# UNITARY REPRESENTATIONS DEFINED BY BOUNDARY CONDITIONS-THE CASE OF $\mathfrak{s l}(2, R)$ 

BY

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## § 1. Introduction

Let $\mathfrak{g}$ be a Lie algebra over $R$, the field of real numbers, and $\sigma$, a $\mathfrak{g}$-module in a Hilbert space $\mathcal{H}$. If the domain of $\sigma$ is dense, one can define an adjoint module $\sigma^{+}$in $\boldsymbol{\mathcal { H }}$ such that

$$
(\sigma(a) f, g)=\left(f, \sigma^{\dagger}\left(a^{\dagger}\right) g\right)
$$

for all $f \in \mathcal{D}(\sigma), g \in \mathcal{D}\left(\sigma^{\dagger}\right), a \in \mathcal{U}[g]$, (see Appendix A for notation and details). The module $\sigma$ is said to be symmetric or (infinitesimally) unitary if $\sigma \subset \sigma^{\dagger}$ and self-adjoint if $\sigma=\sigma^{\dagger}$. The importance of self-adjointness comes from the fact that $d T$ is a self-adjoint module (see Appendix A). Here $T$ is a unitary representation of the simply connected group corresponding to $\mathfrak{g}$, and $d T$ is the usual $\mathfrak{g}$-module with the set of $C^{\infty}$-vectors of $T$ as its domain. Calling a g-module exact if it is equal to $d T$ for some $T$, a natural problem would be to determine all exact extensions of a given symmetric $g$-module. The theory here is analogous to the theory of self-adjoint extensions of a single unbounded symmetric operator. In fact if $\operatorname{dim} g=1$, it is well known that $g$-module is exact if and only if it is self-adjoint. For the general case, self-adjointness is necessary but not sufficient for exact-

[^0]ness. (See Appendix A.) However there are many interesting cases where self-adjointness alone is enough to assure integrability to the group. In this paper one such module for $\mathfrak{R l}(2, \boldsymbol{R})$ is studied in detail, and all its self-adjoint extensions are obtained. Since the extensions are determined by boundary conditions, it is natural to consider the corresponding group representations as being defined by the boundary conditions. An interesting feature is that all non-trivial unitary irreducible representations are obtained by determining all self-adjoint extensions of the module $\sigma_{\lambda}$. Although the representations of $G$ have been known for a long time this way of deriving them appears to be new. It is an interesting problem to find modules similar to $\sigma_{\lambda}$ for other groups as well. This question will be pursued in future papers.

A brief description of the contents follows. Generalities about $\mathfrak{g}$-modules and some basic results which are used repeatedly are collected together in Appendix A. In Section 3, the basic homomorphism $\varrho_{2}$ of $\mathfrak{ß l}(2, \boldsymbol{R})$ into differential operators and the modules $\sigma_{\lambda}, \sigma_{\lambda}$ are defined. All self-adjoint extensions are determined in Theorems 1, 2 and Lemma 10. They are shown to be integrable (Theorem 3). Their unitary equivalence classes are identified in Theorem 4. Theorems 5 and 6 describe the set of $C^{\infty}$ vectors, $K$-eigenbasis and the group operators for the representations $T_{\lambda}^{+}$, which correspond to self-adjoint extensions of $\sigma_{\lambda}^{t}$. Theorems 7, 8 and 10 do the same for the representations $T_{\lambda, \delta, \delta^{\prime}}$ (here $\delta, \delta^{\prime}$ parametrize self-adjoint extensions of $\sigma_{\lambda}$ ). Theorem 9 is an auxiliary result which determines all representations of $G$ in $L^{2}(R)$ with a given restriction to a parabolic subgroup. Theorem 11 gives the intertwining operator between the unitary principal series and the representations $T_{\lambda, \delta, \delta^{\prime}}$, when $\lambda$ is imaginary.

The methods and results of this paper will be used to obtain explicit decomposition of (1) the tensor product of two discrete series representations of $\mathfrak{g l}(2, \boldsymbol{R})([12])$, (2) the Weil representation associated to a quadratic form [13]. In fact the present paper grew out of an attempt to get such an explicit decomposition. In this connection we refer to [11] for another approach to the same problem.

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## § 2. Preliminaries on $\mathfrak{S l}_{2}$

Let $G$ denote the simply connected Lie group with Lie algebra $\mathfrak{G l ( 2 , \boldsymbol { R } ) \text { . Let } X , H , Y / 2 )}$ denote the standard basis of $\mathfrak{\xi l}(\mathbf{2}, \boldsymbol{R})$, i.e.

$$
X=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right), \quad H=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad Y=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right)
$$

Then $[H, X]=2 X,[H, Y]=-2 Y$, and $[X, Y]=H$. Write $\varkappa(\theta)=\exp \theta(X-Y), h(t)=\exp t H$, and $u(s)=\exp s X$. Here we have written $\exp$, rather than $\exp _{G}$ for the exponential map of $\mathfrak{g l}(2, \boldsymbol{R})$ into $G$. Put $K=\{x(\theta) \mid \theta \in R\}, A=\{h(t) \mid t \in R\}$ and $N=\{u(s) \mid s \in R\}$. Then $K, A, N$ are closed subgroups of $G$ and $G=K \cdot A \cdot N$. Write

$$
\begin{equation*}
w=\exp \frac{\pi}{2}(X-Y), \quad \gamma=\exp \pi(X-Y)=w^{2} \tag{1}
\end{equation*}
$$

then the center $Z(G)$ is the cyclic group $=\left\{\gamma^{n} \mid n \in \mathbf{Z}\right\}$. Also

$$
\begin{equation*}
\operatorname{Ad} w \cdot H=-H, \quad \operatorname{Ad} w \cdot X=-Y, \quad \operatorname{Ad} w \cdot Y=-X \tag{2}
\end{equation*}
$$

Let $\mathcal{U}\left[\mathfrak{S l}_{2}\right]$ denote the universal enveloping algebra of $\mathfrak{g l}(2, \mathbf{C})$. The following basis of $\mathfrak{s l}(2, \mathbb{C})$ will be frequently used

$$
\begin{equation*}
X^{\prime}=(i H+X+Y) / 2, \quad H^{\prime}=i(X-Y), \quad Y^{\prime}=(-i H+X+Y) / 2 \tag{3}
\end{equation*}
$$

Then $\left\{X^{\prime}, H^{\prime}, Y^{\prime}\right\}$ is another Lie triple. Let

$$
\begin{equation*}
\Omega=(H+1)^{2}+4 Y X=(H-1)^{2}+4 X Y \tag{4}
\end{equation*}
$$

Then $\Omega$ generates the center of $\mathcal{U}\left[\mathfrak{S H}_{2}\right]$.
Let $\mathcal{E}(G)$ denote the set of equivalence classes of irreducible unitary representations of G. These are all known (Bargman [1], Kunze-Stein [6], Pukhansky [10] and Sally [15]). In any irreducible representation the center of the group and the center of the algebra $\mathcal{U}\left[\mathfrak{H}_{2}\right]$ are mapped into scalars and they can be used to parametrize them. It is known that Spec $H^{\prime}$ and Spec $\Omega$ determine the unitary equivalence classes. Then the points of $\mathcal{E}(G)$ can be parametrized as follows:

1. $\omega(\eta, \lambda)$ where $\lambda^{2}$ and $\eta$ are real and $|\operatorname{Re} \lambda|+|\eta| \leqslant 1$, with inequality holding if $\lambda$ is real. These representations are characterized by Spec $H^{\prime}=\eta+2 Z, \operatorname{Spec} \Omega=\lambda^{2}$ and Spec $\gamma=e^{-i \eta \pi}$. Also $\omega(\eta,-\lambda)=\omega(\eta, \lambda), \omega(\eta+2, \lambda)=\omega(\eta, \lambda)$.
2. $\omega^{ \pm}(\lambda)$, where $\lambda$ is real and $\lambda+1>0$. Here Spec $H^{\prime}= \pm(\lambda+1+2 N)$, Spec $\Omega=\lambda^{2}$ and Spec $\gamma=\exp (\mp i(\lambda+1) \pi)$.
3. $\omega(0,1)$, the class of the trivial representation.

The class $\omega(\eta, \lambda)$ is known as the principal series when $\lambda$ is purely imaginary, and as the complementary series when $\lambda$ is real. The classes $\omega^{ \pm}(\lambda)$ are the so-called discrete series. It is then known that representations of the class $\omega^{ \pm}(\lambda)$ are integrable (or matrix entries belong to $\left.L^{1}(G / Z(G))\right)$ when $\lambda>1$ and are square integrable when $\lambda>0$. It may be of some interest to note that although $\omega^{+}(\lambda)$ and $\omega^{-}(\lambda)$ are not unitarily equivalent, they are physically (viewpoint of physicists) equivalent since there is an anti-unitary isomorphism between representations of the two classes.
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## § 3. The homomorphism $Q_{\lambda}$ and the modules $\sigma_{\lambda}$

For each $\lambda \in \mathbf{C}$, define the following differential operators on $R^{\prime}=R \backslash\{0\}$ :

$$
\begin{gather*}
\varrho_{\lambda}(H)=2 t \partial_{t}+1, \quad \varrho_{\lambda}(X)=-i t \\
\varrho_{\lambda}(X)=-i\left(t \partial_{t}^{2}+\partial_{t}-\lambda^{2} / 4 t\right) \tag{5}
\end{gather*}
$$

Then one checks easily that these operators satisfy the same commutation rules as the Lie triple $\{X, H, Y\}$. Thus $\varrho_{\lambda}$ extends as an algebra homomorphism of $\mathcal{U}\left[\mathfrak{I X}_{2}\right]$ into differential operators on $R^{\prime}$. A simple calculation shows that

$$
\varrho_{\lambda}(\Omega)=\lambda^{2}
$$

The formal transpose and adjoint with respect to Lebesgue measure on $R$ is easily checked to be

$$
{ }^{t} \varrho_{\lambda}(H)=-\varrho_{\lambda}(H), \quad \varrho_{\lambda}(X)=\varrho_{\lambda}(X), \quad{ }^{t} \varrho_{\lambda}(Y)=\varrho_{\lambda}(Y)
$$

In particular, if $\lambda^{2}$ is real the operators $\varrho_{\lambda}(Z)^{*}=-\varrho_{\lambda}(Z)$ for each $Z \in \mathbb{Z} l(2, R)$, i.e. they are formally skew adjoint. The natural domain of $\varrho_{\lambda}$ is $C^{\infty}\left(R^{\prime}\right)$. In this part we shall determine all sub-modules of $\varrho_{\lambda}$ in $L^{2}\left(R^{\prime}\right)$ which are integrable to a group representation. We also discuss simultaneously the $g$-modules on $R_{ \pm}$, defined by $\sigma_{\lambda}^{4}=\left(C^{\infty}\left(R_{ \pm}\right)\right.$, $\left.\varrho_{\lambda}\right)$. Define

$$
\begin{equation*}
\sigma_{\lambda}=\left(C_{c}^{\infty}\left(R^{\prime}\right), \varrho_{\lambda}\right), \quad \sigma_{\lambda}^{+}=\left(C_{c}^{\infty}\left(R_{ \pm}\right), \varrho_{\lambda}\right) . \tag{6}
\end{equation*}
$$

Then $\sigma_{\lambda}$ is a $g$-module in $L^{2}(R)$ and its adjoint module is described in the following lemma.
Lemma 1. The domain of the adjoint module $\mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)=\left\{f \in C_{c}^{\infty}\left(R^{\prime}\right) \mid \varrho_{\lambda}(a) f \in L^{2}\right.$, for all a in $\left.U\left[\mathfrak{I J}_{2}\right]\right\}$. Moreover, $\varrho_{\lambda}^{\dagger}(a) f=\varrho_{\lambda}(a) f$ for all $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$. If $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$, then for each $t_{0}>0$

$$
\sup _{|t|>t_{0}}\left|t^{m}\left(t \partial_{t}\right)^{n} f\right|<\infty
$$

for all $m, n \in \mathbf{N}$. Similar results hold for the modules $\sigma_{\lambda}^{\dagger}$ in $L^{2}\left(R_{ \pm}\right)$.
Proof. The description of $\mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ is just Weyl's lemma. (See Appendix A.) Let $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$. Then

$$
\begin{aligned}
|f(s)| & \leqslant \int_{s}^{\infty}\left|\partial_{t} f\right| d t \leqslant s^{-1 / 2}\left\|t \partial_{t} f\right\|_{2} \\
& \leqslant \frac{1}{2} s^{-1 / 2}\left\|\varrho_{\lambda}(H-1) f\right\|_{2}
\end{aligned}
$$

The lemma follows by replacing $f$ by $\varrho_{\lambda}\left(X^{m}(H-1)^{n}\right) f$.

## § 4. Eigenfunctions of $\boldsymbol{\varrho}_{\lambda}\left(\boldsymbol{H}^{\prime}\right)$

Suppose $\varrho_{\lambda}\left(H^{\prime}\right) f=\xi f$ for some $\xi \in \mathbb{C}$. Then $f$ satisfies the Tricomi's equation

$$
\left\{\partial_{t}^{2}+\frac{1}{t} \partial_{t}-1+\xi / t-\lambda^{2} / 4 t^{2}\right\} f=0
$$

If we put $f=2|t|^{-1 / 2} g(2 t)$ (see [3], p. 251) then $g$ satisfies the Whittaker's equation

$$
\left[\partial_{t}^{2}-1 / 4+\varkappa / t+\left(1-\mu^{2}\right) / t^{2}\right] g=0
$$

where $x=\xi / 2, \mu=\lambda / 2$. On any connected interval $W_{x, \mu}(t), W_{-x, \mu}(-t)$ are a basis of solutions of the Whittaker's equation, where $W_{\varkappa, \mu}$ is the Whittaker's function. Also $W_{\varkappa, \mu}(t)=$ $O\left(t^{\kappa / 2} e^{-t / 2}\right)$ as $t \rightarrow \infty$, so that a solution which is in $L^{2}$ near $\infty$ has to be a multiple of $W_{x, \mu}(t)$. Thus we introduce the function

$$
\begin{equation*}
L_{\varkappa, \mu}(t)=(2 t)^{-1 / 2} W_{\varkappa, \mu}(2 t), \quad t>0 . \tag{7}
\end{equation*}
$$

Then

$$
f=\left\{\begin{array}{l}
c_{1} L_{\xi / 2, \lambda_{2} 2}(t), \quad t>0 \\
c_{2} L_{-\xi / 2, \lambda_{2} / 2}(-t), \quad t<0
\end{array}\right.
$$

for suitable constants $c_{1}, c_{2}$. The function $f$ will be in $L^{2}$ when $\xi$ and $\lambda$ are such that $L_{\xi / 2 . \lambda i 2}(t)$ is in $L^{2}$ near 0 . Now the point $t=0$ is a regular singular point of the Tricomi's equation with indices $\pm \lambda / 2$. Now $L_{\xi / 2, \lambda / 2}(t)=e^{-t}(2 t)^{\lambda / 2} \Psi(a, c ; 2 t)$ with $a=(1-\xi+\lambda) / 2$ and $c=1+\lambda, \Psi$ being Tricomi's function. ([3], p. 255, Vol. I.) From fractional power series expansion of $\Psi$ near 0 ( $[3]$, p. 257, Vol. I) one obtains those of $L$ and the results are summarized in the following

Lemma 2. There exist convergent power series $P_{j}(t)$ depending on $\lambda$ such that for $t>0$

$$
L_{\xi / 2, \lambda / 2}(t)=\left\{\begin{array}{l}
t^{\lambda / 2} P_{1}(t)+t^{-\lambda / 2} P_{2}(t), \quad \text { if } \quad \lambda \notin \mathbf{Z} \\
t^{\lambda / 2}\left(P_{1}(t) \ln t+P_{\mathbf{3}}(t)\right\}+t^{-\lambda / 2} P_{2}(t) \quad \text { for } \quad \lambda \in \mathbf{Z}, \quad \lambda \neq 0 .
\end{array}\right.
$$

Moreover, define

$$
c(\xi, \lambda)=\Gamma(-\lambda) 2^{\lambda / 2}\{\Gamma((1-\xi-\lambda) / 2)\}^{-1}
$$

for $\lambda \notin \mathbf{N}$, and $c(\xi, n)$ as the residue of $c(\xi, \lambda)$ at $\lambda=n$ (for example $\left.c(\xi, 0)=-\{\Gamma((1-\xi) / 2)\}^{-1}\right)$. Then $P_{1}(0)=c(\xi, \lambda)$ for all $\lambda$, and $P_{2}(0)=c(\xi,-\lambda)$ for all $\lambda \notin-N$. When $\lambda=0$ we have

$$
L_{\xi / 2,0}(t)=\left\{\begin{array}{l}
Q_{1}(t), \quad \xi-1 \in 2 \mathbf{N} \\
Q_{2}(t) \ln |t|+Q_{3}(t), \quad \text { if } \quad \xi-1 \oplus 2 \mathbf{N}
\end{array}\right.
$$

where $Q_{1}(t)$ is a polynomial, $Q_{1}(0)=(-1)^{(\xi-1) / 2}$ and $Q_{2}, Q_{3}$ are convergent power series with $Q_{2}(0)=c(\xi, 0)$ and $Q_{3}(0)=c(\xi, 0) d(\xi)$, where $d(\xi)=\psi((1-\xi) / 2)-2 \psi(0)+\ln 2, \psi$ is the logarith. mic derivative of the $\Gamma$-function.

From the above lemma, the following corollary is immediate.
Corollary 1. If $|\operatorname{Re} \lambda|<1$, then $L_{\xi / 2, \lambda / 2} \in L^{2}\left(R_{+}\right)$for all $\xi \in \mathbb{C}$. If $\operatorname{Re} \lambda \geqslant 1$, then $L_{\xi / 2, \lambda / 2} \in L^{2}\left(R_{+}\right)$if and only if $c(\xi,-\lambda)=0$ or if and only if $\xi \in 1+\lambda+2 \mathbf{N}$.

