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# On the irreducibility and inequivalence of unitary representations of gauge groups 

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## Introduction

Let $M$ be a connected, smooth, orientable, paracompact, $n$ dimensional manifold. Let $\langle, \quad\rangle$ be a Riemannian structure and $\omega$ a volume form on $M$. Let $U$ be a compact, semi-simple Lie group and let $G$ be the group of all smooth mappings of $M$ into $U$ that equal the identity outside of a compact set with pointwise multiplication. Let $B$ denote an $A \mathrm{~d}(U)$-invariant inner product on $\mathfrak{u}$, the Lie algebra of $U$. Then certain "non-local" unitary representations, $T=T_{\langle,\rangle, B, \omega}$, of $G$ were introduced in $[I]$ for $U=S U(2)$ and in [GGV, I, II], [AKT] for general $U$ (see §5).
Our main results on these representations are
(1) If $n \geqslant 3$ then $T$ is irreducible.
(2) Let $n=2$. Write $\omega=\varrho \mathrm{Vol}_{\langle,\rangle}\left(\mathrm{Vol}_{\langle,\rangle}\right.$a volume element of $M$ with respect to $\langle\rangle$,$) . Fix \mathfrak{h}$ a maximal abelian subalgebra of $\mathfrak{u}$ and let $\|\cdots\|_{B}$ denote the norm on $\mathfrak{b} *$ corresponding to $B$. If $\|\alpha\|_{B}>(8 \pi|\varrho(x)|)^{1 / 2}$ for $x \in M$ and all roots $\alpha$ of $\mathfrak{u}$ relative to $\mathfrak{h}$ then $T$ is irreducible.

Whenever one has $\langle\rangle, B,, \omega$ as above then one has an inner product $(,)_{\langle,\rangle, B, \omega}$ on, $\Omega_{c}^{1}(M ; \mathfrak{u})$, the compactly supported smooth one forms on $M$ with values in $\mathfrak{u}$ (see $\S 6$ ).
(3) If $\operatorname{dim} M=2$ then assume that $\langle\rangle, B,, \omega$ satisfy the condition in (2). If $\operatorname{dim} M \geqslant 2$ then $T_{\langle,\rangle, B, \omega}$ and $T_{\langle,\rangle_{\rangle}, B_{1}, \omega_{1}}$ are either equivalent or disjoint. They are equivalent if and only if $(,)_{\langle,\rangle, B, \omega}=(,)_{\langle,\rangle_{1}, B_{1}, \omega_{1}}$. Condition (2) can be made to hold by varying $\varrho$ and or $B$. (1), (3) for $n \geqslant 5$ and $U=S U(2)$ and $\omega, \omega_{1}$ the Riemannian volume elements are due to [I] (as are the main "algebraic" ideas in the proof of the general case). (1), (2), (3) without any conditions are asserted for $n \geqslant 2$ and $M$ non-compact in [GGV, I]. But that paper is severely flawed. In [GGV, II] there is a proof of (1), (3) for $n \geqslant 4$ if $M$ is the interior of a compact manifold with boundary although it is not clear if this condition is necessary to their proof, since finite
volume and non-compact seem to be enough) for the Riemannian volume elements. In [AKT], (1) is proved for $n \geqslant 3$ as is (2) and a slightly weaker form of (3) for $\|\alpha\|_{B}>(32 \pi|\varrho(x)|)^{1 / 2}, x \in M$ for the same class of manifolds as in [GGV, II] (since they refer to this paper for the details of the proof of irreducibility). They also indicate the likelihood that the 32 can be replaced by 8 . We note that slight modifications of the original argument are necessary in the case of compact $M$. Thus the new results in this paper involve establishing the validity of (1), (2), (3) for general manifolds, replacing a 32 by an 8 and a stronger criterion for disjointness.

In this paper, the first four sections contain technical results on Gaussian measures. The representation theory is in Sections 5 and 6. We suggest on first reading that the reader start with Section 6 and refer back to the necessary preliminaries.

As indicated above the main line of the proof of irreducibility is contained in [I] and [GGV, I]. The differences involve precise results on uniform mutual singularity of measures on spaces of distributions (our results can be found in $\S 3$ and $\S 4)$. We give a complete proof of the "algebraic" aspects of the proof of irreducibility in $\S 5,6$ for several reasons. One is that [GGV, I, II] and [AKT] make use of undocumented "well known" results on direct integrals (which are essentially proved in §5). Secondly, there is a rather subtle argument regarding singularity of convolutions in [AKT], Lemma 3.2 , for the case when $n=2$ that we don't understand (this of course, is not meant to imply that it is wrong). We avoid this argument (which also appears in [GGV, II], however there seems to be no problem with it if $n \geqslant 3$ ). Thirdly, the details of our argument are necessary in order to prove (3). Fourthly, we fix a minor error in [GGV, II]. Finally, our proof should be accessible to novices to quantum field theory and probability theory.

## 1. Gaussian measures

Let $V$ be a locally convex, separable, topological vector space over $\mathbb{R}$. Let (, ) be a continuous, positive definite, symmetric, bilinear form (inner product for short) on $V$. Let $H$ be the Hilbert space completion of $V$ with respect to (, ). If $W$ is a topological vector space then we use the notation $W^{\prime}$ for the space of all continuous linear functionals on $W$ endowed with the weak topology.

If $W$ is a finite dimensional subspace of $V$ and if $\Omega$ is a Borel set in $W^{\prime}$ then we set $Z_{W, \Omega}=\left\{\lambda \in V^{\prime}: \lambda_{\mid W} \in \Omega\right\} . Z_{W, \Omega}$ is called a cylinder set. Let $A_{W}$ be the isomorphism of $W$ onto $W^{\prime}$ given by $A_{W}(v)(u)=(u, v)$. Let $\mathrm{d}_{W} x$ denote the

Lebesgue measure on $W$ corresponding to an orthonormal basis of $W$ relative to (, ). If $\operatorname{dim} W=n$ then we set (cf. [GV], IV, 3.1)

$$
\mu\left(Z_{W, \Omega}\right)=(2 \pi)^{-n / 2} \int_{A_{W}^{-1}(\Omega)} \mathrm{e}^{-(x, x) / 2} \mathrm{~d}_{W} x
$$

Let $\mathscr{B}=\mathscr{B}\left(V^{\prime}\right)$ denote the $\sigma$-algebra of sets generated by the cylinder sets. We assume that $\mu=\mu_{(,)}$has a countably additive extension to $\mathscr{B}$. In this case $\mu$ is a probability measure which is called the Gaussian measure associated with (, ).

The following simple lemma will be used often in this paper.
Lemma 1.1. Let $v \in V$ be such that $(v, v)=1$. Then if $r \geqslant 1$

$$
\mu\left(\left\{\lambda \in V^{\prime}:|\lambda(v)| \leqslant r\right\}\right) \geqslant 1-\mathrm{e}^{-r^{2} / 2}
$$

Proof. By definition, the measure of the indicated set is

$$
\begin{aligned}
& (2 \pi)^{-1 / 2} \int_{-r}^{r} \mathrm{e}^{-x^{2} / 2} \mathrm{~d} x=(1 / \Gamma(1 / 2)) \int_{0}^{r^{2} / 2} x^{-1 / 2} \mathrm{e}^{-x} \mathrm{~d} x \\
& \quad=1-(1 / \Gamma(1 / 2)) \int_{r^{2} / 2}^{\infty} x^{-1 / 2} \mathrm{e}^{-x} \mathrm{~d} x \geqslant 1-\left(2^{1 / 2} / \Gamma(1 / 2)\right) \int_{r^{2} / 2}^{\infty} \mathrm{e}^{-x} \mathrm{~d} x \\
& \quad=1-(1 / \pi)^{1 / 2} \mathrm{e}^{-r^{2} / 2} \geqslant 1-\mathrm{e}^{-r^{2} / 2}
\end{aligned}
$$

We now record two results. The first will be used later. The second is standard but it gives a simple instance of the technique that we will use to prove singularity of measures.

Lemma 1.2. Let $S, T \subset V^{\prime}$. Suppose that for each $\varepsilon>0$ there exists $X_{\varepsilon} \in \mathscr{B}$ such that $\mu\left(X_{\varepsilon}\right) \geqslant 1-\varepsilon, X_{\varepsilon}+\lambda=X_{\varepsilon}$ for $\lambda \in T$ and $\mu\left(X_{\varepsilon}+\lambda\right) \leqslant \varepsilon$ for $\lambda \in S$. Then there exists $Y \in \mathscr{B}$ such that $\mu(Y)=1, Y+\lambda=Y$ for $\lambda \in T$ and $\mu(Y+\lambda)=0$ for $\lambda \in S$.

Proof. Set $Z_{n}=X_{2-n}$. Put $Y=\bigcap_{k=1}^{\infty} \bigcup_{n=k}^{\infty} Z_{n}$.
If $h \in H$ then we define $\lambda_{h} \in V^{\prime}$ by $\lambda_{h}(v)=(v, h)$. If $\lambda \in V^{\prime}$ then we say $\lambda \in H$ if $\lambda=\lambda_{h}$ for some $h \in H$.

Lemma 1.3. If $\lambda \in V^{\prime} . \lambda \notin H$ then there exists $X \in \mathscr{B}$ such that $\mu(X)=1$ and $\mu(X+\lambda)=0$.

Proof. Since $\lambda \notin H$ there exists for each $n \geqslant 1, v_{n} \in V$ such that $\lambda\left(v_{n}\right) \geqslant n$ and $\left(v_{n}, v_{n}\right)=1$. Set $Z_{n}=\left\{\xi \in V^{\prime}:\left|\xi\left(v_{n}\right)\right|<n / 2\right\}$. Then $Z_{n}+\lambda \subset V^{\prime}-Z_{n}$. So Lemma 1.1 implies that

$$
\mu\left(Z_{n}\right) \geqslant 1-\mathrm{e}^{-n^{2} / 8}
$$

and

$$
\mu\left(Z_{n}+\lambda\right) \leqslant \mu\left(V^{\prime}-Z_{n}\right) \leqslant \mathrm{e}^{-n^{2} / 8} .
$$

Now apply the previous lemma with $S=\{\lambda\}, T$ the empty set.

## 2. Some observations about the first Sobelev space

We first prove a simple lemma which will be adequate to prove our results on Gaussian measures for $n \geqslant 3$. If $\eta$ is a 1 -form on $\mathbb{R}^{n}, \eta=\Sigma \eta_{i} \mathrm{~d} x_{i}$ then set

$$
\|\eta\|^{2}=\int\left(\sum\left(\eta_{i}\right)^{2}\right) \mathrm{d} x .
$$

We use the notation $B(x ; r)$ for the open $r$-ball with center $x$ in $\mathbb{R}^{n} .\|\cdots\|$ will denote the usual Hilbert space norm on $\mathbb{R}^{n}$.

Lemma 2.1. Assume that $n \geqslant 3$. Then there exists $C(n)>0$ depending only on $n$ such that for each $0<\varepsilon<1 / 2$ there exists $f_{\varepsilon} \in C_{c}^{\infty}(B(0 ; 1))$ with
(1) $\left\|d f_{\varepsilon}\right\|=1$,
(2) $f_{\varepsilon}(x) \geqslant \varepsilon^{-n / 2+1} C(n)$ for $\|x\| \leqslant \varepsilon$.

Proof. Let $h \in C^{\infty}(\mathbb{R})$ be such that $h(x)=1$ for $x \leqslant 1$ and $h(x)=0$ for $x \geqslant 2$. Set (as usual) $r(x)=\|x\|$. Put $\varphi_{\varepsilon}(x)=h(r(x) / \varepsilon)$. Then $\varphi_{\varepsilon} \in C_{c}^{\infty}(B(0 ; 1))$ if $0<\varepsilon<1 / 2$ and $\mathrm{d} \varphi_{\varepsilon}=h^{\prime}(r / \varepsilon) \mathrm{d} r / \varepsilon$. Let $\Omega_{n}$ denote the volume of the $n-1$ dimensional Euclidean sphere. Then

$$
\left\|\mathrm{d} \varphi_{\varepsilon}\right\|^{2}=\left(\Omega_{n} / \varepsilon^{2}\right) \int_{0}^{\infty} r^{n-1} h^{\prime}(r / \varepsilon)^{2} \mathrm{~d} r=\varepsilon^{n-2} \Omega_{n} \int_{0}^{\infty} r^{n-1} h^{\prime}(r)^{2} \mathrm{~d} r .
$$

Set

$$
C(n)^{-2}=\Omega_{n} \int_{0}^{\infty} r^{n-1} h^{\prime}(r)^{2} \mathrm{~d} r
$$

and put $f_{\varepsilon}=\varepsilon^{-n / 2+1} C(n) \varphi_{\varepsilon}$.

We now prove an analogous result for $n=2$. In this case one can show (as was pointed out to us by Roger Nussbaum) the estimates are best possible.

Lemma 2.2. Let $n=2$. Given $0<C<1$ there exists for each $0<\varepsilon<1$, $f_{\varepsilon, C} \in C_{c}^{\infty}(B(0 ; 1))$ such that
(1) $\left\|\mathrm{d} f_{\varepsilon, C}\right\|=1$,
(2) $f_{\varepsilon, C}(x) \geqslant\left(C /(2 \pi)^{1 / 2}\right)|\log \varepsilon|^{1 / 2}$ for $\|x\| \leqslant \varepsilon$.

