# Decay Preserving Operators and Stability of the Essential Spectrum 

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#### Abstract

We establish criteria for the stability of the essential spectrum for unbounded operators acting in Banach modules. The applications cover operators acting on sections of vector fiber bundles over non-smooth manifolds or locally compact abelian groups, in particular differential operators of any order with complex measurable coefficients on $\mathbb{R}^{n}$, singular Dirac operators, and Laplace-Beltrami operators on Riemannian manifolds with measurable metrics.

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## 1 Introduction

The main purpose of this paper is to establish criteria which ensure that the difference of the resolvents of two operators is compact. In order to simplify later statements, we use the following defi nition (our notations are quite standard; we recall however the most important ones at the end of this section).

Definition 1.1 Let $A$ and $B$ be two closed operators acting in a Banach space $\mathscr{H}$. We say that $B$ is a compact perturbation of $A$ if there is $z \in \rho(A) \cap \rho(B)$ such that $(A-z)^{-1}-(B-z)^{-1}$ is a compact operator.

Under the conditions of this defi nition the difference $(A-z)^{-1}-(B-z)^{-1}$ is a compact operator for all $z \in \rho(A) \cap \rho(B)$. In particular, if $B$ is a compact perturbation of $A$, then $A$ and $B$ have the same essential spectrum, and this for any reasonable defi nition of the essential spectrum, see [GW]. To be precise, in this paper we defi ne the essential spectrum of $A$ as the set of points $\lambda \in \mathbb{C}$ such that $A-\lambda$ is not Fredholm.

We shall describe now a standard and simple, although quite powerful, method of proving that $B$ is a compact perturbation of $A$. Note that we are interested in situations where $A$ and $B$ are differential (or pseudo-differential) operators with complex measurable coeffi cients which differ little on a neighborhood of infi nity. An important point in such situations is that one has not much information about the domains of the operators. However, one often knows explicitly a generalized version of the "quadratic form domain" of the operator. Since we want to consider operators of any order (in particular Dirac operators) we shall work in the following framework, which goes beyond the theory of accretive forms.

Let $\mathscr{G}, \mathscr{H}, \mathscr{K}$ be reflexive Banach spaces such that $\mathscr{G} \subset \mathscr{H} \subset \mathscr{K}$ continuously and densely. We are interested in operators in $\mathscr{H}$ constructed according to the following procedure: let $A_{0}, B_{0}$ be continuous bijective maps $\mathscr{G} \rightarrow \mathscr{K}$ and let $A, B$ be their restrictions to $A_{0}^{-1} \mathscr{H}$ and $B_{0}^{-1} \mathscr{H}$. These are closed densely defi ned operators in $\mathscr{H}$ and $z=0 \in \rho(A) \cap \rho(B)$. Then in $\mathcal{B}(\mathscr{K}, \mathscr{G})$ we have

$$
\begin{equation*}
A_{0}^{-1}-B_{0}^{-1}=A_{0}^{-1}\left(B_{0}-A_{0}\right) B_{0}^{-1} \tag{1.1}
\end{equation*}
$$

In particular, we get in $\mathcal{B}(\mathscr{H})$

$$
\begin{equation*}
A^{-1}-B^{-1}=A_{0}^{-1}\left(B_{0}-A_{0}\right) B^{-1} \tag{1.2}
\end{equation*}
$$

We get the simplest compactness criterion: if $A_{0}-B_{0}: \mathscr{G} \rightarrow \mathscr{K}$ is compact, then $B$ is a compact perturbation of $A$. But in this case we have more: the operator $A_{0}^{-1}-$
$B_{0}^{-1}: \mathscr{K} \rightarrow \mathscr{G}$ is also compact, and this can not happen if $A_{0}, B_{0}$ are differential operators with distinct principal part (cf. below). This also excludes singular lower order perturbations, e.g. Coulomb potentials in the Dirac case.

