi_{n}\right)$ the vectors of coordinates on $E^{\perp}$ and $V^{*}$ respectively. The automorphisms $\rho^{*}\left(\sigma_{i}(\dot{h})\right)$ can then be regarded as $n \times n$-matrices. Putting $\theta_{i}: U_{i} \times E^{\perp} \rightarrow V^{*}, \theta_{i}(\hat{h}, w):=\sigma_{i}(\hat{h})^{*} \cdot w, 1 \leq i \leq q$, we have $\left(X_{1}, \ldots, X_{m}\right.$ a basis of the Lie algebra $\mathfrak{h}$ of $H$ ):

$$
\begin{gathered}
\frac{\partial}{\partial w_{j}} F \circ \theta_{i}(\dot{h}, w)=\sum_{l=1}^{n} \frac{\partial F}{\partial \xi_{l}}\left(\sigma_{i}(\dot{h})^{*} \cdot w\right) \cdot \rho^{*}\left(\sigma_{i}(\dot{h})\right)_{l j}, \quad 1 \leq j \leq n-d, \\
X_{k}\left(F \circ \theta_{i}\right)(\dot{h}, w):=\left.\frac{d}{d t} F \circ \theta_{i}\left(\exp \left(t X_{k}\right) \dot{h}, w\right)\right|_{t=0} \\
=\left.\sum_{l=1}^{n} \frac{\partial F}{\partial \xi_{l}}\left(\sigma_{i}(\dot{h})^{*} \cdot w\right) \cdot \frac{d}{d t}\left[\rho^{*}\left(\sigma_{i}\left(\exp \left(t X_{k}\right) \dot{h}\right)\right) w\right]_{l}\right|_{t=0}, \quad 1 \leq k \leq m .
\end{gathered}
$$

We consider these equations as a linear system in $\frac{\partial F}{\partial \xi_{1}}, \ldots, \frac{\partial F}{\partial \xi_{n}}$. For every $\left(\dot{h}_{0}, w_{0}\right) \in$ $\dot{K} \times\left(E^{\perp} \backslash\{0\}\right)$ there is an index $i$ such that $\dot{h}_{0} \in U_{i}$, and we can choose $n$ of the above equations such that the determinant of the coefficients of this part of the system does not vanish $\left(\theta_{i}\right.$ is on $U_{i} \times\left(E^{\perp} \backslash\{0\}\right)$ a submersion $)$. Since this determinant is a homogeneous polynomial in $w$, it remains non-zero on a neighborhood of $\left(\dot{h}_{0}, w_{0}\right)$ of the form $O \times \mathbb{R}^{*} \cdot W$, where $O$ is a neighborhood of $\dot{h}_{0}$ in $U_{i}$ and $W$ one of $\frac{w_{0}}{\left\|w_{0}\right\|}$ in the sphere $S_{E^{\perp}}$ of radius 1 in $E^{\perp}$. Using Cramer's rule, we obtain:

$$
\frac{\partial F}{\partial \xi_{l}}\left(\sigma_{i}(\dot{h})^{*} \cdot w\right)=\sum_{j=1}^{n-d} f_{j}(\dot{h}, w) \frac{\partial}{\partial w_{j}} F \circ \theta_{i}(\dot{h}, w)+\sum_{k=1}^{m} g_{k}(\dot{h}, w) X_{k}\left(F \circ \theta_{i}\right)(\dot{h}, w),
$$

where $f_{j}, g_{k}$ are smooth homogeneous rational functions in $w\left(f_{j}\right.$ of degree $0, g_{k}$ of degree -1 ) (take $f_{j}, g_{k}$ to be zero, if a term does not appear). Repeating the procedure with $\frac{\partial F}{\partial \xi_{l}}$ instead of $F$, we obtain an expression for the second derivatives of $F$, etc. Eventually this leads to an expression of the operator $D$ in coordinates $(\dot{h}, w)$ :

$$
D F\left(\sigma_{i}(\dot{h})^{*} \cdot w\right)=\sum_{j} e_{j}(\dot{h}, w) T_{j} P_{j}\left(\frac{\partial}{\partial w_{1}}, \ldots, \frac{\partial}{\partial w_{n-d}}\right) F \circ \theta_{i}(\dot{h}, w)
$$

$$
\begin{equation*}
=\sum_{j} e_{j}(\dot{h}, w) T_{j} P_{j}\left(\frac{\partial}{\partial w_{1}}, \ldots, \frac{\partial}{\partial w_{n-d}}\right) \tilde{\varphi}\left(\dot{h}, \sigma_{i}(\dot{h})^{*} \cdot w\right) . \tag{14}
\end{equation*}
$$

Therein $T_{j}$ is a differential operator on ${ }^{H} / H_{E}, P_{j}\left(\frac{\partial}{\partial w_{1}}, \ldots, \frac{\partial}{\partial w_{n-d}}\right)$ a polynomial in $\frac{\partial}{\partial w_{1}}, \ldots, \frac{\partial}{\partial w_{n-d}}$, and $e_{j}$ smooth rational functions, homogeneous in $w$ of degree $\leq 0$.

Relation (13) is a property for $w$ in a neighborhood of infinity. For $\|w\| \leq 1$ we actually have $\left\|\sigma_{i}(\dot{h})^{*} \cdot w\right\| \leq\left\|\rho^{*}\left(\sigma_{i}(\dot{h})\right)\right\|$, and

$$
A_{i}:=\sup _{\dot{h} \in U_{i}}\left\|\rho^{*}\left(\sigma_{i}(\dot{h})\right)\right\|<\infty
$$

since $U_{i}$ is relatively compact and $\sigma_{i}$ a local cross section on $\bar{U}_{i}$. The region of $V^{*}$ we obtain by $\sigma_{i}(\dot{h})^{*} \cdot w(1 \leq i \leq q)$ from $\dot{h} \in \dot{K},\|w\| \geq 1$, contains $\left\{\xi \in V^{*} \mid\|\xi\|>\max _{1 \leq i \leq q} A_{i}\right\}$. Therefore we can restrict ourselves to $\|w\| \geq 1$ in (13).

Equation (14) holds for $(\dot{h}, w) \in O \times \mathbb{R}^{*} \cdot W$. If we take $O$ and $W$ to be compact, all functions $e_{j}$ are bounded on $O \times(\mathbb{R} \backslash]-1,1[) \cdot W$, and (13) holds for $(\dot{h}, w) \in O \times(\mathbb{R} \backslash]-1,1[) \cdot W$, because $\tilde{\varphi}$ is a rapidly decreasing function (cf. section 2). Due to the compactness of $K$ and $S_{E \perp}$, we now only need finitely many sets of the form $O \times(\mathbb{R} \backslash]-1,1[) \cdot W$ to obtain a full neighborhood of the infinity in $V^{*}$. The proof is thereby completed.