Corollary 2. Let $V_{\lambda}(\xi)\left(V_{\lambda}^{ \pm}(\xi)\right)$ denote the linear space of eigenfunctions of $\varrho_{\lambda}\left(H^{\prime}\right)$ in $L^{2}(R)\left(L^{2}\left(R_{ \pm}\right)\right)$for the eigen-value $\xi$. Then
(i) $V_{\lambda}(\xi)=V_{\lambda}^{+}(\xi)+V_{\lambda}^{-}(\xi)$,
(ii) $\operatorname{dim} V_{\lambda}^{ \pm}(\xi)=1$ for all $\xi \in \mathbf{C}$ if $|\operatorname{Re} \lambda|<1$ and
(iii) for $\lambda \geqslant 1, \operatorname{dim} V_{\lambda}^{ \pm}(\xi)=1$ if $\xi \in \pm(1+\lambda+2 N), \operatorname{dim} V_{\lambda}^{ \pm}(\xi)=0$ otherwise.

From the above corollary $V_{\lambda}^{ \pm}( \pm i)=0$ when $\lambda \geqslant 1$, and so we have

Corollary 3. The symmetric operators $\sigma_{\lambda}\left(H^{\prime}\right), \sigma_{\lambda}^{+}\left(H^{\prime}\right)$ are essentially self-adjoint in $L^{2}(R)$ and $L^{2}\left(R_{ \pm}\right)$respectively when $\lambda$ is real and $\geqslant 1$.

The situation is very different when $|\operatorname{Re} \lambda|<1$. In the following paragraphs we shall find all $\mathfrak{\omega l}_{2}$-modules $\sigma$ such that $\sigma_{\lambda} \subset \sigma \subset \sigma_{\lambda}^{\dagger}$ for which $\sigma\left(H^{\prime}\right)$ is essentially self-adjoint.

## § 5. Boundary forms

For each $a \in \mathcal{U}$ and $f, g \in D\left(\sigma_{\lambda}^{\ddagger}\right)$, let $B_{\lambda}$ denote the boundary form of the module $\sigma_{\lambda}$, i.e.

$$
B_{\lambda}(a: f: g)=\left(\varrho_{\lambda}(a) f, g\right)-\left(f, \varrho_{\lambda}\left(a^{\dagger}\right) g\right)
$$

(see the appendix for the identities satisfied by $B_{\lambda}$ ). If $Z=t_{1} H+t_{2} X$, with $t_{1}, t_{2}$ real, then it is easy to check that none of the eigenfunctions of $\varrho_{\lambda}(Z)$ for real eigenvalues are in $L^{2}$ and so the operators $\sigma_{\lambda}(Z)$ are essentially skew adjoint. Thus $B_{\lambda}(Z: f: g) \equiv 0$ and so from the identities satisfied by $B_{\lambda}$ (see Lemma A.2) it follows that

$$
B_{\lambda}(a: f: g)=0
$$

for all $f, g \in \mathcal{D}\left(\sigma_{\lambda}^{+}\right)$and $a$ of the form $H^{m} X^{n}$. Thus it is sufficient to study $B_{\lambda}\left(H^{\prime n}: f: g\right)$. Now

$$
\varrho_{\lambda}\left(H^{\prime}\right)=-\partial_{t} \circ t \circ \partial_{t}+t+\lambda^{2} / 4 t
$$

so that

$$
\begin{align*}
(d / d t) t W(f, g) & =\left(\varrho_{\lambda}\left(H^{\prime}\right) f\right) g-f \varrho_{\lambda}\left(H^{\prime}\right) g \\
& =-i\left\{\left(\varrho_{\lambda}(Y) f\right) g-f \varrho_{\lambda}(Y) g\right\} \tag{8}
\end{align*}
$$

where $W(f, g)=f \partial_{t} g-g \partial_{t} f$ is the Wronskian. Thus

$$
\begin{equation*}
B_{\lambda}\left(H^{\prime}: f: \bar{g}\right)=-i B_{\lambda}(Y: f: \bar{g})=-\left\{\left.t W(f, g)\right|^{+}\right\} . \tag{9}
\end{equation*}
$$

Here $\left.\varphi\right|^{+}=\lim _{t \rightarrow 0+} \varphi(t)-\lim _{t \rightarrow 0^{-}} \varphi(t)$.

## § 6. Boundary values

We assume that $\lambda^{2}$ is real and $|\operatorname{Re} \lambda|<1$. To consider all different cases simultaneously we use the following device. Define

$$
a_{\lambda}(t)=\left\{\begin{array}{l}
|t|^{-\lambda / 2}, \lambda \neq 0  \tag{10}\\
\ln |t|, \lambda=0
\end{array}\right.
$$

and define

$$
\begin{aligned}
A_{1}^{ \pm}(f) & =\lim _{t \rightarrow 0 \pm}\left\{t W\left(f, a_{\lambda}(t)\right)\right\} \\
A_{2}^{ \pm}(f) & =\lim _{t \rightarrow 0 \pm}\left\{t W\left(f,|t|^{\lambda / 2}\right)\right\}
\end{aligned}
$$

Lemma 3. For $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right), A_{j}^{ \pm}(f), j=1,2$, exist and

$$
\begin{aligned}
& A_{1}^{+}(f)=( \pm i) \int_{R_{ \pm}} a_{\lambda}(t) \varrho_{\lambda}(Y) f d t \\
& A_{2}^{+}(f)=( \pm i) \int_{R_{ \pm}}|t|^{\alpha / 2} \varrho_{\lambda}(Y) f d t .
\end{aligned}
$$

Moreover, for a suitable choice of $\varphi^{ \pm} \in \mathscr{S}\left(R^{\prime}\right)$, we have

$$
\begin{aligned}
& A_{1}^{ \pm}(f)=i B\left(Y: f: a_{\lambda}(t) \varphi^{ \pm}\right) \\
& A_{2}^{+}(f)=i B\left(Y: f:|t|^{\lambda / 2} \varphi^{ \pm}\right)
\end{aligned}
$$

Proof. Note that $\varrho_{\lambda}(Y) a_{\lambda}(t)=0, \varrho_{\lambda}(Y)|t|^{\lambda / 2}=0$, the first two formulae follow from (8) and the fact that $f$ is rapidly decreasing at $\infty$ (cf. Lemma 1). The last two follow from (9).

The following class of functions is somewhat more convenient to work with than $\mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$. For any open subset $U$ of $R$ denote $\mathscr{S}(U)$ as the Schwartz space of $U$, i.e. $\mathscr{S}(U)=$ $\left\{f \in C^{\infty}(U)\left|\sup _{U}\right| t^{m} \partial_{t}^{n} f \mid<\infty\right.$ for all $\left.m, n \in \mathbf{N}\right\}$. If $f \in \mathscr{S}(U)$, then $f$ and its derivatives have limits as $t$ approaches a boundary point of $U$. It is somewhat more convenient to work with classes of functions in $\mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ which have asymptotic expansions at the boundary. For this purpose we introduce the class $\boldsymbol{X}_{\lambda}$.

Let $\mathscr{X}_{\lambda}$ denote the class of functions $f \in C^{\infty 0}\left(R^{\prime}\right)$ such that
(i) for each $\delta>0, \sup \left\{\left|t^{m} \partial_{t}^{n} f\right|:|t|>\delta\right\}<\infty$
(ii) there exists $\delta>0$ and functions $f_{1} f_{2} \in \mathscr{P}\left(R^{\prime}\right)$ such that for $0<|t|<\delta$

$$
\begin{equation*}
f(t)=|t|^{\lambda / 2} f_{1}(t)+a_{\lambda}(t) f_{2}(t) \tag{11}
\end{equation*}
$$

Also write $\mathscr{X}_{\lambda}^{t}=|t|^{\lambda / 2} \mathscr{S}\left(R_{ \pm}\right)$.
Lemma 4. (i) $\boldsymbol{X}_{\lambda} \subset L^{2}(R)$ and $\mathfrak{X}_{\lambda}^{ \pm} \subset L^{2}\left(R_{ \pm}\right)$. (ii) $\varrho_{\lambda}(Z) \mathcal{X}_{\lambda} \subset \mathfrak{X}_{\lambda}$, and $\varrho_{\lambda}(Z) \mathfrak{X}_{\lambda}^{ \pm} \subset \boldsymbol{X}_{\lambda}^{ \pm}$ for all $Z \in \Xi I(2, R)$. In particular, $\mathcal{X}_{2} \subset \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ and $\mathcal{X}_{\lambda}^{ \pm} \subset \mathcal{D}\left(\left(\sigma_{\lambda}^{+}\right)^{\dagger}\right)$.

Proof. Only (ii) needs checking. This follows from the following. Let $\alpha= \pm \lambda / 2$. Put $\tilde{\varrho}_{\lambda}=|t|^{-\alpha / 2} \circ \varrho_{\lambda} \circ|t|^{\alpha / 2}$. Then $\tilde{\varrho}_{\lambda}(H)=2 t \partial_{t}+1+\alpha, \tilde{\rho}_{\lambda}(X)=-i t$, and $\tilde{\varrho}_{\lambda_{1}}(Y)=-i\left(t \partial_{t}^{2}+(1+\alpha) \partial_{t}\right)$. This proves (ii) for $\lambda \neq 0$. For $\lambda=0$, define $\varrho_{0}=(\ln |t|)^{-1} \rho \varrho_{0} \circ(\ln |t|)$, then $\tilde{\varrho}_{0}(H)=$ $2 t \partial_{i}+1-(\ln |t|)^{-1}, \tilde{\varrho}_{0}(X)=-i t$ and $\tilde{\varrho}_{0}(Y)=-i\left(t \partial_{t}^{2}+\partial_{t}+2(\ln |t|)^{-1} \partial_{t}\right)$. Thus $\varrho_{0}(Z) \mathcal{X}_{0} \subset \mathcal{X}_{0}$ for all $Z \in \mathscr{L}(2, R)$. The rest is clear.

The following lemma is easily checked, by direct calculation.
Lemma 5. Let $|\operatorname{Re} \lambda|<1$. Suppose that $f=|t|^{\lambda / 2} f_{1}+a_{\lambda}(t) f_{2}, g=|t|^{\lambda / 2} g_{1}+a_{\lambda}(t) g_{2}$ with $f_{j}, g_{j} \in \mathscr{S}\left(R^{\prime}\right)$, then

$$
\{t W(f, g)\}_{-}^{+}=\left.c_{\lambda}\left(f_{1} g_{2}-g_{1} f_{2}\right)\right|_{-} ^{ \pm}
$$

where $c_{\lambda}$ is the constant $\equiv t W\left(|t|^{\lambda^{2}}, a_{\lambda}(t)\right),=-\lambda$ if $\lambda \neq 0$ and $=1$ if $\lambda=0$. In particular, $B_{\lambda}\left(H^{\prime}: f: \bar{g}\right)=-c_{\lambda}\left(f_{1} g_{2} \sim f_{2} g_{1}\right)_{-}^{+}$.

Lemma 6. Assume $|\operatorname{Re} \lambda|<1$. Let $f \in \mathcal{X}_{\lambda}$ have the local expansion $f=|t|^{\lambda / 2} f_{1}+a_{\lambda}(t) f_{2}$ near 0. Then

$$
A_{2}^{+}(f)=-c_{\lambda} f_{2}(0 \pm), \quad A_{1}^{\prime}(f)=c_{\lambda} f_{1}(0 \pm) .
$$

Moreover, for all $f, g \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$

$$
c_{\lambda} B_{\lambda}\left(H^{\prime}: f: g\right)=A_{1}^{+}(f) A_{2}^{+}(\bar{g})-A_{2}^{+}(f) A_{1}^{+}(\bar{g})-A_{1}^{-}(f) A_{2}^{-}(\bar{g})+A_{2}^{-}(f) A_{1}^{-}(\bar{g}) .
$$

Proof. The first part of the lemma follows from the definition and the fact $c_{\lambda} \equiv$ $t W\left(|t|^{\lambda / 2}, a_{\lambda}\right)$. Next the formula for $B_{\lambda}\left(H^{\prime}: f: g\right)$ follows from the previous lemma when $f, g \in \boldsymbol{X}_{\lambda}$. To prove it in the general case note that the eigenspaces $V_{\lambda}(\xi) \subset \mathcal{X}_{\lambda}$ and from known results about the adjoint of a symmetric operator, it follows that $\mathcal{D}\left(\sigma_{\lambda}\left(H^{\prime}\right)^{*}\right)=$ $\mathcal{D}\left(\mathrm{Cl}_{\lambda}\left(H^{\prime}\right)\right)+V_{\lambda}(i)+V_{\lambda}(-i)$. Thus, given $\varphi, \psi \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$, it follows that there exist $\varphi_{1}, \psi_{1} \in \mathcal{X}_{\lambda}$ such that $\varphi-\varphi_{1}, \psi-\psi_{1} \in \mathcal{D}\left(\mathrm{Cl}_{\lambda}\left(H^{\prime}\right)\right)$. But $B_{\lambda}\left(H^{\prime}: \varphi, \psi\right)=B_{\lambda}\left(H^{\prime}: \varphi_{1}: \psi_{1}\right)$. Finally observe that $A_{j}^{+}(f)=0$ for all $f \in C_{c}^{\infty}\left(R^{\prime}\right)$ and thus $A_{j}^{ \pm} \equiv 0$ on $\mathcal{D}\left(\mathrm{Cl}_{\lambda}\left(H^{\prime}\right)\right) \cap \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ by Lemma 3. The formula for $B_{\lambda}\left(H^{\prime}: f: g\right)$ thus follows from the corresponsing formula for $f, g \in \mathcal{X}_{\lambda}$.

Lemma 7. Define the sesquilinear form $F_{\lambda}(x, y)$ on $\mathbf{C}^{4}$ as follows:

$$
F_{\lambda}(x, y)= \begin{cases}x_{1} \bar{y}_{1}-x_{2} \bar{y}_{2}-x_{3} \bar{y}_{3}+x_{4} \bar{y}_{4}, & \lambda^{2}<0 \\ x_{1} \bar{y}_{2}-x_{2} \bar{y}_{1}-x_{3} \bar{y}_{4}+x_{4} \bar{y}_{3}, & \text { if } \lambda^{2} \geqslant 0,|\lambda|<1\end{cases}
$$

Define $\Lambda_{0}: \mathcal{D}\left(\sigma_{\lambda}^{+}\right) \rightarrow \mathbf{C}^{4}$ by setting

$$
\Lambda_{0}(f)=\left(A_{1}^{+}(f), A_{2}^{+}(f), A_{1}^{-}(f), A_{2}^{-}(f)\right)
$$

then $\Lambda_{0}$ maps onto and

$$
\begin{equation*}
c_{\lambda} B_{\lambda}\left(H^{\prime}: f: g\right)=F_{\lambda}\left(\Lambda_{0}(f), \Lambda_{0}(g)\right) \tag{12}
\end{equation*}
$$

for all $f, g \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$.
Proof. This follows from the previous lemma if you note the following. If $\lambda^{2}<0$, then $A_{1}^{ \pm}(\bar{g})=\overline{A_{2}^{ \pm}(g)}$, and if $\lambda^{2} \geqslant 0,|\operatorname{Re} \lambda|<1$ then $A_{1}^{ \pm}(\bar{g})=\overline{A_{1}^{ \pm}(g)}, A_{2}^{ \pm}(\bar{g})=\overline{A_{2}^{ \pm}(g)}$. From the formulas for $A_{j}^{\dagger}$ in Lemma 6, if follows that $\Lambda_{0}$ maps $\boldsymbol{X}_{\lambda}$ onto $\mathbf{C}^{4}$.

Lemma 8. (i) $\Lambda_{0}\left(\varrho_{\lambda}(X) f\right)=0$ for all $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ and (ii) there exists a matrix $M_{\lambda} \in G L(\mathbf{4}, \mathbf{C})$ such that $\Lambda_{0}\left(\varrho_{\lambda}(H) f\right)=M_{\lambda} \cdot \Lambda_{0}(f)$. Also $M_{\lambda}=\operatorname{diag}(1+\lambda, 1-\lambda, 1+\lambda, 1-\lambda)$ if $\lambda \neq 0$ and

$$
M_{0}=\left(\begin{array}{rrrr}
1 & -2 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & -2 \\
0 & 0 & 0 & 1
\end{array}\right)
$$

Proof. To show that $A_{j}\left(\varrho_{\lambda}(X) f\right)=0$, it is sufficient to show that $\lim _{t \rightarrow 0} t W(t f, \psi)=0$, where $\psi=|t|^{\lambda / 2}$ or $a_{\lambda}(t)$. Now $t W(t f, \psi)=-t f \psi+t^{2} W(f, \psi)$. Now $t^{2} W(f, \psi) \rightarrow 0$ as $t \rightarrow 0$. Now $f \in D\left(\sigma_{\lambda}^{+}\right)$implies that $f=O\left(|t|^{-1 / 2}\right)$, and thus $t f \psi=O\left(|t|^{1-\lambda / 2}\right)$ if $\lambda \neq 0$, and $O\left(|t|^{1 / 2} \ln |t|\right)$ if $\lambda=0$. Thus in all cases $\lim t f \psi=0$. This proves that $\Lambda_{0}\left(\varrho_{\lambda}(X) f\right)=0$.