Before we give the proof we recall some well known (or easily proved) sophomore calculus results. If $f$ is a continuous function on $\mathbb{R}^{2}$ then we say that $\partial f / \partial x=u$ and $\partial f / \partial y=v$ in $L^{2}$ if $u$ and $v$ are square integrable and whenever $g \in C_{c}^{\infty}\left(\mathbb{R}^{2}\right)$

$$
\begin{aligned}
& \int f(x, y) \frac{\partial g}{\partial x}(x, y) \mathrm{d} x \mathrm{~d} y=-\int u(x, y) g(x, y) \mathrm{d} x \mathrm{~d} y \\
& \int f(x, y) \frac{\partial g}{\partial y}(x, y) \mathrm{d} x \mathrm{~d} y=-\int v(x, y) g(x, y) \mathrm{d} x \mathrm{~d} y
\end{aligned}
$$

Under this condition we write $\mathrm{d} f=u \mathrm{~d} x+v \mathrm{~d} y$ and we say that $\mathrm{d} f$ exists in $L^{2}$.

The following assertion is an easy calculation using Stoke's theorem.
(1) Let $h$ be a piecewise smooth function on $\mathbb{R}$ with $\operatorname{supp} h \subset(-\infty, a]$, $a<\infty$, such that $h$ is constant in a neighborhood of 0 . Set $f(x)=h(r(x))$. Then $\mathrm{d} f$ exists in $L^{2}$ and $\mathrm{d} f=h^{\prime}(r) \mathrm{d} r$.

$$
\text { If } f \in L^{1}\left(\mathbb{R}^{2}\right) \text { and if } g \in L^{2}\left(\mathbb{R}^{2}\right) \text { then we set (as usual) }
$$

$$
f * g(z)=\int f(u) g(z-u) \mathrm{d} u=\int f(z-u) g(u) \mathrm{d} u
$$

The following result is also standard.
(2) Let $f$ be continuous on $\mathbb{R}^{2}$ with $\operatorname{supp} f \subset B(0 ; 1-\eta)$ for some $0<\eta<1$ and suppose that $\mathrm{d} f$ exists in $L^{2}$. Then given $\varepsilon>0$ there exists $0<\delta<\eta$ and $\varphi \in C_{c}^{\infty}(B(0 ; \delta))$ such that

$$
\begin{aligned}
& \|\mathrm{d}(\varphi * f)\| \leqslant\|\mathrm{d} f\| \\
& \sup _{x \in B(0 ; 1)}|\varphi * f(x)-f(x)|<\varepsilon
\end{aligned}
$$

Indeed, the first inequality is true for any non-negative $\varphi$ with $L^{1}$-norm one without any assumption on $\delta$, and the second is an easy consequence of uniform continuity.

We now give the proof of Lemma 2.2. For each $0<\delta<1-\varepsilon$ define $h_{\varepsilon, \delta}$ by

$$
h_{\varepsilon, \delta}(x)=\left\{\begin{array}{cc}
1 & x \leqslant \varepsilon \\
\frac{\log x-\log (1-\delta)}{\log \varepsilon-\log (1-\delta)} & \varepsilon \leqslant x \leqslant 1-\delta \\
0 & 1-\delta \leqslant x
\end{array}\right.
$$

Set $\varphi_{\varepsilon, \delta}(z)=h_{\varepsilon, \delta}(r(z))$ for $z \in \mathbb{R}^{2}$. If we apply (1) above we have

$$
\begin{aligned}
\left\|\mathrm{d} \varphi_{\varepsilon, \delta}\right\|^{2} & =\frac{2 \pi}{(\log \varepsilon-\log (1-\delta))^{2}} \int_{\varepsilon}^{1-\delta} \mathrm{d} r / r \\
& =\frac{2 \pi}{|\log \varepsilon|\left|1-\frac{\log (1-\delta)}{\log \varepsilon}\right|}
\end{aligned}
$$

Set $C(\varepsilon, \delta)^{2}=2 \pi /|1-(\log (1-\delta)) / \log \varepsilon|$. Let $\mu>0$ be given. Then (2) above implies that there exists $u \in C_{c}^{\infty}(B(0 ; \eta)), 0<\eta<\delta$, with $u \geqslant 0$ and $\int u=1$ such that

$$
\left\|\mathrm{d}\left(u * \varphi_{\varepsilon, \delta}\right)\right\| \leqslant C(\varepsilon, \delta) /|\log \varepsilon|^{1 / 2}
$$

and

$$
u * \varphi_{\varepsilon, \delta}(x) \geqslant 1-\mu \text { for }\|x\| \leqslant \varepsilon .
$$

Put $g_{\varepsilon, \delta, \mu}=u * \varphi_{\varepsilon, \delta} /\left\|\mathrm{d}\left(u * \varphi_{\varepsilon, \delta}\right)\right\|$. Then $\left\|\mathrm{d} g_{\varepsilon, \delta, \mu}\right\|=1$ and

$$
g_{\varepsilon, \delta, \mu}(x) \geqslant(1-\mu) C(\varepsilon, \delta)^{-1}|\log \varepsilon|^{1 / 2} \quad \text { for } \quad\|x\| \leqslant \varepsilon .
$$

Now $\lim _{\delta, \mu \rightarrow 0}(1-\mu) C(\varepsilon, \delta)^{-1}=(2 \pi)^{-1 / 2}$. So we can take $f_{\varepsilon, C}=g_{\varepsilon, \delta, \mu}$ for $\delta$ and $\mu$ sufficiently small.

The following simple covering Lemma is sufficient for our purposes.

Lemma 2.3. Let $n \geqslant 2$. If $0<\varepsilon<1 / 4$ then there exist $z_{1}, z_{2}, \ldots, z_{N} \in$ $B(0 ; 5 / 4)$ with $N \leqslant(2 \sqrt{n}+1)^{n} / \varepsilon^{n}$ such that $B(0 ; 1) \subset U_{j=1}^{N} B\left(z_{j} ; \varepsilon\right)$.

Note. $B\left(z_{j} ; 1\right) \subset B(0 ; 5 / 2)$.

Proof. This is standard. For each $x \in B(0 ; 1)$ let $m_{i} \in Z$ be the unique element that satisfies $m_{i} \leqslant x_{i} \sqrt{n} / \varepsilon<m_{i}+1$. Put $m(x)=\left(m_{1}, \ldots, m_{n}\right)$. Then $\|x-\varepsilon m(x) / \sqrt{n}\|<\varepsilon$. Also $\left|m_{i}\right|<\sqrt{n} / \varepsilon+1$. Take the $z_{j}$ to be an enumeration of the set $\{\varepsilon m(x) / \sqrt{n}: x \in B(0 ; 1)\}$. Clearly, there are at most $2(\sqrt{n} / \varepsilon+3)^{n}$ such points.

## 3. Singularity of translates of Gaussian measures

Let $M$ be a smooth, orientable, paracompact, connected manifold of dimension $n$. Let $W$ be a finite dimensional vector space over $\mathbb{R}$. If $M$ is noncompact set $V=C_{c}^{\infty}(M ; W)$, the smooth compactly supported functions on $M$ with values in $W$. If $M$ is compact then we fix once and for all a base point $x_{0} \in M$ and set $V=\left\{f \in C^{\infty}(M ; W): f\left(x_{0}\right)=0\right\}$. If $K$ is a compact set in $M$ set $C_{K}^{\infty}(M ; W)$ equal to the smooth functions from $M$ to $W$ that equal 0 outside of $K$. We endow $C_{K}^{\infty}(M ; W)$ with the topology of uniform convergence with all derivatives. If $M$ is non-compact then we endow $V$ with the "union topology" (cf. [GV], p. 330). If $M$ is compact then we use the topology of uniform convergence with all derivatives. Then $V$ is either a nuclear space or a "union of nuclear spaces" (cf. [GV], p. 330). Let $\Omega^{1}(M ; W)$ denote the space of all 1-forms on $M$ with values in $W$. Let $\langle,>$ be a Riemannian structure on $M, B$ an inner product on $W$ and $\omega$ a volume form on $M$. If $\alpha, \beta \in \Omega^{1}(M ; W)$ then $\alpha_{x}, \beta_{x} \in \operatorname{Hom}_{\mathbf{R}}\left(T M_{x}, W\right)$. We write $(\alpha, \beta)_{x}$ for the Hilbert-Schmidt inner product of $\alpha_{x}$ with $\beta_{x}$. That is, $(\alpha, \beta)_{x}=\operatorname{Tr}\left(\beta_{x}^{*} \alpha_{x}\right)$. If $f, g \in V$ then we set

$$
(f, g)=\int_{\mathrm{M}}(\mathrm{~d} f, \mathrm{~d} g)_{x} \omega
$$

Let $\mu$ denote the corresponding Gaussian measure on $V^{\prime}(\S 1)$. Then $\mu$ is countably additive ([GV], Theorem 6, p. 333).

If $v \in W^{\prime}$ and if $x \in M\left(x \neq x_{0}\right.$ if $M$ is compact $)$ then we set $v_{x}(f)=v(f(x))$ for $f \in V$. Then $v_{x} \in V^{\prime}$ and $v_{x}=0$ if and only if $v=0$. We set $\|v\|_{B}=\sup _{B(w, w)=1}|v(w)|$. We also write $\omega=\varrho \mathrm{Vol}_{\langle,\rangle}$as in the introduction. If it is necessary to indicate the dependence of (, ) and $\mu$ on $\langle$,$\rangle ,$ $B, \omega$ then we write (,$)_{\langle,\rangle, B, \omega}$ and $\mu_{\langle,\rangle, B, \omega}$.

Proposition 3.1. Let $v_{1}, \ldots, v_{d} \in W^{\prime}-(0)$. Let $\langle,\rangle_{i}, B_{i}, \omega_{i}$ be as above for $i=1$, 2. Set $\mu_{i}=\mu_{\langle,\rangle_{i}, B_{i}, \omega_{i}}$. If $n \geqslant 3$ there is no additional condition. If $n=2$ we assume

$$
\left\|v_{i}\right\|_{B_{1}}>\sqrt{8 \pi\left|\varrho_{1}(x)\right|} \quad \text { for } \quad x \in M
$$

Then if $n \geqslant 2$ there exists $X \in \mathscr{B}$ such that $\mu_{i}(X)=1$ and $\mu_{2}\left(X+\Sigma\left(v_{i}\right)_{x_{1}}\right)=$ 0 for $x_{1}, \ldots, x_{d} \in M\left(M-x_{0}\right.$ if $M$ is compact $)$ with $x_{i} \neq x_{j}$ if $i \neq j$.

We will derive this result from a lemma which will also be used in $\S 5$.
Lemma 3.2. Let $\mu_{i}, i=1.2$ be as in Proposition 3.1. Assume that $v \in W^{*}-\{0\}$ and if $\operatorname{dim} M=2$ that $\|v\|_{B_{1_{1}}}>\left(8 \pi \varrho \varrho_{1}(x) \mid\right)^{1 / 2}$ for $x \in M$. If $U$ is an open subset of $M$ (not containing $x_{0}$ if $M$ is compact) then there exists $Y \in \mathscr{B}$ such that $\mu_{I}(Y)=1, Y+\varphi_{y}=Y$ for $y \in M-U, \varphi \in W^{*}$, and $\mu_{2}\left(Y+v_{x}\right)=0$ for $x \in U$.

We first show that Lemma 3.2 implies Proposition 3.1. So assume it. Set $M^{d}=\left\{\left(x_{1}, \ldots, x_{d}\right) \in x^{d} M: x_{i} \neq x_{j}\right.$ if $i \neq j\left(x_{i} \neq x_{0}\right.$ if $M$ is compact $\left.)\right\}$.
(1) If $x \in M^{d}$ then there exists an open neighborhood, $U_{x}$, of $x$ in $M^{d}$ and $Y_{x} \in \mathscr{B}$ such that $\mu_{1}\left(Y_{x}\right)=1$ and $\mu_{2}\left(Y_{x}+\Sigma\left(v_{j}\right)_{y}\right)=0$ for $y \in U_{x}$.
Let us show how (1) implies the proposition. Then we will use Lemma 3.2 to prove (1). Clearly, $M^{d}$ is separable. There is therefore a countable subcovering $\left\{U_{x_{j}}\right\}$ of $\left\{U_{x}\right\}_{x \in M^{d}}$. Set $U_{j}=U_{x j}$ and $Y_{j}=Y_{x,}$. Then

$$
\mu_{1}\left(Y_{t}\right)=1 \text { and } \mu_{2}\left(Y_{i}+\sum_{J}\left(v_{j}\right)_{y_{j}}\right)=0 \text { for } y \in U_{i} .
$$

Take $X=\cap Y_{i}$.
We now derive (1). Let $x=\left(x_{1}, \ldots, x_{d}\right)$. Let $W_{i}, i=1, \ldots, d$ be an open neighborhood of $x_{i}$ such that $W_{i} \cap W_{J}=\varnothing$ if $i \neq j$ and $x_{0} \notin W_{i}$ if $M$ is compact. Let $U_{x}=W_{1} \times \ldots \times W_{d}$. Let $Y_{x}$ be the " $Y$ " of Lemma 3.2 for $U=W_{1}$ and $v=v_{1}$. Then $\mu_{1}\left(Y_{x}\right)=1$. If $y=\left(y_{1}, \ldots, y_{d}\right) \in U_{x}$ then $y_{j} \in$ $M-W_{1}$ for $j \geqslant 2$. Thus $Y_{x}+\Sigma_{j \geqslant 2}\left(v_{j}\right)_{y_{j}}=Y_{x}$. Thus $\mu_{2}\left(Y_{x}+\Sigma_{j \geqslant 1}\left(v_{j}\right)_{y}\right)=$ 0 for $y \in U_{x}$.

We are left with the proof of Lemma 3.2. As above the following "local" assertion implies the lemma.
(2) If $x \in U$ then there exist a neightborhood, $U_{x}$, of $x$ in $U$ and $Y_{x} \in \mathscr{B}$ such that $\mu_{1}\left(Y_{x}\right)=1, Y_{x}+\varphi_{y}=Y_{x}$ for $\varphi \in W^{*}, y \in M-U$ and $\mu_{2}\left(Y_{x}+v_{y}\right)=0$ if $y \in U_{x}$.
We are left with the proof of (2). Since the proof in the case of $n \geqslant 3$ fairly simple and contains most of the essential ideas for the more delicate case of $n=2$, we will now give the complete proof for the case $n \geqslant 3$. It would be worthwhile to read this even if there is only interest in the case $n=2$ which we will prove in the next section.