The advantage of the preceding criterion is that no knowledge of the domains $\mathcal{D}(A), \mathcal{D}(B)$ is needed. To avoid the mentioned disadvantages, one may assume that one of the operators is more regular than the second one, so that the functions in its domain are, at least locally, slightly better than those from $\mathscr{G}$. Note that $\mathcal{D}(B)$ when equipped with the graph topology is such that $\mathcal{D}(B) \subset \mathscr{G}$ continuously and densely and we get a second compactness criterion by asking that $A_{0}-B_{0}: \mathcal{D}(B) \rightarrow \mathscr{K}$ be compact. This time again we get more than needed, because not only $B$ is a compact perturbation of $A$, but also $A_{0}^{-1}-B_{0}^{-1}: \mathscr{H} \rightarrow \mathscr{G}$ is compact. However, perturbations of the principal part of a differential operator are allowed and also much more singular perturbations of the lower order terms, cf. [N1] for the Dirac case.

In this paper we are interested in situations where we have really no information concerning the domains of $A$ and $B$ (besides the fact that they are subspaces of $\mathscr{G}$ ). The case when $A, B$ are second order elliptic operators with measurable complex coeffi cients acting in $\mathscr{H}=L^{2}\left(\mathbb{R}^{n}\right)$ has been studied by Ouhabaz and Stollmann in [OS] and, as far as we know, this is the only paper where the "unperturbed" operator is not smooth. Their approach consists in proving that the difference $A^{-k}-B^{-k}$ is compact for some $k \geq 2$ (which implies the compactness of $A^{-1}-B^{-1}$ ). In order to prove this, they take advantage of the fact that $\mathcal{D}\left(A^{k}\right)$ is a subset of the Sobolev space $W^{1, p}$ for some $p>2$, which means that we have a certain gain of local regularity. Of course, $L^{p}$ techniques from the theory of partial differential equations are required for their methods to work.

We shall explain now in the most elementary situation the main ideas of our approach to these questions. Let $\mathscr{H}=L^{2}(\mathbb{R})$ and $P=-i \frac{d}{d x}$. We consider operators of the form $A_{0}=P a P+V$ and $B_{0}=P b P+W$ where $a, b$ are bounded operators on $\mathscr{H}$ such that $\operatorname{Re} a$ and $\operatorname{Re} b$ are bounded below by strictly positive numbers. $V$ and $W$ are assumed to be continuous operators $\mathscr{H}^{1} \rightarrow \mathscr{H}^{-1}$, where $\mathscr{H}^{s}$ are Sobolev spaces associated to $\mathscr{H}$. Then $A_{0}, B_{0} \in \mathcal{B}\left(\mathscr{H}^{1}, \mathscr{H}^{-1}\right)$ and we put some conditions on $V, W$ which ensure that $A_{0}, B_{0}$ are invertible (e.g. we could include the constant $z$ in them). Thus we are in the preceding abstract framework with $\mathscr{G}=\mathscr{H}^{1}$ and $\mathscr{K}=\mathscr{H}^{-1} \equiv \mathscr{G}^{*}$. Then from (1.2) we get

$$
\begin{equation*}
A^{-1}-B^{-1}=A_{0}^{-1} P(b-a) P B^{-1}+A_{0}^{-1}(W-V) B^{-1} \tag{1.3}
\end{equation*}
$$

Let $R$ be the first term on the right hand side and let us see how we could prove that it is a compact operator on $\mathscr{H}$. Note that the second term should be easier to treat since we expect $V$ and $W$ to be operators of order less than 2 .

We have $R \mathscr{H} \subset \mathscr{H}^{1}$, so we can write $R=\psi(P) R_{1}$ for some $\psi \in B_{0}(\mathbb{R})$ (bounded Borel function which tends to zero at infi nity) and $R_{1} \in \mathcal{B}(\mathscr{H})$. This is just half of the conditions needed for compactness, in fact $R$ will be compact if and only if one can also fi nd $\varphi \in B_{0}(\mathbb{R})$ and $R_{2} \in \mathcal{B}(\mathscr{H})$ such that $R=\varphi(Q) R_{2}$, where $\varphi(Q)$ is the operator of multiplication by $\varphi$. Of course, the only factor which can help to get such a decay is $b-a$. So let us suppose that we can write $b-a=\xi(Q) U$ for some $\xi \in B_{0}(\mathbb{R})$ and a bounded operator $U$ on $\mathscr{H}$. We denote $S=A_{0}^{-1} P$ and note that this
is a bounded operator on $\mathscr{H}$, because $P: \mathscr{H} \rightarrow \mathscr{H}^{-1}$ and $A_{0}^{-1}: \mathscr{H}^{-1} \rightarrow \mathscr{H}^{1}$ are bounded. Then $R=S \xi(Q) U P B^{-1}$ and $U P B^{-1} \in \mathcal{B}(\mathscr{H})$, hence $R$ will be compact if the operator $S \in \mathcal{B}(\mathscr{H})$ has the following property: for each $\xi \in B_{0}(\mathbb{R})$ there are $\varphi \in B_{0}(\mathbb{R})$ and $T \in \mathcal{B}(\mathscr{H})$ such that $S \xi(Q)=\varphi(Q) T$.