## 6. The image in the case $d<n-1$

Richter and Kurusa have shown that in the classical case ( $H=S O(n), V=\mathbb{R}^{n}$ ) and for $d<n-1$ the range-characterizing partial differential equations (which are equivalent to $d \nu(T) \varphi=0$ for $T \in \mathcal{K}^{1}\left(s o(n) \times{ }_{s} \mathbb{R}^{n}\right)$, see [14, page 72 ff$]$ ) suffice to characterize the range $\mathcal{S}\left(\mathbb{R}^{n}\right)^{\wedge}$ (the latter author derives directly the moment conditions of Helgason) [11], [13]. Gonzalez has given a group theoretical proof [6]. The fact is also true in the general situation, provided that the set of planes in $\mathcal{E}$ (actually in $\mathcal{E}^{*}$ ) is sufficiently rich.

Definition 6.1. A collection $\mathcal{U}$ of $k$-dimensional linear subspaces of $V^{*}$ will be called fundamental system if every homogeneous polynomial on $V^{*}$ is completely determined by its restriction on the union of all elements of $\mathcal{U}$.

For instance, a set of hyperplanes $(k=n-1)$ is a fundamental system iff it contains infinitely many elements.

We proceed to the following assumption (denoted in the sequel by $(G)$ ) on the action of the group $H$ :

There is a dense set $A \subseteq V^{*} \backslash\{0\}$ with the following property: For every $\xi \in A$ the planes of the form $h^{*} \cdot E^{\perp}$ through $\xi$ form a fundamental system.

The set of $\xi \in V^{*} \backslash\{0\}$ for which the statement in $(G)$ holds is $\rho^{*}(H)$ invariant. Therefore we only have to establish $(G)$ for a dense set in $E^{\perp} \backslash\{0\}$. This makes things much easier in concrete situations.

Proposition 6.2. Assumption $(G)$ implies: $\mathcal{S}(V)^{\wedge}=\mathcal{S}_{D}(\mathcal{E})$.
Proof. Let $\varphi \in \mathcal{S}_{D}(\mathcal{E}), \tilde{\varphi}$ the Fourier transform. We denote by $D_{2}^{m} \tilde{\varphi}\left(\dot{h}, \xi_{0}\right)$ the $m$-th derivative of $\xi \mapsto \tilde{\varphi}(h, \xi)$ (as a function on $h^{*} \cdot E^{\perp}$ ) at the point $\xi_{0} \in h^{*} \cdot E^{\perp}$. The proof will be given by establishing the moment conditions for $\varphi$. These are equivalent to

$$
D_{2}^{m} \tilde{\varphi}(\dot{h}, 0) \cdot \xi^{\otimes^{m}}=(-i)^{m} P_{m}(\xi), \quad \xi \in h^{*} \cdot E^{\perp}
$$

$P_{m}$ homogeneous polynomials of degree $m \geq 0$ on $V^{*}$ (cf. proposition 4.4).
Let $Y_{j}(\xi), 1 \leq j \leq l$, be a basis of $\mathfrak{h}_{\xi}$ for $\xi$ in a dense open subset of $V^{*} \backslash\{0\}$ (notation as in section 5). It follows from the considerations in section 5 that for $1 \leq j \leq l$ the function

$$
(t, s) \longmapsto \tilde{\varphi}\left(\exp \left(-t Y_{j}(\xi)\right) \dot{h}, s \xi\right)
$$

does not depend on $t$, nor does the $m$-th partial derivative with respect to $s$ at $s=0$ :

$$
D_{2}^{m} \tilde{\varphi}\left(\exp \left(-t Y_{j}(\xi)\right) \dot{h}, 0\right) \cdot \xi^{\otimes^{m}}
$$

Following the reasoning of section 5 , there exists a function $F_{m} \in C^{\infty}\left(V^{*} \backslash\{0\}\right)$ such that

$$
\begin{equation*}
D_{2}^{m} \tilde{\varphi}(\dot{h}, 0) \cdot \xi^{\otimes^{m}}=F_{m}(\xi) \tag{15}
\end{equation*}
$$

for $\xi \in h^{*} \cdot E^{\perp}, \xi \neq 0$. Let $\xi_{0} \in A, A$ as in $(G)$. The map

$$
\xi \longmapsto D^{k} F_{m}\left(\xi_{0}\right) \cdot \xi^{\otimes^{k}}, \quad k>m,
$$

( $D^{k}$ : the $k$-th derivative) is a homogeneous polynomial, whose restriction on every plane $h^{*} \cdot E^{\perp}$ containing $\xi_{0}$ is equal to zero by (15). The fact $\xi_{0} \in A$ implies $D^{k} F_{m}\left(\xi_{0}\right) \cdot \xi^{\otimes^{k}} \equiv 0$, and by lemma $5.6 D^{k} F_{m}\left(\xi_{0}\right)=0$. Since $A$ is dense, we obtain $D^{k} F_{m}=0$ for $k>m$. Therefore $F_{m}$ equals a polynomial $(-i)^{m} P_{m}$ of degree $\leq m$ on $V^{*} \backslash\{0\}$ and can be extended to it on $V^{*}$. By continuity,

$$
D_{2}^{m} \tilde{\varphi}(\dot{h}, 0) \cdot \xi^{\otimes^{m}}=(-i)^{m} P_{m}(\xi), \quad \xi \in h^{*} \cdot E^{\perp}
$$

Hereby we conclude that $P_{m}$ is homogeneous of degree $m$, and the proof is completed.

However, assumption $(G)$ is not necessary for a relation of the form $\mathcal{S}(V)^{\wedge}=\mathcal{S}_{D}(\mathcal{E})$, as the last example in section 9 shows.

## 7. The Radon transform $(d=n-1)$

In this section we focus our attention on the special case $d=n-1$. As it will come out, this leads to the classical Radon transform. Different situations can only appear, if we remove some of our conditions on the group action, or if we completely change the setting and allow for $\mathcal{E}$ sets of planes which do not admit a transitive group action. For the latter possibility the reader is referred to [2].

Let $H$ be a Lie group, $\rho$ a representation of $H$ on a (finite-dimensional) vector space $V$. The contragredient representation will be denoted by $\rho^{*}$. We fix a linear subspace $E \subseteq V$ of codimension 1 and denote by $E^{\perp} \subseteq V^{*}$ its (onedimensional) orthogonal complement. We make the following assumptions:

1. $\rho(H)$ consists of unimodular automorphisms of $V$, and the same holds for $\rho\left(H_{E}\right)$ (the isotropy subgroup of $E$ ) in the case of $E$.
2. The map from $H \times E^{\perp}$ into $V^{*},(h, w) \mapsto \rho^{*}(h) w$, is a surjective submersion. ${ }^{6}$