Let $x \in U$. Let $\left(U_{1}, \Psi\right)$ be a chart for $U$ such that $x \in U_{1}$ and
(i) $\Psi\left(U_{1}\right)=B(0 ; 3)$,
(ii) $\Psi(x)=0$.

Set $V_{1}=(\Psi)^{-1}(B(0 ; 5 / 2))$. On $\mathbb{R}^{n}$ we use the usual Riemannian structure and Lebesgue measure for the volume element. We will write $\|\cdot \cdot\|_{i}$ for the pointwise norm on Hom (TM*, $W$ ) corresponding to $\langle,\rangle_{i}$ and $B_{i}$. We fix an arbitrary inner product, $B$, on $W^{*}$ and write $\|\|$ for the norm on $\left(\mathbb{R}^{n}\right)^{*} \otimes W^{*}$ corresponding to the usual inner product tensored with $B$. (A) There exist constants $D_{1}, D_{2}>0$ such that if $f \in C_{c}^{\infty}\left(V_{1} ; W\right)$ then

$$
\begin{aligned}
& D_{1} \int_{B(0 ; 3)}\left\|\mathrm{d}\left(f \cdot \Psi^{-1}\right)(x)\right\|^{2} \mathrm{~d} x \leqslant \int_{M}\|\mathrm{~d} f\|_{i, x}^{2} \omega_{i} \\
& \leqslant D_{2} \int_{B(0 ; 3)}\left\|\mathrm{d}(f \cdot \Psi)^{-1}(x)\right\|^{2} \mathrm{~d} x
\end{aligned}
$$

Let $\left\{y_{1}, \ldots, y_{n}\right\}$ be the local coordinates on $U_{1}$ corresponding to $\Psi$. Set

$$
g_{p, q, i}(z)=\left\langle\mathrm{d} y_{p}, \mathrm{~d} y_{q}\right\rangle_{i, x}, \quad z \in U_{1} .
$$

Then there exist $\alpha_{i}, \beta_{i}>0$ such that if $z \in$ Closure $\left(V_{1}\right)$ then

$$
\alpha_{i}\langle y, y\rangle \leqslant \sum g_{p, q, i}(z) y_{p} y_{q} \leqslant \beta_{i}\langle y, y\rangle, \quad y \in \mathbb{R}^{n} .
$$

Also, $\omega_{i \mid U_{1}}=u_{i} \mathrm{~d} y_{1} \Lambda \mathrm{~d} y_{2} \Lambda \ldots \Lambda \mathrm{~d} y_{n}$. There exist $\gamma_{i}, \delta_{i}>0$ such that

$$
\gamma_{i} \leqslant\left|u_{i}(z)\right| \leqslant \delta_{i}, \quad z \in \text { Closure }\left(V_{1}\right)
$$

There exist $M_{i}, m_{i}>0$ such that

$$
m_{i} B(v, v) \leqslant B_{i}(v, v) \leqslant M_{i} B(v, v), \quad v \in W^{*}
$$

Take $D_{1}=\min \left\{\gamma_{i} \alpha_{i} m_{i}\right\}, D_{2}=\max \left\{\delta_{i} \beta_{i} M_{i}\right\}$. (A) now follows.
Note. We will also use this result in our proof in the case $n=2$.
Set $U_{x}=\Psi^{-1}(B(0 ; 1))$.
(B) There exist $E, F>0$ such that if $1 / 4>\varepsilon>0$ there exists an open covering $W_{i}, i=1, \ldots, N \leqslant E \varepsilon^{-n}$, of $U_{x}$ with $W_{i} \subset V_{1}$ and there exist $f_{i, \varepsilon} \in C_{c}^{\infty}\left(V_{1}\right)$ such that
(a) $\left\|\mathrm{d} f_{i, \varepsilon}\right\|=1$ and
(b) $f_{i, \varepsilon}(z) \geqslant F \varepsilon^{-n / 2+1} \quad$ for $\quad z \in W_{i}$.

Indeed, let $f_{\varepsilon}$ be as in Lemma 2.1 and $z_{1}, \ldots, z_{N}$ be as in Lemma 2.3. Set $W_{i}=\Psi^{-1}\left(B\left(z_{i} ; \varepsilon\right)\right)$. Put $g_{i, \varepsilon}(y)=f_{\varepsilon}\left(\Psi(y)-z_{i}\right)$ if $y \in V_{1}$ and 0 otherwise. Then $\left\|\mathrm{d} g_{i, \varepsilon}\right\|_{1} \leqslant\left(D_{2}\right)^{1 / 2}$ by (A) and $g_{i, \varepsilon}(y) \geqslant C(n) \varepsilon^{-n / 2+1}$. Set $f_{i, \varepsilon}=$ $g_{i, \varepsilon} /\left\|\mathrm{d} g_{i, \varepsilon}\right\|_{1}, F=C(n) /\left(D_{2}\right)^{1 / 2}$ and $E=\left(2 n^{1 / 2}+1\right)$. This proves (B).

Let $w \in W$ be such that $B_{1}(w, w)=1$ and $v(w)=\|v\|_{B_{1}}$. Set $\alpha_{i, \varepsilon}=f_{i, \varepsilon} \otimes w$. Then $\left\|\mathrm{d} \alpha_{i, \varepsilon}\right\|=1$. If $u \in W_{j}$ then
$(v)_{u}\left(\alpha_{j, \varepsilon}\right)=f_{j, \varepsilon}(u)\|v\|_{B_{1}} \geqslant C \varepsilon^{-n / 2+1}$
with $C=F\|v\|_{B_{1}}$.
Set $Z_{j, \varepsilon}=\left\{\lambda \in V^{\prime}:\left|\lambda\left(\alpha_{j, \varepsilon}\right)\right|<C \varepsilon^{-n / 2+1} / 2\right\}$. Then Lemma 1.1 implies that

$$
\mu_{1}\left(Z_{j, \varepsilon}\right) \geqslant 1-\mathrm{e}^{-C^{2} \varepsilon^{-n+2 / 8}}
$$

If $y \notin U$ then $\alpha_{j, \varepsilon}(y)=0$ so $Z_{j, \varepsilon}+\varphi_{y}=Z_{j, \varepsilon}$ for $\varphi \in W^{*}$. Also, if $u \in U_{x}$ then

$$
Z_{j, \varepsilon}+v_{u} \subset\left\{\lambda \in V^{\prime}: \lambda\left(\alpha_{j, \varepsilon} /\left\|\alpha_{j, \varepsilon}\right\|_{2}\right) \geqslant C \varepsilon^{-n / 2+1} / 2\left\|\alpha_{j, \varepsilon}\right\|_{2}\right\} .
$$

Hence Lemma 1.1 implies that

$$
\mu_{2}\left(Z_{J, \varepsilon}+v_{u}\right) \leqslant \exp \left(-C^{2} \varepsilon^{-n+2} / 8\left\|\mathrm{~d} \alpha_{ر, \varepsilon}\right\|_{2}^{2}\right)
$$

(A) implies that $\left\|\alpha_{j, \varepsilon}\right\|_{2} \leqslant\left(D_{2} / D_{1}\right)^{1 / 2}$. Set $\xi=C^{2} D_{1} / 8 D_{2}$. Take $Z_{\varepsilon}=\bigcap_{\jmath} Z_{j, \varepsilon}$. Then

$$
\begin{aligned}
& \mu_{1}\left(Z_{\varepsilon}\right) \geqslant 1-E \varepsilon^{-n} \mathrm{e}^{-C^{2} \varepsilon^{-n+2 / 8}} \\
& \mu_{2}\left(Z_{\varepsilon}+v_{\mathrm{u}}\right) \leqslant \mathrm{e}^{-\xi \varepsilon^{-n+2}}
\end{aligned}
$$

if $u \in U_{x}$ and $Z_{\varepsilon}+\varphi_{y}=Z_{\varepsilon}$ if $\varphi \in W^{*}$ and $y \in M-U$. Since $\lim _{\varepsilon \rightarrow 0} \varepsilon^{-n} \mathrm{e}^{-C^{2} \varepsilon^{-n+2 / 8}}=0$ if $n \geqslant 3$, this implies that given $\varepsilon>0$ there exists $Y_{\varepsilon} \in \mathscr{B}$ such that $\mu_{1}\left(Y_{\varepsilon}\right)=1, Y_{\varepsilon}+\varphi_{y}=Y_{\varepsilon}$ for $\varphi \in W^{*}, y \in M-U$ and $\mu_{2}\left(Y_{\varepsilon}+v_{u}\right)=0$ for $u \in U_{x}$. This combined with the argument in the proof of Lemma 1.2 completes the proof of (2) in the case $n \geqslant 3$.

## 4. The proof of Lemma 3.2 for $n=2$

In this section we assume that $n=2$, otherwise the assumptions and notation are as in §3. To simplify notation, we denote $\langle,\rangle_{1}, B_{1}, \omega_{1}$ by the same symbols without the subscripts. We will be using some elementary Riemannian geometry. For this we refer to [H], Chapter 1. Let $\mathrm{d}(x, y)$ be the Riemannian distance on $M$. If $m \in M$ then we denote by $\|v\|_{m}$ the norm of
$v \in T M_{m}$ relative to $\langle,\rangle_{m}$. Put $B(m ; x ; r)=\left\{v \in T M_{m}:\|v-x\|_{m}<r\right\}$ for $x \in T M_{m}$. If $x \in M$ we set $B_{x}(r)=\{y \in M: \mathrm{d}(x, y)<m\}$. For $m \in M$, let $\exp _{m}$ be the (geodesic) exponential map of $M$ at $m$. For each $m \in M$ let $s_{m}>0$ be such that $\exp _{m}: B\left(m ; 0 ; s_{m}\right) \rightarrow B_{m}\left(s_{m}\right)$ is a surjective diffeomorphism.

If we identify $T\left(T M_{m}\right)_{v}$ for $v \in T M_{m}$ with $T M_{m}$ in the canonical way then $\mathrm{d}\left(\exp _{m}\right)_{0}=I$ for $m \in M$. This implies that for each $\delta>0$ there exists $0<\eta_{m}(\delta)<s_{m}$ such that if $f \in C_{c}^{\infty}\left(B_{m}\left(\eta_{m}(\delta)\right)\right.$ then

$$
\begin{equation*}
\int_{M}\|\mathrm{~d} f\|^{2} \omega \leqslant(1+\delta)|\varrho(m)| \int_{B\left(m ; ; ; \eta_{m}(\delta)\right)}\left\|\mathrm{d}\left(f \circ \exp _{m}\right)\right\|_{m}^{2} \mathrm{~d}_{m} x . \tag{*}
\end{equation*}
$$

Here $\mathrm{d}_{m} x$ is the Lebesgue measure on $T M_{m}$ corresponding to an orthonormal basis relative to $\langle,\rangle_{m}$.

Fix $x \in U$. Choose $s>0$ such that $B_{x}(3 s) \subset U$. Put $r(\delta)=$ $\min \left\{s / 3, \eta_{x}(\delta) / 3\right\}$. Set $U_{1}=B_{x}(3 r(\delta)), V_{1}=B_{x}(5 r(\delta) / 2), W_{1}=B_{x}(r(\delta))$. We now prove the assertion of (2), §3, for $U_{x}=W_{1}$ if $\delta$ is chosen to be sufficiently small. Let $0<\zeta<1$, and let $f_{\varepsilon, 1-\zeta}$ be as in lemma 2.2. Set $u_{\varepsilon, \zeta}(x)=f_{\varepsilon, 1-\zeta}(x / r(\delta))$. Then $\left\|u_{\varepsilon, \zeta}\right\|=1$ (here we are using the norms as in §2) and

$$
u_{\varepsilon, \zeta}(x) \geqslant(1-\zeta)|(\log \varepsilon) / 2 \pi|^{1 / 2}, \quad\|x\| \leqslant \varepsilon r(\delta)
$$

Let $z_{1}, \ldots, z_{N}\left(N \leqslant(1+2 \sqrt{2})^{2} \varepsilon^{-2}=E \varepsilon^{-2}\right)$ be as in Lemma 2.3 for $\varepsilon>0$. Put $Z_{i}(\delta)=\exp _{x}\left(B\left(x ; r(\delta) z_{i} ; r(\delta) \varepsilon\right)\right)$. Define $\xi_{i, \varepsilon, \delta, \zeta} \in C_{c}^{\infty}\left(V_{1}\right)$ by $\xi_{i, \varepsilon, \delta, \zeta}(y)=0$ if $y \notin V_{1}$ and $\xi_{i, \varepsilon, \delta, \zeta}\left(\exp _{x}(y)\right)=u_{\varepsilon, \zeta}\left(y-r(\delta) z_{i}\right)$ for $y \in$ $B\left(x ; r(\delta) z_{i} ; r(\delta) \varepsilon\right)$. Put $f_{i, \varepsilon, \delta, \zeta}=\xi_{i, \varepsilon, \delta, \zeta} /\left\|\mathrm{d} \xi_{j, \varepsilon, \delta, \xi}\right\|$. Then (*) implies that

$$
\begin{equation*}
f_{i, \varepsilon, \delta, \zeta}(z) \geqslant(1-\zeta)|(\log \varepsilon) / 2 \pi(1+\delta) \varrho(x)|^{1 / 2} \text { for } z \in Z_{i}(\delta) \tag{**}
\end{equation*}
$$

Let $w \in W$ be such that $v(w)=\|v\|_{B}$ and $B(w, w)=1$. Put

$$
\alpha_{i, \varepsilon, \delta, \zeta}=f_{i, \varepsilon, \delta, \zeta} \otimes w
$$