To prove (ii) we use the identity (see Lemma A.2) satisfied by boundary forms. Thus $2 B_{\lambda}(Y: f: g)=B_{\lambda}(Y H-H Y: f: g)=B_{\lambda}\left(Y: \varrho_{\lambda}(H) f: g\right)+B_{\lambda}\left(Y: f: \varrho_{\lambda}(H) g\right)$ since (see section 5), $B_{\lambda}(H: \cdot: \cdot) \equiv 0$. Thus $B_{\lambda}\left(Y: \varrho_{\lambda}(H) f: g\right)=B_{\lambda}\left(Y: f: \varrho_{\lambda}(2-H) g\right)$. Now let $\varphi \in C_{c}^{\infty}(R)$ such that $\varphi=1$ around 0 , then

$$
\varrho_{\lambda}(2-H)|t|^{\lambda / 2} \varphi=\left(-2 t \partial_{t}\right)|t|^{\lambda_{i} / 2} \varphi=(1-\lambda)|t|^{\lambda / 2}
$$

in a neighborhood of $0 . \quad A_{2}^{ \pm}\left(\varrho_{\lambda}(H) f\right)=i B_{\lambda}\left(Y: \underline{\varrho}_{\lambda}(H) f:|t|^{\lambda / 2} \varphi\right)=(1-\lambda) A_{2}^{ \pm}(f)$. Again $\varrho_{\lambda}(2-H)\left(a_{\lambda}(t) \varphi\right)=(1+\lambda)|t|^{-\lambda / 2}$ around 0 if $\lambda \neq 0$, and $=\ln |t|-2$ near 0 if $\lambda=0$. Thus $A_{1}^{ \pm}\left(\varrho_{\lambda}(H) f\right)=(1+\lambda) A_{1}^{ \pm}(f)$ if $\lambda \neq 0,=A_{1}^{ \pm}(f)-2 A_{2}^{ \pm}(f)$ if $\lambda=0$. Thus $\Lambda_{0}\left(\varrho_{\lambda}(H) f\right)=M_{\lambda} \cdot \Lambda_{0}(f)$, where $M_{\lambda}$ is given in the lemma. This completes the proof.

Lemma 9. Let the boundary values $\Lambda_{m}$ be defined by $\Lambda_{m}(f)=\Lambda_{0}\left(\varrho_{2}(Y)^{m} f\right)$. Then we have
(i) $\Lambda_{m} \circ \varrho_{\lambda}(X)=-m\left(M_{\lambda}+m-1\right) \cdot \Lambda_{m-1}$
(ii) $\Lambda_{m} \circ \varrho_{\lambda}(H)=\left(M_{\lambda}+2\right) \cdot \Lambda_{m}$
(iii) $\Lambda_{m} \circ \varrho_{\lambda}(Y)=\Lambda_{m+1}$
(iv) $c_{\lambda} B_{\lambda}\left(H^{\prime m}: f: g\right)=\sum_{r=0}^{m-1} F_{\lambda}\left(\Lambda_{r}(f), \Lambda_{m-r-1}(g)\right)$
(v) Let $v_{m} \in \mathbb{C}^{4}$ be arbitrary. Then there exists an $f \in \boldsymbol{X}_{\lambda}$ such that $\Lambda_{m}(f)=v_{m}, m \in \mathbf{N}$.

Proof. The first statement follows from the identity $Y^{m} X=X Y^{m}-m(H+m-1) Y^{m-1}$ in $\mathcal{U}\left[\mathfrak{E l}_{2}\right]$. The statements (ii) and (iii) are obvious. The part (iv) follows from the identity (iii) of Lemma A.2, satisfied by boundary forms and Lemma 7. Finally let $f \in \mathcal{X}_{\lambda}$ and suppose that

$$
f=|t|^{\lambda / 2} f_{1}+a_{\lambda} f_{2} \text { near } 0, \text { with } f_{1}, f_{2} \in \mathscr{S}\left(R^{\prime}\right)
$$

Put $\varrho_{\lambda}(Y)^{m} f=|t|^{\lambda / 2} f_{1, m}+a_{\lambda} f_{2, m}$. Then we have the formula

$$
|t|^{-\alpha / 2} \circ \varrho_{\lambda}(Y) \circ|t|^{\alpha / 2}=-i\left(t \partial_{t}^{2}+(1+\alpha) \partial_{t}\right)
$$

One checks easily by induction that

$$
\left\{\left(t \partial_{t}^{2}+(1+a) \partial_{t}\right)\right\}^{m}=\left\{\prod_{j=1}^{m-1}\left(t \partial_{t}+j+a\right)\right\} \partial_{t}^{m}=D_{m, a}, \text { say }
$$

Then $f_{1, m}=(-i)^{m} D_{m, \lambda / 2} f_{1}$, so that $f_{1, m}( \pm 0)=(-i)^{m}(1+\alpha)_{m}\left(\partial_{t}^{m} f_{1}\right)( \pm 0)$. If $\lambda \neq 0$, we have similarly $f_{2, m}( \pm 0)=(-1)^{m}(1+\alpha)_{m} \partial_{t}^{m} f_{2}( \pm 0)$. If $v_{m}=\Lambda_{m}(f)$, then

$$
v_{m}=c_{\lambda}\left(f_{1, m}(0+), f_{2, m}(0+), f_{1, m}(0-), f_{2, m}(0-)\right)
$$

Put

$$
w_{m}=\left(\partial_{t}^{m} f_{1}(0+), \partial_{t}^{m} f_{2}(0+), \partial_{t}^{m} f_{1}(0-), \partial_{t}^{m} f_{2}(0-)\right)
$$

Then

$$
w_{m}=i^{m}\left\{(1+\alpha)_{m}\right\}^{-1} c_{\lambda}^{-1} \cdot v_{m} .
$$

By Borel's theorem ( $[6], \mathrm{p} .30$ ), one can choose $f_{1}, f_{2} \in \mathscr{S}\left(R^{\prime}\right)$ with values $w_{m}$ for the derivatives at 0 . Thus there exist $f \in \mathcal{X}_{\lambda}$, such that $\Lambda_{m}(f)=v_{m}$. In the case $\lambda=0$

$$
\binom{f_{1, m}}{f_{2, m}}=\left(\begin{array}{ll}
t \partial_{t}+1 & 3 \\
0 & t \partial_{t}+1
\end{array}\right)\left(\begin{array}{ll}
-i \partial_{t} & f_{1, m-1} \\
-i \partial_{t} & f_{2, m-1}
\end{array}\right)
$$

One checks by induction on $m$

$$
\binom{f_{1, m}}{f_{2, m}}=\left(\begin{array}{ll}
t \partial_{t}+1 & 2 \\
0 & t \partial_{t}+1
\end{array}\right) \cdots\left(\begin{array}{ll}
t \partial_{t}+m & 2 \\
0 & t \partial_{t}+m
\end{array}\right)\binom{\left(-i \partial_{t}\right)^{m} f_{1}}{\left(-i \partial_{t}\right)^{m} f_{2}}
$$

From which one gets easily

$$
\begin{aligned}
& f_{1, m}(0 \pm)=m!\left\{\left(-i \partial_{t}\right)^{m} f_{1}(0 \pm)+2\left(1+\frac{1}{2}+\ldots+1 / m\right)\left(-i \partial_{t}\right)^{m} f_{2}(0 \pm)\right\} \\
& f_{2, m}(0 \pm)=m!(-i)^{m}\left(\partial_{t}^{m} f_{2}\right)(0 \pm)
\end{aligned}
$$

An argument similar to the case $\lambda \neq 0$, now gives that there exists $f \in \mathcal{X}_{\lambda}$ such that $\Lambda_{m}(f)=v_{m}$, for all $m$.

## § 7. Self-adjoint extensions

With these preparations we can now obtain all the self-adjoint extensions of $\sigma_{2}$.
Theorem 1. (i) Let $\sigma$ be an $\mathfrak{s l}_{2}$-module such that $\sigma_{\lambda} \subset \sigma \subset \sigma_{\lambda}^{\dagger}$. Then $\sigma_{\lambda} \subset \sigma^{\dagger} \subset \sigma_{\lambda}^{\dagger}$. Let $E(\sigma)$ denote the subspace of $\mathrm{C}^{4}$ defined by $E(\sigma)=\left\{\Lambda_{0}(f) \mid f \in D(\sigma)\right\}$. Then $M_{\lambda} \cdot E(\sigma)=E(\sigma)$ and $\mathcal{D}\left(\sigma^{\dagger}\right)=\left\{f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right) \mid \Lambda_{m}(f) \in E(\sigma)^{\perp}\right.$ for all $\left.m \in \mathbf{N}\right\}$.
(ii) Conversely let $E \subset \mathrm{C}^{4}$ be such that $M_{\lambda} \cdot E=E$. Let a $\mathfrak{ß l}_{2}$-module $\sigma \subset \sigma_{\lambda}^{\dagger}$ be defined by $\mathcal{D}(\sigma)=\left\{f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right\rangle \mid \Lambda_{m}(f) \in E\right.$, for all $\left.m \in \mathbf{N}\right\}$. Then $E(\sigma)=E$ and $\sigma^{+\pi}=\sigma$. In particular, the map $\sigma \rightarrow E(\sigma)$ is a bijection of self-adjoint $\mathfrak{ß l}_{2}$-modules $\sigma$ such that $\sigma_{\lambda} \subset \sigma$ and subspaces $E$ such that (a) $M_{\lambda} \cdot E=E$ and (b) $E=E^{\perp}$. Here the orthogonal complement is with respect to the form $F_{\lambda}$ introduced in Lemma 7.

Proof. Since $\Lambda_{0}\left(\varrho_{\lambda}(H) f\right)=M_{\lambda} \cdot \Lambda_{0}(f)$ and $M_{\lambda}$ is invertible, it follows that $M_{\lambda} \cdot E(\sigma)=E(\sigma)$. From the relation $\Lambda_{m} \varrho_{\lambda}(X)=-m\left(M_{\lambda}+m-1\right) \Lambda_{m-1}$, it follows by induction on $m$, that $\Lambda_{m}(D(\sigma))=E(\sigma)$ for all $m$. Now $\sigma_{\lambda} \subset \sigma \subset \sigma_{\lambda}^{\dagger}$ implies that $\sigma_{\lambda} \subset \sigma^{\dagger} \subset \sigma_{\lambda}^{\dagger}$ and $\mathcal{D}\left(\sigma^{\dagger}\right)=$ $\left\{f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right) \mid B_{\lambda}(a: g: f)=0\right.$ for all $a \in \mathcal{U}$ and $\left.g \in \mathcal{D}(\sigma)\right\}$. Since $B_{\lambda}\left(H^{m} X^{n}: g: f\right)=0$ for all $g$ and $f$ it follows that $f \in \mathcal{D}\left(\sigma^{\dagger}\right)$ if and only if $B_{\lambda}\left(H^{\prime m}: g: f\right)=0$ for all $m$. From part (iv) of previous lemma, this is equivalent to

$$
\sum_{r=0}^{n-1} F_{\lambda}\left(\Lambda_{r}(g), \Lambda_{n-r-1}(f)\right)=0
$$

for all $n \in \mathbb{N}$, and $g \in \mathcal{D}(\sigma)$. Since the range $\Lambda_{m}(\mathcal{D}(\sigma))=E(\sigma)$, it follows by induction on $n$ that the above identity holds if and only if $\Lambda_{m}(f) \in E(\sigma)^{\perp}$, for all $m \in \mathbf{N}$. That $E\left(\sigma^{\dagger}\right)=E(\sigma)^{\perp}$ follows from the statement $(\mathrm{v})$ of the previous lemma.

To prove (ii), note first that $\mathcal{D}(\sigma)$ is invariant under $\varrho_{\lambda}(Z)$, for all $Z \in \mathfrak{I l}(2, R)$. In fact this follows from the properties (i)-(iii) of the previous lemma. Since $\Lambda_{m}(f)=0$ for $f \in \mathcal{D}\left(\sigma_{\lambda}\right)$, it follows that $\sigma_{\lambda} \subset \sigma \subset \sigma_{\lambda}^{\ddagger}$. Thus $\sigma$ is a well defined $\mathfrak{I}_{2}$-module. From the first part it follows that $D\left(\sigma^{\dagger \dagger}\right)=\left\{f \in D\left(\sigma_{\lambda}^{\dagger}\right) \mid \Lambda_{m}(f) \in E\left(\sigma^{\dagger}\right)^{\perp}\right.$, for all $\left.m\right\}$. Since $E\left(\sigma^{\dagger}\right)=E^{\perp}$ and $E^{\perp \perp}=E$, it follows that $\sigma=\sigma^{+\dagger}$. The rest is clear.

Corollary. $\mathcal{D}\left(\sigma_{\lambda}^{\dagger \dagger}\right)=\left\{f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right) \mid \Lambda_{m}(f)=0\right.$ for all $\left.m \in \mathbf{N}\right\}$.
Proof. For $E\left(\sigma_{\lambda}^{\dagger}\right)=C^{4}$ and thus $E\left(\sigma_{\lambda}^{\dagger \dagger}\right)=\{0\}$, the corollary follows from this.
The next lemma describes all such subspaces.

Lemma 10. Let $\lambda^{2}<1$. For each $\lambda$, the following is a complete list of all subspaces $E$ such that $M_{\lambda} \cdot E=E$ and $E=E^{\perp}$.

Case 1. $\lambda^{2}<0$. In this case $E$ is of the form $E_{\delta, \delta^{\prime}}=\mathbf{C}\left(e_{1}+\delta e_{3}\right)+\mathbf{C}\left(e_{2}+\delta^{\prime} e_{4}\right)$, where $\delta, \delta^{\prime} \in \mathbf{C}$ are such that $|\delta|=\left|\delta^{\prime}\right|=1$.

Case 2. $0<\lambda^{2}<1$. There are two classes. Case $2 a . E$ is of the form $E_{\delta, \delta^{\prime}}=\mathbf{C}\left(e_{1}=\delta e_{3}\right)+$ $\mathrm{C}\left(e_{2}+\delta^{\prime} \epsilon_{4}\right)$, where $\delta^{\prime} \delta=1$. Case $2 b . E$ is one of the following $E_{13}=\mathbf{C} e_{1}+\mathbf{C} e_{3}, E_{14}=\mathbf{C} e_{1}+\mathbf{C} e_{4}$, $E_{23}=\mathbf{C} e_{2}+\mathbf{C} e_{3}$ and $E_{24}=\mathbf{C} e_{2}+\mathbf{C} e_{4}$.

Case 3. $\lambda=0$. There are two classes. Case $3 a . E$ is of the form

$$
E_{\delta, \delta^{\prime}}=\left\{\left(x_{1}, x_{2}, x_{3}, x_{4}\right) \in \mathrm{C} \mid x_{4}-\delta x_{2}=0 \quad \text { and } \quad \delta x_{1}+\delta^{\prime} x_{2}-x_{3}=0\right\},
$$

where $|\delta|=1, \delta^{\prime} \delta$ is real. Case $3 b . E=E_{13}$.
Proof. In the statement of the Lemma, $e_{j}$ is the standard basis of $\mathbf{C}^{4}$. Put $W_{1}=\mathbf{C} e_{1}+\mathbf{C} e_{3}$ and $W_{2}=\mathbf{C} e_{2}+\mathbf{C} e_{4}$. We consider each case separately.

Case 1. $\lambda^{2}<0$. Here the form $F_{\lambda}$ (see Lemma 7) is symmetric, and $M_{\lambda}$ is a diagonal matrix, with $W_{1}, W_{2}$ being the eigenspaces. Thus $E=E \cap W_{1}=E \cap W_{2}$. In this case $W_{1}^{\perp}=W_{2}$ so that we must have $\operatorname{dim} E \cap W_{1}=1, j=1,2$. Let $v_{1}=a_{1} e_{1}+a_{3} e_{3}, v_{2}=a_{2} e_{2}+a_{4} e_{4}$ be a basis of $E$. Then $F_{\lambda}\left(v_{1}, v_{1}\right)=\left|a_{1}\right|^{2}-\left|a_{3}\right|^{2}=0 ; F_{\lambda}\left(v_{2}, v_{2}\right)=-\left|a_{2}\right|^{2}+\left|a_{4}\right|^{2}=0$. Thus $E$ is of the form $E_{\delta, \delta^{\prime}}$. One checks that $\left(E_{\delta, \delta^{\prime}}\right)^{\perp}=E_{\delta, \delta^{\prime}}$.

Case 2. $\lambda^{2}>0$. In this case the form $F_{\lambda}$ is symplectic, but $W_{f}$ are still the eigenspaces of $M_{\lambda}$. Consider first the case $\operatorname{dim} E \cap W_{1}=1$. Suppose $v_{1}=a_{1} e_{1}+a_{3} e_{3}, v_{2}=a_{2} e_{2}+a_{4} e_{4}$, $E=\mathbf{C} v_{1}+C v_{2}$. Then $B\left(v_{1}, v_{2}\right)=a_{1} \bar{a}_{2}-a_{3} \bar{a}_{4}=0$. Now consider the case where none of the $a_{j}$ 's are zero. In this case $E$ is of the form $E_{\delta, \delta^{\prime}}=\mathbf{C}\left(e_{1}+\delta e_{3}\right)+\mathbf{C}\left(e_{2}+\delta^{\prime} e_{4}\right)$ with $\delta, \delta^{\prime}$ satisfying $\delta \delta^{\prime}=1$. One checks that $E^{\perp}=E$.

Case 2b. $\operatorname{dim} E \cap W_{1}=1$. But one of $a_{f}$ 's is zero. Suppose $a_{1}=0$. Then $a_{3} \neq 0$, so that $a_{4}=0$. Thus in this case $E=\mathbf{C} e_{2}+\mathbf{C} e_{3}$. Similarly you get the other possibilities listed.

Case 3. In this case $M_{0}$ (see Lemma 8) is unipotent. Write $M_{0}=\exp (-2 C)$, then $C$ is nilpotent and $C^{2}=0$. Also $W_{1}=\mathbf{C} e_{1}+\mathbf{C} e_{3}=\left\{v \in \mathbf{C}^{4} \mid C v=0\right\}$.

Case $3 a . \operatorname{dim} E \cap W_{1}=1$. Then $E$ has a basis $v_{1}, v_{2}$ such that $C v_{1}=0, C v_{2}=v_{1}$. Thus $v_{1}=a_{1} e_{1}+a_{3} e_{3}$ and $\dot{v}_{2}=\Sigma b, e_{j}$. Then $b_{2}=a_{1}$, and $b_{4}=a_{3} . F_{0}\left(v_{1}, v_{2}\right)=a_{1} b_{2}-a_{3} b_{4}=\left|a_{1}\right|^{2}-$ $\left|a_{3}\right|^{2}=0$. $F_{0}\left(v_{2}, v_{2}\right)=b_{1} b_{2}-b_{2} b_{1}-b_{3} b_{4}+b_{4} b_{3}=0$. Thus we may suppose $E=\mathbf{C}\left(e_{1}+\delta e_{3}\right)+$ $\mathbf{C}\left(\delta^{\prime} e_{3}+e_{2}+\delta e_{4}\right)$, where $\delta, \delta^{\prime} \in \mathbf{C}$. Then $M_{0} E=E$. The condition $E^{\perp}=E$ gives that $|\delta|=1$ and $\delta \delta^{\prime}$ is real. Thus in this case $E=E_{\delta, \delta^{\prime}}$.