We put $h \cdot u:=\rho(h) u, h^{*} \cdot w:=\rho^{*}(h) w$. Let $\|\cdot\|$ be a norm on $V^{*}, S_{V^{*}}$ the sphere of radius 1 . The group $H$ acts transitively on $S_{V^{*}}$ :

$$
h \cdot w:=\frac{h^{*} \cdot w}{\left\|h^{*} \cdot w\right\|} .
$$

Let $w_{0} \in S_{V^{*}} \cap E^{\perp}, H_{w_{0}}^{S}$ the isotropy subgroup of $w_{0}$. Since $H_{E}$ operates by unimodular automorphisms on $E^{\perp}$, we obtain $\stackrel{o}{H}_{E} \subseteq H_{w_{0}}^{S} \subseteq H_{E}$ ( $\stackrel{o}{H}_{E}$ denoting the connected component of $H_{E}$ ) and $H_{w_{0}}^{S}=H_{w_{0}}$ (the isotropy subgroup in the contragredient action). The transitivity of the action on $S_{V^{*}}$ implies $H_{\xi}^{S}=H_{\xi}$ for every $\xi \in S_{V^{*}}$. Thus, ${ }^{H} / H_{\xi} \cong S_{V^{*}}$ for $\xi \in V^{*} \backslash\{0\}$. The orbit $\rho^{*}(H) \xi$ is bounded because ${ }^{H} / H_{\xi}$ is compact. Since this holds for every $\xi, \rho^{*}(H)$ itself is bounded (in the vector space $L\left(V^{*}\right)$ of all endomorphisms of $V^{*}$ ). Therefore, $\operatorname{det}\left[\rho^{*}(H)\right]$ is a bounded subgroup of $\mathbb{R}^{*}$, so $\operatorname{det}\left[\rho^{*}(H)\right] \subseteq\{-1,1\}$. This relation also holds for the closure $\overline{\rho^{*}(H)}$ in $L\left(V^{*}\right)$. Thus, $\overline{\rho^{*}(H)}$ is a compact subgroup of $G L\left(V^{*}\right)$ and hence leaves a positive definite inner product invariant. With the help of it, $V^{*}$ can be identified with $V^{* *} \cong V$, and we have $\rho^{*}=\rho$. Eventually, $V \cong V^{*}$ is identified with $\mathbb{R}^{n}$ with the help of an orthonormal basis. We obtain:
$\rho(H)$ is a group of orthogonal matrices and operates transitively on the sphere $S^{n-1}$. It holds $\rho=\rho^{*}$, and $E \subseteq \mathbb{R}^{n}$ is a subspace of codimension 1 .

This is the situation in the classical case.
As for the range of $\mathcal{S}(V)$ under $f \mapsto \hat{f}$, it is completely characterized by the moment conditions (see [9, page 100]):

Proposition 7.1. In the case $d=n-1$ we have: $\mathcal{S}(V)^{\wedge}=\mathcal{S}_{M}(\mathcal{E})$.
Proof. Since $\xi \in h \cdot E^{\perp}, \xi \neq 0$, completely determines $\dot{h}=h H_{E}\left(h H_{E} \Leftrightarrow\right.$ $h \cdot E=\xi^{\perp}$ ), we deduce at once the result of the first step in section 5. The differential operators are not needed for the second and third step. ${ }^{7}$

## 8. The range of $C_{c}^{\infty}(V)$

The characterization of the range of $C_{c}^{\infty}(V)$, the set of $C^{\infty}$-functions with compact support, under the $d$-plane transform will be done with the help of a modified version of the Paley-Wiener theorem, due to Helgason [8], which is recalled first. Condition (ii) in the beginning of section 5 can be dropped throughout this section.

Let $\mathcal{D}_{A}\left(\mathbb{R}^{n}\right)$ denote the class of $C^{\infty}$-functions on $\mathbb{R}^{n}$ with support in $\{x||x| \leq A\}$.

[^4]Theorem. The Fourier transform $f \mapsto \tilde{f}$ maps $\mathcal{D}_{A}\left(\mathbb{R}^{n}\right)$ onto the set of functions $\tilde{f}(\lambda \omega)=\varphi(\lambda, \omega) \in C^{\infty}\left(\mathbb{R} \times S^{n-1}\right)$ satisfying:
(i) For each $\omega$, the function $\lambda \mapsto \varphi(\lambda, \omega)$ extends to a holomorphic function on $\mathbb{C}$ such that

$$
\sup _{\lambda, \omega}\left|\varphi(\lambda, \omega) \lambda^{N} e^{-A|I m \lambda|}\right|<\infty
$$

for each integer $N \in \mathbb{Z}_{+}$.
(ii) For each $k \in \mathbb{Z}_{+}$and every vector $a=\left(a_{1}, \ldots, a_{n}\right) \in \mathbb{C}^{n}$ satisfying $\sum_{i=1}^{n} a_{i}^{2}=$ 0 the function

$$
\lambda \mapsto \lambda^{-k} \int_{S^{n-1}} \varphi(\lambda, \omega)\langle a, \omega\rangle^{k} d \omega
$$

is even and holomorphic on $\mathbb{C}$.
For a proof we refer to [9, page 23ff].
Before we state the main theorem, we need some facts about spherical harmonics.

Proposition . The eigenspaces of the Laplacian on $S^{n-1}$ are of the form

$$
E_{k}=\operatorname{span}\left\{f_{a, k} \mid f_{a, k}(\omega)=\langle a, \omega\rangle^{k}, \omega \in S^{n-1}\right\},
$$

where $a=\left(a_{1}, \ldots, a_{n}\right) \in \mathbb{C}^{n}$ satisfies $\sum_{i=1}^{n} a_{i}^{2}=0$, and $k \in \mathbb{Z}_{+}$. The subspaces $E_{k}$ are pairwise orthogonal in $L^{2}\left(S^{n-1}\right)$. If $P_{k}$ denotes the set of restrictions on $S^{n-1}$ of the homogeneous polynomials of degree $k$ on $\mathbb{R}^{n}$, it holds

$$
P_{k}=E_{k}+E_{k-2}+\ldots+E_{k-2 m}, \quad m=\left[\frac{k}{2}\right] .
$$

For a proof see [9, page 17 ff$]$.
Let $\langle\cdot, \cdot\rangle$ be a (positive definite) inner product on $V$ such that $E$ and $U$ are orthogonal. $S_{V}$ stands for the sphere of radius 1 in the induced norm. For $h \in H$ we put

$$
\tau(h):=\left\|\rho\left(h^{-1}\right)\right\|,
$$

$\rho$ being the representation of $H$ on $V$ and $\|\cdot\|$ the operator norm. The map $\tau$ is continuous. For $u \in U,\|u\| \geq \tau(h), h \cdot(u+E)$ has at most one point with $S_{V}$ in common. If $\operatorname{supp} f \subseteq\{x \in V \mid\|x\| \leq A\}$, we therefore have $\hat{f}(\dot{h}, h \cdot u)=0$ for $\|u\| \geq A \tau(h)$.

Theorem 8.1. $\varphi \in \mathcal{S}_{M}(\mathcal{E})$ is the d-plane transform of a $C^{\infty}$-function $f$ of compact support on $V$ iff there exists a continuous function $c$ on $H$ such that for every $h \in H \quad \varphi(\dot{h}, h \cdot u)=0$ for $\|u\| \geq c(h)$.