Then $\left\|\mathrm{d} \alpha_{i, \varepsilon, \delta, \zeta}\right\|=1$ and if $u \in Z_{i}(\delta)$ then

$$
\begin{aligned}
& v_{u}\left(\alpha_{i, \varepsilon, \delta, \xi}\right)=\mid v \|_{B} f_{i, \varepsilon, \delta, \zeta}(u) \\
& \quad \geqslant\|v\|_{B}(1-\zeta)|(\log \varepsilon) / 2 \pi(1+\delta) \varrho(x)|^{1 / 2}
\end{aligned}
$$

Our assumption on $v$ implies that $\|v\|_{B}>(8 \pi|\varrho(x)|)^{1 / 2}$. Thus we can choose $\zeta$ and $\delta$ so small that

$$
\|v\|_{B}(1-\zeta)|(\log \varepsilon) 2 \pi(1+\delta) \varrho(x)|^{1 / 2} \geqslant(1+\gamma)(2+\gamma)|\log \varepsilon|^{1 / 2}
$$

for some $\gamma>0$. Fix these values of $\delta$ and $\zeta$. Set $\alpha_{j, \varepsilon}=\alpha_{j, \varepsilon, \delta, \zeta}$ and $U_{x}=W_{1}$. Set $Z_{j, \varepsilon}=\left\{\lambda \in V^{\prime}:\left|\lambda\left(\alpha_{j, \varepsilon}\right)\right|<(2+\gamma)|\log \varepsilon|^{1 / 2}\right\}$. Then lemma 1.1 implies that

$$
\mu\left(Z_{j, \varepsilon}\right) \geqslant 1-\exp \left(\left((2+\gamma)^{2} \log \varepsilon\right) / 2\right)
$$

$$
\geqslant 1-\varepsilon^{2+\gamma}
$$

If $u \in M-U$ and $\varphi \in W$ then $\varphi_{u}\left(\alpha_{j, \varepsilon}\right)=0$ so $Z_{j, \varepsilon}+\varphi_{u}=Z_{j, \varepsilon}$. Also if $u \in$ $Z_{j}$ then $Z_{j, \varepsilon}+(v)_{u} \subset\left\{\lambda \in V^{\prime}: \lambda\left(\alpha_{j, \varepsilon} /\left\|\mathrm{d} \alpha_{j, \varepsilon}\right\|_{2}\right) \geqslant \gamma(2+\gamma)|\log \varepsilon|^{1 / 2} /\left\|\mathrm{d} \alpha_{j, \varepsilon}\right\|_{2}\right\}$. Thus, if we set $\xi_{j}=\left(\gamma(2+\gamma) /\left\|\mathrm{d} \alpha_{j, \varepsilon}\right\|_{2}\right)^{2} / 2$ then Lemma 1.1 implies that

$$
\mu_{2}\left(Z_{j, \varepsilon}+v_{u}\right) \leqslant \mathrm{e}^{-\xi,|\log \varepsilon|}=\varepsilon^{\xi} .
$$

Now $\S 3$ (A) implies that there exists a constant $D>0$ independent of $j, \varepsilon$ such that $\left\|\mathrm{d} \alpha_{j, \varepsilon}\right\|<D$. Thus if we set $\xi=\gamma / 2 D^{2}$ then $\xi_{j}>\xi$. Hence if $u \in Z_{j}$ then

$$
\mu_{2}\left(Z_{j, \varepsilon}+v_{u}\right) \leqslant \varepsilon^{\xi}
$$

Put $Z_{\varepsilon}=\bigcap_{,} Z_{j, \varepsilon}$. Then

$$
\mu\left(Z_{\varepsilon}\right) \geqslant 1-E \varepsilon^{-2} \varepsilon^{2+\gamma}
$$

$Z_{\varepsilon}+\varphi_{u}=Z_{\varepsilon}$ for $\varphi \in W^{*}, u \in M-U$ and

$$
\mu_{2}\left(Z_{\varepsilon}+v_{y}\right) \leqslant \varepsilon^{\xi} \quad \text { and } \quad u \in U_{x} .
$$

Thus given $\varepsilon>0$ there exists $Y_{\varepsilon} \in \mathscr{B}$ such that $\mu\left(Y_{\varepsilon}\right) \geqslant 1-\varepsilon$, $Y_{\varepsilon}+\varphi_{y}=Y_{\varepsilon}$ for $y \in M-U, \varphi \in W^{*}$ and $\mu_{2}\left(Y_{\varepsilon}+v_{u}\right) \leqslant \varepsilon$ for $u \in U_{x}$. The result now follows from the argument in the proof of Lemma 1.2.

## 5. Some representation theory

As in $\S 3$, let $M$ be a smooth, paracompact, connected, orientable manifold. Let ( $W$, (, )) be a finite dimensional real inner product space. Fix a Riemannian structure, $\langle$,$\rangle and a volume element, \omega$, on $M$. If $M$ is compact then fix a base point, $x_{0}$, set $V=C_{c}^{\infty}(M ; W)$ if $M$ is non-compact and $V=\left\{f \in C^{\infty}(M ; W): f\left(x_{0}\right)=0\right\}$ if $M$ is compact. We set $Q(f, g)=$ ( $\mathrm{d} f, \mathrm{~d} g$ ) as in $\S 3$. Let $\mu$ be the Gaussian measure on $V^{\prime}$ corresponding
to $Q$. We define a unitary representation, $S$, of $V$ on $L^{2}\left(V^{\prime}, \mu\right)$ by $S(v) f(\lambda)=\mathrm{e}^{\mathrm{i}(v)} f(\lambda)$.

Let $v_{1}, v_{2}, \ldots$ be a fixed sequence of non-zero (not necessarily distinct) elements of $W^{\prime}$. If $\operatorname{dim} M=2$ then we assume that $\left\|v_{i}\right\|>(8 \pi|\varrho(x)|)^{1 / 2}$ for $x \in M$ and all $i$. If $I=\left(i_{1}, \ldots, i_{d}\right)$ define for $x=\left(x_{1}, \ldots, x_{d}\right) \in M^{d}$ (see §3), $\Psi_{I}(x)(f)=\Sigma_{j} v_{i j}\left(f\left(x_{j}\right)\right)$, for $f \in V$. Then $\Psi_{I}$ defines a continuous mapping of $M^{d}$ into $V^{\prime}$ (with the weak topology). It is easily seen that $\Psi\left(M^{d}\right) \in \mathscr{B}$. If $E$ is a Hilbert vector bundle over $M^{d+k}$ for $d \geqslant 0, k \geqslant 0$ and if $\omega_{1}$ is a volume form on $M$ then we set $L^{2}\left(E, \omega_{1}\right)$ equal to the space of square integrable cross sections of $E$ (here we use the product measure $\omega_{1}^{d+k}$ on $\left.M^{d+k}\right)$. If $f \in L^{2}\left(E, \omega_{1}\right), v \in V$ and $I=\left(i_{1}, \ldots, i_{d}\right)$ then we set

$$
\sigma_{I, E}(v) f(x)=\mathrm{e}^{i \Psi_{I}\left(x_{1}, \ldots, x_{d}\right)(v)} f(x)
$$

Let $\langle,\rangle_{1}, B_{1}, \omega_{1}$ be respectively a Riemannian structure on $M$, an inner product on $W$ and a volume form on $M$. Let $Q_{1}$ be the inner product defined as above on $V$ using $\langle,\rangle_{1}, B_{1}$ and $\omega_{1}$ in place of $\langle\rangle, B,, \omega$. Let $\mu_{1}$ be the corresponding Gaussian measure on $V^{\prime}$.

Lemma 5.1. Let $d>0$. If $C$ is a bounded linear operator from $L^{2}\left(V^{\prime}, \mu\right)$ to $L^{2}\left(V^{\prime}, \mu_{1}\right) \hat{\otimes} L^{2}\left(E, \omega_{1}\right)$ such that $C S(v)=S(v) \otimes \sigma_{I, E}(v) C$ for all $v \in V$ then $C=0$.

Proof. We write $\Psi_{I}(x)=\Psi_{I}\left(x_{1}, \ldots, x_{d}\right)$. Set $\Omega(\lambda)=1$ for all $\lambda \in V^{\prime}$. We note that Closure $(\operatorname{span}(S(v) \Omega))=L^{2}\left(V^{\prime}, \mu\right)$. Indeed, $\operatorname{span}\left\{\mathrm{e}^{i \cdot(v)}: v \in V\right\}$ is dense in $L^{2}\left(V^{\prime}, \mu\right)$ (cf. [Gu, §7.2]). This implies that $C=0$ if and only if $C \Omega=0$. Let $f=C \Omega$. We assume that $f \neq 0$ and derive a contradiction. Then we can look upon $f$ as a function on $V^{\prime}$ with values in $L^{2}(E)$. Thus we can write $f(\lambda, x) \in E_{x}$. Let

$$
D=\operatorname{Closure}\left(\operatorname{span}\left\{\left(S(v) \otimes \sigma_{I, E}(v)\right) f: v \in V\right\}\right)
$$

Then $C$ is a continuous linear map of $L^{2}\left(V^{\prime}, \mu\right)$ onto $D$. On $V^{\prime} \times M^{d+k}$ we put the product measure, $\mu \times \omega_{1}^{d+k}$. Let $\Gamma(\lambda, x)=\lambda+\Psi_{I}(x) \in V^{\prime}$ for $\lambda \in V^{\prime}, x \in M^{d+k}$. Then $\Gamma$ is continous. We define a measure, $\gamma$, on $V^{\prime}$ as follows:

$$
\gamma(X)=\int_{\Gamma^{-1}(X)}\|f(\lambda, x)\|_{x}^{2} \mathrm{~d} \mu_{1}(\lambda) \omega_{1}^{d+k}
$$

Thus

$$
\begin{equation*}
\gamma(X)=\int_{X \times M^{d+k}}\left\|f\left(\lambda-\Psi_{I}(x), x\right)\right\|_{x}^{2} \mathrm{~d} \mu_{1}\left(\lambda-\Psi_{I}(x)\right) \omega_{1}^{d+k} . \tag{*}
\end{equation*}
$$

Then the representation of $V$ on $D$ given by the restriction of $S \otimes \sigma_{I, E}$ to $D$ is equivalent with the representation, $\beta$, of $V$ on $L^{2}\left(V^{\prime}, \gamma\right)$ with $\beta(v) \varphi(\lambda)=\mathrm{e}^{\mathrm{i} \lambda(v)} \varphi(\lambda)$. Thus $C$ induces a continuous linear map, $C_{1}$, of $L^{2}\left(V^{\prime}, \mu\right)$ into $L^{2}\left(V^{\prime}, \gamma\right)$ such that $C_{1} S(v)=\beta(v) C_{1}, v \in V$ and $C_{1} \Omega=\Omega_{1}$ ( $\Omega_{1}$ is the constant function 1 on $V^{\prime}$ looked upon as an element of $L^{2}\left(V^{\prime}, \gamma\right)$ ). Now

$$
\left(C_{1} S(v) \Omega, \Omega_{1}\right)=\left(\beta(v) \Omega_{1}, \Omega_{1}\right)=\int_{V^{\prime}} \mathrm{e}^{i \lambda(v)} \mathrm{d} \gamma(\lambda)
$$

On the other hand

$$
\left(C_{1} S(v) \Omega, \Omega_{1}\right)=\left(S(v) \Omega, C_{1}^{*} \Omega_{1}\right)
$$

Set $C_{1}^{*} \Omega_{1}=h$. Then

$$
\left(C_{1} S(v) \Omega, \Omega_{1}\right)=\int_{V^{\prime}} \mathrm{e}^{i \lambda(v)} \overline{h(\lambda)} \mathrm{d} \mu(\lambda)
$$

Since both $\gamma$ and $\overline{h(\lambda)} \mathrm{d} \mu(\lambda)$ are cylinder set measures (finite valued $\sigma$-additive measures on $\mathscr{B}$ ), we see that this implies that $\mathrm{d} \gamma=\overline{h(\lambda)} \mathrm{d} \mu(\lambda)$.

Proposition 3.1 implies that there exists $X \in \mathscr{B}$ such that $\mu(X)=1$ and $\mu_{1}\left(X+\Psi_{I}(x)\right)=0$ for all $x \in M^{d+k}$. Thus $\gamma(X)=\int_{x} \overline{h(\lambda)} \mathrm{d} \mu(\lambda)=$ $f_{V^{\prime}} \overline{h(\lambda)} \mathrm{d} \mu(\lambda)=\left(\Omega_{1}, \Omega_{1}\right)=\|f\|^{2}$. On the other hand $\gamma(X) \stackrel{J_{x}}{=} 0$ by $(*)$. This is a contradiction, so the lemma follows.

We now assume that $v_{1}, \ldots, v_{r}$ are distinct and satisfy the hypothesis above. Let $E_{1}, E_{2}, \ldots, E_{r}$ be Hermitian vector bundles over $M$. We define an action of $V$ on each $E_{i}$ by $\sigma_{j}(v)_{x\left(E_{j}\right)_{x}}=\mathrm{e}^{i \nu_{j}\left(v\left(x_{j}\right)\right)} I$. Let $E=E_{1} \oplus E_{2} \oplus \ldots \oplus E_{r}$ with action $\oplus \sigma_{i}=\sigma$. Let $\otimes^{d} E$ be the $d$-fold tensor product bundle over $M^{d}$ with the corresponding tensor product action of $V, \otimes^{d} \sigma$. Let $S_{d}$, the symmetric group on $d$ letters, act on $M^{d}$ by permuting the coordinates. We also let $s \in S_{d}$ act on $\otimes^{d} E$ by $e_{1} \otimes \ldots \otimes e_{d}$ over $\left(x_{1}, \ldots, x_{d}\right)$ goes to $s\left(e_{1} \otimes \ldots \otimes e_{d}\right)=e_{s 1} \otimes \ldots \otimes e_{s d}$ over $\left(x_{s 1}, \ldots, x_{s d}\right)$. Then $\otimes^{d} \sigma(v)$ commutes with the action of $S_{d}$.