Case $3 b$. If $\operatorname{dim} E \cap W_{1}=2$, then $E=\mathbf{C} e_{1}+C e_{3}$, is the only solution in this case.
Definition. Let $\sigma_{\lambda, \delta, \delta^{\prime}}$ be the $\mathfrak{I l}_{2}$-module which is self-adjoint and is defined by the boundary conditions $\sigma_{\lambda} \subset \sigma_{\lambda, \delta, \delta} \subset \sigma_{\lambda}^{\dagger}$ and

$$
\begin{equation*}
\mathcal{D}\left(\sigma_{\lambda, \delta, \delta^{\prime}}\right)=\left\{f \in \mathcal{D}\left(\sigma_{\lambda}^{+}\right) \mid \Lambda_{m}(f) \in E_{\delta, \delta^{\prime}}, \text { for all } m \in \mathbf{N}\right\} \tag{13}
\end{equation*}
$$

Here $E_{\delta, \delta^{\prime}}$ is defined in the above lemma, and $\delta, \delta^{\prime}$ satisfy the appropriate conditions (depending on $\lambda$ ) stated there.

Remark. In this connection note that the module $\sigma_{\lambda}$ depends only on $\lambda^{2}$, while $\sigma_{\lambda, \delta, \delta}$ depends on $\lambda$. In fact we have $(\lambda \neq 0)$,

$$
\begin{equation*}
\sigma_{-\lambda, \delta, \delta^{\prime}}=\sigma_{\lambda, \delta^{\prime}, \delta} \tag{14}
\end{equation*}
$$

This may be seen as follows. Writing $A_{1}^{+}(f: \lambda)$ for $A_{1}^{+}(f)$ to denote its dependence on $\lambda$, it is clear that $A_{1}^{ \pm}(f:-\lambda)=A_{2}^{+}(f: \lambda)$, if $\lambda \neq 0$. From this it follows that $\Lambda_{0}(f:-\lambda) \in E_{\dot{\delta}, \delta^{\prime}}$ if and only if $\Lambda_{0}(f: \lambda) \in E_{\delta^{\prime}, \delta}$, proving (14).

Remark. Let $V$ denote the unitary operator $V f=c_{1} f$ if $t>0,=c_{2} f$ if $t<0$, where $\left|c_{1}\right|=\left|c_{2}\right|=1$. Let $\sigma=V \circ \sigma_{\lambda, \delta, \delta} \circ V^{-1}$. Then it is easily checked that

$$
\begin{equation*}
\sigma=\sigma_{\lambda, \delta c_{2} / c_{1}, \delta^{\delta} c: / c_{1}} \tag{15}
\end{equation*}
$$

Theorem 2. (i) For each $\lambda$ real and $\lambda+1>0$, there exists a unique self-adjoint $\mathfrak{s l}_{2}$ module $\mu_{\lambda}^{+}$in $L^{2}\left(R_{ \pm}\right)$such $\mathcal{D}\left(\mu_{\lambda}^{+}\right) \supset \mathcal{X}_{\lambda}^{+}$.
(ii) For $\lambda=0$, and for $\lambda \geqslant 1, \sigma_{\lambda}^{ \pm}$has a unique self-adjoint extension and for $-1<\lambda<1$, $\lambda \neq 0, \sigma_{\lambda}^{ \pm}$has exactly two self-adjoint extensions namely $\mu_{\lambda}^{ \pm}, \mu_{-\lambda}^{+}$.

Proof. Case 1. $\lambda \geqslant 1$. In this case $\sigma_{\lambda}^{ \pm}(Z)$ is essentially self-adjoint, for all $Z \in \mathfrak{l l}(2, \boldsymbol{R})$ (Cor. 3 of Lemma 2). Hence if $\mu_{\lambda}^{t}$ is a self-adjoint extension of $\sigma_{\lambda}^{t}$, then by Lemma A. 5 $\left(\sigma_{\lambda}^{t}\right)^{\dagger}=\mu_{\lambda}^{t}$, and $\mu_{\lambda}^{t}$ is the unique self-adjoint extension of $\sigma_{\lambda}^{t}$.

Case 2. $-1<\lambda<1$. Consider the self-adjoint module $\sigma$ in $L^{2}(R)$ such that $\sigma_{\lambda} \subset \sigma$ and $E(\sigma)=E_{18}$. Then it is clear that $\mathcal{D}(\sigma)=\mathcal{D}(\sigma) \cap L^{2}\left(R_{+}\right)+\mathcal{D}(\sigma) \cap L^{2}\left(R_{-}\right)$where we consider
$L^{2}\left(R_{ \pm}\right)$as subspaces of $L^{2}(R)$. Define $\sigma^{ \pm}$as $\mathbb{B l}_{2}$-modules in $L^{2}\left(R_{ \pm}\right), \mathcal{D}\left(\sigma^{ \pm}\right)=\mathcal{D}(\sigma) \cap L^{2}\left(R_{ \pm}\right)$. Then $\sigma$ self-adjoint means that $\sigma^{ \pm}$are self-adjoint in $L^{2}\left(R_{ \pm}\right)$. From the definition of $E(\sigma)$, it follows that $\mathcal{D}\left(\sigma^{ \pm}\right) \supset \boldsymbol{X}_{\lambda}^{ \pm}$. Define $\mu_{\lambda}^{ \pm}=\sigma^{ \pm}$. Then $\mu_{\lambda}^{ \pm}$are self-adjoint and $\mathcal{D}\left(\mu_{\lambda}^{ \pm}\right) \supset \boldsymbol{X}_{\lambda}^{ \pm}$. To prove uniqueness suppose $\mu^{\prime}$ is another self-adjoint module in $L^{2}\left(R_{+}\right)$such that $\mathcal{D}\left(\mu^{\prime}\right) \supset \mathcal{X}_{\lambda}^{+}$. Consider $\sigma^{\prime}=\mu^{\prime}+\mu_{\lambda}^{-}$. Then $\sigma^{\prime}$ is self-adjoint and $E(\sigma) \supset E_{18}$. Therefore $E\left(\sigma^{\prime}\right)=$ $E_{13}$ and $\sigma=\sigma^{\prime}$, implying $\mu^{\prime}=\mu_{\lambda}^{+}$. The other cases are handled similarly.

Lemma 11. Let $\sigma_{\lambda} \subset \sigma \subset \sigma_{\lambda}^{\dagger}$ and $\sigma^{\dagger}=\sigma$. Then the operator $\sigma\left(H^{\prime}\right)$ is essentially self-adjoint.

Proof. Now $\sigma=\sigma^{\dagger}$ implies that $\sigma\left(H^{\prime}\right)$ is a symmetric operator. Also $\sigma\left(H^{\prime}\right)^{*} \subset \sigma_{\lambda}\left(H^{\prime}\right)^{*}$. Suppose $\xi$ is not real and $\sigma\left(H^{\prime}\right)^{*} f=\xi f$. Then $\sigma_{\lambda}\left(H^{\prime}\right)^{*} f=\xi f$ and so $\varrho_{\lambda}\left(H^{\prime}\right) f=\xi f$, i.e., $f \in \mathcal{X}_{\lambda}$ (see Lemma 2). Also $f \in \mathcal{D}\left(\sigma\left(H^{\prime}\right)^{*}\right)$. Thus $\left(\varrho_{\lambda}\left(H^{\prime}\right) g, f\right)-\left(g, \varrho_{\lambda}\left(H^{\prime}\right) f\right)=0$ for all $g \in \mathcal{D}(\sigma)$ or $B_{\lambda}\left(H^{\prime}: g: f\right)=0$ for all $g \in \mathcal{D}(\sigma)$. Thus $F_{\lambda}\left(\Lambda_{0}(g), \Lambda_{0}(f)\right)=0$, for all $g \in D(\sigma)$. This implies that $\Lambda_{0}(f) \in E(\sigma)^{\perp}=E(\sigma)$, since $\sigma$ is self-adjoint. Finally, by the identity satisfied by boundary forms, it follows that

$$
B_{\lambda}\left(H^{\prime n}: g: f\right)=\sum_{r=0}^{n-1} B\left(H^{\prime}: \varrho\left(H^{\prime}\right)^{r} g: \xi^{n-r-1} f\right)=0
$$

for all $g \in \mathcal{D}(\sigma)$, and all $n$. Using formula (iv) of Lemma 9 it follows by induction on $n$, that $\Lambda_{n}(f) \in E(\sigma)$ for all $n \in \mathbf{N}$. Thus $f \in \mathcal{D}\left(\sigma^{+}\right)=\mathcal{D}(\sigma)$. This contradicts symmetry of $\sigma\left(H^{\prime}\right)$, since $\xi$ is not real. Thus $f=0$. This completes the proof.

Lemma 12. Let $\sigma_{\lambda} \subset \sigma \subset \sigma_{\lambda}^{\dagger}$ and $\sigma^{\dagger}=\sigma$. Then the Spec $\sigma\left(H^{\prime}\right)$ is discrete. Moreover, all the eigenfunctions of $\mathrm{Cl} \sigma\left(H^{\prime}\right)$ are actually in $D(\sigma)$.

Proof. Let $R_{\zeta}$ denote the resolvent of the closure of the operator $\sigma\left(H^{\prime}\right)$, i.e., $R_{\zeta}=$ $\left(\mathrm{Cl} \sigma\left(H^{\prime}\right)-\zeta\right)^{-1}$. We shall show that for each $f \in C_{c}^{\infty}\left(R^{\prime}\right), R_{\zeta} f$ is a meromorphic function of $\zeta$. This implies, by a well known formula for the spectral measure in terms of the resolvent, that the spectrum is discrete.

Let $\zeta \in \mathbf{C}, \zeta$ not real. Two linearly independent solutions of $\varrho_{\lambda}\left(H^{\prime}\right) f=\zeta f$ may be chosen as follows. The Whittaker's function has an analytic continuation to a domain containing the upper half-plane and $R^{\prime}$, and so we may define $y_{1}=(2 z)^{-1 / 2} W_{x, \mu}(2 z), y_{2}=$ $(2 z)^{-1 / 2} W_{-x . \mu}(-z)$ where $x=\zeta / 2, \mu=\lambda / 2$. Then $y_{1}, y_{2}$ are solutions of the equation $\varrho_{\lambda}\left(H^{\prime}\right) f=$ $\zeta$. Also from the formula for the Wronskian of Whittaker's functions, it follows that $W\left(y_{1}, y_{2}\right)=z^{-1} e^{-i t \pi / 2}$. Thus, $t W\left(y_{1}, y_{2}\right)=e^{-i 5 \pi / 2}$ is a constant. Define $K_{\xi}(t, s)=$ $e^{-i e \pi / 2} y_{1}(t) y_{2}(s)$ if $s<t$ and $e^{i t \pi / 2} y_{2}(t) y_{1}(s)$ if $s>t$. Then it follows from standard methods in differential equations that the function $g_{1}=K_{\xi} f: t \rightarrow \int_{R^{\prime}} K_{\xi}(t, s) f(s) d s$ is a solution of $\left(\varrho_{\lambda}\left(H^{\prime}\right)-\zeta\right) g_{1}=f$, in $R^{\prime}$, whenever $f \in C_{c}\left(R^{\prime}\right)$. Actually $K_{\xi} f \in \mathcal{X}_{\lambda}$ when $f \in C_{c}^{\infty}\left(R^{\prime}\right)$, since $y_{1}\left(y_{2}\right)$
is rapidly decreasing as $t \rightarrow \infty$ (as $t \rightarrow-\infty$ ). If $R_{\zeta} f=g$, then it is clear that ( $\left.\rho_{\lambda}\left(H^{\prime}\right)-\zeta\right)\left(g-g_{1}\right)=$ 0 . Thus $g-g_{1}$ is an $L^{2}$-eigenfunction and so $g-g_{1} \in \mathcal{X}_{\lambda}$. Thus $g-g_{1}=b_{1}(f) y_{1}$ if $t>0,=b_{2}(f) y_{2}$ if $t<0$, where $b_{1}, b_{2}$ are constants depending on $f$. Now $g \in \mathcal{X}_{2}$ and $g \in \mathcal{D}\left(\mathrm{Cl} \sigma\left(H^{\prime}\right)\right)$ means that $B_{\lambda}\left(H^{\prime}: \varphi: g\right)=0$ for all $\varphi \in \mathcal{D}(\sigma)$. Thus the boundary condition to be satisfied by $g$ is that $\Lambda_{0}(g) \in E(\sigma)$. Or the constants are to be determined from the condition $\Lambda_{0}\left(K_{\zeta} f\right)+$ $b_{1}(f) \Lambda_{0}\left(y_{1}^{+}\right)+b_{2}(f) \Lambda_{0}\left(y_{2}^{-}\right) \in E(\sigma)$ where $y_{1}^{+}=y_{1}$ if $t>0$, and $=0$ if $t<0$. Since $\Lambda_{0}\left(K_{\zeta} f\right), \Lambda_{0}\left(y_{1}^{+}\right)$ and $\Lambda_{0}\left(y_{2}^{-}\right)$are meromorphic in $\zeta$ (in fact, see Lemma 2 , they involve only $\Gamma$-functions), it follows that $b_{1}(f), b_{2}(f)$ depend meromorphically on $\zeta$. Thus the function $\zeta \rightarrow\left(R_{\zeta} f, \varphi\right)$ is a meromorphic function of $\zeta$. This proves then the spectrum is discrete. Since $L^{2}$-eigenfunctions of the operator $\varrho_{\lambda}\left(H^{\prime}\right)$ are in $\mathcal{X}_{\lambda}$, it follows that, if $\sigma\left(H^{\prime}\right)^{*} f=\xi f$, then $f \in \mathcal{X}_{\lambda}$ and $B_{\lambda}\left(H^{\prime}: \varphi: f\right)=0$ for all $\varphi \in \mathcal{D}(\sigma)$. But this implies that $B_{\lambda}\left(H^{\prime m}: \varphi: f\right)=0$ for all $m$, and so $f \in \mathcal{D}\left(\sigma^{\dagger}\right)=\mathcal{D}(\sigma)$. This completes the proof.

In the above lemma we are dealing with the case $\lambda^{2}$ real and $<1$. If $\lambda$ is real and $\geqslant 1$, then we have already seen in Corollary 3 of Lemma 2 that $\sigma_{\lambda}^{ \pm}\left(H^{\prime}\right)$ is essentially self-adjoint in $L^{2}\left(R_{ \pm}\right)$. The proof of the above lemma actually gives the following for this case.

Lemma 13. Let $\lambda \geqslant 1$. Then spectrum of $\sigma_{\lambda}^{ \pm}\left(H^{\prime}\right)$ is discrete. All the eigenfunctions of $\mathrm{Cl} \sigma_{\lambda}^{ \pm}\left(H^{\prime}\right)$ are in $\boldsymbol{X}_{\lambda}^{ \pm}$.

Proof. In this case $R_{\zeta}=\left\{\mathrm{Cl} \sigma_{\lambda}^{ \pm}\left(H^{\prime}\right)-\zeta\right\}^{-1}$, and as in the previous lemma, we have $R_{\zeta} f=K_{\zeta} f$ in this case, for all $f \in C_{c}^{\infty}\left(R_{ \pm}\right)$. Thus spec $\sigma_{\lambda}\left(H^{\prime}\right)$ is discrete.

Combining the previous discussion with Nelson's theorem, we have
Theorem 3. (i) Let $\sigma$ be a self-adjoint $\mathfrak{B l}_{2}$-module in $L^{2}(R)\left(\right.$ in $\left.L^{2}\left(R_{ \pm}\right)\right)$such that $\sigma_{\lambda} \subset \sigma\left(\sigma_{\lambda}^{\dagger} \subset \sigma\right)$, then there exists a unique unitary representation $T$ of the simply connected Lie group of $\mathfrak{l}(2, \boldsymbol{R})$ in $L^{2}(R)\left(\right.$ in $\left.L^{2}\left(R_{ \pm}\right)\right)$such that $d T=\sigma$.
(ii) Let the representations $T_{\lambda, 8, \delta^{\circ}}, T_{\lambda}^{ \pm}$be defined by $d T_{\lambda, \delta, \delta^{\prime}}=\sigma_{\lambda, 8, \delta^{\prime}}$ and $d T_{\lambda}^{+i}=\mu_{\lambda}^{+}$. Then $T_{\lambda, \delta, \delta^{\prime}}$ and $T_{\lambda}^{ \pm}$are all irreducible.
(iii) If $T$ is the unitary representation such that $E(d T)=E_{13}$, then $T=T_{\lambda}^{+} \oplus T_{\lambda}$, if $\lambda \geqslant 0$. Similar results hold for other subspaces listed in Lemma 10.

Proof. Let $\Delta=H^{2}+(X+Y)^{2}+(X-Y)^{2}$. Then $\Delta=\Omega-1-2 H^{\prime 2}$. Then $\varrho_{\lambda}(\Delta)=\lambda^{2}-1-$ $2 \varrho_{\lambda}\left(H^{\prime}\right)^{2}$. From Lemmas 11-13, it follows that $\sigma\left(H^{\prime}\right)$ is essentially self-adjoint and has discrete spectrum. Moreover, all the eigenfunctions belong to $D(\sigma)$. From this it is clear that $\sigma\left(H^{\prime 2}\right)$ is also essentially self-adjoint. So Nelson's theorem now gives (i). The case $\sigma \supset \sigma_{\lambda}^{ \pm}$is discussed similarly.