Proof. We only have to show that the condition is sufficient. We calculate:

$$
\tilde{\varphi}(\dot{h}, \xi)=\int_{U} \varphi(\dot{h}, h \cdot u) e^{-i\langle h \cdot u, \xi\rangle} d u=\int_{U} \varphi(\dot{h}, h \cdot u) \sum_{m=0}^{\infty} \frac{(-i)^{m}}{m!}\langle h \cdot u, \xi\rangle^{m} d u
$$

$$
=\sum_{m=0}^{\infty} \frac{(-i)^{m}}{m!} \int_{U} \varphi(\dot{h}, h \cdot u)\langle h \cdot u, \xi\rangle^{m} d u
$$

(the region of integration is bounded)

$$
=\sum_{m=0}^{\infty} \frac{(-i)^{m}}{m!} P_{m}(\xi)
$$

by the moment conditions. The last series converges for all $\xi \in h^{*} \cdot E^{\perp}$, and, since it is a power series, for all $\xi \in V^{*}$. It defines an analytic function $F$.

The spaces $V$ and $V^{*}$ can be identified with the help of the inner product. We choose an orthonormal basis and identify both spaces with $\mathbb{R}^{n}$.

$$
F(\xi)=\sum_{m=0}^{\infty} \frac{(-i)^{m}}{m!} P_{m}(\xi)
$$

is then an analytic function on $\mathbb{R}^{n}$ and can be extended to a holomorphic function on $\mathbb{C}^{n}$ by the above power series.

Let $F(\lambda \omega)=: \psi(\lambda, \omega) \in C^{\infty}\left(\mathbb{R} \times S^{n-1}\right)$. We show that $\psi$ satisfies the two conditions of the theorem in the beginning.

Ad (i): It is enough to verify the relation for even exponents $N$. We calculate $\left(\omega=h^{*} \cdot w \in S^{n-1}, w \in U\right)$ :

$$
\begin{gathered}
\left|\psi(\lambda, \omega) \lambda^{2 N}\right|=\left|\tilde{\varphi}(\dot{h}, \lambda \omega) \lambda^{2 N}\right|=\left|\int_{U} \varphi(\dot{h}, h \cdot u) \lambda^{2 N} e^{-i\left\langle h \cdot u, \lambda h^{*} \cdot w\right\rangle} d u\right| \\
=\left|\int_{U} \varphi(\dot{h}, h \cdot u) \lambda^{2 N} e^{-i \lambda\langle u, w\rangle} d u\right| \\
=\left|\int_{U} \varphi(\dot{h}, h \cdot u)\left(-\sum_{k=1}^{n-d} \frac{\partial^{2}}{\partial u_{k}^{2}}\right)^{N} e^{-i \lambda\langle u, w\rangle} \cdot\|w\|^{-2 N} d u\right|
\end{gathered}
$$

(taking such an orthonormal basis that the first $n-d$ vectors belong to $U$ )

$$
\begin{aligned}
&=\|w\|^{-2 N}\left|\int_{U}\left(-\sum_{k=1}^{n-d} \frac{\partial^{2}}{\partial u_{k}^{2}}\right)^{N} \varphi(\dot{h}, h \cdot u) e^{-i \lambda\langle u, w\rangle} d u\right| \\
& \leq\|w\|^{-2 N} \int_{U}\left|\left(\sum_{k=1}^{n-d} \frac{\partial^{2}}{\partial u_{k}^{2}}\right)^{N} \varphi(\dot{h}, h \cdot u)\right| d u \cdot e^{c(h)\|w\| I m \lambda \mid} .
\end{aligned}
$$

For $\tau^{*}(h):=\left\|\rho^{*}\left(h^{-1}\right)\right\|$ we have: $\frac{1}{\tau^{*}\left(h^{-1}\right)} \leq\|w\| \leq \tau^{*}(h)$, so

$$
\left|\psi(\lambda, \omega) \lambda^{2 N}\right| \leq \tau^{*}\left(h^{-1}\right)^{2 N} \int_{U}\left|\left(\sum_{k=1}^{n-d} \frac{\partial^{2}}{\partial u_{k}^{2}}\right)^{N} \varphi(\dot{h}, h \cdot u)\right| d u \cdot e^{c(h) \tau^{*}(h)|I m \lambda|}
$$

Let $K \subseteq H$ be a compact subset with the property of lemma 5.4. Putting

$$
C_{N}:=\sup _{h \in K} \tau^{*}\left(h^{-1}\right)^{2 N} \int_{U}\left|\left(\sum_{k=1}^{n-d} \frac{\partial^{2}}{\partial u_{k}^{2}}\right)^{N} \varphi(\dot{h}, h \cdot u)\right| d u, \quad A:=\sup _{h \in K} c(h) \tau^{*}(h)
$$

we obtain the desired relation

$$
\left|\psi(\lambda, \omega) \lambda^{2 N}\right| \leq C_{N} e^{A|I m \lambda|}
$$

Ad (ii): Since $\sum_{m=0}^{\infty} \frac{(-i)^{m}}{m!} P_{m}(\xi)$ converges uniformly on compact sets, we have:

$$
\begin{gathered}
\int_{S^{n-1}} \psi(\lambda, \omega)\langle a, \omega\rangle^{k} d \omega=\int_{S^{n-1}} F(\lambda \omega)\langle a, \omega\rangle^{k} d \omega \\
=\int_{S^{n-1}} \sum_{m=0}^{\infty} \frac{(-i)^{m}}{m!} P_{m}(\lambda \omega)\langle a, \omega\rangle^{k} d \omega \\
=\sum_{m=0}^{\infty} \int_{S^{n-1}} \frac{(-i)^{m}}{m!} P_{m}(\lambda \omega)\langle a, \omega\rangle^{k} d \omega=\sum_{m=0}^{\infty} \frac{(-i \lambda)^{m}}{m!} \int_{S^{n-1}} P_{m}(\omega)\langle a, \omega\rangle^{k} d \omega .
\end{gathered}
$$

The proposition above implies that $\int_{S^{n-1}} P_{m}(\omega)\langle a, \omega\rangle^{k} d \omega$ vanishes, unless $m-k \in$ $\mathbb{Z}_{+}$. This verifies (ii).

Concluding, there exists a function $f \in C_{c}^{\infty}(V)$ such that

$$
\tilde{f}(\lambda \omega)=\psi(\lambda, \omega)=F(\lambda \omega)=\tilde{\varphi}(\dot{h}, \lambda \omega), \quad \omega \in h^{*} \cdot U
$$

Since on the other hand $\tilde{f}(\lambda \omega)=\tilde{\hat{f}}(\dot{h}, \lambda \omega)$, we obtain $\varphi=\hat{f}$ by the inverse Fourier transform.

Remark . The application of the modified version of the Paley-Wiener theorem in the above proof shows that the support of $f$ is contained in the ball of radius $A$, which depends on the function $c(h)$ and the compact set $K$. In the classical case the theorem reduces to theorem 2.10 in [9, page 109], since $H$ is compact. Furthermore, for functions $f$ of the Schwartz class, we obtain Helgason's support theorem [9, page 105] by taking $c(h) \equiv A$.

## 9. Examples

1. We begin by a counterexample, which shows that condition (i) in the beginning of section 5 cannot be weakened substantially.