By our definition of $M^{d}$ the action of $S_{d}$ is free. We therefore have a manifold $N^{d}=S_{d} \backslash M^{d}$. Let $\pi$ be the canonical projection of $M^{d}$ onto $N^{d}$. We note that $\omega^{d}$ "pushes down" to a measure on $N^{d}$. We write $L^{2}\left(\otimes^{d} E\right)$ for $L^{2}\left(\otimes^{d} E, \omega^{d}\right)$. Let $H^{d}$ be the space of all $f \in L^{2}\left(\otimes^{d} E\right)$ such that $s f\left(x_{1}, \ldots, x_{d}\right)=f\left(x_{s 1}, \ldots, x_{s d}\right)$. We define a representation, $\tau$, of $V$ on by $\tau(v) f(x)=\otimes^{d} \sigma(v)_{x} f(x)$.

The following result is undoubtedly a very special case of a well known result that is true for totally discontinuous actions of discrete groups. We include a proof since it is short.

Lemma 5.2. There exists an open subset, $F^{d}$, of $M^{d}$ such that $\pi$ is injective on $F^{d}$ and $N^{d}-\pi F^{d}$ has measure 0.

Proof. Whitney has shown [W] that we may assume that $M$ is a closed analytic submanifold of $\mathbb{R}^{N}$ for some large $N$. Choose a non-constant real analytic function, $f$, on $M$. Set $F^{d}=\left\{\left(x_{1}, \ldots, x_{d}\right): f\left(x_{1}\right)>f\left(x_{2}\right)>\ldots>f\left(x_{d}\right)\right\}$. Then clearly $s F^{d} \cap F^{d}$ is empty if $s \neq 1$. Set $f_{i j}\left(x_{1}, \ldots, x_{d}\right)=f\left(x_{i}\right)-$ $f\left(x_{j}\right)$ for $i \neq j$. Then $f_{i j}$ is real analytic and non-constant on $M^{d}$ for $i \neq j$. Now, the complement to $\bigcup_{s \in S_{d}} s F^{d}$ is $\bigcup_{i \neq j}\left\{x \in M^{d}: f_{i j}(x)=0\right\}$. Since the zero set of a non-constant real analytic function has measure 0 relative to any volume form, the result follows.

We will "abuse notation" and think of $\pi$ as projecting onto $F^{d}$. Also $H^{d}$ is, under these identifications, just $L^{2}\left(\left.\otimes^{d} E\right|_{F^{d}}\right)$ with the same action of $V$. Set $\left.\otimes^{d} E\right|_{F^{d}}=E^{d}$. Then $E^{d}$ splits into a direct sum $\oplus E_{I}^{d}$ over $I=\left(i_{1}, \ldots, i_{d}\right)$, $1 \leqslant i_{j} \leqslant r$ and

$$
\left.\otimes^{d} \sigma(v)\right|_{\left(E_{I}^{d}\right)_{x}}=\mathrm{e}^{i \Psi_{I}(x)(v)} I
$$

Let $\tau_{I, E}$ be the representation of $V$ on $L^{2}\left(E_{I}^{d}\right)$ given by $\tau_{I, E}(v) f(x)=$ $\tau(v) f(x)$.

Lemma 5.3. Let $U \subset N^{d}$ be open. Let $C$ be a continous linear operator on $L^{2}\left(V^{\prime}, \mu\right) \otimes L^{2}\left(E_{I}^{d}\right)$ such that $C\left(S(v) \otimes \tau_{I, E}(v)\right)=\left(S(v) \otimes \tau_{I, E}(v)\right) C$ for all $v \in V$. Then $C\left(L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid U}^{d}\right)\right) \subset L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid U}^{d}\right)$.

Proof. Let $Z=N^{d}-U$. Then

$$
L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I}^{d}\right)=L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid U}^{d}\right) \oplus L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid Z}^{d}\right)
$$

a direct sum of invariant subspaces under $S \otimes \tau_{I, E}$. Thus we must show that if $C$ is a continuous linear operator from $L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid U}^{d}\right)$ to $L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid Z}^{d}\right)$ such that $C\left(S(v) \otimes \tau_{I, E}(v)\right)=\left(S(v) \otimes \tau_{I, E}(v)\right) C$ for all $v \in V$ then $C=0$. We first reduce this result to a special case. Let $x \in U$ and let $y \in F^{d}$ be such that $\pi(y)=x$. Then there exist open neighborhoods $W_{1}, \ldots, W_{d}$ of $y_{1}, \ldots, y_{d}$ such that $W_{i} \cap W_{j}=\varnothing$ if $i \neq j$ and $W_{x}=$ $W_{1} \times W_{2} \times \ldots \times W_{d} \subset F^{d}$. Now $\bigcup_{x} \pi W_{x}=U$. A countable number of the $\pi W_{x}$ cover. Let $P_{x}$ be the projection of $L^{2}\left(E_{I \mid U}^{d}\right)$ onto $L^{2}\left(E_{I \mid \pi W_{x}}^{d}\right)$ given by multiplication by the characteristic function of $\pi W_{x}$. If $C\left(I \otimes P_{x}\right)=0$ for all
$x \in \pi^{-1} U$ then $C=0$. We may thus assume that $U=W_{1} \times W_{2} \times \ldots \times W_{d}$ as above.

For this special case we need the following simple observation:
(1) $N^{d}-U=\pi\left(\left(\bigcup_{i} \times{ }^{d}\left(M-W_{i}\right)\right) \cap M^{d}\right)$.

Set $Z_{i}=\pi\left(\left(\times{ }^{d}\left(M-W_{i}\right)\right) \cap M^{d}\right)$. Let $Q_{i}$ be the projection of $L^{2}\left(\left.E\right|_{Z}\right)$ onto $L^{2}\left(E_{I \mid Z_{i}}^{d}\right)$ given by multiplication by the characteristic function of $Z_{i}$. Then it is enough to prove that $C_{i}=\left(I \otimes Q_{i}\right) C=0$ for all $i=1, \ldots, d$. So assume that $C_{i} \neq 0$. We now follow precisely the same line of argument as in Lemma 5.1. Let $f \in L^{2}\left(E_{I \mid U}^{d}\right)$ be such that $g=C_{i} f \neq 0$. As above define

$$
\sigma_{1}(X)=\int_{X \times U}\left\|f\left(\lambda-\Psi_{I}(x), x\right)\right\|_{x}^{2} \mathrm{~d} \mu(\lambda) \omega^{d}
$$

and

$$
\begin{equation*}
\sigma_{2}(X)=\int_{X \times Z_{t}}\left\|g\left(\lambda-\Psi_{I}(x), x\right)\right\|_{x}^{2} \mathrm{~d} \mu(\lambda) \omega^{d} \tag{**}
\end{equation*}
$$

Let $\varrho_{i}$ be the representation of $V$ on $L^{2}\left(V^{\prime}, \sigma_{i}\right)$ given by $\left(\varrho_{i}(v) u\right)(\lambda)=$ $\mathrm{e}^{t(v)} u(\lambda)$. Then $C_{t}$ induces a continuous linear map, $L$, of $L^{2}\left(V^{\prime}, \sigma_{1}\right)$ to $L^{2}\left(V^{\prime}, \sigma_{2}\right)$ such that if $\Omega_{t}$ is the function identically equal to 1 on $V^{\prime}$ as an element of $L^{2}\left(V^{\prime}, \sigma_{t}\right)$ then $L \Omega_{1}=\Omega_{2}$ and $L \varrho_{1}(v)=\varrho_{2}(v) L$ f or $v \in V$. As above, this implies that there exists $h \in L^{2}\left(V^{\prime}, \sigma_{1}\right)$ such that $d \sigma_{2}=h d \sigma_{1}$.

Lemma 3.2 implies that there exists $X \in \mathscr{B}$ such that $\mu\left(X+\Psi_{I}(x)\right)=0$ for $x \in U$ and $\mu\left(X+\Psi_{I}(x)\right)=1$ for $x \in Z_{i}$ (take $X$ to be the $X$ in Lemma 3.2 corresponding to $W_{i}$ and $\left.v_{t}\right)$. Thus $\sigma_{1}(X)=0$ and $\sigma_{2}(X)=\|g\|^{2}$. On the other hand, $(*)$ implies that $\sigma_{2}(X)=0$. We have derived our desired contradiction. The lemma now follows.

If $H$ is a Hilbert space then we denote by End $H$ the space of all bounded linear operators of $H$ to $H$ with the strong topology (the topology defined by the semi-norms $\|T v\|, v \in H$ ). If $S \subset$ End $H$ then we set Comm ( $S$ ) = $\{A \in$ End $H: A T=T A$ for all $T \in S\}$. The Von Neumann density Theorem asserts that if $S$ is a subalgebra of End $H$, containing $I$, such that if $T \in S$ then $T^{*} \in S$ then $\operatorname{Comm}(\operatorname{Comm}(S))=\operatorname{Closure}(S)(c f .[D]$, p. 42).

Lemma 5.4. Let the notation be as in the previous lemma. If $v \in V$ then $S(v) \otimes I \subset \operatorname{Closure}\left(\operatorname{Span}\left\{S(v) \otimes \tau_{I, E}(v): v \in V\right)\right.$ ).

Proof. If $X$ is a Borel set in $N^{d}$ let $\chi_{X}$ denote the characteristic function of $X$. If $f \in L^{\infty}\left(N^{d}\right)$ let $M_{f}$ be the operator of multiplication by $f$ on $L^{2}(E)$. Then

$$
\text { Closure }\left(\operatorname{span}\left\{M_{\chi u}: U \text { open in } N^{d}\right\}\right) \supset\left\{M_{f}: f \in L^{\infty}\left(N^{d}\right)\right\}
$$

Now the previous Lemma implies that

$$
\operatorname{Comm}\left(\operatorname{Comm}\left(\operatorname{span}\left\{S(v) \otimes \tau_{I, E}(v): v \in V\right\}\right)\right) \supset\left\{M_{\chi U}: U \text { open in } N^{d}\right\}
$$

Thus the density theorem implies that
Closure $\left(\operatorname{span}\left\{S(v) \otimes \tau_{I, E}(v): v \in V\right\}\right)=A$
contains the operators $I \otimes M_{f}, f \in L^{\infty}\left(N^{d}\right)$. Since $I \otimes \tau_{I, E}(v)$ is such an operator, we see that $A$ contains the operators

$$
\left(I \otimes \tau_{I, E}(-v)\right)\left(S(v) \otimes \tau_{I, E}(v)\right)=S(v) \otimes I \quad \text { for } \quad v \in V
$$

For the next lemma we assume in addition that if $v_{i} \neq v_{j}$ and if $\left(v_{i}, v_{j}\right)>0$ then $v_{i}-v_{j} \in\left\{v_{1}, \ldots, v_{r}\right\}$. We note that this implies that if $v_{i}-v_{j} \neq 0$ then $\left\|v_{i}-v_{j}\right\| \geqslant \min \left\{\left\|v_{k}\right\|: k=1, \ldots, r\right\}$. Indeed, if $\left(v_{i}, v_{j}\right) \leqslant 0$ then $\left\|v_{i}-v_{j}\right\|^{2} \geqslant \mid v_{i}\left\|^{2}+\right\| v_{j} \|^{2}$. If $\left(v_{i}, v_{j}\right)>0$ and $v_{i} \neq v_{j}$ then $v_{i}-v_{j}=v_{k}$ for some $k$.

LEMMA 5.5. Let $d, d^{\prime} \geqslant 0$ and suppose that there exists a continuous non-zero linear map of $L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I}^{d}\right)$ into $L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{J}^{d^{\prime}}\right)$ such that $C\left(S(v) \otimes \tau_{I, E}(v)\right)=\left(S(v) \otimes \tau_{J, E}(v)\right) C$ for all $v \in V$. Then $d=d^{\prime}$ and $I=J$.

Proof. Suppose that $d \neq d^{\prime}$. Then by replacing $C$ by $C^{*}$, if necessary, we may assume that $d>d^{\prime}$. If we argue as above it is enough to prove that $C$ is 0 on

$$
L^{2}\left(V^{\prime}, \mu\right) \hat{\otimes} L^{2}\left(E_{I \mid W_{1} \times \cdots \times W_{d}}^{d}\right)
$$

for $W_{1} \times \cdots \times W_{d} \subset F^{d}$ and $W_{i} \cap W_{j}=\Phi$ if $i \neq j$. We thus replace $C$ by its restriction to this space. If $Q_{i}$ is the usual projection of $L^{2}\left(E_{J}^{d^{\prime}}\right)$ onto $L^{2}\left(\left.E_{J}^{d^{\prime}}\right|_{\left(\times d^{d}\left(M-W_{i}\right)\right) \cap F^{d^{d}}}\right)$ then Lemma 3.2 combined with the argument in Lemma 5.3 implies that $\left(I \otimes Q_{i}\right) C=0$.

Since $d>d^{\prime}$ it is easily seen that $\bigcup_{i=1}^{d} \times{ }^{d^{\prime}}\left(M-W_{i}\right)=\times{ }^{d} M$. This proves the result for $d \neq d^{\prime}$. So assume that $d=d^{\prime}$. In this case we may argue as in the proof of Lemmas 5.2 and 5.3 to see that if $\varphi \in L^{\infty}\left(N^{d}\right)$ then $C\left(I \otimes M_{\varphi}\right)=\left(I \otimes M_{\varphi}\right) C$. This implies that if we define a new action, $\gamma$, on $L^{2}\left(E_{J}^{d}\right)$ given by

$$
\gamma(v) f(x)=\mathrm{e}^{\mathrm{l}\left(\Psi_{f}(x)-\Psi_{J}(x)\right)(v)} f(x)
$$

then $C(S(v) \otimes I)=(S(v) \otimes \gamma(v)) C$.