An operator $S$ with the property specifi ed above will be called decay preserving. Thus we see that the compactness of $R$ follows from the fact that $S$ preserves decay and our main point is that it is easy to check this property under very general assumptions on $A$, cf. Corollary 2.21 and Proposition 4.10 for abstract criteria, Lemmas 4.14, 8.2 and 7.1 and Theorems 6.5 and 6.6 for more concrete examples. Note that the perturbative technique described in Proposition 3.7 shows that in many cases it suffi ces to prove the decay preserving property only for operators with smooth coeffi cients (cf. Lemma 7.1).

An abstract formulation of the ideas described above (see Proposition 2.16) allows one to treat situations of a very general nature, like pseudo-differential operators on fi nite dimensional vector spaces over a local (for example $p$-adic) fi eld, in particular differential operators of arbitrary order with irregular coeffi cients on $\mathbb{R}^{n}$, the Laplace operator on manifolds with locally $L^{\infty}$ Riemannian metrics, and operators acting on sections of vector bundles over locally compact spaces. Sections $4,5,7$ and 8 are devoted to such applications. We stress once again that, in the applications to differential operators, we are interested only in situations where the coeffi cients are not smooth and the lower order terms are singular.

Plan of the paper: In Section 2 we introduce an algebraic formalism which allows us to treat in a unifi ed and simple way operators which have an algebraically complicated structure, e.g. operators acting on sections of vector fi ber bundles over a locally compact space. The class of decay improving (or vanishing at infi nity) operators is defi ned through an a priori given algebra of operators on a Banach space $\mathscr{H}$, that we call multiplier algebra of $\mathscr{H}$, and this allows us to defi ne the notion of decay preserving operator in a natural and general context, that of Banach modules. Several examples of multiplier algebras are given Subsections 2.1, 4 and 6 . We stress that Section 2 is only an accumulation of defi nitions and straightforward consequences.

We mention that this algebraic framework allows one to study differential operators in $L^{p}$ or more general Banach spaces. Since these extensions are rather obvious and the examples are not particularly interesting, we shall not consider explicitly such situations.

Section 3 contains several abstract compactness criteria which formalize in the context of Banach modules the ideas involved in the example discussed above.

In Subsection 4.4 we give our fi rst concrete applications of the abstract theory: we consider "hypoelliptic" operators on abelian groups and treat as an example the Dirac operators on $\mathbb{R}^{n}$. In Section 5 we discuss operators in divergence form on $\mathbb{R}^{n}$, hence of order $2 m$ with $m \geq 1$ integer, with coeffi cients of a rather general form (they do not have to be functions, for example).

In Section 7 we present several results concerning the case when the coeffi cients of the operator $A-B$ vanish at infi nity only in some weak sense. This question has been studied before, for example in [He, LV, OS, We]. We present the notion of weakly

[^0]vanishing at infinity functions in terms of filters fi ner than the Fr'echet filter, a natural idea in our context being to extend the standard notion of neighborhood of infi nity.