Let

$$
H=\left\{\left.\left(\begin{array}{ccc}
1 & a & c \\
0 & 1 & b \\
0 & 0 & 1
\end{array}\right) \right\rvert\, a, b, c \in \mathbb{R}\right\}
$$

$V=\mathbb{R}^{3}, E=\langle(0,0,1)\rangle, U=\langle(1,0,0),(0,1,0)\rangle \cong E^{\perp}$ via the standard scalar product. The assumptions we have made in the general case are all satisfied except of (i), section 5 . The vector $(0,0,1)$ is not contained in the image (which is dense in $\mathbb{R}^{3}$ ) of the mapping

$$
\left(a, b, c, w_{1}, w_{2}\right) \longmapsto\left(\begin{array}{ccc}
1 & a & c \\
0 & 1 & b \\
0 & 0 & 1
\end{array}\right)^{*} \cdot\left(\begin{array}{c}
w_{1} \\
w_{2} \\
0
\end{array}\right)=\left(\begin{array}{c}
w_{1} \\
-a w_{1}+w_{2} \\
(a b-c) w_{1}-b w_{2}
\end{array}\right) .
$$

Let $g \in C_{c}^{\infty}(\mathbb{R}), g$ nowhere zero on $[-1,1]$, and define $F: \mathbb{R}^{3} \rightarrow \mathbb{R}$ by

$$
F\left(\xi_{1}, \xi_{2}, \xi_{3}\right):=g\left(\xi_{1}^{2}+\xi_{2}^{2}+\xi_{3}^{2}\right) \cdot g\left(\frac{1}{\xi_{1}^{2}+\xi_{2}^{2}+\xi_{3}^{2}}\right) \cdot \frac{1}{\xi_{1}^{2}+\xi_{2}^{2}}
$$

Then $\psi(\dot{h}, \xi):=F(\xi), \xi \in h^{*} \cdot E^{\perp}$, belongs to $\mathcal{S}\left(\mathcal{E}^{*}\right)$ and hence there exists $\varphi \in$ $\mathcal{S}(\mathcal{E})$ such that $\tilde{\varphi}=\psi$. Due to the construction of $\psi$ it holds $\varphi \in \mathcal{S}_{D}(\mathcal{E}) \cap \mathcal{S}_{M}(\mathcal{E})$. However, there is no $f \in \mathcal{S}\left(\mathbb{R}^{3}\right)$ such that $\hat{f}=\varphi$, because otherwise we would have $\tilde{f}=F$, which is a contradiction, since $F$ does not extend to a $C^{\infty}$-function on $\mathbb{R}^{3}$.
2. The action of $H=S O(n)$ on $V=\mathbb{R}^{n}$ (by matrix multiplication) leads to the classical $d$-plane transform. For the range of $\mathcal{S}\left(\mathbb{R}^{n}\right)$ the reader is referred to [13] and [6].
3. $H=S O(p, n-p)$ operates on $V=\mathbb{R}^{n}$ by matrix multiplication. We introduce the following basis on the Lie algebra $s o(p, n-p)\left(E_{i j}\right.$ denotes the matrix $\left.\left(\delta_{i k} \delta_{j l}\right)_{k, l}\right)$ :

$$
\begin{aligned}
X_{i j} & :=E_{i j}-E_{j i} \quad \text { for } \quad 1 \leq i<j \leq p, p+1 \leq i<j \leq n, \\
Y_{i j} & :=E_{i j}+E_{j i} \text { for } 1 \leq i \leq p, p+1 \leq j \leq n .
\end{aligned}
$$

We agree to write $X_{j i}$ for $-X_{i j}$ and $Y_{j i}$ for $Y_{i j}$ (in the above range of the indices). $Z_{1}, \ldots, Z_{n}$ will denote the standard basis of $\mathbb{R}^{n}$.

The Lie algebra so $(n) \times \mathbb{R}^{n}$ is isomorphic to the Lie algebra of all $(n+$ 1) $\times(n+1)$-matrices of the form

$$
\left(\begin{array}{cc}
T & Z \\
0 & 0
\end{array}\right), T \in \operatorname{so}(n), Z \in \mathbb{R}^{n} .
$$

Similarly, $s o(p, n-p) \times{ }_{s} \mathbb{R}^{n}$ is isomorphic to the Lie algebra of the $(n+1) \times(n+1)-$ matrices of the form

$$
\left(\begin{array}{ccc}
T_{1} & X & U_{1} \\
X^{T} & T_{2} & U_{2} \\
0 & 0 & 0
\end{array}\right)
$$

$T_{1} \in s o(p), T_{2} \in s o(n-p), U_{1} \in \mathbb{R}^{p}, U_{2} \in \mathbb{R}^{n-p}$ and $X$ an arbitrary $p \times(n-p)$ matrix. The complexification of $s o(n) \times{ }_{s} \mathbb{R}^{n}$ is isomorphic to the complexification of $s o(p, n-p) \times s \mathbb{R}^{n}$ by the mapping

$$
\left(\begin{array}{rcl}
X_{1} & X_{2} & U_{1} \\
-X_{2}^{T} & X_{3} & U_{2} \\
0 & 0 & 0
\end{array}\right) \longmapsto\left(\begin{array}{rcl}
X_{1} & i X_{2} & U_{1} \\
i X_{2}^{T} & X_{3} & -i U_{2} \\
0 & 0 & 0
\end{array}\right),
$$

where $X_{1} \in \operatorname{so}(p, \mathbb{C}), X_{3} \in \operatorname{so}(n-p, \mathbb{C}), U_{1} \in \mathbb{C}^{p}, U_{2} \in \mathbb{C}^{n-p}$ and $X_{2}$ is an arbitrary $p \times(n-p)$-matrix. With the help of this isomorphism we obtain from [14, page 62 ff$]$ a basis $V_{i j l}, 1 \leq i<j<l \leq n$, of $\mathcal{K}^{1}\left(s o(p, n-p) \times{ }_{s} \mathbb{R}^{n}\right)$ :

$$
V_{i j l}=\left\{\begin{array}{c}
X_{i j} Z_{l}+X_{j l} Z_{i}+X_{l i} Z_{j} \text { for } l \leq p \text { or } i \geq p+1 \\
X_{i j} Z_{l}-Y_{j l} Z_{i}+Y_{l i} Z_{j} \text { for } j \leq p<l \\
Y_{i j} Z_{l}+X_{j l} Z_{i}-Y_{l i} Z_{j} \text { for } i \leq p<j
\end{array} .\right.
$$

The dual space of $V=\mathbb{R}^{n}$ will be identified with $\mathbb{R}^{n}$ via

$$
B(x, y):=\sum_{i=1}^{p} x_{i} y_{i}-\sum_{i=p+1}^{n} x_{i} y_{i} .
$$

Since $H$ leaves $B$ invariant, $\rho=\rho^{*}$.
Let $E:=\left\langle Z_{k}, Z_{k+1}, \ldots, Z_{k+d-1}\right\rangle, 1<k \leq k+d-1<n$, that is, $1<k<n+1-d$ (in particular, $d<n-1$ ), $U:=E^{\perp}$. Conditions (i) and (ii) in the beginning of section 5 are satisfied.