Our new hypothesis implies that if $v_{i_{k}}-v_{j_{k}} \neq 0$ and if $\operatorname{dim} M=2$ then $\left\|v_{t_{k}}-v_{j k}\right\|>(8 \pi|\varrho(x)|)^{1 / 2}$ for $x \in M$. So if $I \neq J$ then Lemma 5.1 implies that $C=0$.

## 6. Unitary representations of groups of smooth mappings

Let $U$ be a compact Lie group. We write 1 for the identity element of $U$. Let $M$ be as in $\S 3$. Let $\langle,>$ be a Riemannian structure on $M$ and let $\omega$ be a volume element for $M$. Let $G=\left\{g \in C^{\infty}(M ; U): g(m)=1\right.$ outside of a compact set $\}$. If $K$ is a compact subset of $M$ then we set $G_{K}=\{g \in G$ : $g(m)=1$ if $m \notin K\}$. We endow $G_{K}$ with the topology of uniform convergence with all derivatives and look upon $G$ as $\bigcup_{K} G_{K}$. Then $G$ is a topological group under pointwise multiplication.

Let $\mathfrak{u}$ be the Lie algebra of $U$ which we identify with $T U_{1}$ (the tangent space at 1) as usual. Fix $B$, an $\operatorname{Ad}(U)$-invariant inner product on $\mathfrak{u}$. If $x$, $y \in U$ then set $R(y) x=x y^{-1}$. If $W$ is a finite dimensional vector space over $\mathbb{R}$ then let $\Omega^{1}(M, W)$ denote the space of all smooth 1 -forms on $M$ with values in $W$. If $K$ is a compact subset of $M$ then we set $\Omega_{K}^{1}(M ; W)$ equal to the space of all $\eta \in \Omega^{1}(M ; W)$ such that $\eta_{x}=0$ for $x \notin K$. We endow $\Omega_{K}^{1}(M ; W)$ with the topology of uniform convergence with all derivatives. We set $\Omega_{c}^{1}(M ; W)=\bigcup_{K} \Omega_{K}^{1}(M ; W)$ with the corresponding union topology.

Before we introduce the main results of this paper let us record a result which we feel is necessary in the course of their proof. The argument below is based on a suggestion of A . Borel.

Lemma 6.1. $d\left(C_{c}^{\infty}(M ; W)\right)$ is closed in $\Omega_{c}^{1}(M ; W)$.

Proof. It is enough to prove this result for $W=\mathbb{R}$. Let $N$ be a connected, paracompact, orientable, smooth $n$-dimensional manifold. Set $Z_{c}^{k}(N)=$ $\left\{\eta \in \Omega_{c}^{k}(N): \mathrm{d} \eta=0\right\}$ and put $B_{c}^{k}(N)=\mathrm{d} \Omega_{c}^{k-1}(N)$. Then $H_{c}^{k}(N)=$ $Z_{c}^{k}(N) / B_{c}^{k}(N)$ is called the $k$-th (de Rham) cohomology of $M$ with compact support. If $\eta \in \Omega_{c}^{k}(N)$ and if $v \in \Omega^{n-k}(N)$ then set

$$
(\eta \mid v)=\int_{N} \eta \Lambda v .
$$

Then $(\mathrm{d} \eta \mid v)=(-1)^{k-1}(\eta \mid \mathrm{d} v)$ for $\eta \in \Omega_{c}^{k-1}(N), v \in \Omega^{n-k}(N)$ and if we set $Z^{k}(N)=\left\{\eta \in \Omega^{k}(N): \mathrm{d} \eta=0\right\}, B^{k}(N)=\mathrm{d} \Omega^{k-1}(N), H^{k}(N)=Z^{k}(N) / B^{k}(N)$ then ([deR; §22, §23]) (|) induces a nondegenerate pairing of $H_{c}^{k}(N)$ with
$H^{n-k}(N)$ (i.e. $\left(\eta \mid H^{n-k}(N)\right)=0$ implies $\eta=0$ ). Thus if $\operatorname{dim} H^{n-1}(N)<\infty$ then $\operatorname{dim} H_{c}^{1}(N)<\infty$.
(1) If $\operatorname{dim} H^{n-1}(N)<\infty$ then $\mathrm{d}_{c}^{\infty}(N)$ is closed in $\Omega_{c}^{1}(N)$.

Indeed, $Z_{c}^{1}(N)$ is clearly closed in $\Omega_{c}^{1}(N)$. Let $Z$ be a finite dimensional subspace of $Z_{c}^{1}(N)$ such that $Z_{c}^{1}(N)=\mathrm{d} C_{c}^{\infty}(N) \oplus Z$. If $N$ is compact, then choose $x_{0} \in N$ and let $V=\left\{f \in C^{\infty}(N): f\left(x_{0}\right)=0\right\}$. If $N$ is non-compact then set $V=C_{c}^{\infty}(N)$. We set $A(v, z)=\mathrm{d} v+z$ for $v \in V, z \in Z$. Then

$$
A: V \times Z \rightarrow Z_{c}^{1}(N)
$$

is continuous and bijective. Since $V \times Z$ is an $L F$ space this implies that $A^{-1}$ is continuous. Thus $A(V \times\{0\})=\mathrm{d} C_{c}^{\infty}(N)$ is closed.

We now return to $M$. Let $\left\{U_{i}\right\}$ be a covering of $M$ such that all non-empty finite intersections of the $U_{i}$ are contractible (e.g., take a covering by convex neighborhoods relative to 〈, >). For $m=1,2, \ldots$, define $N_{m}=\bigcup_{i \leqslant m} U_{i}$. Then $\operatorname{dim} H^{n-1}\left(N_{m}\right)<\infty$ for all $m$. If $K$ is a compact subset of $M$ then $K \subset N_{m}$ for some $m$. Thus (1) implies that $\mathrm{d} C_{K}^{\infty}(M)$ is closed. This completes the proof.

We now return to the situation at the beginning of this section. Let $\beta$ : $G \rightarrow \Omega_{c}^{1}(M ; \mathfrak{u})$ be defined by

$$
\beta(g)_{x}=\mathrm{d} R(g(x))_{g(x)} \mathrm{d} g_{x}
$$

(i.e., $\beta(g)=(\mathrm{d} g) g^{-1}$ ). We set $V=\Omega_{c}^{1}(M ; \mathfrak{u})$ and define a representation $\pi$ of $G$ on $V$ by $(\sigma(g) \eta)_{x}=A \mathrm{~d}(g(x))\left(\eta_{x}(v)\right), g \in G, x \in M, v \in T M_{x}$. Then $\beta(x y)=\beta(x)+\sigma(x) \beta(y)$. We will sometimes write $A \mathrm{~d}(g) \eta$ for $\sigma(g) \eta$.

For $\eta, v \in V$, let

$$
(\eta, v)=\int_{M}\left(\eta_{x}, v_{x}\right)_{x} \omega
$$

Here we use the inner product on $T M_{x}^{*} \otimes \mathfrak{u}$ corresponding to $\langle,>$ and $B$. It will sometimes be necessary to write

$$
(,)=(, \quad)_{\langle,\rangle, B, \omega}
$$

Let $H$ denote the Hilbert space completion of $V$ with respect to (, ). Let $\mu$ be the Gaussian measure corresponding to this inner product (§1). Then $\mu$ is countably additive (cf. [GV], Theorem 6, p. 332). We note that if $g \in G$ then $\sigma(g)$ extends to a unitary operator on $H$ and $(\sigma, H)$ is a unitary representation of $G$.

If $f \in L^{2}\left(V^{\prime}, \mu\right)$ and if $g \in G$ then we set $T(g) f(\lambda)=\mathrm{e}^{i \lambda(\beta(g))} f(\lambda \cdot \sigma(g))$. Then $\left(T, L^{2}\left(V^{\prime}, \mu\right)\right.$ ) is a (strongly continuous) unitary representation of $G$. We will also write $T=T_{\langle,\rangle, B, \omega}$.

Let $\mathfrak{h}$ be a maximal abelian subalgebra of $\mathfrak{u}$. Let $\Delta$ be the root system of $\mathfrak{u}_{\mathbf{c}}$ with respect to $\mathfrak{h}_{\mathbb{C}}$ (here sub $\mathbb{C}$ indicates complexification). If $\alpha \in \Delta$ then $\left.\alpha\right|_{\mathfrak{h}}=i \tilde{\alpha}$ with $\tilde{\alpha} \in \mathfrak{h}^{\prime}$. Let $\|\alpha\|_{B}$ be the norm of $\tilde{\alpha}$ relative to $\left.B\right|_{\mathfrak{h}}$. The main theorem is

Theorem 6.2. Assume that $U$ is semi-simple. If $\operatorname{dim} M=2$ then we assume that $\|\alpha\|_{B}>(8 \pi|\varrho(x)|)^{1 / 2}$, for $x \in M\left(\omega=\varrho \operatorname{Vol}_{\langle,\rangle}\right)$and $\alpha \in \Delta$. If $\operatorname{dim} M \geqslant 2$ then $T_{\langle,\rangle, B, \omega}$ is irreducible. Let $\langle,\rangle_{1}, B_{1}, \omega_{1}$ be an arbitrary triple as above and let $\mu_{1}$ be the corresponding Gaussian measure on $V^{\prime}$. If $C$ is a non-zero bounded operator from $L^{2}\left(V^{\prime}, \mu\right)$ to $L^{2}\left(V^{\prime}, \mu_{1}\right)$ such that $C T_{\langle,\rangle, B, \omega}(g)=$ $T_{\langle,\rangle_{1}, B_{1}, \omega_{1}}(g) C$ for $g \in G$ then

$$
(,)_{\langle,\rangle, B, \omega}=(,)_{\langle,\rangle_{1}, B_{1}, \omega_{1}} .
$$

The proof of this result will involve more notation and concepts. For the moment we assume that $V$ is a locally convex, separable, topological vector space over $\mathbb{R}$. Let $(, \quad)$ be an inner product on $V$ and let $\mu$ be the corresponding Gaussian measure on $V^{\prime}$. We assume that $V_{1}, V_{2}$ are closed subspaces of $V$ such that $V=V_{1} \oplus V_{2}$ and that $\left(V_{1}, V_{2}\right)=0$. Let $\mu_{i}$ be the corresponding Gaussian measure on $V_{t}^{\prime}$ for $i=1,2$. We assume that $\mu_{i}$ is countably additive for $i=1,2$. We identify $V_{1}^{\prime}$ with $\left\{\lambda \in V^{\prime}: \lambda\left(V_{2}\right)=0\right\}$. Then $V^{\prime}=V_{1}^{\prime} \oplus V_{2}^{\prime}$. It is easily seen that $\mu=\mu_{1} \times \mu_{2}$ (product measure). Thus Fubini's theorem implies that the map $S$ from $L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) \hat{\otimes} L^{2}\left(V_{2}^{\prime}, \mu_{2}\right)$ to $L^{2}(V, \mu)$ given by $S(f \otimes g)(\lambda+v)=f(\lambda) g(v)$ is a unitary isomorphism. Here $\hat{\otimes}$ denotes completed tensor product.

We will also make use of the Fock space. Let $\mathscr{F}(H)=\oplus_{n \geqslant 0} \hat{S}^{n}(H)_{\mathbf{c}}$. Here $\hat{S}^{n}(H)$ is the completed $n$-fold symmetric power of $H$ (the inner product is defined by $\left.\left(v^{n}, v^{n}\right)=\|v\|^{2 n}\right)$. The subscript $\mathbb{C}$ will always indicate complexification with the Hermitian extension if there is an inner product. If $h \in H$ then set

$$
\operatorname{EXP} h=\sum_{n \geqslant 0} h^{n} / \sqrt{n!}
$$

Here $h^{0}=0$ and $h^{n}=h \otimes h \otimes \ldots \otimes h, n$-times. Then (EXP $h_{1}$, EXP $h_{2}$ ) $=$ $\mathrm{e}^{\left(h_{1}, h_{2}\right)}$. The next results that we will be describing can be found in [Gu], §2.1, §7.2. Span $\{\operatorname{EXP} v: v \in V\}$ is dense in $\mathscr{F}(H)$. We set for $v \in V$,
$\mathrm{e}_{v}(\lambda)=\exp (i \lambda(v)+(v, v) / 2)$. Then $\operatorname{span}\left\{\mathrm{e}_{v}: v \in V\right\}$ is dense in $L^{2}\left(V^{\prime}, \mu\right)$. Furthermore $\{\operatorname{EXP}(v): v \in V\}$ and $\left\{\mathrm{e}_{v}: v \in V\right\}$ are linearly independent sets in their respective spaces. Since $\left(\mathrm{e}_{v}, \mathrm{e}_{w}\right)=\mathrm{e}^{(v, w)}$ for $v, w \in V$ we can define a natural isometry $F_{V}: L^{2}\left(V^{\prime}, \mu\right) \rightarrow \mathscr{F}(H)$ by $F_{V}\left(\mathrm{e}_{v}\right)=$ EXP $v$.

If $V=V_{1} \oplus V_{2}$ as above then $\mathscr{F}(H)=\mathscr{F}\left(H_{1}\right) \hat{\otimes} \mathscr{F}\left(H_{2}\right)$ and $F_{V}\left(L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) 1_{V_{2}^{\prime}}\right)=$ Closure $\left(\operatorname{span}\left\{\operatorname{EXP} v: v \in V_{1}\right\}\right)$. At this point we will explain and fix an error in [GGV, II]. In that paper they look at a situation such as this and consider the space $Q=\left\{f \in L^{2}\left(V^{\prime}, \mu\right): f(\cdot+\lambda)=f(\cdot)\right.$ for $\left.\lambda \in V_{2}^{\prime}\right\}$. We assert that if $V_{2} \neq H_{2}$ then $Q$ is not defined. Indeed, let $\lambda \in V_{2}^{\prime}-H_{2}$. Then there exists $X \in \mathscr{B}$ such that $\mu(X)=1, \mu(X-\lambda)=0$ (Lemma 1.2). Let $\chi_{X}$ be the characteristic function of $X$. If $f \in L^{2}\left(V^{\prime}, \mu\right)$ then $\chi_{X} f=f$. But $\chi_{X}(\cdot+\lambda)=\chi_{X-\lambda}=0$ as an element of $L^{2}\left(V^{\prime}, \mu\right)$. Thus if " $f \in Q$ " then $\chi_{X+\lambda} f=f$ so $f=0$ in $L^{2}$. We replace this nonsense with the following result.