If $X$ is a locally compact space, it is usual to defi ne the filter of neighborhoods of infi nity as the family of subsets of $X$ with relatively compact complement; we shall call this the Fréchet filter. If $\mathcal{F}$ is a filter on $X$ fi ner than the Fr'echet filter then a function $\varphi: X \rightarrow \mathbb{C}$ such that $\lim _{\mathcal{F}} \varphi=0$ can naturally be thought as convergent to zero at infi nity in a generalized sense (recall that $\lim _{\mathcal{F}} \varphi=0$ means that for each $\varepsilon>0$ the set of points $x$ such that $|\varphi(x)|<\varepsilon$ belongs to $\mathcal{F}$ ). In Subsection 6 we consider three such fi lters and describe corresponding classes of decay preserving operators in Theorem 6.1, Proposition 6.7 and Theorem 6.5.

Theorem 6.5 is a consequence of a factorization theorem that we prove in Section 9 and which involves interesting tools from the modern theory of Banach spaces. In fact, Theorem 9.7, the main result of Section 9, is a version of the "strong factorization theorem" of B. Maurey (see Theorem 9.1) which does not seem to be covered by the results existing in the literature. We also use Maurey's theorem directly to prove some of our main results, for example Theorems 8.7 and 8.8 which depend on Theorem 6.1.

Theorem 7.4 is one of the main applications of our formalism: we prove a compactness result for operators of order $2 m$ in divergence form assuming that the difference between their coeffi cients vanishes at infi nity in a weak sense. Such results were known before only in the case $m=1$, see especially Theorem 2.1 in [OS]. We assume that the coeffi cients of the higher order terms are bounded, thus their Theorem 3.1 is not covered unless we add an implicit assumption, as is done in [OS] (or in our Theorems 8.7 and 8.8). In fact, our main abstract compactness result Theorem 3.2 is stated such as to apply to situations when the coeffi cients of the principal part of the operators are locally unbounded, as in [Ba1, Ba2], but we have not developed this idea here.

Perturbations of the Laplace operator on a Riemannian manifold with locally $L^{\infty}$ metric are considered in Section 8. We introduce and study an abstract model of this situation which fi ts very naturally in our algebraic framework and covers the case of Lipschitz manifolds with measurable metrics. We consider in more detail the case when the manifold is $C^{1}$ (but the metric is only locally $L^{\infty}$ ) and establish stability of the essential spectrum under certain perturbations of the metric, see Theorems 8.4, 8.7 and 8.8. We also consider, in an abstract setting and without going into technical details, the Laplace operator acting on differential forms.

In an Appendix we collect some general facts concerning operators acting in scales of spaces which are often used without comment in the rest of the paper.

Notations: If $\mathscr{G}$ and $\mathscr{H}$ are Banach spaces then $\mathcal{B}(\mathscr{G}, \mathscr{H})$ is the space of bounded linear operators $\mathscr{G} \rightarrow \mathscr{H}$, the subspace of compact operators is denoted $\mathcal{K}(\mathscr{G}, \mathscr{H})$, and we set $\mathcal{B}(\mathscr{H})=\mathcal{B}(\mathscr{H}, \mathscr{H})$ and $\mathcal{K}(\mathscr{H})=\mathcal{K}(\mathscr{H}, \mathscr{H})$. The domain and the resolvent set of an operator $S$ will be denoted by $\mathcal{D}(S)$ and $\rho(S)$ respectively. The norm of a Banach space $\mathscr{G}$ is denoted by $\|\cdot\|_{\mathscr{G}}$ and we omit the index if the space plays a central rôle. The adjoint space (space of antilinear continuous forms) of a Banach space $\mathscr{G}$ is denoted $\mathscr{G}^{*}$ and if $u \in \mathscr{G}$ and $v \in \mathscr{G}^{*}$ then we set $v(u)=\langle u, v\rangle$. The embedding $\mathscr{G} \subset \mathscr{G}^{* *}$ is realized by defi ning $\langle v, u\rangle=\overline{\langle u, v\rangle}$.