We follow [13] and introduce a system of coordinates on a neighborhood of $E$ in $\mathcal{E}$. Let $G(d, n)$ denote the set of $d$-planes in $\mathbb{R}^{n}, G_{d, n} \subseteq G(d, n)$ the subset of $d$-dimensional subspaces. For every multiindex $j=\left(j_{1}, \ldots, j_{d}\right)$, $1 \leq j_{1}<\ldots<j_{d} \leq n$, let

$$
U_{j}:=\left\{\begin{array}{l|c}
\sigma \in G_{d, n} & \begin{array}{c}
\text { the orthogonal projection of } \sigma \text { on the subspace } \\
\text { spanned by } Z_{j_{1}}, \ldots, Z_{j_{d}} \text { is bijective }
\end{array}
\end{array}\right\}
$$

Let $1 \leq j_{d+1}<\ldots<j_{n} \leq n$ be such that $\left\{j_{d+1}, \ldots, j_{n}\right\} \cup\left\{j_{1}, \ldots, j_{d}\right\}=$ $\{1, \ldots, n\}$. Let $\sigma \in U_{j}, v_{j_{\alpha}} \in \sigma$ be the inverse image of $Z_{j_{\alpha}}$ under the orthogonal projection onto $\operatorname{span}\left\{Z_{j_{1}}, \ldots, Z_{j_{d}}\right\} \quad(1 \leq \alpha \leq d)$. It holds:

$$
V_{j_{\alpha}}=Z_{j_{\alpha}}+\sum_{\kappa=d+1}^{n} \sigma_{j_{\kappa} j_{\alpha}} Z_{j_{\kappa}}
$$

with real coefficients $\sigma_{j_{\kappa} j_{\alpha}}$. In what follows, the range of $\alpha, \kappa$ will always be $[1, d],[d+1, n]$ respectively. We define the $n \times d$-matrix $\Sigma$ by taking $v_{j_{1}}, \ldots, v_{j_{d}}$ for its columns (expressed in the basis $Z_{1}, \ldots, Z_{n}$ ). The open set $U_{j}$ can now be parametrized by $\Sigma$, that is, by $\left(\sigma_{j_{\kappa} j_{\alpha}}\right)_{\kappa, \alpha}$. Since $\cup_{j} U_{j}=G_{d, n}$, we obtain in this way an atlas for $G_{d, n}$.

The space $G(d, n)$ is a vector bundle over $G_{d, n}$. Let $\pi: G(d, n) \rightarrow G_{d, n}$ denote the projection. We take $\xi \in G(d, n)$ and put $\sigma:=\pi(\xi) \in G_{d, n}$. If $\sigma \in U_{j},\{y\}:=\xi \cap\left\langle Z_{j_{d+1}}, \ldots, Z_{j_{n}}\right\rangle$, it holds: $\xi=y+\sigma$. We agree to write $y=\sum_{\kappa=d+1}^{n} y_{j_{\kappa}} Z_{j_{\kappa}}$ and parametrize $\pi^{-1}\left(U_{j}\right)$ by the vector

$$
\begin{equation*}
\left(\left(\sigma_{j_{\kappa} j_{\alpha}}\right)_{\kappa, \alpha}, y_{j_{d+1}}, \ldots, y_{j_{n}}\right) \in \mathbb{R}^{d(n-d)} \times \mathbb{R}^{n-d} \tag{16}
\end{equation*}
$$

However, not every $d$-plane in $\pi^{-1}\left(U_{j}\right)$ can be found in $\mathcal{E}$. In fact, $\pi(\mathcal{E})$ consists precisely of the subspaces $\sigma$, where the signature of $B$ equals its signature on $E$. We conclude that $\mathcal{E} \subseteq G(d, n)$ is open, so we can restrict the above parametrization to $\pi^{-1}\left(U_{j}\right) \cap \mathcal{E}$ (and obtain in this way an atlas of $\mathcal{E}$ ).

The action of $S O(p, n-p)$ on $\mathbb{R}^{n}$ possesses property $(G)$ of section 6 , and consequently the differential operators $d \nu\left(V_{i j l}\right)$ characterize the range of $\mathcal{S}\left(\mathbb{R}^{n}\right)$. However, the computation of these operators in the coordinates (16) is very lengthy and tedious (see [15, page 36-38]). Their second order terms are linear combinations of

$$
\mathcal{T}_{j_{\kappa} j_{\alpha} j_{\lambda}}:=\frac{\partial^{2}}{\partial \sigma_{j_{\kappa} j_{\alpha}} \partial y_{j_{\lambda}}}-\frac{\partial^{2}}{\partial \sigma_{j_{\lambda} j_{\alpha}} \partial y_{j_{\kappa}}},
$$

their first order ones consist of derivatives in the $y$-coordinates, and the coefficients of all terms only depend on $\sigma_{j_{\kappa} j_{\alpha}}$. Nevertheless, we can prove the following result, which is analogous to that of Richter for the classical case [14, page 73ff]:

Proposition 9.1. A function $\varphi \in \mathcal{S}(\mathcal{E})$ satisfies $d \nu\left(V_{i j l}\right) \varphi=0$ for $1 \leq i<j<$ $l \leq n$ iff for every multiindex $j=\left(j_{1}, \ldots, j_{d}\right)$ as above, $\varphi_{0}$ satisfies $\mathcal{T}_{j_{\kappa} j_{\alpha} j_{\lambda}} \varphi_{0}=0$ for $1 \leq \alpha \leq d, d+1 \leq \kappa, \lambda \leq n$ on $\pi^{-1}\left(U_{j}\right) \cap \mathcal{E}$, where $\varphi_{0}$ is a modified function, defined with the help of the $n \times d$-matrix $\Sigma$ :

$$
\varphi_{0}(\Sigma, y):=\left|\operatorname{det} \Sigma^{T} I_{p, n-p} \Sigma\right|^{-\frac{1}{2}} \varphi(\Sigma, y)
$$

( $I_{p, n-p}$ denoting the matrix of $B$ in the basis $Z_{1}, \ldots, Z_{n}$ ).
Proof. - The condition is necessary: If $\varphi \in \mathcal{S}_{D}(\mathcal{E})$, there exists $f \in \mathcal{S}\left(\mathbb{R}^{n}\right)$ such that $\hat{f}=\varphi$. By looking at the integration measure on the various $d$-planes $f$ is integrated on, we conclude that

$$
\varphi_{0}(\Sigma, y)=\int_{\mathbb{R}^{d}} f(\Sigma x+y) d x_{1} \ldots d x_{d}, \quad x=\left(x_{1}, \ldots, x_{d}\right)^{T} .
$$

The relations $\mathcal{T}_{j_{\kappa} j_{\alpha} j_{\lambda}} \varphi_{0}=0$ now follow immediately.