To prove (ii) note that if $T$ is any unitary representation such that $\sigma_{\lambda} \subset d T$, then
$T(\exp s H) f: t \rightarrow e^{s} f\left(e^{2 s t}\right)$, and $T(\exp s X) f: t \rightarrow e^{-i s t} f(t)$, Any bounded operator which commutes with these two one parameter groups must be scalars on each of the subspaces $L^{2}\left(R_{ \pm}\right)$. If the bounded operator commutes with $T$, then $A \mathcal{D}(d T) \subset \mathcal{D}(d T)$ also. Thus if the boundary condition $E(\sigma)$ relates the boundary values on $R_{+}$and $R_{-}$, the two scalars on $R_{ \pm}$must coincide. This proves that $T_{\lambda, \delta, \delta^{\prime}}$ is irreducible. Similarly $T_{\lambda}^{ \pm}$is always irreducible. The rest is clear.

Theorem 4. (i) Spec $d T_{\lambda}^{+} \cdot\left(H^{\prime}\right)= \pm(\lambda+1+2 N)$ for all $\lambda$ real $\lambda+1>0$. The unitary equivalence class of $T_{\lambda}^{ \pm}$is $\omega^{ \pm}(\lambda)$.
(ii) Spec $d T_{\lambda, \delta, \delta^{\prime}}\left(H^{\prime}\right)=\xi+2 \mathrm{Z}$ where $\xi, \delta, \delta^{\prime}$ and $\lambda$ are related as follows: Case 1. Suppose $\lambda \neq 0$, then

$$
\frac{\cos \pi(\xi-\lambda) / 2}{\cot \pi(\xi+\lambda) / 2}=\frac{\delta^{\prime}}{\delta} \quad \text { or } \quad \frac{\sin \pi(\mu+\lambda) / 2}{\sin \pi(\mu-\lambda) / 2}=e^{i \pi \xi}
$$

where $\delta / \delta^{\prime}=e^{-i \mu \pi}$. Case 2. $\lambda=0$. In this case $\pi \tan \pi \xi / 2=\delta^{\prime} / \delta$. In each of these cases $\xi$ in the spectrum can be chosen uniquely so that $|\xi| \leqslant 1$ and if $\lambda$ is real, then $|\lambda|+|\xi|<1$. The unitary equivalence class of $T_{\lambda, \delta, \delta^{\prime}}$ is $\omega(\xi, \lambda)$.

Proof. (i) In this case $\xi \in \operatorname{Spec} d T_{\lambda}^{+}\left(H^{\prime}\right)$ if and only if $L_{\xi / 2, \lambda / 2} \in \mathcal{X}_{\lambda}^{+}$. From Lemma 2 this happens if and only if $c(\xi,-\lambda)=0$, i.e. iff $\xi \in \lambda+1+2 \mathbf{N}$. A similar argument works for the case $T_{\lambda}^{-}$.

Proof of (ii). In this case $\xi \in \operatorname{Spec} d T_{\lambda, \delta, \delta^{\prime}}\left(H^{\prime}\right)$ if and only if there exists an $f \in \mathcal{X}_{\lambda}$, $\varrho_{\lambda}\left(H^{\prime}\right) f=\xi f$ such that $\Lambda_{0}(f) \in E_{\delta, \delta^{\prime}}$. Now $f=\alpha L_{\xi / 2, \delta / 2}(t)$ for $t>0,=\beta L_{\xi / 2, \lambda / 2}(-t)$ if $t<0$. Then we have two cases.

Case 1. Suppose $\lambda^{2}<1, \lambda \neq 0$. Then from Lemma 2, we have the following

$$
\Lambda_{0}(f)=(\alpha c(\xi, \lambda), \alpha c(\xi,-\lambda), \beta c(-\xi, \lambda), \beta c(-\xi,-\lambda))
$$

Thus $\xi$ belongs to the spectrum if and only if $\beta c(-\xi, \lambda)=\delta \alpha c(\xi, \lambda)$ and $\beta c(-\xi, \lambda)=$ $\delta^{\prime} \alpha c(\xi,-\lambda)$ since $c(\xi, \lambda), c(-\xi,-\lambda)$ cannot both be zero simultaneously it follows that both $\alpha$ and $\beta \neq 0$. (Note $\left|\delta \delta^{\prime}\right|>0$.) Thus $\xi$ belongs to the spectrum if and only if

$$
\frac{c(-\xi, \lambda)}{c(\xi, \lambda)} \frac{c(\xi,-\lambda)}{c(-\xi,-\lambda)}=\frac{\delta}{\delta^{\prime}}
$$

Now

$$
\begin{aligned}
c(\xi, \lambda) c(-\xi,-\lambda) & =\Gamma(-\lambda) \Gamma(\lambda)\{\Gamma(1-\xi-\lambda) / 2) \Gamma((1+\xi+\lambda) / 2)\}^{-1} \\
& =\left(\pi^{-\lambda}\right) \Gamma(-\lambda) \Gamma(\lambda) \cos \pi(\xi+\lambda) / 2
\end{aligned}
$$

Thus we have

$$
\frac{\cos \pi(\xi-\lambda) / 2}{\cos \pi(\xi+\lambda) / 2}=\frac{\delta}{\delta^{\prime}} .
$$

If you put $\delta / \delta^{\prime}=e^{-i \mu \pi}$ and simplify this expression above, we get the equivalent formulation given in the theorem. It is then easy to check that there exists a $\xi$ in the spectrum such that $|\operatorname{Re} \lambda|+|\xi|<1$, when $\lambda$ is real.

Case 2. Suppose $\lambda=0$. In this case we have from Lemma 2 that

$$
\begin{aligned}
\Lambda_{0}(f) & =\left(\alpha(-1)^{(\xi-1) 2}, 0, \beta c(-\xi, 0) d(-\xi), \beta c(-\xi, 0)\right) \quad \text { or } \\
& =\left(\alpha c(\xi, 0) d(\xi), \alpha c(\xi, 0),(-1)^{(1+\xi) 2}, 0\right) \quad \text { or } \\
& =(\alpha c(\xi, 0) d(\xi), \alpha c(\xi, 0), \beta c(-\xi, 0) d(-\xi), \beta c(-\xi, 0))
\end{aligned}
$$

where the first expression holds if $\xi-1 \in 2 \mathbf{N}$, the second one holds if $-(\xi+1) \in 2 \mathbf{N}$, and the last one is valid if $\xi \notin \pm(1+2 N)$. One checks easily from the definition of $E_{\delta, \delta^{\prime}}$ (see Lemma 10 ) that in the first two cases $\Lambda_{0}(f) \notin E_{\delta, \delta^{\prime}}$. Thus $\xi \notin \pm(1+2 N)$ and $\Lambda_{0}(f) \in E_{\delta, \delta^{\prime}}$ implies

$$
\begin{gathered}
\beta c(-\xi, 0)=\delta \alpha c(\xi, 0) \\
\delta \alpha c(\xi, 0) d(\xi)+\delta^{\prime} \alpha c(\xi, 0)-\beta c(-\xi, 0) d(-\xi)=0
\end{gathered}
$$

Note that $a, \beta$ cannot both be zero. If $\alpha=0$, then $c(-\xi, 0)=0$, which is impossible since $\xi \notin \pm(1+2 N)$. Thus both $\alpha \neq 0, \beta \neq 0$. Thus we have

$$
d(\xi)-d(-\xi)=\delta^{\prime} / \delta
$$

Now $d(\xi)-d(-\xi)=\psi((1-\xi) / 2)-\psi((1+\xi) / 2)=\pi \tan \pi \xi / 2$, since $\psi$ is the logarithmic derivative of the $\Gamma$-function. Thus $\pi \tan (\pi \xi) / 2=\delta^{\prime} / \delta$. The rest is clear.

## § 8. Bases of eigenfunctions-the discrete series

We next obtain a basis of $K$-finite vectors for each of the representations $T_{\lambda}^{+}$. We begin with

Lemma 14. Let $L_{\xi / 2 . \lambda / 2}(t)=(2 t)^{-1 / 2} W_{\xi / 2, \lambda / 2}(2 t)(t>0)$ be the eigenfunctions of $\varrho_{\xi}\left(H^{\prime}\right)$ introduced earlier. Then
(i) $\varrho_{\lambda}\left(H^{\prime}\right) L_{\xi / 2, \lambda / 2}=\xi L_{\xi / 2, \lambda / 2}$
(ii) $\varrho_{\lambda}\left(X^{\prime}\right) L_{\xi / 2, \lambda / 2}=-i L_{(\xi+2) / 2, \lambda / 2}$
(iii) $\varrho_{\lambda}\left(Y^{\prime}\right) L_{\xi / 2, \lambda / 2}=-a(a-c+1) L_{(\xi-2) / 2, \lambda / 2}$
where $a=(1-\xi+\lambda) / 2, c=1+\lambda, X^{\prime}$ and $Y^{\prime}$ are defined by (3).

Proof. Let $\Psi(a, c ; x)$ denote Tricomi's confluent hypergeometric function. Then the following identities for $\Psi$ are known (see [3], p. 258)

$$
\begin{gathered}
\left(x \partial_{x}-x+c-a\right) \Psi=-\Psi(a-1, c ; x) \\
\left(x \partial_{x}+a\right) \Psi=a(a-c+1) \Psi(a+1, c ; x) .
\end{gathered}
$$

Since $L_{\xi / 2, \lambda / 2}(t)=(2 t)^{(c-1) / 2} e^{-t} \Psi(a, c ; 2 t)$, the lemma follows from the identities satisfied by $\Psi$.
Theorem 5. (i) Let $\lambda$ be real and $\lambda+1>0$. Define $\psi_{\xi}^{+}(t)=\{c(\xi, \lambda)\}^{-1} L_{\xi / 2, \lambda / 2}(t), t>0$. Then $\psi_{\xi}^{+}, \xi \in \lambda+1+2 \mathbf{N}$ is a $K$-eigenbasis in $L^{2}\left(R_{+}\right)$for the representation $T_{\lambda}^{+}$. Another expression for $\psi_{\xi}^{+}$is the following

$$
\begin{equation*}
\psi_{\xi}^{+}(t)=\left\{\binom{n+\lambda}{n}\right\}^{-1} e^{-t} t^{\lambda / 2} L_{n}^{(\lambda)}(2 t) \tag{13}
\end{equation*}
$$

where $\xi=\lambda+1+2 n$. Also

$$
\left(\psi_{\xi}^{+}, \psi_{\xi}^{+}\right)=\left\{\binom{n+\lambda}{\lambda}\right\}^{-1} \Gamma(\lambda+1) 2^{-(\lambda+1)} .
$$

(ii) If $\mathscr{S}\left(R_{+}\right)$denotes the Schwartz space of $R_{+}$then $\mathcal{D}\left(d T_{\lambda}^{+}\right)=\left\{t^{\lambda^{2} 2} f \mid f \in \mathscr{S}\left(R_{+}\right)\right\}$.
(iii) If $J$ is the anti-unitary isomorphism $J f: t \rightarrow \overline{f(-t)}, f \in L^{2}\left(R_{-}\right)$, then $T_{\lambda}^{-}=J^{-1} \circ T_{\lambda}^{+} \circ J$.

Proof. This theorem can be proved independently of the earlier development. The expression for $\psi_{\xi}^{+}$in terms of the Laguerre polynomials follows from the identity $\Psi(-n, 1+\lambda ; x)=n!(-1)^{n} L_{n}^{(\lambda)}(x)$ (see [3], p. 268). One could deduce this directly, since the differential equation $\varrho_{\lambda}\left(H^{\prime}\right) f=\xi f$ reduces to that of Laguerre polynomials, by putting $f=e^{-t} t^{\lambda / 2} g$. To prove (ii) note that $t^{\gamma / 2} \mathscr{S}\left(R_{+}\right) \subset \mathcal{D}\left(\left(\sigma_{\lambda}^{+}\right)^{\dagger}\right)=\mathcal{D}\left(d T_{\lambda}^{+}\right)$. Let $f$ be a $C^{\infty}$-vector. Then we know (see Lemma 1) $f \in C^{\infty}\left(R_{+}\right)$and $\varrho_{\lambda}(a) f \in L^{2}\left(R_{+}\right)$for all a $\in \mathcal{U}$. From Lemma 1 it follows that $f$ is rapidly decreasing at $\infty$. Now let $f=\Sigma a_{\xi} \psi_{\xi}^{+}$be the eigenfunction expansion of $f$. It follows then that $a_{\xi} \xi^{k}=\left(\varrho_{\lambda}\left(H^{\prime}\right)^{k} j, \psi_{\xi}^{+}\right)$. Now ( $\left.\psi_{\xi}^{+}, \psi_{\xi}^{+}\right)=O\left(n^{-\lambda}\right)$ if $\xi=\lambda+1+2 n$. Thus $a_{\xi}=O\left(n^{-k}\right)$ for every $k$. Now it is known that $\left|L_{n}^{(\lambda)}(x)\right| \leqslant C n^{\mu}$ for all $x, 0<x<1$ (see Szegö [16], p. 176), where $\mu=\operatorname{Max}(1 / 2-1 / 4, \lambda / 2)$. Also (d/dx) $L_{n}^{(\lambda)}=-L_{n-1}^{(\lambda)}$ ([16], p. 101). Thus the series

$$
\Sigma a_{\xi}(d / d t)^{r} L_{n}^{(\lambda)}(2 t)
$$

converges absolutely and uniformly in $(0,1 / 2)$. Then the function $\varphi(t)=\Sigma a_{\xi} L_{n}^{(\lambda)}(2 t) \in C^{\infty}[0,1 / 2)$ and $f=t^{\lambda / 2} e^{-t} \varphi$ in $0<t<1 / 2$, proving $f \in t^{\lambda / 2} \mathscr{S}\left(R_{+}\right)$.

For the part (iii), it is easy to check that $\sigma_{\lambda}^{-}=J^{-1} \circ \sigma_{\lambda}^{+} \circ J$, and since representations $T$, such that $\mathcal{D}(d T) \supset \boldsymbol{X}_{\lambda}^{ \pm}$are unique, the result follows.

The following theorem is known. We state it and sketch a proof since it fits in naturally
with the development here, and we will need it.for another paper. (Foor $S L(2, R)$ see Kunze and Stein [6], Vilenkin [17], and for simply connected covering group of $S L(2, R)$ see Sally [15].)

Theorem 6. The unitary representations $T_{\lambda}^{+}$of the simply connected group $G$ of $S L(2, R)$ may be described by the formulae
(i) $T_{\lambda}^{+}\left(h_{s}\right) f: t \rightarrow e^{s} f\left(e^{2 s} t\right)$
(ii) $T_{\lambda}^{+}\left(u_{s}\right) f: t \rightarrow e^{-i s t} f(t)$
(iii) $T_{\lambda}^{+}(w) f=e^{-i(\lambda+1) \pi / 2} H_{\lambda} f$, where $H_{\lambda}$ is Hankel transform

$$
\boldsymbol{H}_{\lambda} f=\underset{a \rightarrow \infty}{\operatorname{lin} . \mathrm{m}} \cdot \int_{0}^{a} f(s) J_{\lambda}\left(2(s t)^{1 / 2}\right) d s
$$

Proof. The first two statements are clear. The last one can be proved in several ways. It is known that the Hankel transform is a unitary operator and self-reciprocal or $H_{\lambda}^{2}=$ identity. It is thus sufficient to check that $H_{\lambda} \psi_{\xi}^{+}=\exp (-i \xi \pi) \psi_{\xi}^{+}$. But this follows from a known integral formula. (See [4], p. 42, No. (3).) One could also prove it by observing that the operator $\sigma_{\lambda}(Y)$ with domain $t^{\lambda / 2} \mathscr{S}\left(R_{+}\right)$is essentially skew-adjoint, and the operator $H_{\lambda}$ is really the spectral map (or 'diagonalizing' operator) for $\sigma_{\lambda}(Y)$. In other words $H_{\lambda} \circ \sigma_{\lambda}(Y) \circ H_{\lambda}^{-1}=\sigma_{\lambda}(X)$ (see Dunford and Schwartz [2], p. 1535). Since $T_{\lambda}^{+}(w)$ is also a spectral map, it follows that the operator $T_{\lambda}^{+}(w) \circ H_{\lambda}^{+1}$ commutes with $\sigma_{\lambda}(X) ; \sigma_{\lambda}(Y)$ and hence with $\sigma_{\lambda}(Z)$ for all $Z \in \mathfrak{Z l}(2, R)$. Thus the operator $T_{\lambda}^{+}(w) \circ H_{\lambda}^{-1}$ is a scalar. The scalar can be evaluated by evaluating $T_{\lambda}^{+}(w) f$ and $H_{\lambda} \cdot f$, for $f=e^{-t} t^{\lambda / 2}=\psi_{\lambda^{+}}^{+}$. We omit the details.

Remark 1. It is easy to calculate the matrix entry

$$
\left(T_{\lambda}^{+}\left(h_{t}\right) \psi_{\lambda+1}^{+}, \psi_{\lambda+1}^{+}\right)=\int_{0}^{\infty} e^{-y \cosh t} y^{\lambda} d y=\Gamma(\lambda+1) /(\cosh t)^{\lambda+1}
$$

Since any element of $G$ can be written in the form $x\left(\theta_{1}\right) h(t) x\left(\theta_{2}\right)$ with $t \geqslant 0$, and the Haar measure in this decomposition is $\left|e^{2 t}-e^{-2 t}\right| d \theta_{1} d t d \theta_{2}$, it follows that $T_{\lambda}^{+} \in L^{1}(G / Z)$ if $\lambda>1$ and $T_{\lambda}^{+} \in L^{2}$ if $\lambda>0$. These are well known.