If $\mathscr{G}, \mathscr{H}, \mathscr{K}$ are Banach spaces such that $\mathscr{G} \subset \mathscr{H}$ continuously and densely and $\mathscr{H} \subset \mathscr{K}$ continuously then we we have a natural continuous embedding $\mathcal{B}(\mathscr{H}) \hookrightarrow$
$\mathcal{B}(\mathscr{G}, \mathscr{K})$ that will be used without comment later on.
A Friedrichs couple $(\mathscr{G}, \mathscr{H})$ is a pair of Hilbert spaces $\mathscr{G}, \mathscr{H}$ together with a continuous dense embedding $\mathscr{G} \subset \mathscr{H}$. The Gelfand triplet associated to it is obtained by identifying $\mathscr{H}=\mathscr{H}^{*}$ with the help of the Riesz isomorphism and then taking the adjoint of the inclusion map $\mathscr{G} \rightarrow \mathscr{H}$. Thus we get $\mathscr{G} \subset \mathscr{H} \subset \mathscr{G}^{*}$ with continuous and dense embeddings. Now if $u \in \mathscr{G}$ and $v \in \mathscr{H} \subset \mathscr{G}^{*}$ then $\langle u, v\rangle$ is the scalar product in $\mathscr{H}$ of $u$ and $v$ and also the action of the functional $v$ on $u$. As noted above, we have $\mathcal{B}(\mathscr{H}) \subset \mathcal{B}\left(\mathscr{G}, \mathscr{G}^{*}\right)$.

If $X$ is a locally compact topological space then $B(X)$ is the $C^{*}$-algebra of bounded Borel complex functions on $X$, with norm $\sup _{x \in X}|\varphi(x)|$, and $B_{0}(X)$ is the subalgebra consisting of functions which tend to zero at infi nity. Then $C(X), C_{\mathrm{b}}(X), C_{0}(X)$ and $C_{\mathrm{c}}(X)$ are the spaces of complex functions on $X$ which are continuous, continuous and bounded, continuous and convergent to zero at infi nity, and continuous with compact support respectively. We denote $\chi_{S}$ the characteristic function of a set $S \subset X$.

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## 2 Banach modules and decay preserving operators

### 2.1 Banach modules

We use the terminology of [FD] but with some abbreviations, e.g. a morphism is a linear multiplicative map between two algebras, and a *-morphism is a morphism between two $*$-algebras which commutes with the involutions. We recall that an approximate unit in a Banach algebra $\mathcal{M}$ is a net $\left\{J_{\alpha}\right\}$ in $\mathcal{M}$ such that $\left\|J_{\alpha}\right\| \leq C$ for some constant $C$ and all $\alpha$ and $\lim _{\alpha}\left\|J_{\alpha} M-M\right\|=\lim _{\alpha}\left\|M J_{\alpha}-M\right\|=0$ for all $M \in \mathcal{M}$. An approximate unit exists if and only if there is a number $C$ such that for each $\varepsilon>0$ and for each fi nite set $\mathcal{F} \subset \mathcal{M}$ there is $J \in \mathcal{M}$ with $\|J\| \leq C$ and $\|J M-M\| \leq$ $\varepsilon,\|M J-M\| \leq \varepsilon$ for all $M \in \mathcal{F}$. It is well known that any $C^{*}$-algebra has an approximate unit. If $\mathscr{H}$ is a Banach space, we shall say that a Banach subalgebra $\mathcal{M}$ of $\mathcal{B}(\mathscr{H})$ is non-degenerate if the linear subspace of $\mathscr{H}$ generated by the elements $M u$, with $M \in \mathcal{M}$ and $u \in \mathscr{H}$, is dense in $\mathscr{H}$.

In view of its importance in our paper, we state below the Cohen-Hewitt factorization theorem [FD, Ch. V-9.2].

Theorem 2.1 Let $\mathcal{C}$ be a Banach algebra with an approximate unit, let $\mathscr{E}$ be a Banach space, and let $Q: \mathcal{C} \rightarrow \mathcal{B}(\mathscr{E})$ be a continuous morphism. Denote $\mathscr{E}_{0}$ the closed linear subspace of $\mathscr{E}$ generated by the elements of the form $Q(\varphi) v$ with $\varphi \in \mathcal{C}$ and $v \in \mathscr{E}$. Then for each $u \in \mathscr{E}_{0}$ there are $\varphi \in \mathcal{C}$ and $v \in \mathscr{E}$ such that $u=Q(\varphi) v$.

Now we introduce the framework in which we shall work.