- The condition is sufficient: We compute the Fourier transform of $\varphi$ $(A \in S O(p, n-p), u \in U)$ :

$$
\tilde{\varphi}(A \cdot E, A u)=\int_{U} \varphi\left(A \cdot E, A u^{\prime}\right) e^{-i\left\langle A u, A u^{\prime}\right\rangle} d u^{\prime}=\int_{U} \varphi(\Sigma, y) e^{-i\langle A u, y\rangle} d u^{\prime}
$$

in the notation of (16),

$$
\begin{gathered}
=\left|\operatorname{det} \Sigma^{T} I_{p, n-p} \Sigma\right|^{-\frac{1}{2}} \int_{\mathbb{R}^{n-d}} \varphi(\Sigma, y) e^{-i\langle A u, y\rangle} d y_{j_{d+1}} \ldots d y_{j_{n}} \\
=\int_{\mathbb{R}^{n-d}} \varphi_{0}(\Sigma, y) e^{-i\langle A u, y\rangle} d y_{j_{d+1}} \ldots d y_{j_{n}} .
\end{gathered}
$$

In a completely analogous way to that of Richter [13, page 252-254] it can now be shown that $\tilde{\varphi}(A \cdot E, A u)$ only depends on $A u$, so there exists a function $F \in C^{\infty}\left(\mathbb{R}^{n} \backslash\{0\}\right)$ such that $\tilde{\varphi}(A \cdot E, A u)=F(A u)$. The proof now follows by remark 2 to lemma 5.2 and the subsequent computation of the Fourier transforms.
4. The group of orthogonal matrices

$$
H=\left\{\left.\left(\begin{array}{cccc}
x_{1} & -x_{2} & -x_{3} & -x_{4} \\
x_{2} & x_{1} & x_{4} & -x_{3} \\
x_{3} & -x_{4} & x_{1} & x_{2} \\
x_{4} & x_{3} & -x_{2} & x_{1}
\end{array}\right) \right\rvert\, x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+x_{4}^{2}=1\right\}<S O(4)
$$

operates on $V=\mathbb{R}^{4}$, which is identified with $\left(\mathbb{R}^{4}\right)^{*}$ via the canonical scalar product. We have $\rho=\rho^{*}$.

The choice $E=\left\langle Z_{1}\right\rangle$ (notation as above) shows how indispensable assumption (ii), section 5 , is for our considerations. The other assumptions are all satisfied. For $\xi \neq 0$ we have $H_{\xi}=\{e\}$. Thus, there should be no need of differential operators. But since every point $\xi$ lies on more than one hyperplanes of the form $h \cdot E^{\perp}$, the moment conditions fail to characterize the range of $\mathcal{S}\left(\mathbb{R}^{4}\right)$. A counterexample can be constructed as in [6, page 605].

Now let $E=\left\langle Z_{1}, Z_{2}\right\rangle$. All assumptions are satisfied. The action does not have property $(G)$ of section 6 . We have instead the extreme situation where two planes of the form $h \cdot E^{\perp}$ either coincide or are complementary. Here we cannot dispense with the moment conditions. Otherwise, $\mathcal{S}\left(\mathbb{R}^{4}\right)^{\wedge}=\mathcal{S}(\mathcal{E})$ would hold, but this is contradicted by the following counterexample:

Let $g \in C_{c}^{\infty}\left(\mathbb{R}^{4}\right), g(0) \neq 0$. Let $\left(x_{1}, x_{2}, x_{3}, x_{4}\right)$ denote the matrix in $H$, which has this vector in its first column. $x_{3}^{2}+x_{4}^{2}$ is constant on the coset $\left(x_{1}, x_{2}, x_{3}, x_{4}\right) H_{E}$. For $\xi \in\left(x_{1}, x_{2}, x_{3}, x_{4}\right) \cdot E^{\perp}, \xi \neq 0$, it holds:

$$
x_{3}^{2}+x_{4}^{2}=\frac{\xi_{1}^{2}+\xi_{2}^{2}}{\xi_{1}^{2}+\xi_{2}^{2}+\xi_{3}^{2}+\xi_{4}^{2}} .
$$

Now take $\varphi \in \mathcal{S}(\mathcal{E})$ such that

$$
\tilde{\varphi}\left(\left(x_{1}, x_{2}, x_{3}, x_{4}\right) H_{E}, \xi\right)=\left(x_{3}^{2}+x_{4}^{2}\right) g(\xi)=\frac{\left(\xi_{1}^{2}+\xi_{2}^{2}\right) g(\xi)}{\xi_{1}^{2}+\xi_{2}^{2}+\xi_{3}^{2}+\xi_{4}^{2}}=: F(\xi)
$$

for $\xi \in\left(x_{1}, x_{2}, x_{3}, x_{4}\right) \cdot E^{\perp}, \xi \neq 0$. Obviously, $\varphi$ is not an element of $\mathcal{S}(V)^{\wedge}$, because $F$ cannot be continuously extended to 0 .

Next, $H$ is being enriched by the central matrices of the form

$$
\left(\begin{array}{ccc}
\cos t & -\sin t &  \tag{17}\\
\sin t & \cos t & \\
& 0 & \cos t \\
\hline & -\sin t \\
0 & \sin t & \cos t
\end{array}\right)
$$

We study the interesting case $E=\left\langle Z_{2}, Z_{3}\right\rangle$. The matrices

$$
\begin{gathered}
X_{1}=\left(\begin{array}{ccc}
0 & 1 & \\
-1 & 0 & \\
Y=\left(\begin{array}{ccc}
0 & 1 & 0 \\
-1 & 0 & 0
\end{array}\right), \quad X_{2}=\left(\begin{array}{ccc}
0 & & \\
& 0 & 1 \\
& -1 & 0
\end{array}\right), \quad Z=\left(\begin{array}{cccc}
0 & 0 & 1 \\
0 & -1 & 0 & -1
\end{array}\right) \\
-1 & 0 & 0
\end{array}\right)
\end{gathered}
$$

form a basis of the Lie algebra of $H$. Now $\mathcal{K}^{1}\left(\mathfrak{h} \times{ }_{s} \mathbb{R}^{4}\right)$ is spanned by

$$
P=X_{1}\left(Z_{3}^{2}+Z_{4}^{2}\right)+X_{2}\left(Z_{1}^{2}+Z_{2}^{2}\right)+Y\left(Z_{1} Z_{4}-Z_{2} Z_{3}\right)-Z\left(Z_{1} Z_{3}+Z_{2} Z_{4}\right) .
$$

For $w=w_{1} Z_{1}+w_{2} Z_{4} \in E^{\perp} \backslash\{0\}$ we have: $\mathfrak{h}_{w}=\left\langle w_{2}^{2} X_{1}+w_{1}^{2} X_{2}+w_{1} w_{2} Y\right\rangle$. Therefore, every 2-plane $h \cdot E^{\perp}$ through $w$ has the form $\exp \left[t\left(w_{2}^{2} X_{1}+w_{1}^{2} X_{2}+\right.\right.$ $\left.\left.w_{1} w_{2} Y\right)\right] \cdot E^{\perp}, t \in \mathbb{R}$. Since $\bar{w}:=w_{1} Z_{2}-w_{2} Z_{3}$ is orthogonal to all these planes $\left(\mathfrak{h}_{\bar{w}}=\mathfrak{h}_{w}\right.$ and $\left.\bar{w} \perp E^{\perp}\right)$, the action does not have property $(G)$. Nevertheless, we can still dispense with the moment conditions. The latter are satisfied if the functions $F_{m}$ in (15) are polynomials. So we proceed as follows.

Let $F \in C^{\infty}\left(\mathbb{R}^{4} \backslash\{0\}\right)$ have the property that its restriction to every plane $h \cdot E^{\perp}$ is a homogeneous polynomial of degree $m$. We shall show that $F$ is a homogeneous polynomial of degree $m$.