Remark 2. Let $P=Z A N$ be the minimal parabolic subgroup of $G$. One has the Bruhat decomposition $G=P \cup N w P . G$ is then generated by $P$ and $w$ and the relations satisfied by $w$ are $w^{2}=\gamma, w h_{t} w^{-1}=h_{-t}$ and for $t \neq 0$,

$$
\begin{equation*}
w u_{t} w=u_{-1 / t} w u_{-t} h_{1 \mathrm{n}|t|} \cdot \gamma^{m(t)} \tag{16}
\end{equation*}
$$

where the integer $m(t)=1$ if $t>0$ and $=0$ if $t<0$. This can be cheeked directly for $S L(2, R)$. However, for the simply connected covering group, the third relation above is not so 14-772905 Acta Mathematica 139. Imprimé le 30 Décembre 1977
obvious. One could presumably use Bargmann's parametrization of $G$ to verify this. Another method would be to use the representations $T_{\lambda}^{+}$for this purpose. For example, let $t>0$, define

$$
F(y, t)=\left\{T_{\lambda}^{+}\left(w u_{t} w\right) \psi_{\lambda+1}^{+}\right\}(y) \cdot e^{t \pi(\lambda+1)}
$$

Then the identity for $w u_{t} w$ gives the following functional equation

$$
e^{i \pi(\lambda+1)(m-1 / 2)} F(y, t)=t^{-1} e^{i y / t} F\left(y / t^{2},-1 / t\right)
$$

The integer $m$ can be determined from the above identity. In fact

$$
F(y, t)=\int_{0}^{\infty} e^{-x(1+i t)} x^{\lambda / 2} J_{\lambda}\left(2(x y)^{1 / 2}\right) d x
$$

and from formula (10), ([4], p. 29), it follows that

$$
F(y, t)=y^{\lambda / 2}(1+i t)^{-(\lambda+1)} \exp (-y / 1+i t)
$$

From this the value of the integer $m$ is easily calculated to be $=1$. Thus $m(t)=1$ for $t>0$. Also it is easy to check directly that $m(t)+m(-t)=1$. Thus $m(t)=0$ for $t<0$, proving the identity completely.

Remark 3. Another formula for $T_{\lambda}^{+}(w)$ is known. (See Sally [15], for details.) It may be described as follows. Let $M: L^{2}\left(R^{+}\right) \rightarrow L^{2}(R)$ be the unitary isomorphism given by the Mellin transform

$$
\begin{equation*}
M f: x \rightarrow \frac{1}{\sqrt{2 \pi}} \int_{0}^{\infty} t^{-(t+i x)} f(t) d t \tag{17}
\end{equation*}
$$

for $f \in C_{c}\left(R_{+}\right)$. Then $T_{\lambda}^{+}(w)=W_{\lambda} \cdot V$, where $V$ is the operator

$$
V t: t \rightarrow \frac{1}{t} f\left(\frac{1}{t}\right)
$$

and $M W_{\lambda} M^{-1}$ is the operator of multiplication by the function

$$
e^{-\operatorname{tin}(\lambda+1) / 2} \Gamma\left(\frac{1+\lambda-2 i x}{2}\right) / \Gamma\left(\frac{1+\lambda+2 i x}{2}\right) .
$$

This can be checked easily from the following facts.
(i) $V T_{\lambda}^{+}\left(h_{s}\right) V^{-1}=T_{\lambda}^{+}\left(h_{-s}\right)$ and thus $W_{\lambda}$ commutes with $h_{s}$ for all $s$.
(ii) $M W_{\lambda} M^{-1}$ is thus a multiplication operator. Finally the multiplier function can be evaluated by using the fact that $e^{-t} t^{\lambda / 2}$ is an eigenfunction for $T_{\lambda}^{+}(w)$. We omit the details.

## § 9. The representations $T_{\lambda, \delta, \delta^{\prime}}$

The representations $T_{\lambda, \delta, \delta^{\prime}}$ were defined infinitesimally by the condition $d T_{\lambda, 8, \delta^{\prime}}=$ $\sigma_{\lambda, \delta, \delta^{\prime}}$, i.e.

$$
\mathcal{D}\left(d T_{\lambda, \delta, \delta^{\prime}}\right)=\left\{f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right) \mid \Lambda_{m}(f) \in E_{\delta, \delta^{\prime}}, \text { for all } m \in \mathbb{N}\right\}
$$

Here $\lambda^{2}$ is real and <1. $E_{\delta, \delta}$ is defined in Lemma 10. A $K$-eigenbasis for the representation is given in the following.

Lemma 15. For each $\xi \in \operatorname{Spec} d T_{\lambda, \delta . \delta^{\prime}}\left(H^{\prime}\right)$, define

$$
\psi_{\lambda . \xi}(t)=\left\{\begin{array}{l}
\{c(\xi, \lambda)\}^{-1} L_{\xi / 2, \lambda / 2}(t), \quad t>0 \\
\delta\{c(-\xi, \lambda)\}^{-1} L_{-\xi / 2 . \lambda / 2}(-t), \quad t<0
\end{array}\right.
$$

Then $\psi_{\lambda, 5}$ is an eigenbasis of $d T_{\lambda 8, \delta^{\prime}}\left(H^{\prime}\right)$. Moreover, the following formulae hold:
(i) $d T_{\lambda, \delta, \delta^{\prime}}\left(H^{\prime}\right) \psi_{\lambda, \xi}=\xi \psi_{\lambda, \xi}$,
(ii) $d T_{\lambda, \delta, \delta^{\prime}}\left(X^{\prime}\right) \psi_{\lambda, \xi}=\{i(1+\xi+\lambda) / 2\} \psi_{\lambda, \xi+2}$,
(iii) $d T_{\lambda, \delta, \delta^{\prime}}\left(Y^{\prime}\right) \psi_{\lambda . \xi}=\{-i(1-\xi+\lambda) / 2\} \psi_{\lambda . \xi-2}$.

Proof. Since $\xi$ is an eigenvalue, the eigenfunction $f_{\xi}$ is of the form $c_{1} L_{\xi / 2, \lambda / 2}(t), t>0$ and $c_{2} L_{-\xi / 2, \lambda / 2}(-t)$, for $t<0$. The constants $c_{1}, c_{2}$ are to be determined from the condition $\Lambda_{0}(f) \in E_{\delta, \delta^{\circ}}$. From the local expansion of $L_{\xi / 2, \lambda / 2}$ it is easy to check that $\Lambda_{0}\left(\psi_{\lambda, \xi}\right) \in E_{\delta, \delta^{\prime}}$. Since the multiplicities are one, it follows that $\left\{\psi_{\lambda, 6}\right\}$ is an eigenbasis. The formulae (i)-(iii) follow from Lemma 14.

Remark. The basis $\psi_{\lambda . \xi}$ is not orthonormal. If $\lambda=i \nu, \nu$ real $\neq 0$, then the identity $\left(X^{\prime} \cdot \psi_{\lambda, \xi}, \psi_{\lambda, \xi+2}\right)=-\left(\psi_{\lambda, \xi}, Y^{\prime} \cdot \psi_{\lambda, \xi+2}\right)$ gives $\left(\psi_{\lambda, \xi+2}, \psi_{\lambda, \xi+2}\right)=\left(\psi_{\lambda, \xi}, \psi_{\lambda, \xi}\right)$, for all $\xi$. Thus all the functions $\psi_{\lambda, \xi}$ have the same norm. Using the formula (40) on page 409 of [4], one can show that in this case

$$
\begin{equation*}
\left(\psi_{\lambda_{0} \xi}, \psi_{\lambda_{, \xi} \xi}\right)=\frac{\nu \tan h(\nu \pi / 2) \sec ^{2}(\pi \xi / 2)}{1+\tan ^{2}(\pi \xi / 2) \tan h^{2}(\nu \pi / 2)} \tag{19}
\end{equation*}
$$

The norm of $\psi_{\lambda . \xi}$ can be evaluated for other values of $\lambda$, but it is no longer independent of $\xi$.
We next describe $C^{\infty}$-vectors of the representation.
Theorem 7. Suppose $\lambda \neq 0$. Then $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ if and only if
(i) $f \in C^{\infty}\left(R^{\prime}\right)$ and $\sup _{|t|>t_{0}}\left|t^{m} \partial_{t}^{n} f\right|<\infty$ for all $t_{0}>0$ and $m, n \in \mathbb{N}$.
(ii) Let $\alpha= \pm \lambda$. Then $\lim _{t \rightarrow 0 \pm}|t|^{\alpha / 2} \varrho_{\lambda}(a) f$ exists for all a belonging to the right ideal $(H-1-\alpha) U+X U$

Moreover, $f \in \mathcal{D}\left(d T_{\lambda, \delta, \delta^{\prime}}\right)$ or is a $C^{\infty}$-vector for the representation if and only if

$$
\begin{equation*}
\lim _{t \rightarrow 0-}|t|^{\alpha / 2} \varrho_{\lambda}(a) f=\delta_{\alpha} \lim _{t \rightarrow 0+}|t|^{\alpha / 2} \varrho_{\lambda}(a) f \tag{20}
\end{equation*}
$$

for all $a \in(H-1-\alpha) \mathcal{U}+X \mathcal{U}, \alpha= \pm \lambda$, where $\delta_{\alpha}=\delta_{i}$ for $\alpha=\lambda, \delta_{\alpha}=\delta^{\prime}$ for $\alpha=-\lambda$.
Proof. Suppose $f \in \mathcal{D}\left(\sigma_{\lambda}^{t}\right)$. Then $f$ satisfies (i) by Lemma 1. Moreover, $f(t)=\mathrm{O}\left(|t|^{-1 / 2}\right)$ as $t \rightarrow 0$ so that $|t|^{\alpha / 2} \varrho_{\lambda}(a) f=0$ for all $a \in X U$. Next

$$
\begin{aligned}
A_{i_{1}^{\prime}}^{ \pm}(f) & =\lim _{t \rightarrow 0 \pm} t W\left(f,|t|^{-\lambda / 2}\right) \\
& =\lim _{t \rightarrow 0 \pm}(-1 / 2)|t|^{-\lambda / 2} \underline{\varrho}_{\lambda}(H-1+\lambda) f .
\end{aligned}
$$

Similarly

$$
A_{2}^{ \pm}(f)=(-1 / 2) \lim _{t \rightarrow 0 . \pm}|t|^{\lambda / 2} \varrho_{\lambda}(H-1-\lambda) f .
$$

Thus if $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$, then (ii) holds. Conversely suppose $f$ satisfies (i) and (ii). Then $\varrho_{\lambda}(H-1-\alpha) f \in L^{2}$ for $\alpha= \pm \lambda$ and thus $f \in L^{2}$. Thus $\varrho_{\lambda}(a) f \in L^{2}$ for all $a \in \mathcal{U}$, or $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$. Next note that if $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$, then

$$
\lim _{t \rightarrow 0 \pm}|t|^{\alpha / 2} \varrho_{\lambda}(a) f=0
$$

if $a \in X . U$. On the other hand, the result $\Lambda_{0} \circ \varrho_{\lambda}(H)=M_{\lambda} \cdot \Lambda_{0}$ of Lemma 8, implies that

$$
\lim _{t \rightarrow 0_{ \pm}}|t|^{\alpha / 2} \varrho_{\lambda}\left\{(H-1-\alpha)^{2} a\right\} t=0
$$

for all $a \in \mathcal{U}$. Thus when $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$, the condition (20) is satisfied when a is of the form $(H-1-\alpha)^{2} \mathcal{U}+X \mathcal{U}$. Now $E_{\delta, \delta^{\prime}}=\left\{\left(x_{1}, x_{2}, x_{3}, x_{4}\right) \in \mathbb{C}^{4} \mid x_{3}=\delta x_{1}, x_{4}=\delta^{\prime} x_{2}\right\}$. Also $\Lambda_{m}=\Lambda_{0} 0 \varrho_{\lambda}(Y)^{m}$, where $\Lambda_{0}=\left(A_{1}^{+}, A_{2}^{+}, A_{1}^{-}, A_{2}^{-}\right)$. Thus if $f \in \mathcal{D}\left(\sigma_{\lambda}^{+}\right)$, then $f$ is a $C^{\infty}$. vector if and only if it satisfies $\Lambda_{m}(f) \in E_{\delta, d^{\prime}}$, this is equivalent to the condition (20).

Theorem 8. Suppose $\lambda=0$. Then $f \in \mathcal{D}\left(\sigma_{\lambda}^{\dagger}\right)$ if and only if
(i) $f \in C^{\infty}\left(R^{\prime}\right)$ and $\sup _{i t t>t_{0}}\left|t^{m} \partial_{t}^{n} f\right|<\infty$ for all $t_{0}>0$ and $m, n \in \mathbf{N}$.
(ii) $\lim _{t \rightarrow 0 \pm} \varrho_{0}(a) f$ exists for all $a \in(H-1) \mathcal{U}+X \boldsymbol{U}$.
(iii) $\lim _{t \rightarrow 0 \pm}\left\{2-\ln |t| \varrho_{0}(H-1)\right\} \varrho_{0}(a) f$ exists for all $a \in \mathcal{U}$. Moreover, $f$ is a $C^{\infty}$-vector for $T_{0, \delta, \delta}$ if and only if
(iv) $\lim _{t \rightarrow 0-} \varrho_{0}(a) f=\delta \lim _{t \rightarrow 0+} \varrho_{0}(a) f$ for all $a \in(H-1) \mathcal{U}+X \mathcal{U}$ and
(v) $\lim _{t \rightarrow 0-}\left\{2-\ln |t| \varrho_{0}(H-1)\right\} \varrho_{0}(a) f$
$=\delta \lim _{t \rightarrow 0^{+}}\left\{2-\ln |t| \varrho_{0}(H-1)\right\} \varrho_{0}(a) f-\delta^{\prime} \lim _{t \rightarrow 0^{+}} \varrho_{0}\{(H-1) a\}$ for all $a \in \mathcal{U}$.
This theorem is proved the same way as the previous one. We omit the details.

Theorem 9. Let $\pi$ be a continuous (not necessarily unitary) representation of $G$ in $L^{2}(R)$ such that
(i) $\pi\left(h_{s}\right) f: t \rightarrow e^{s} f\left(e^{2 s} t\right)$
(ii) $\pi\left(u_{s}\right) f: t \rightarrow e^{-t \cdot s t} f(t)$ for all $s \in R$. Then $C_{c}^{\infty}\left(R^{\prime}\right) \subset \mathcal{D}(d \pi)$ and there exist complex constants $\lambda, \mu$ such that

$$
d \pi\left|C_{\mathrm{c}}^{\infty}\left(R_{+}\right) \supset \sigma_{\lambda}^{\ddagger}, \quad d \pi\right| C_{\mathrm{c}}^{\infty}\left(R_{-}\right) \supset \sigma_{\mu}^{-} .
$$

In particular, if $\pi$ is unitary, then there are only two possibilities: (a) $\pi$ is irreducible and $\pi=T_{\lambda, \delta, \delta^{\prime}}$ for some $\lambda, \delta, \delta^{\prime} ;(\mathrm{b}) \pi$ is reducible and $\pi=T_{\lambda}^{+} \oplus T_{\mu}^{-}$, for some $\lambda, \mu$.

Proof. We first show that $D_{0}=\mathcal{D}(d \pi)$-the set of $C^{\infty}$-vectors of $\pi$ is an $\mathscr{S}(R)$-module; i.e., if $f \in \mathscr{S}(R), \varphi \in \mathcal{D}_{0}$, then $f \varphi \in \mathcal{D}_{0}$. Let $\pi_{1}=\pi \mid A N$. Then one checks by Weyl's lemma (see Appendix A) that $D\left(d \pi_{1}\right)=\left\{f \in C^{\infty}\left(R^{\prime}\right) \mid t^{m}\left(t \partial_{t}\right)^{n} f \in L^{2}(R)\right.$, for all $\left.m, n \in \mathbf{N}\right\}$, and $d \pi_{1}(H) f=$ $\left(2 t \partial_{t}+1\right) f$, and $\left(d \pi_{k}\right)(X) f=-i t f$. Clearly $\mathcal{D}_{0} \subset \mathcal{D}\left(d \pi_{1}\right) \subset C^{\infty}\left(R^{\prime}\right)$. Next $\mathcal{D}_{0}$ is a complete locally convex vector space in the semi-norms, $f \rightarrow\|d \pi \cdot(a) f\|_{2}, a \in \mathcal{U}$. Moreover, $\pi(x) \mathcal{D}_{0} \subseteq \mathcal{D}_{0}$ and $\pi \mid \mathcal{D}_{0}$ is a continuous representation of $G$. If $\varphi \in \mathcal{D}_{0}, g \in L^{1}(R)$, and if

$$
\int|g(s)| v\left(\pi\left(u_{s}\right) \cdot \varphi\right) d s<\infty
$$

for all continuous semi-norms $\boldsymbol{\nu}$ on $\mathcal{D}_{\mathbf{0}}$, then $\hat{g} \varphi=\int g(s) \pi\left(u_{s}\right) \varphi d s \in \mathcal{D}_{\mathbf{0}}$. Now suppose $a \in \mathcal{U}$, $\operatorname{deg} a \leqslant r$, and $\nu(f)=\|d \pi(a) \cdot f\|_{2}, f \in \mathcal{D}_{0}$. From the properties of universal enveloping algebras, it follows that if $a_{1}, a_{2}, \ldots, a_{m} \in \mathcal{U}$ is a basis of the subspace of elements of degree $\leqslant r$, then there exist polynomials $p_{1}, \ldots, p_{m}$ in $s$ such that

$$
\operatorname{Ad} u^{-1} \cdot a=\Sigma p_{j}(s) a_{f} .
$$

Thus

$$
\nu\left(\pi\left(u_{s}\right) \varphi\right) \leqslant \Sigma\left|p_{j}(s)\right|\left\|d \pi\left(a_{j}\right) \varphi\right\|
$$

It thus follows that if $\int|g(s)||s|^{r} d s<\infty$ for all $r \in \mathbf{N}$, then $\hat{g} \varphi \in \mathcal{D}_{0}$, where $\hat{g}$ denotes the Fourier transform of $g$. In particular, $\mathscr{S}(R) \varphi \subset \mathcal{D}_{\mathbf{0}}$. In particular, $C_{c}^{\infty}\left(R^{\prime}\right) \mathcal{D}_{\mathbf{0}} \subset \mathcal{D}_{\mathbf{0}}$.