Definition 2.2 A Banach module is a couple $(\mathscr{H}, \mathcal{M})$ consisting of a Banach space $\mathscr{H}$ and a non-degenerate Banach subalgebra $\mathcal{M}$ of $\mathcal{B}(\mathscr{H})$ which has an approximate unit. If $\mathscr{H}$ is a Hilbert space and $\mathcal{M}$ is a $C^{*}$-algebra of operators on $\mathscr{H}$, we say that $\mathscr{H}$ is a Hilbert module.

We shall adopt the usual abus de language and say that $\mathscr{H}$ is a Banach module (over $\mathcal{M})$. The distinguished subalgebra $\mathcal{M}$ will be called multiplier algebra of $\mathscr{H}$ and, when required by the clarity of the presentation, we shall denote it $\mathcal{M}(\mathscr{H})$. We are only interested in the case when $\mathcal{M}$ does not have a unit: the operators from $\mathcal{M}$ are the prototype of decay improving (or vanishing at infi nity) operators, and the identity operator cannot have such a property. Note that it is implicit in Defi nition 2.2 that if $\mathscr{H}$ is a Hilbert module then its adjoint space $\mathscr{H}^{*}$ is identified with $\mathscr{H}$ with the help of the Riesz isomorphism.

If $\left\{J_{\alpha}\right\}$ is an approximate unit of $\mathcal{M}$, then the density in $\mathscr{H}$ of the linear subspace generated by the elements $M u$ is equivalent to

$$
\begin{equation*}
\lim _{\alpha}\left\|J_{\alpha} u-u\right\|=0 \quad \text { for all } \quad u \in \mathscr{H} \tag{2.4}
\end{equation*}
$$

But much more is true:

$$
\begin{equation*}
u \in \mathscr{H} \Rightarrow u=M v \text { for some } M \in \mathcal{M} \text { and } v \in \mathscr{H} . \tag{2.5}
\end{equation*}
$$

This follows from the Cohen-Hewitt theorem, see Theorem 2.1. By using (2.4) we could avoid any reference to this result in our later arguments; this would make them more elementary but less simple. From Theorem 2.1 we also get:

Lemma 2.3 Assume that $\mathcal{A}$ is a Banach algebra with approximate unit and that a morphism $\Phi: \mathcal{A} \rightarrow \mathcal{M}(\mathscr{H})$ with dense image is given. Then each $u \in \mathscr{H}$ can be written as $u=A v$ where $A \in \Phi(\mathcal{A})$ and $v \in \mathscr{H}$.

Example 2.4 The simplest example of Banach module is the following. Let $X$ be a locally compact non-compact topological space and let $\mathscr{H}$ be a Banach space. We say that $\mathscr{H}$ is a Banach $X$-module if a continuous morphism $Q: C_{0}(X) \rightarrow \mathcal{B}(\mathscr{H})$ has been given such that the linear subspace generated by the vectors of the form $Q(\varphi) u$, with $\varphi \in C_{0}(X)$ and $u \in \mathscr{H}$, is dense in $\mathscr{H}$. If $\mathscr{H}$ is a Hilbert space and $Q$ is a $*$-morphism, we say that $\mathscr{H}$ is a Hilbert $X$-module. We shall use the notation $\varphi(Q) \equiv Q(\varphi)$. The Banach module structure on $\mathscr{H}$ is defi ned by the closure $\mathcal{M}$ in $\mathcal{B}(\mathscr{H})$ of the set of operators of the form $\varphi(Q)$ with $\varphi \in C_{0}(X)$. In the case of a Hilbert $X$-module the closure is not needed and we get a Hilbert module structure (because a $*$-morphism between two $C^{*}$-algebras is continuous and its range is a $C^{*}$ algebra). Banach $X$-modules appear naturally in differential geometry as spaces of sections of vector fi ber bundles over a manifold $X$, and this is the point of interest for us.