Lemma 6.3. Let $V=V_{1} \oplus V_{2}=V_{3} \oplus V_{4}$ be two decompositions of $V$ as above with closures in $H, H_{i}, i=1,2,3,4$. Suppose that $H_{1} \cap H_{3}=(0)$. Then $L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) 1_{V_{2}^{\prime}} \cap L^{2}\left(V_{3}^{\prime}, \mu_{3}\right) 1_{V_{4}^{\prime}}=\mathbb{C} 1_{V^{\prime}}$. Here the $\mu_{i}$ are the Gaussian measures corresponding to $\left.()\right|_{,V_{i}}$ respectively and we assume that the properties above of $\left(V_{1}, V_{2}\right)$ are satisfied by $\left(V_{3}, V_{4}\right)$.

Proof. By the above we must show that $\left(\oplus \hat{S}^{n}\left(H_{1}\right)\right) \cap\left(\oplus \hat{S}^{n}\left(H_{3}\right)\right)=\mathbb{C} 1$. So suppose that $a=\Sigma a_{n} \in\left(\oplus \hat{S}^{n}\left(H_{1}\right)\right) \cap\left(\oplus \hat{S}^{n}\left(H_{3}\right)\right)$. Then comparing homogeneity, we see that $a_{n} \in \hat{S}^{n}\left(H_{1}\right) \cap \hat{S}^{n}\left(H_{3}\right)$ for all $n$. Thus we need only show that $\hat{S}^{n}\left(H_{1}\right) \cap \hat{S}^{n}\left(H_{3}\right)=(0)$ for $n>0$. If $n=1$ this just says that $H_{1} \cap H_{3}=(0)$ as assumed. So assume, inductively, the desired result for $n(\geqslant 1)$. If $a \in H$ then define $\partial_{a}: \hat{S}^{n+1}(H) \rightarrow \hat{S}^{n}(H)$ by $\partial_{a} x^{n+1}=$ $(n+1)(x, a) x^{n}$. Then $\partial_{a}$ defines a bounded operator. Furthermore, if $x \in \hat{S}^{n+1}\left(H_{1}\right)$ and $\partial_{a} x=0$ for all $a \in H_{1}$ then $x=0$. Now $\partial_{a}\left(\hat{S}^{n+1}\left(H_{i}\right)\right) \subset$ $\hat{S}^{n}\left(H_{i}\right)$ for $i=1,2$ and $a \in H_{1}$. Thus the inductive hypothesis implies that $\partial_{a}\left(\hat{S}^{n+1}\left(H_{1}\right) \cap \hat{S}^{n+1}\left(H_{3}\right)\right)=(0)$ hence $\hat{S}^{n+1}\left(H_{1}\right) \cap \hat{S}^{n+1}\left(H_{3}\right)=(0)$.

We now begin the proof of Theorem 6.2. Let $\mathfrak{h}$ be a maximal abelian subalgebra of $\mathfrak{u}$. Let $\mathfrak{h}^{\perp}=\{X \in \mathfrak{u}: B(\mathfrak{h}, X)=0\}$. Put $V_{1}=\Omega_{c}^{1}(M ; \mathfrak{h})$, and $V_{2}=\Omega_{c}^{1}\left(M ; \mathfrak{h}^{\perp}\right)$. Then $V=V_{1} \oplus V_{2}$ as above. We fix a base point $x_{0}$ (as usual) if $M$ is compact, and set $A=\left\{f \in C^{\infty}(M ; \mathfrak{h})\right.$ : $\left.f\left(x_{0}\right)=0\right\}$ if $M$ is compact and $A=C_{c}^{\infty}(M ; \mathfrak{h})$ otherwise. Endow $A$ with the topology given as in $\S 3,4$. We look upon $A$ as an abelian topological group. Define $W(a)=$ $T(\exp a)$ for $a \in A$. Then $W$ defines a (strongly continuous) unitary representation of $A$ of $L^{2}\left(V^{\prime}, \mu\right)$.

As above $L^{2}\left(V^{\prime}, \mu\right)=L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) \otimes L^{2}\left(V_{2}^{\prime}, \mu_{2}\right)$. Under this identification the action $W$ is given as follows

$$
W(a) f(\lambda, \xi)=\mathrm{e}^{i \lambda(\mathrm{~d} a)} f(\lambda, \xi \cdot \sigma(\exp a)), \quad \lambda \in V_{1}^{\prime}, \quad \xi \in V_{2}^{\prime}, \quad a \in A
$$

Thus if we define $W_{1}(a) f(\lambda)=\mathrm{e}^{i \lambda(d a)} f(\lambda)$ for $f \in L^{2}\left(V_{1}^{\prime}, \mu_{1}\right)$ and $W_{2}(a) f(\xi)=$ $f(\xi \cdot \sigma(\exp a))$ for $f \in L^{2}\left(V_{2}^{\prime}, \mu_{2}\right)$ then under the above identification $W=W_{1} \otimes W_{2}$. We now analyze each of the representations $W_{i}, i=1,2$.

Let $H_{i}$ denote the Hilbert space completion of $V_{i}, i=1,2$. Let $F_{i}=F_{V_{i}}$ : $L^{2}\left(V_{i}^{\prime}, \mu_{i}\right) \rightarrow \mathscr{F}\left(H_{i}\right)$ (defined as above). We note that $W_{2}(a) \mathrm{e}_{h}=\mathrm{e}_{\sigma(\operatorname{expa} a)}$ for $h \in V_{2}$. Thus

$$
F_{2} \cdot W_{2}(a)=\left(\oplus_{m \geqslant 0} S^{m}(\sigma)(\exp a)\right) \cdot F_{2}, \quad a \in A
$$

Here $S^{m}(\sigma)$ is the representation on $\hat{S}^{m}\left(H_{2}\right)$ corresponding to $\left.\sigma(a)\right|_{H_{2}}$.
We have $\left(\mathfrak{h}^{\perp}\right)_{\mathbb{C}}=\oplus_{\alpha \in \Delta}\left(\mathfrak{u}_{\mathbf{C}}\right)_{\alpha}$, an orthogonal direct sum relative to the Hermitian extension of $B$ to $\mathfrak{u}_{\mathbf{C}}$ (here, as usual, $\left(\mathfrak{u}_{\mathbf{C}}\right)_{\alpha}$, is the $\alpha$-root space). Let $\alpha_{1}, \ldots, \alpha_{r}$ be an enumeration of $\Delta$. Thus the complexification of the vector bundle $\operatorname{Hom}\left(T M, \mathfrak{h}^{\perp}\right)$ splits into a direct sum $E_{1} \oplus E_{2} \oplus \ldots \oplus E_{r}$ where $E_{i}$ is the complexification of $T M^{*}$ and $A$ acts on $E_{j}$ via $\beta_{J}(v)_{x} u=\mathrm{e}^{i \tilde{\delta}_{J}(a(x))} u$ for $u \in\left(E_{J}\right)_{x}$. Let $\beta$ denote the action $\beta_{1} \oplus \ldots \oplus \beta_{r}$ on $E_{1} \oplus \ldots \oplus E_{r}=E$. Then $H_{2}=L^{2}(E)$. Here we are using the notation in $\S 5$. Thus $\hat{\otimes}^{d} H_{2}=L^{2}\left(\otimes^{d} E\right)$ where $\otimes^{d} E$ is a vector bundle over $M^{d}$. The action $\otimes^{d} \sigma$ goes over to $\otimes^{d} \beta$ as in $\S 5$. Under this identification the action of $S_{d}$ on $\hat{\otimes}^{d} H_{2}$ given by $s\left(v_{1} \otimes \ldots \otimes v_{d}\right)=v_{s 1} \otimes \ldots \otimes v_{s d}$ corresponds to $\left(s^{-1} f\right)\left(x_{1}, \ldots, x_{d}\right)=s^{-1} f\left(x_{s 1}, \ldots x_{s d}\right)$. Let $F^{d}$ be as in Lemma 5.2. Thus, as in $\S 5, \hat{S}^{d}\left(H_{2}\right)=L^{2}\left(\left.\otimes^{d} E\right|_{F^{d}}\right)$. We write (as in $\left.\S 5\right)\left.\otimes^{d} E\right|_{F^{d}}=E^{d}$. Then $E^{d}=\oplus E_{I}^{d}$ an orthogonal direct sum over $I=\left(i_{1}, \ldots, i_{d}\right), 1 \leqslant i \leqslant r$ (here we have replaced the $v_{j}$ with $\tilde{\alpha}_{j}$ ). The conditions on the roots in Theorem 6.2 imply that the $\tilde{\alpha}_{j}$ satisfy all of the conditions on the $v_{j}$. Indeed, it is standard that if $\alpha, \tau$ are roots and if $\alpha \neq \tau$, and $B(\tilde{\alpha}, \tilde{\tau})>0$ then $\alpha-\tau$ is a root. We therefore see that $\mathscr{F}\left(H_{2}\right) \simeq \mathbb{C} 1 \oplus \oplus_{d>0} \oplus_{I} L^{2}\left(E_{I}^{d}\right)$ where the action on each $L^{2}\left(E_{I}^{d}\right)$ is given by $\tau_{I, E}$ as in $\S 5$.

We now analyze $W_{1}$. Lemma 6.1 implies that $\mathrm{d} A$ is closed in $V_{1}$. Let $Z=$ $\left\{h \in H_{1}:(\mathrm{d} A, h)=0\right\}$. Then $H_{1}=\operatorname{Closure}(\mathrm{d} A) \oplus Z$. Define $Q(a, b)=$ $(\mathrm{d} a, \mathrm{~d} b), a, b \in A$. Then $Q$ defines a continuous inner product on $A$. Let $\mu_{Q}$ be the corresponding Gaussian measure on $A^{\prime}$. We define a representation of $A$ on $L^{2}\left(A^{\prime}, \mu_{Q}\right), S$, by $S(a) f(\lambda)=\mathrm{e}^{i \lambda(a)} f(\lambda)$. We note that as a representation of $A, L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) \simeq L^{2}\left((\mathrm{~d} A)^{\prime}, \mu^{\prime}\right) \otimes \mathscr{F}(Z)$ with $\mu^{\prime}$, the Gaussian measure on $(\mathrm{d} A)^{\prime}$ induced by $($,$) restricted to \mathrm{d} A$ and the action of $A$ on
$L^{2}\left((\mathrm{~d} A)^{\prime}, \mu^{\prime}\right)$ is given by $\xi(a) f(\lambda)=\mathrm{e}^{\lambda /(\mathrm{d} a)} f(\lambda)$. Since $\mathrm{d}: A \rightarrow \mathrm{~d} A$ is continuous, linear and bijective the closed graph theorem implies that it is a topological isomorphism. The pullback of $\mu^{\prime}$ is $\mu_{Q}$. Thus as a representation of $A,\left(W_{1}, L^{2}\left(V_{1}^{\prime}, \mu_{1}\right)\right)$ is equivalent with $\left(S \otimes I, L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z)\right)$. Let us recapitulate our analysis in the following result.

Lemma 6.4. $\left(W, L^{2}\left(V^{\prime}, \mu\right)\right)$ is unitarily equivalent with the direct sum of $\left(S \otimes I, L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z)\right)$ and

$$
\oplus_{d>0} \oplus_{I}\left(S \otimes I \otimes \tau_{I, E}, L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z) \otimes L^{2}\left(E_{I}^{d}\right)\right)
$$

Furthermore, the unitary equivalence, $F$, can be chosen so that $F\left(L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) 1_{V_{2}^{\prime}}\right)=L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z)$.

We use this decomposition to prove the theorem.
(1) Closure $\left(\operatorname{span}\left\{T(g) 1_{V^{\prime}}: g \in G\right\}\right)=L^{2}\left(V^{\prime}, \mu\right)$.