Since $\pi\left(h_{t}\right) \mathcal{D}_{0} \subset \mathcal{D}_{0}$, it follows that for each $t_{0} \in R^{\prime}$, there exists $\varphi \in \mathcal{D}_{0}$ such that $\varphi\left(t_{0}\right) \neq 0$, and hence a $f_{t_{0}} \in C_{c}^{\infty}\left(R^{\prime}\right) \cap \mathcal{D}_{0}$, such that $t_{t_{0}}=1$ in a neighborhood of $t_{0}$. This implies easily that $C_{c}^{\infty}\left(R^{\prime}\right) \subset \mathcal{D}_{0}$. Let $D=d \pi(Y)$. It will be shown next that $D$ is a second order differential operator. In fact $\pi\left(u_{s}^{-1}\right) D \pi\left(u_{s}\right) \varphi=d \pi\left(e^{-s a d X} Y\right) \varphi=d \pi\left(Y \pm s H-s^{2} X\right) \varphi=\varphi_{1}-s \varphi_{2}-$ $s^{2} \varphi_{3}$, where $\varphi_{1}=D \varphi, \varphi_{2}=d \pi(H) \varphi, \varphi_{3}=d \pi(X) \varphi$. Thus

$$
D \pi\left(u_{s}\right) \varphi=e^{-i s t} \varphi_{1}-s e^{-i s t} \varphi_{2}-s^{2} e^{-i s t} \varphi_{3}
$$

Let

$$
f(t)=\int g(s) e^{-i s t} d s, \text { with } g \in \mathscr{S}(R)
$$

Then

$$
D(t \varphi)=\int g(s) D \pi\left(u_{s}\right) \varphi d s=f \varphi_{1}-i\left(\partial_{t} f\right) \varphi_{2}+\left(\partial_{t}^{2} f\right) \varphi_{3}
$$

Now suppose $f \in C_{c}^{\infty}\left(R^{\prime}\right)$ and $\varphi \in C_{c}^{\infty}\left(R^{\prime}\right)$ with $\varphi=1$ on supp $f$. Then $f \varphi=f, \varphi_{2}=1$ on supp $f$, $\varphi_{3}=-i t$ on supp $f$. Thus

$$
D f=f \varphi_{1}-\left(i \partial_{t} f\right)-i t\left(\partial_{t}^{2} f\right)
$$

where $f \varphi_{1}=f D \varphi$, for any $\varphi$ in $C_{c}^{\infty}\left(R^{\prime}\right)$ equal to 1 on $\operatorname{supp} f$. From this it follows that $D \mid C_{c}^{\infty}\left(R^{\prime}\right)$ is a second order differential operator of the form $-i\left(t \partial_{t}^{2}+\partial_{t}+\psi\right)$. Now the commutation rule $[H, Y]=-2 Y$, gives, that there exist constants $c_{1}, c_{2}$ such that $\psi=c_{1} / t$ for $t>0,=c_{2} / t$ for $t<0$. Thus $d \pi\left|C_{c}^{\infty}\left(R_{+}\right) \supset \sigma_{\lambda}^{+}, d \pi\right| C_{c}^{\infty}\left(R_{-}\right) \supset \sigma_{\mu}^{-}$for some $\lambda, \mu$. Unitarity of $\pi$ implies $\lambda^{2}, \mu^{2}$ are real. If $\pi$ is irreducible, then $d \pi(\Omega)=$ constant, so that $\lambda=\mu$ in this case. In this case $d T_{\lambda, 8,8}$ are the only irreducible self-adjoint extensions of $\sigma_{\lambda}$. Thus $\pi=T_{\lambda, \delta, \delta^{\prime}}$ for some $\lambda, \delta, \delta^{\prime}$. The second case is proved similarly.

For the representations $T_{\lambda, \delta, \delta}$, an analogue of Theorem 6 can be given via the twocomponent Mellin transform. (See also Remark 3 following Theorem 6.) (In this connection see Sally [15], Vilenkin [17], for another approach.) Let $M: L^{2}\left(R_{+}\right) \rightarrow L^{2}(R)$ be the Mellin transform defined earlier (17). Define $M^{\prime}: L^{2}\left(R_{-}\right) \rightarrow L^{2}(R)$ by $M^{\prime} f=M f^{0}, f^{0}(t)=f(-t)$. Put $M_{2} f=\left(M f, M^{\prime} f\right)$. Then $M_{2}$ gives a unitary isomorphism of $L^{2}\left(R^{\prime}\right)$ with $L^{2}(R) \otimes \mathbf{C}^{2}$ and is called the two-component Mellin transform. Let $A(x)=\left(a_{i j}(x)\right)_{1 \leqslant 1,1 \leqslant 2}$ be a unitary matrix valued function on $R$. Define the operator $I \otimes A$ on $L^{2}(R) \otimes \mathbf{C}^{2}$ as follows:

$$
I \otimes A\binom{f_{1}}{f_{2}}=\binom{a_{11}(x) f_{1}(x)+a_{12}(x) f_{2}(x)}{a_{21}(x) f_{1}(x)+a_{22}(x) f_{2}(x)} .
$$

We call $I \otimes A$ a matrix multiplication operator.

Theorem 10. Let $V$ be the unitary operator $V f: t \rightarrow|t|^{-1} f(1 / t)$, in $L^{2}(R)$. Put $W_{2, \delta, \delta^{\prime}}=$ $T_{\lambda, 8 . \delta \cdot(w) V . T h e n}$

$$
M_{2} \circ W_{\lambda, \delta, \delta^{\prime}} \circ M_{2}^{-1}=I \otimes A_{\lambda, \delta, \delta^{\prime}}
$$

for a suitable matrix-function $A_{\lambda, \delta, 8^{\circ}}$, and its matrix entries have explicit formulas in terms of the gamma and hypergeometric functions.

Proof. For the sake of brevity we use the module notation and write $h_{s}$ instead of the operator $T_{\lambda, \delta, \delta}\left(h_{s}\right)$. One checks easily that $V h_{s} V^{-1}=h_{-s}$ and $w h_{s} w^{-1}=h_{-s}$ in $G$. So it follows that the operator $W_{\lambda, \delta, \delta^{\prime}}$ commutes with the operator $h_{s}$. However, $M_{2} h_{s} M_{2}^{-1}$ is just the operator of multiplication by $\exp (-2 i s x)$. Thus it follows from known results that $W_{\lambda, \delta, \delta^{\prime}}$ is a matrix multiplication operator. To determine this matrix we use that $\psi_{\lambda .5}$ are eigenfunctions. Note $M_{2} V M_{2}^{-1}: \varphi \rightarrow \varphi^{0}$ for $\varphi \in L^{2}(R) \otimes \mathbf{C}^{2}, \varphi^{0}(x)=\varphi(-x)$. Now put

$$
g_{\lambda, \xi}(x)=\frac{1}{\sqrt{2 \pi}} \int_{0}^{\infty} t^{-(1 / 2+t x)} \psi_{\lambda, \xi}(t) d t=M \psi_{\lambda, \xi} .
$$

Then

$$
M_{2} \psi_{\lambda, \xi}=\binom{g_{\lambda, \xi}(x)}{\delta g_{\lambda,-\xi}(x)}
$$

Since $\psi_{\lambda . \xi}$ is an eigenfunction for the eigenvalue $\exp \left(-\frac{1}{2} i \pi \xi\right)$, it follows that

$$
A\binom{g_{\lambda_{k} \xi}(-x)}{\delta g_{\lambda_{0}-\xi}(-x)}=e^{-i \pi \xi / 2}\binom{g_{\lambda_{, \xi}}(x)}{\delta g_{\lambda_{,-\xi}}(x)}
$$

for all $\xi$ in the spectrum. Now it is easy to check that if

$$
A\binom{\alpha_{1}}{\alpha_{2}}=\binom{\beta_{1}}{\beta_{2}}, \quad A\binom{\alpha_{1}^{\prime}}{\alpha_{2}^{\prime}}=\binom{\beta_{1}^{\prime}}{\beta_{2}^{\prime}}
$$

then

$$
A=\left(\begin{array}{ll}
\beta_{1} & \beta_{1}^{\prime} \\
\beta_{2} & \beta_{2}^{\prime}
\end{array}\right)\left(\begin{array}{ll}
\alpha_{1} & \alpha_{1}^{\prime} \\
\alpha_{2} & \alpha_{2}^{\prime}
\end{array}\right)^{-1}
$$

Let $\xi, \xi^{\prime}$ be in the spectrum of $H^{\prime}$ and $\xi \neq \xi^{\prime}$. Then we have, using the above formulas,

$$
\left(\begin{array}{ll}
1, & 0 \\
0, & 1 / \delta
\end{array}\right) A\left(\begin{array}{ll}
1, & 0 \\
0, & \delta
\end{array}\right)=Q(x)\left(\begin{array}{ll}
e^{i \pi \xi / 2} & 0 \\
0 & e^{-i \pi \xi^{\prime} / 2}
\end{array}\right) Q(-x)^{-1}
$$

where

$$
Q(x)=\left(\begin{array}{ll}
g_{\lambda, \xi}(x) & g_{\lambda, \xi^{\prime}}(x) \\
g_{\lambda,-\xi}(x) & g_{\lambda,-\xi^{\prime}}(x)
\end{array}\right) .
$$

From the formula for Mellin transform of $W_{\varkappa, \mu}(x)$ ([4], p. 337), it follows that

$$
g_{\lambda .5}(x)=\frac{2^{b}}{\sqrt{2 \pi} \Gamma(-\lambda)} \frac{\Gamma(a) \Gamma(b) \Gamma(c-a)}{\Gamma(c)} F\left(a, b ; c: \frac{1}{2}\right)
$$

where $a=(1+\lambda-2 i x) / 2, b=(1-\lambda-2 i x) / 2, c=1-i x-\xi / 2$. This completes the proof.
Finally we mention here the intertwining operator connecting $T_{\lambda, \delta, \delta^{\prime}}$ with the unitary principal series. Let $\eta \otimes e^{\lambda}$ denote the representation $\gamma^{m} h_{t} u_{s} \rightarrow e^{i \eta m \pi-\lambda t}$ of $P$ and $U_{\eta . \lambda}=$
$\operatorname{Ind}_{P \uparrow G} \eta \otimes e^{\lambda}$. The following realization of $U_{\eta . \mathrm{A}}$ in $L^{2}(R)$ is well known. It comes from using $G=P \cup N w P$ and using

$$
f \rightarrow \int_{R} f\left(u_{t} w P\right) d t
$$

as the quasi-invariant measure in $G / P$. In fact, we have

$$
\begin{equation*}
U_{\eta, \lambda}(x) f: t \rightarrow|-c t+a|^{-(1-\lambda)} e^{i \eta m\left(x^{-1}, t\right) \pi} f\left(\frac{d t-b}{-c t+a}\right) \tag{21}
\end{equation*}
$$

where $\sigma(x)=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$, and $x \rightarrow \sigma(x)$ is the covering homomorphism and the integer valued function $m(x, t)$ is defined as follows:

$$
x \cdot u_{t} \cdot w\left\{\begin{array}{l}
\in P, \quad \text { if } \quad c t+d=0 \\
=u_{(a t+b))(c t+d)} w p(x, t), \quad \text { if } \quad c t+d \neq 0 .
\end{array}\right.
$$

The element $p(x, t) \in P$ and $m(x, t)$ is defined uniquely by

$$
\begin{equation*}
p(x, t)=u_{-c(c t+d)} h_{\ln |c t+d|} \gamma^{m(x, t)} . \tag{23}
\end{equation*}
$$

We have already seen that $m(w, t)=1$ if $t>0,=0$ if $t<0$ and $m(x, t)$ for arbitrary $x$ can be found by using the identities

$$
m(x \cdot y, t)=m(y, t)+m(x, \sigma(y) \cdot t) \quad \text { and } \quad m\left(\gamma^{\tau} y, t\right)=r \quad \text { if } \quad y \in A N
$$

Let $\mathcal{I}$ denote the Fourier transform and let $F_{\lambda}$ denote the operator $F_{\lambda}=|t|^{\lambda / 2} \circ \mathcal{F}$, $F_{\lambda}^{-1}=\mathcal{F}^{-1} \circ|t|^{-\lambda / 2}$. Then $F_{\lambda}$ is unitary for $\lambda \in i R$ and for $f \in C_{c}(R)$

$$
F_{\lambda} f: t \rightarrow|t|^{\lambda / 2} \frac{1}{\sqrt{2 \pi}} \int_{R} e^{-i y t} f(y) d y
$$

and

$$
F_{\lambda}^{-1} f: x \rightarrow \frac{1}{\sqrt{2 \pi}} \int_{R} e^{i x t} f(t)|t|^{-\lambda / 2} d t
$$

We then have

Theorem 11. Let $\lambda \in i R,-1<\eta \leqslant 1$ and suppose $(\lambda, \eta) \neq(0,1)$. Then

$$
T_{\lambda, \delta, \delta}=F_{\lambda} \circ U_{\eta, \lambda} \circ F_{\lambda}^{-1}
$$

where $\delta=-1, \delta^{\prime}=-\{\cos \pi(\eta+\lambda) / 2\} /\{\cos (\eta-\lambda) / 2\}$ if $\lambda \neq 0$, and $\delta=-1, \delta^{\prime}=-\pi \tan (\pi \eta / 2)$ if $\lambda=0$.

Proof. Let $S=F_{\lambda} \circ U_{\eta, \lambda} \circ F_{\lambda}^{-1}$. Then direct calculation shows that $S\left(h_{s}\right) f: t \rightarrow e^{8} f\left(e^{2 s} t\right)$ and $S\left(u_{s}\right) f: t \rightarrow e^{-i s t} f(t)$. Also $U_{\eta, \lambda}$ is known to be irreducible. Thus by Theorem 9 , it follows
that there exists $\delta, \delta^{\prime}$ such that $T_{\lambda, \delta, \delta^{\prime}}=S=F_{\lambda} \circ U_{\eta, \lambda} \circ F_{\lambda}^{-1}$. We have to find $\delta$ and $\delta^{\prime}$. These are obtained by comparing $K$ eigenfunctions in both. First of all Spec $d U_{\eta, \lambda}\left(H^{\prime}\right)=\eta+2 \mathbf{Z}$. Comparing it with spectrum of $d T_{\lambda, \delta, \delta^{\prime}}\left(H^{\prime}\right)$ we have (see Theorem 4)

$$
\delta^{\prime} / \delta=\{\cos \pi(\eta+\lambda) / 2\} /\{\cos \pi(\eta-\lambda) / 2\} \quad \text { if } \quad \lambda \neq 0 \quad \text { and }=\pi \tan \pi \eta / 2, \quad \text { if } \lambda=0
$$

Next consider the representation $d U_{\eta . \lambda}$. From the formulas for the one-parameter groups $U_{\eta, \lambda}(\exp s Z)$, where $Z=H$ or $X$ or $Y$, it is not difficult to show that $C_{c}^{\infty}(R) \subset$ $\mathcal{D}\left(d U_{\eta, \lambda}\right)$ and $d U_{\eta, \lambda} \supset \tau_{\lambda}$ where $\mathcal{D}\left(\tau_{\lambda}\right)=C_{c}^{\infty}(R)$ and

$$
\tau_{\lambda}(H)=-\left(2 t \partial_{t}+1-\lambda\right), \quad \tau_{\lambda}(X)=-\partial_{t}, \quad \text { and } \quad \tau_{\lambda}(Y)=\left(t^{2} \partial_{t}+(1-\lambda) t\right)
$$

By Weyl's lemma $\mathcal{D}\left(\tau_{\lambda}^{\dagger}\right)=\left\{f \in C^{\infty}(R) \mid \tau_{\lambda}(a) f \in L^{2}(R)\right.$, for all $\left.a \in \mathcal{U}\right\}$ and for $f \in \mathcal{D}\left(\tau_{\lambda}^{\dagger}\right)$, $\tau_{\lambda}^{\dagger}(a) f=\tau_{\lambda}(a) f$. Since $U_{\eta, \lambda}$ is unitary, it follows that $d U_{\eta, \lambda} \subset \tau_{\lambda}^{\dagger}$. Thus the eigenfunctions $\left\{d U_{\eta, \lambda}\left(H^{\prime}\right) f=\xi f\right\} \subset\left\{f \in C^{\infty}(R) \mid \tau_{\lambda}\left(H^{\prime}\right) f=\xi f\right\}$. By direct computation, these eigenspaces are one-dimensional and if $f_{\xi}$ is an eigenfunction then

$$
f_{\xi}=(1+i t)^{-(1-\lambda-\xi) / 2}(1-i t)^{-(1-\lambda+\xi) / 2}
$$

for $\xi \in \eta+2 Z$. Thus there exists a constant $\beta_{\xi}$ such that $F_{\lambda} f_{\xi}=\beta_{\xi} \psi_{\lambda, \xi}$, where $\psi_{\lambda, \xi}$ is the eigenfunction introduced in Lemma 15. From formula (12) on page 119 of [4], it follows that

$$
F_{\lambda} f_{\xi}=\left\{\begin{array}{c}
-\sqrt{2 \pi} 2^{\lambda / 2}\{\Gamma((1-\lambda+\xi) / 2)\}^{-1} L_{\xi / 2, \lambda / 2}(t), \quad t>0 \\
\sqrt{2 \pi} 2^{\lambda / 2}\{\Gamma((1-\lambda-\xi) / 2)\}^{-1} L_{-\xi / 2, \lambda / 2}(-t), \quad t<0 .
\end{array}\right.
$$

Comparing this with the formula for $\psi_{\lambda, \xi}$, one gets that $\delta=-1$ always. This completes the proof.

## Appendix A. On extensions of symmetric $\boldsymbol{g}$-modules

Let $\mathcal{U}$ be an associative algebra over $\mathbf{C}$ and $\tau$ an involutory, conjugate linear antiautomorphism of $\mathcal{U}$. We write $a^{\boldsymbol{\tau}}$ instead of $\tau(a)$ for $a \in \mathcal{U}$. A $\boldsymbol{U}$-module $\sigma$ consists of a complex vector space $\mathcal{D}(\sigma)$, called the domain of $\sigma$, and linear operators $\sigma(a): \mathcal{D}(\sigma) \rightarrow \mathcal{D}(\sigma)$, such that $a \rightarrow \sigma(a)$ is a homomorphism of $\mathcal{U}$. If $\sigma_{1}, \sigma_{2}$ are two such modules, we say $\sigma_{1} \subset \sigma_{2}$ if $\mathcal{D}\left(\sigma_{1}\right) \subset \mathcal{D}\left(\sigma_{2}\right)$ and $\sigma_{1}(a)=\sigma_{2}(a) \mid \mathcal{D}\left(\sigma_{1}\right)$, for all $a \in \mathcal{U}$. Let $\mathcal{H}$ be a complex separable Hilbert space. We shall be mainly concerned with $\mathcal{U}$-modules in $\boldsymbol{\mathcal { H }}$, i.e., $\sigma$ such that $\mathcal{D}(\sigma) \subset \mathcal{H}$. In the following we write $(\cdot, \cdot)$ for the inner product on $\mathcal{H}$ and use standard terminology for operators in a Hilbert space.