Remark 2.5 In the case of a Banach $X$-module, Lemma 2.3 gives: each $u \in \mathscr{H}$ can be written as $u=\psi(Q) v$ with $\psi \in C_{0}(X)$ and $v \in \mathscr{H}$. In particular, we deduce that the morphism $Q$ has an extension, also denoted $Q$, to a unital continuous
morphism of $C_{\mathrm{b}}(X)$ into $\mathcal{B}(\mathscr{H})$ which is uniquely determined by the following strong continuity property: if $\left\{\varphi_{n}\right\}$ is a bounded sequence in $C_{\mathrm{b}}(X)$ such that $\varphi_{n} \rightarrow \varphi$ locally uniformly, then $\varphi_{n}(Q) \rightarrow \varphi(Q)$ strongly on $\mathscr{H}$. Indeed, we can defi ne $\varphi(Q) u=$ $(\varphi \psi)(Q) v$ for each $\varphi \in C_{\mathrm{b}}(X)$; then if $e_{\alpha}$ is an approximate unit for $C_{0}(X)$ with $\left\|e_{\alpha}\right\| \leq 1$ we get $\varphi(Q) u=\lim \left(\varphi e_{\alpha}\right)(Q) u$ hence the defi nition is independent of the factorization of $u$ and $\|\varphi(Q)\| \leq\|Q\| \sup |\varphi|$.

Remark 2.6 If $\mathscr{H}$ is a Hilbert $X$-module one can extend the morphism even further: $Q$ canonically extends to $a *$-morphism $\varphi \mapsto \varphi(Q)$ of $B(X)$ into $\mathcal{B}(\mathscr{H})$ such that ${ }^{2}$ : if $\left\{\varphi_{n}\right\}$ is a bounded sequence in $B(X)$ and $\lim _{n \rightarrow \infty} \varphi_{n}(x)=\varphi(x)$ for all $x \in X$, then $\mathrm{s}-\lim _{n} \varphi_{n}(Q)=\varphi(Q)$. This follows from standard integration theory see [Be, Lo]. In particular, a separable Hilbert $X$-module is essentially a direct integral of Hilbert spaces over $X$, see [Di, II.6.2], but we shall not need this fact.

The class of $X$-modules is more general than it appears at first sight. Indeed, if $\mathcal{C}$ is an abelian $C^{*}$-algebra then one has a canonical identifi cation $\mathcal{C} \equiv C_{0}(\mathcal{X})$ where $\mathcal{X}$ is the spectrum of $\mathcal{C}$. However, the space $\mathcal{X}$ is in general rather complicated so it is not really useful to take it into account. In particular, this happens in the following class of examples of interest in applications (see Section 7).

Example 2.7 Let $X$ be a set and $\mathcal{F}$ a fi lter on $X$. Let us choose a $C^{*}$-algebra $\mathcal{C}$ of bounded complex functions on $X$ (with the sup norm) and then let $\mathcal{C}_{0}$ be the set of $\varphi \in \mathcal{C}$ such that $\lim _{\mathcal{F}} \varphi=0$. Then $\mathcal{C}_{0}$ is a $C^{*}$-algebra and its spectrum $\mathcal{X}$ contains $X$ but is much larger than $X$ in general.

Let us say that a Banach module structure defi ned by a Banach algebra $\mathcal{N}$ on $\mathscr{H}$ is finer than that defi ned by $\mathcal{M}$ if $\mathcal{M} \subset \mathcal{N}$. In the next example we show that, by using the same idea as in Example 2.7, one can defi ne on each $X$-module new Banach module structures fi ner than the initial one. In Section 7 we shall consider the question of the stability of the essential spectrum in situations of this type, when the perturbation vanishes at infi nity in a weak sense.

Example 2.8 Let $\mathscr{H}$ be a Hilbert $X$-module over a locally compact non-compact topological space $X$ and let $\mathcal{F}$ be a filter on $X$ fi ner than the $\mathrm{Fr}^{\prime}$ echet filter. We extend the morphism $Q$ to all of $B(X)$ as explained in Remark 2.6 and observe that we get a fi ner Hilbert module structure on $\mathscr{H}$ by taking $\left\{\varphi(Q) \mid \lim _{\mathcal{F}} \varphi=0\right\}$ as multiplier algebra. One can proceed similarly in the case of a Banach $X$-module, it suffi ces to replace $B(X)$ by $C_{\mathrm{b}}(X)$.