At first, we establish this for the hyperplane $\mathbb{R}^{3} \times\{0\}$ in $\mathbb{R}^{4}$.
Let $\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right)$ denote the matrix $\left(x_{1}, x_{2}, x_{3}, x_{4}\right)$ multiplied by the matrix (17). Then $\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right) \cdot E^{\perp}$ lies in $\mathbb{R}^{3} \times\{0\}$ iff

$$
x_{3} \sin t+x_{4} \cos t=x_{1} \cos t+x_{2} \sin t=0 .
$$

In that case,

$$
h \cdot w=\frac{1}{\cos t}\left(\begin{array}{c}
-w_{1} x_{2} \sin 2 t+w_{2} x_{3} \sin 2 t \\
w_{1} x_{2} \cos 2 t-w_{2} x_{3} \cos 2 t \\
w_{1} x_{3}+w_{2} x_{2} \\
0
\end{array}\right)
$$

for $\cos t \neq 0$ and $h \cdot E^{\perp}=\left\langle Z_{2}, Z_{3}\right\rangle$ for $\cos t=0$. We obtain all planes through $(0,0,1)$ in $\mathbb{R}^{3}$. Since $F$, restricted to each of them, is a homogeneous polynomial, we can write

$$
F\left(\xi_{1}, \xi_{2}, \xi_{3}, 0\right)=A_{0}\left(\xi_{1}, \xi_{2}\right)+A_{1}\left(\xi_{1}, \xi_{2}\right) \xi_{3}+\ldots+A_{m}\left(\xi_{1}, \xi_{2}\right) \xi_{3}^{m}
$$

If we write down this equation for $m+1$ different values of $\xi_{3} \neq 0$, we obtain a linear system in $A_{0}, \ldots, A_{m}$. By Cramer's rule we deduce that the $A_{i}$ 's are $C^{\infty}$-functions and, therefore, so is $\left.F\right|_{\mathbb{R}^{3} \times\{0\}}$. Now since $\left.F\right|_{\mathbb{R}^{3} \times\{0\}}$ is homogeneous, it is a polynomial. We denote it by $P_{0}$.

If $F\left(\xi_{1}, \xi_{2}, \xi_{3}, \xi_{4}\right) \not \equiv P_{0}\left(\xi_{1}, \xi_{2}, \xi_{3}\right)$, let $0<k_{0} \leq m$ be the highest possible exponent, such that

$$
F_{1}\left(\xi_{1}, \xi_{2}, \xi_{3}, \xi_{4}\right):=\frac{F\left(\xi_{1}, \xi_{2}, \xi_{3}, \xi_{4}\right)-P_{0}\left(\xi_{1}, \xi_{2}, \xi_{3}\right)}{\xi_{4}^{k_{0}}}
$$

is a $C^{\infty}$-function on $\mathbb{R}^{4} \backslash\{0\}$. This function is a homogeneous polynomial of the smaller degree $m-k_{0}$ on every plane $h \cdot E^{\perp}$ that is not contained in $\mathbb{R}^{3} \times\{0\}$. This means that every expression of the form

$$
\left(\frac{\partial}{\partial w_{1}}\right)^{m-k_{0}+1-j}\left(\frac{\partial}{\partial w_{2}}\right)^{j} F_{1}(h \cdot w)
$$

vanishes for $h=\left(x_{1}, x_{2}, x_{3}, x_{4}, t\right),\left(x_{3} \sin t+x_{4} \cos t, x_{1} \cos t+x_{2} \sin t\right) \neq(0,0)$. By continuity, it vanishes for every $h$. Consequently, $F_{1}$ is a homogeneous polynomial on every plane $h \cdot E^{\perp} \subseteq \mathbb{R}^{3} \times\{0\}$ too.

We repeat the argument with $F_{1}$ instead of $F$, define a polynomial $P_{1}$, a number $k_{1}$ and a function $F_{2}$ with which we continue in the same manner. After finitely many steps

$$
\left(P_{0}, k_{0}\right) \rightarrow\left(P_{1}, k_{1}\right) \rightarrow \ldots \rightarrow\left(P_{n}, k_{n}\right)
$$

we obtain a function $F_{n+1} \in C^{\infty}\left(\mathbb{R}^{4} \backslash\{0\}\right)$, constant on every plane $h \cdot E^{\perp}$. Since all planes $h \cdot E^{\perp} \in \mathbb{R}^{3} \times\{0\}$ contain $(0,0,1), F_{n+1}$ is constant on $\mathbb{R}^{3} \times\{0\}$. Since every other plane $h \cdot E^{\perp}$ meets $\mathbb{R}^{3} \times\{0\}$, it is globally constant. We obtain:

$$
F\left(\xi_{1}, \xi_{2}, \xi_{3}, \xi_{4}\right)=F_{n+1} \xi_{4}^{\sum_{i=0}^{n} k_{i}}+\sum_{j=0}^{n} P_{n-j}\left(\xi_{1}, \xi_{2}, \xi_{3}\right) \xi_{4}^{\sum_{i=0}^{n-j-1} k_{i}}
$$

$$
\left(\sum_{i=0}^{n} k_{i}=m\right)
$$

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[^0]:    ${ }^{1}$ Richter's definition leads to a much greater space $\mathcal{K}(n)$ (following his notation), and the relation ' $\mathcal{K}(n) \subseteq \operatorname{Ker} \lambda$ ' fails to hold in [14, Proposition 3]. For example, if $P$ belongs to $\mathcal{K}^{M}(n)$, then $P+\mathcal{U}^{M-1}(n) \subseteq \mathcal{K}(n)$, but of course $P+\mathcal{U}^{M-1}(n) \nsubseteq \operatorname{Ker} \lambda$. Furthermore, the use of the operators $Z_{i}$ as $\operatorname{Ad}\left(\exp \left(Z_{i}\right)\right.$ - Id instead of $\operatorname{ad} Z_{i}$ in the definition of $\mathcal{K}^{M}(n)$ makes it much harder to deduce the last equation in the proof of Proposition 3 if $Y$ is not just one of the $Z_{1}, \ldots, Z_{n}$, but a linear combination of them.

[^1]:    ${ }^{2}$ In the first and last example of section 9 it is shown that they are also necessary for a relation like (4).

[^2]:    ${ }^{3}$ This very technical assumption is needed, so as to construct a unique cross section $\sigma$ from the various $\sigma_{i}$. After all, Glaeser's theorem is formulated for open sets in $\mathbb{R}^{p}$ and not in the context of manifolds, as is the case here.
    ${ }^{4}$ After choosing bases we can identify $E^{\perp}$ and $V^{*}$ with $\mathbb{R}^{n-d}$ and $\mathbb{R}^{n}$ respectively.

[^3]:    ${ }^{5}$ This theorem is a special case of theorem I in [16, page 65].

[^4]:    ${ }^{6}$ Notice that condition (ii) in the beginning of section 5 here degenerates.
    ${ }^{7}$ Obviously, we have here $d \nu(Y)=0$ for $Y \in \mathcal{K}^{1}\left(\mathfrak{h} \times{ }_{s} V\right)$.