Definition of the adjoint module. Let $\sigma$ be a densely defined $U$-module, i.e., $\mathcal{D}(\sigma)$ is dense in $\mathcal{H}$. Then a vector $g$ belongs to the domain of the adjoint of $\sigma$ or $g \in \mathcal{D}\left(\sigma^{\tau}\right)$ if
for each $a \in \mathcal{U}$, there is $g_{a}$ such that $(\sigma(a) f, g)=\left(f, g_{a}\right)$ for all $f \in \mathcal{D}(\sigma)$. Then the vector $g_{a}$ is unique and we define $\sigma^{\tau}(a) g=g a^{\tau}$. It is then clear that if $g \in \mathcal{D}\left(\sigma^{\tau}\right), g_{a} \in \mathcal{D}\left(\sigma^{\tau}\right)$ for all $a$ and $\left(g_{a}\right)_{b}=g_{a b}$. Thus $\sigma^{\tau}$ is a $\mathcal{U}$-module called the adjoint of $\sigma$. Also note that $\mathcal{D}\left(\sigma^{\tau}\right)=$ $\cap\left\{\mathcal{D}\left(\sigma(a)^{*}\right) \mid a \in \mathcal{U}\right\}$ and $\sigma^{\tau}(a) \subset\left(\sigma\left(a^{\tau}\right)\right)^{*}$.
$A U$-module $\sigma$ is said to be $\tau$-symmetric (or simply symmetric when it is clear from the context what $\tau$ is meant) if $\sigma \subset \sigma^{\tau}$. This is equivalent to the statement

$$
(\sigma(a) f, g)=\left(f, \sigma\left(a^{\tau}\right) g\right)
$$

for all $f, g \in \mathcal{D}(\sigma)$.
The following lemma is proved the same way as in the case of a single operator.
Lemma A.1. (i) If $\sigma_{1}$ is densely defined and $\sigma_{1} \subset \sigma_{2}$, then $\sigma_{2}^{\tau} \subset \sigma_{1}^{\tau}$,
(ii) If $\sigma$ and $\sigma^{\tau}$ are both densely defined, then $\sigma^{\tau \tau}=\left(\sigma^{\tau}\right)^{\tau}$ exists, $\sigma \subset \sigma^{\tau \tau}$ and $\sigma^{\tau \tau \tau}=\sigma^{\tau}$.

Lemma A.2. Let $\sigma$ be a densely defined, symmetric $\mathcal{U}$-module. Let $B_{\sigma}(a: f: g)=$ $\left(\sigma^{\tau}(a) f, g\right)-\left(f, \sigma^{\tau}\left(a^{\tau}\right) g\right)$ for $a \in \mathcal{U}, f, g \in \mathcal{D}\left(\sigma^{\tau}\right)$. Then the boundary forms $B_{\sigma}$ satisfy the following identities:
(i) $B_{\sigma}(a: f: g)=0$ if either $f$ or $g \in \mathcal{D}(\sigma)$,
(ii) $B_{\sigma}(a b: f: g)=B_{\sigma}\left(a: \sigma^{\tau}(b) f: g\right)+B_{\sigma}\left(b: f: \sigma^{\tau}\left(a^{\tau}\right) g\right)$ and
(iii) $B_{\sigma}\left(a^{n}: f: g\right)=\sum_{r=0}^{n-1} B_{\sigma}\left(a: \sigma_{\tau}(a)^{r} f: \sigma^{\tau}\left(a^{\tau}\right)^{n-r-1} g\right)$. Moreover, $\sigma \subset \sigma^{\tau \tau} \subset \sigma^{\tau}$ and $\mathcal{D}\left(\sigma^{\tau \tau}\right)=$ $\left\{f \in \mathcal{D}\left(\sigma^{\tau}\right) \mid B_{\sigma}(a: f: g)=0\right.$ for all $\left.g \in \mathcal{D}\left(\sigma^{\tau}\right)\right\}$.
(iv) If $\sigma \subset \sigma_{1} \subset \sigma^{\tau}$, then $\sigma \subset \sigma_{1}^{\tau} \subset \sigma^{\tau}$ and $\mathcal{D}\left(\sigma_{1}^{\tau}\right)=\left\{f \in \mathcal{D}\left(\sigma^{\tau}\right) \mid B_{\sigma}(a: f: g)=0\right.$ for all $a \in \mathcal{U}$ and $\left.g \in \mathcal{D}\left(\sigma_{1}\right)\right\}$.

The proofs are all straightforward and are omitted.
There are two prime examples of the pairs $(\mathcal{U}, \tau)$. In the first one, let $\mathcal{U}=\boldsymbol{C}[t]$, the polynomial algebra in one indeterminate $t$, and $\tau$ is defined by $\tau(t)=t$. In this case, a $\mathcal{U}$. module is defined by an operator $A$ with domain $\mathcal{D}$ such that $A D \subset \mathcal{D}$ and $\sigma\left(t^{n}\right)=A^{n}$. In this case, $\mathcal{D}\left(\sigma^{x}\right)=\cap \mathcal{D}\left(A^{* n}\right)=\mathcal{D}^{\infty}\left(A^{*}\right)$, and $\tau$ symmetry coincides with the usual notion of symmetry for a single operator.

The second example arises naturally in representation theory of Lie groups. Let $\mathfrak{g}$ be a Lie algebra over $R$, and $\mathfrak{g}_{c}$ its complexification. Let $\mathcal{U}=\boldsymbol{U}\left(\mathfrak{g}_{c}\right)$ denote the universal enveloping algebra of $\mathfrak{g}_{c}$. There exists a unique conjugate linear involutory anti-automorphism $a \rightarrow a^{\dagger}$ of $U$ onto itself, such that $X^{\dagger}=-X$ for all $X \in g$. When dealing with Lie algebras, the pair $\left(\mathcal{U}\left(g_{c}\right), \dagger\right)$ is the one we shall be concerned with. Since a $\mathfrak{g}$-module extends uniquely to a $\mathcal{U}\left(g_{c}\right)$-module, we shall treat the concepts synonymously.

Let $G$ be the simply connected Lie group with Lie algebra $g$. Let $T$ be a continuous re-
presentation of $G$ in $\mathcal{H}$. We define $d T$ as the $g$-module whose domain is the collection of all $C^{\infty}$-vectors of the representation $T$, and when $v \in \mathcal{D}(d T), d T \cdot(X) v=\left.(d / d t)\right|_{t=0} T(\exp t X) v$.

The following lemma is well known and is stated here only for ease of reference (see Warner [18], chapter 4).

Lemma A.3. (i) $A$ vector $v \in \mathcal{D}(d T)$ if and only if the function $x \rightarrow\left(T(x) v, v^{\prime}\right)$ is a $C^{\infty}$ function on $G$, for each $v^{\prime} \in \mathcal{H}$.
(ii) For each $X \in \mathfrak{g}$, the closure of $d T(X)$ is the infinitesimal generator of the one parameter group $t \rightarrow T(\exp t X)$ and

$$
\mathcal{D}(d T)=\cap\left\{\mathcal{D}\left(A_{j 1} \ldots A_{j r}\right) \mid r=1,2, \ldots ; A_{j}=\mathrm{Cl} d T \cdot\left(X_{j}\right)\right.
$$

where $X_{1}, \ldots, X_{n}$ is a basis of $\left.\mathfrak{g}\right\}$.

Lemma A.4. Let $\sigma$ be a g-module and $T$ a unitary representation such that $\sigma \subset d T$. Assume that $\mathcal{D}(\sigma)$ is dense and $T(x) \mathcal{D}(\sigma) \subset \mathcal{D}(\sigma)$ for all $x \in G$. Then, (i) $\sigma^{\dagger}=d T$; (ii) If $z \in$ cent $\mathcal{U}\left[g_{c}\right]$ and $z^{\dagger}=z$, then $\sigma(z)$ is essentially self-adjoint; (iii) If $X \in \mathfrak{g}$; then $\sigma(X)$ is essentially skew adjoint.

Proof. Now $\sigma \subset d T$ implies that $(d T)^{\dagger} \subset \sigma^{\dagger}$. Since $T$ is unitary, we have $d T \subset(d T)^{\dagger}$. Thus it is sufficient to check that $\mathcal{D}\left(\sigma^{\dagger}\right)=\mathcal{D}(d T)$. Let $f \in \mathcal{D}\left(\sigma^{\dagger}\right)$ and let $g \in \mathcal{H}$ be arbitrary. It is sufficient to verify that the function $x \rightarrow \varphi_{g, f}(x)=(g, T(x) f)$ is $C^{\infty}$. Let $g_{n} \in D(\sigma)$ be such that $g_{n} \rightarrow g$ in $\mathcal{H}$. Then

$$
\varphi_{g_{n} . f}(x ; \eta)=\left(d T\left(\eta^{0}\right) T\left(x^{-1}\right) g_{n}, f\right)=\left(T\left(x^{-1}\right) g_{n}, \sigma\left(\eta^{0}\right)^{*} f\right), \quad \text { since } \quad D(\sigma)
$$

is $G$-invariant. Thus

$$
\varphi_{s_{n}, f}(x ; \eta)=\varphi_{a_{n}, \sigma^{+}(\eta) r}(x)
$$

for all $\eta \in \mathcal{U}[g]$. Thus $\varphi_{g_{n}, f}(x ; \eta)$ converges uniformly on compact subsets to $\varphi_{o . a^{\dagger}(\eta) r}(x)$. Since $\varphi_{o_{n, f}}$ is $C^{\infty}$, it follows that $\varphi_{o . f}$ is also $C^{\infty}$ or $f \in \mathcal{D}(d T)$. The rest of the statements of the lemma are known, cf. [8] or [9].

Corollary 1. $(d T)^{\dagger}=d T$.

Corollary 2. With the same assumptions on $\sigma$ as in Lemma A.4, let $\sigma_{1} \subset \sigma$ be a $\mathfrak{g}$-submodule such that $\mathcal{D}\left(\sigma_{1}\right)$ is dense. Then $\sigma_{1} \subset \sigma \subset d T \subset \sigma_{1}^{\dagger}$ and if $\mathcal{D}=\left\{f \in \mathcal{D}\left(\sigma_{1}^{\dagger}\right) \mid\left(\sigma_{1}^{\dagger}(\eta) g, f\right)=\right.$ $\left(g, \sigma_{1}^{\dagger}\left(\eta^{\dagger}\right) f\right)$ for all $g \in \mathcal{D}(\sigma)$ and all $\left.\eta \in \mathcal{U}\right\}$, then $\mathcal{D}=\mathcal{D}(d T)$.

Proof. The first corollary follows from Lemma A. 4 by taking $\sigma=d T$. Since $\sigma_{1} \subset \sigma \subset d T$, it follows that $(d T)^{\dagger}=d T \subset \sigma_{1}^{\dagger}$. Since $\sigma \subset d T \subset \sigma_{1}^{\dagger}$, it follows that $\mathcal{D}(d t) \subset \mathcal{D}$. On the other hand $\mathcal{D} \subset \mathcal{D}\left(\sigma^{\dagger}\right)$. But $\sigma^{\dagger}=d T$, so that $\mathcal{D}(d t)=\mathcal{D}$.

Definition A.1. A symmetric g -module $\sigma$ is said to be integrable if there exists a continuous unitary representation $T$ of the simply connected Lie group in $\mathcal{H}$ such that $\sigma \subset d T$. It is said to be exact if $\sigma=d T$.

We note that if $\sigma=d T_{1}=d T_{2}$, then $T_{1}=T_{2}$. In fact, $T_{1}(\exp t X)$ and $T_{2}(\exp t X)$ have the same infinitesimal generator, so that $T_{1}(\exp X)=T_{2}(\exp X)$ for all $X \in \mathfrak{g}$. However, if $\sigma \subset d T, d T$ is not in general unique. In this connection we note the following

Lemma A.5. Let $\sigma$ be a densely defined g -module such that $\sigma\left(X_{j}\right)$ is essentially skewadjoint for a basis $X$, of $\mathfrak{g}$. Then $\sigma^{+\dagger}=\sigma^{\dagger}$. In particular, if $\sigma \subset d T$, then $d T=\sigma^{\dagger}$, and $T$ is unique.

Proof. If $\sigma\left(X_{j}\right)$ is essentially skew-adjoint, it follows that the boundary forms $B_{\sigma}\left(X_{j}: f: g\right)=0$ for all $f, g \in \mathcal{D}\left(\sigma^{\dagger}\right)$ and, therefore, $B_{o}(a: f: g)=0$ for all $a \in \mathcal{U}$, from identities satisfied by the boundary forms $B_{\sigma}$ (see Lemma A.2). Thus $\sigma^{\dagger \dagger}=\sigma^{\dagger}$, by the same lemma. Now if $\sigma \subset d T$, then $\sigma \subset d T \subset \sigma^{\dagger}$, and so $\sigma^{\dagger \dagger} \subset d T \subset \sigma^{\dagger}$ or $d T=\sigma^{\dagger}$.

The following is just a reformulation in our notation of a theorem of Nelson ([8]).
Theorem. Let $X_{1}, \ldots, X_{n}$ be a basis of $\mathfrak{g}$ and $\sigma$ a $\mathfrak{g}$-module in $\mathcal{H}$. Then $\sigma=d T$ for some continuous unitary representation $T$ if and only if (i) $\sigma=\sigma^{\dagger}$, and (ii) $\sigma(\Delta)$ is essentially self. adjoint, where $\Delta=X_{1}^{2}+\ldots+X_{n}^{2}$.

Remark. An example of Nelson (see [8], section 11) may be interpreted in our notation as follows: there exists $g$-module $\sigma$, of a two-dimensional abelian Lie algebra such that $\sigma=\sigma^{\dagger}$, but $\sigma$ is not integrable. Whether such examples exist for semi-simple $g$ is not known. In this connection it might be of interest to note that all self-adjoint extensions of the module $\sigma_{\lambda}$ considered in Section 3 are integrable to the group (cf. Theorem 3).

Let $M$ be a $C^{\infty}$-manifold and $\mu$ a $C^{\infty}$-density on $M$, i.e., $\mu:(U, \varphi) \rightarrow \mu_{U . \varphi}$ a map from local charts $(U, \varphi)$ to $C^{\infty}$-functions on $U$, such that $\mu_{U, \varphi}$ transforms like the modulus of the Jacobian under change of coordinates. We also denote by $\mu$, the corresponding measure induced on $M$. Let $M^{\prime} \subset M$ be an open subset of $M$ and $D$ a smooth differential operator on $M^{\prime}$. Then there exists a smooth differential operator ${ }^{t} D$, called the transpose of $D$, such that

$$
\int D t \cdot g d \mu=\int f \cdot{ }^{t} D g d \mu
$$

for all $f, g \in C_{\mathrm{c}}^{\infty}\left(M^{\prime}\right)$. The operator $D^{*} g=\left({ }^{t} D \bar{g}\right)^{-}$is called the formal adjoint of $D$, and if $(f, g)=\int f \bar{g} d \mu$, then $(D f, g)=\left(f, D^{*} g\right)$ for all $f, g \in C_{c}^{\infty}\left(M^{\prime}\right)$.

Let $X \rightarrow \varrho(X)$ be a homomorphism of $\mathfrak{g}$ into differential operators on $M^{\prime}$, such that
(i) $\varrho(X)^{*}=-\varrho(X)$, for all $X \in g$,
(ii) For each $m_{0} \in M^{\prime}$, there exists an $a_{0} \in \mathcal{U}\left(g_{c}\right)$, such that the operator $\varrho\left(a_{0}\right)$ is elliptic in a neighborhood of $m_{0}$.

Let $\sigma=\left(C_{c}^{\infty}\left(M^{\prime}\right), \varrho\right)$. Then $\sigma$ is a symmetric $\mathfrak{g}$-module in $L^{2}(\mu)$.

Weyl's Lemma. With the above notation $\mathcal{D}\left(\sigma^{\dagger}\right)=\left\{f \in C^{\infty}\left(M^{\prime}\right) \mid \varrho(a) f \in L^{2}(\mu)\right.$, for all $\left.a \in \mathcal{U}\left(\mathrm{~g}_{\mathrm{c}}\right)\right\}$. Moreover, $\sigma^{\dagger}(a) f=\varrho(a) f$, for all $f \in \mathcal{D}\left(\sigma^{\dagger}\right)$.

Proof. This lemma is quite classical and we include a sketch of proof, for lack of adequate reference. It follows easily from the regularity theorem for elliptic operators. In fact, suppose $f \in \mathcal{D}\left(\sigma^{\dagger}\right)$. Let $u_{f}$ be the distribution $\varphi \rightarrow(\varphi, f), \varphi \in C_{c}^{\infty}\left(M^{\prime}\right) ;{ }^{t} \varrho\left(a_{0}\right) u_{f}(\varphi)=$ $u_{f}\left(\varrho\left(a_{0}\right) \varphi\right)=\left(\varphi, f_{a_{0}}\right)$. Thus the distribution ${ }^{\dagger} \varrho\left(a_{0}^{n}\right) u_{f}$ is an $L^{2}$-function for all $n$; since $\varrho\left(a_{0}\right)$ is elliptic in a neighborhood of $m_{0} \in M^{\prime}$, it follows from regularity theorem that $u_{f}$ is locally a $C^{\infty}$-function near $m_{0}$. Thus $f \in C^{\infty}\left(M^{\prime}\right)$ and $(\varrho(a) \varphi, f)=\left(\varphi, \varrho\left(a^{\dagger}\right) f\right)$ for all $a \in \mathcal{U}$. This proves the lemma.

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