Indeed, let $C: L^{2}\left(V^{\prime}, \mu\right) \rightarrow L^{2}\left(V^{\prime}, \mu\right)$ be such that $C \cdot W(a)=W(a) \cdot C$ for $a \in A$. Set $C_{1}=F C F^{-1}$. Let $P$ be the orthogonal projection of

$$
L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z) \oplus\left(\oplus_{d>0} \oplus_{I} L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z) \otimes L^{2}\left(E_{I}^{d}\right)\right)
$$

onto $L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z)$ and let $P_{d, I}$ be the orthogonal projection onto

$$
L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z) \otimes L^{2}\left(E_{I}^{d}\right)
$$

Then Lemma 5.5 implies that $P C_{1} P_{d, I}=P_{d, I} C_{1} P=0$ and $P_{d, I} C_{1} P_{d^{\prime}, J}=0$ if $d \neq d^{\prime}$ or if $d=d^{\prime}$ and $I \neq J$. Thus Lemma 5.4 implies that

$$
\text { Closure }\left(\operatorname{span}\left\{F W(a) F^{-1}: a \in A\right\}\right) \supset\{\oplus S(a) \otimes I: a \in A\}
$$

Now,
Closure $(\operatorname{span}\{W(a): a \in A\})=F^{-1}$ Closure $\left(\operatorname{span}\left\{F W(a) F^{-1}: a \in A\right\}\right) F$.
Thus, if we set $v(\eta) f(\lambda)=\mathrm{e}^{i \lambda(\eta)} f(\lambda)$ for $f \in L^{2}\left(V^{\prime}, \mu\right)$ and $\eta \in V$ then

$$
F^{-1}\{\oplus S(a) \otimes I: a \in A\} F=\{v(\mathrm{~d} a): a \in A\} .
$$

If $x \in \mathfrak{u}$ then $x$ is contained in a maximal abelian subalgebra of $\mathfrak{u}$. Thus
Closure $(\operatorname{span}\{T(g): g \in G\}) \supset\left\{v(\mathrm{~d} f): f \in C_{c}^{\infty}(M ; \mathfrak{u})\right\}$.
Now $T(g) v(X) T(g)^{-1}=v(A d(g) X)$ for $g \in G, X \in V$.
We assert that $L=\operatorname{span}\left\{A \mathrm{~d}(g) \mathrm{d} f: g \in G, f \in C_{c}^{\infty}(M ; \mathfrak{u})\right\}$ is dense in $H$. In fact, the map $X \rightarrow A \mathrm{~d}(\exp X) \mathrm{d} f$ is real analytic from $C_{c}^{\infty}(M ; \mathfrak{u})$ to $H$. We
may thus differentiate to find that

$$
[X, \mathrm{~d} Y] \in \operatorname{Closure}\left(\operatorname{span}\left\{A \mathrm{~d}(g) \mathrm{d} f: g \in G, f \in C_{c}^{\infty}(M ; \mathfrak{u})\right\}\right)
$$

for $X, Y \in C_{c}^{\infty}(M ; \mathfrak{u})$. If $x, y \in \mathfrak{u}$ and if $f, g \in C_{c}^{\infty}(M)$ then consider $X=f x$, $Y=g y$. Then $[X, \mathrm{~d} Y]=f \mathrm{~d} g \otimes[x, y]$. Since $u$ is assumed to be semi-simple $[\mathfrak{u}, \mathfrak{u}]=\mathfrak{u}$. Thus if $X_{1}, \ldots, X_{d}$ is an orthonormal basis of $\mathfrak{u}$ then

Closure $\left(\operatorname{span}\left\{A \mathrm{~d}(g) \mathrm{d} f: g \in G, f \in C_{c}^{\infty}(M ; \mathfrak{u})\right\}\right)$
$\supset\left\{\sum \eta_{i} \otimes X_{i}: \eta_{i} \in \operatorname{span}\left\{f \mathrm{~d} g: f, g \in C_{c}^{\infty}(M)\right\}\right\}$.
This latter set is obviously dense in $H$. We conclude that Closure $(\operatorname{span}\{T(g)$ : $g \in G\}$ ) contains $\{v(X): X \in L\}$ with $L \subset V$ a dense subspace of $H$. The above described isomorphism of $L^{2}\left(V^{\prime}, \mu\right)$ with the Fock space on $H$ now implies that $\operatorname{span}\left\{v(X) 1_{V^{\prime}}: X \in L\right\}$ is dense in $L^{2}\left(V^{\prime}, \mu\right)$. This proves (1).
(2) If $C$ is a continuous linear operator on $L^{2}\left(V^{\prime}, \mu\right)$ such that $C \cdot T(g)=T(g) \cdot C$ for all $g \in G$ then $C 1_{V^{\prime}} \in \mathbb{C} 1_{V^{\prime}}$.
Let us show how (2) now implies the first assertion of the Theorem. Let $C$ be an operator as above that commutes with the action of $T(g)$ for $g \in G$. Then $C 1_{V^{\prime}}=c 1_{V^{\prime}}$. This implies that $C$ acts by $c I$ on $\operatorname{span}\left\{T(g) 1_{V^{\prime}}: g \in G\right\}$. (1) implies that this space is dense so $C$ acts by $c I$. Hence $T$ is irreducible.

We now prove (2). The argument in the proof of (1) implies that $C F^{-1}\left(L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z)\right) \subset F^{-1}\left(L^{2}\left(A^{\prime}, \mu_{Q}\right) \otimes \mathscr{F}(Z)\right)$ since in particular $C W(a)=W(a) C$ for $a \in A$. But then $C\left(L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) 1_{V_{2}^{\prime}}\right) \subset L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) 1_{V_{2}^{\prime}}$. To complete the proof of (2), we need the following structural property of $\mathfrak{u}$ (this is where the semi-simplicity of $U$ is used).
(3) There exists a maximal abelian subalgebra $\mathfrak{h}_{1}$ of $\mathfrak{u}$ such that $\mathfrak{h} \cap \mathfrak{h}_{1}=(0)$.
We note that there exists $X \in \mathfrak{u}$ such that if $h \in \mathfrak{h}$ and if $[h, X]=0$ then $h=0$. Indeed, choose $X \in \mathfrak{h}^{\perp}$ such that its projection onto every root space is non-zero. Set $\mathfrak{u}_{1}=\{Y \in \mathfrak{u}:[Y, X]=0\}$. Choose $\mathfrak{h}_{1}$ to be a maximal abelian subalgebra of $\mathfrak{u}_{1}$. If $Y \in \mathfrak{u}$ and if $\left[Y, \mathfrak{h}_{1}\right]=(0)$ then, in particular, $[Y, X]=0$. Thus $Y \in \mathfrak{h}_{1}$. So $\mathfrak{h}_{1}$ is maximal abelian in $\mathfrak{u}$. Since $\mathfrak{u}_{1} \cap \mathfrak{h}=(0)$, $\mathfrak{h}_{1} \cap \mathfrak{h}=(0)$. This proves (3).

Let $V_{3}=\Omega_{c}^{1}\left(M ; \mathfrak{h}_{1}\right)$ and $V_{4}=\Omega_{c}^{1}\left(M ; \mathfrak{h}_{1}^{\perp}\right)$. Let $\mu_{3}, \mu_{4}$ be the corresponding Gaussian measures on $V_{3}^{\prime}$ and $V_{4}^{\prime}$. Then the above argument applied to $\mathfrak{h}_{1}$ instead of $\mathfrak{b}$ implies that $C\left(L^{2}\left(V_{3}^{\prime}, \mu_{3}\right) 1_{V_{4}^{\prime}}\right) \subset L^{2}\left(V_{3}^{\prime}, \mu_{3}\right) 1_{V_{4}^{\prime}}$. Now Lemma 6.3 implies that $\left.L^{2}\left(V_{1}^{\prime}, \mu_{1}\right) 1_{V_{2}^{\prime}}, \mu_{3}\right) 1_{V_{4}^{\prime}}=\mathbb{C} 1_{V^{\prime}}$. This proves (2) and hence completes the proof of the first (irreducibility) part of Theorem 6.2.

We now prove the second assertion. Let $C$ be as in the second part of Theorem 6.2. We use the notation in Lemma 5.1. We also write $V_{\mathfrak{h}}$ for
$\Omega_{c}^{1}(M ; \mathfrak{h})$ and $V_{\mathfrak{h} \perp}$ for $\Omega_{c}^{1}\left(M ; \mathfrak{h}^{\perp}\right)$. We write $\mu_{\mathfrak{h}}$ (resp. $\mu_{1, \mathfrak{b}}$ ) for the Gaussian measures on $V_{\mathrm{h}}^{\prime}$ corresponding to the inner product $(,)_{\langle,\rangle, B, \omega}=($, (respectively, $\left.(,)_{\langle,\rangle_{1}, B_{1}, \omega_{1}}=(,)_{1}\right)$.

Since $1_{V^{\prime}}$ is a cyclic vector for $T_{\langle,\rangle, B, \omega}=T, C 1_{V^{\prime}} \neq 0$. If we argue as in the proof of (2) above using Lemma 5.1 we find that $C\left(L^{2}\left(V_{\mathfrak{h}}^{\prime}, \mu_{\mathfrak{h}}\right) 1_{\left(V_{\mathfrak{h})^{\prime}}\right)}\right) \subset$ $L^{2}\left(V_{\mathfrak{h}}^{\prime}, \mu_{1, \mathfrak{h}}\right) 1_{\left(V_{\mathfrak{h}},\right)}$ for all maximal abelian subalgebras $\mathfrak{h}$ of $\mathfrak{u}$. Thus we apply Lemma 6.3 and (3) we find that $C 1_{V^{\prime}} \subset \mathbb{C}_{V^{\prime}}$. Set $\Omega$ (respectively, $\Omega_{1}$ ) equal to $1_{V^{\prime}}$ as an element of $L^{2}\left(V^{\prime}, \mu\right)$ (respectively, $L^{2}\left(V^{\prime}, \mu_{1}\right)$ ). Then we assume that $C \Omega=\Omega_{1}$.

Since $T$ is irreducible, $C^{*} C$ is a multiple of the identity. So we may assume that $C^{*} \Omega_{1}=\Omega$. This implies that

$$
\langle T(g) \Omega, \Omega\rangle_{L^{2}\left(V^{\prime}, \mu\right)}=\left\langle T_{\langle,\rangle, B_{1}, \omega_{1}}(g) \Omega_{1}, \Omega_{1}\right\rangle_{L^{2}\left(V^{\prime}, \mu_{1}\right)}
$$

for all $g \in G$. The obvious calculation of the left and right side of this equation implies that

$$
\mathrm{e}^{(\beta(g), \beta(g)) / 2}=\mathrm{e}^{(\beta(g), \beta(g))_{1} / 2} \quad \text { for all } \quad g \in G
$$

Now, if $X \in C_{c}^{\infty}(M ; \mathfrak{u})$ then $\beta(\exp t X)=t \mathrm{~d} X+\mathrm{O}\left(t^{2}\right)$. We therefore conclude that

$$
\begin{equation*}
(\mathrm{d} X, \mathrm{~d} X)=(\mathrm{d} X, \mathrm{~d} X)_{1} \quad \text { for } \quad X \in C_{c}^{\infty}(M, \mathfrak{u}) \tag{*}
\end{equation*}
$$

We first show that $(*)$ implies that $B_{1}=t B$ for some $t>0$. Indeed, let $X_{1}, \ldots, X_{d}$ be an orthonormal basis of $\mathfrak{u}$ relative to $B$ such that $B_{1}\left(X_{i}, X_{j}\right)=\lambda_{i} \delta_{i, j}$. Set $X=\Sigma f_{i} \otimes X_{i}$ with $f_{i} \in C_{c}^{\infty}(M)$. Then (*) implies that

$$
\sum \int_{M}\left\langle\mathrm{~d} f_{i}, \mathrm{~d} f_{i}\right\rangle \omega=\sum \lambda_{i} \int_{M}\left\langle\mathrm{~d} f_{i}, \mathrm{~d} f_{i}\right\rangle_{1} \omega_{1}
$$

Since this is true for all such $f_{i}$, it is clear that all the $\lambda_{i}$ are equal to (say) $t$. If we change $\omega_{1}$ to $t \omega_{1}$ we may thus assume that $B=B_{1}$.

The second part of Theorem 6.2 now follows from
Lemma 6.5. If

$$
\int_{M}\langle\mathrm{~d} f, \mathrm{~d} f\rangle \omega=\int_{M}\langle\mathrm{~d} f, \mathrm{~d} f\rangle_{1} \omega_{1}
$$

for all $f \in C_{c}^{\infty}(M)$ then $\langle,\rangle_{x} \omega_{x}=\langle,\rangle_{1, x} \omega_{1, x}$ for all $x \in M$.

Proof. Let $x \in M$. By taking local coordinates, we may assume that $M=\mathbb{R}^{n}$ and $x=0$. We fix the usual inner product, (, ), on $\left(\mathbb{R}^{n}\right)^{*}$. Then $\langle v, v\rangle_{y}=$ $(G(y) v, v)$ and $\langle v, v\rangle_{1, v}=\left(G_{1}(y) v, v\right)$. We may assume that $G(0)=I$. Also $\omega=u \mathrm{~d} y$ and $\omega_{1}=u_{1} \mathrm{~d} y$. We must prove that $u(0) I=u_{1}(0) G_{1}(0)$. Let $\varphi \in$ $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ be arbitrary. For $\varepsilon>0$, let $\varphi_{\varepsilon}(y)=\varphi(y / \varepsilon)$. Then $\mathrm{d} \varphi_{\varepsilon}(y)=$ $\varepsilon^{-1} \mathrm{~d} \varphi(y / \varepsilon)$. A direct calculation yields

$$
\begin{aligned}
& \int_{\mathbf{R}^{n}}\left(G(y) \mathrm{d} \varphi_{\varepsilon}(y), \mathrm{d} \varphi_{\varepsilon}(y)\right) u(y) \mathrm{d} y \\
& \quad=\varepsilon^{n-2} \int(G(\varepsilon y) \mathrm{d} \varphi(y), \mathrm{d} \varphi(y)) u(\varepsilon y) \mathrm{d} y
\end{aligned}
$$

If we divide this by $\varepsilon^{n-2}$ and take the limit as $\varepsilon \rightarrow 0$ we find that

$$
\int(G(0) \mathrm{d} \varphi(y), \mathrm{d} \varphi(y)) u(0) \mathrm{d} y=\int\left(G_{1}(0) \mathrm{d} \varphi(y), \mathrm{d} \varphi(y)\right) u_{1}(0) \mathrm{d} y
$$

for all $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. We may choose an orthonormal basis of $\left(\mathbb{R}^{n}\right)^{*}$ such that $G_{1}(0)$ is diagonal with entries $\xi_{1}, \ldots, \xi_{n}$. Set $v_{i}=u_{1}(0) \xi_{i} / u(0)$. Set

$$
D=\sum \frac{\partial^{2}}{\partial x_{i}^{2}} \quad \text { and } \quad D_{v}=\sum v_{i} \frac{\partial^{2}}{\partial x_{i}^{2}}
$$

Then one has

$$
\int \varphi D \varphi \mathrm{~d} x=\int \varphi D_{v} \varphi \mathrm{~d} x
$$

for all $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. This implies that $D=D_{v}$. Hence $v_{i}=1$ for all $i$. This is the content of the lemma.

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