We now give an example of a non-topological nature.
Example 2.9 Let $(X, \mu)$ be a measure space with $\mu(X)=\infty$. We defi ne the class of functions which "vanish at infi nity" as follows. Let us say that a set $F \subset X$ is of cofinite measure if its complement $F^{\mathrm{c}}$ is of fi nite (exterior) measure. The family of

[^1]sets of cofi nite measure is clearly a fi lter $\mathcal{F}_{\mathcal{\mu}}$. If $\varphi$ is a function on $X$ then $\lim _{\mathcal{F}_{\mu}} \varphi=0$ means that for each $\varepsilon>0$ the set where $|\varphi(x)| \geq \varepsilon$ is of fi nite measure. We denote $B_{\mu}(X)$ the $C^{*}$-subalgebra of $L^{\infty}(X)$ consisting of functions such that $\lim _{\mathcal{F}_{\mu}} \varphi=0$. Let $\mathcal{N}_{\mu}$ be the set of (equivalence classes of) Borel subsets of fi nite measure of $X$. Then $\left\{\chi_{N}\right\}_{N \in \mathcal{N}_{\mu}}$ is an approximate unit of $B_{\mu}(X)$ because for each $\varphi \in B_{\mu}(X)$ and each $\varepsilon>0$ we have $N=\{x| | \varphi(x) \mid \geq \varepsilon\} \in \mathcal{N}_{\mu}$ and ess-sup $\left|\varphi-\chi_{N} \varphi\right| \leq \varepsilon$. Now it is clear that $L^{2}(X)$ and, more generally, any direct integral of Hilbert spaces over $X$, has a natural Hilbert module structure with $B_{\mu}(X)$ as multiplier algebra.

If $\mathscr{H}$ is a Banach module and the Banach space $\mathscr{H}$ is reflexive we say that $\mathscr{H}$ is a reflexive Banach module. In this case the adjoint Banach space $\mathscr{H}^{*}$ is equipped with a canonical Banach module structure, its multiplier algebra being $\mathcal{M}\left(\mathscr{H}^{*}\right):=$ $\left\{A^{*} \mid A \in \mathcal{M}(\mathscr{H})\right\}$. This is a closed subalgebra of $\mathcal{B}\left(\mathscr{H}^{*}\right)$ which clearly has an approximate unit and the linear subspace generated by the elements of the form $A^{*} v$, with $A \in \mathcal{M}(\mathscr{H})$ and $v \in \mathscr{H}^{*}$, is weak ${ }^{*}$-dense, hence dense, in $\mathscr{H}^{*}$. Indeed, if $u \in \mathscr{H}$ and $\left\langle u, A^{*} v\right\rangle=0$ for all such $A, v$ then $A u=0$ for all $A \in \mathcal{M}(\mathscr{H})$ hence $u=0$ because of (2.4).

Example 2.10 For each real number $s$ let $\mathscr{H}^{s}:=\mathscr{H}^{s}\left(\mathbb{R}^{n}\right)$ be the Hilbert space of distributions $u$ on $\mathbb{R}^{n}$ such that $\|u\|_{s}^{2}:=\int\left(1+|k|^{2}\right)^{s}|\widehat{u}(k)|^{2} d k<\infty$, where $\widehat{u}$ is the Fourier transform of $u$. This is the usual Sobolev space of order $s$ on $\mathbb{R}^{n}$. The algebra $\mathscr{S}$ of Schwartz test functions on $\mathbb{R}^{n}$ is naturally embedded in $\mathcal{B}\left(\mathscr{H}^{s}\right)$, a function $\varphi \in \mathscr{S}$ being identifi ed with the operator of multiplication by $\varphi$ on $\mathscr{H}^{s}$. If we denote by $\mathcal{M}^{s}$ the closure of $\mathscr{S}$ in $\mathcal{B}\left(\mathscr{H}^{s}\right)$, then clearly $\left(\mathscr{H}^{s}, \mathcal{M}^{s}\right)$ is a Banach module and this Banach module is a Hilbert module if and only if $s=0$. The module adjoint to $\left(\mathscr{H}^{s}, \mathcal{M}^{s}\right)$ is identifi ed with $\left(\mathscr{H}^{-s}, \mathcal{M}^{-s}\right)$. Note that $\mathcal{M}^{s}$ can be realized as a subalgebra of $\mathcal{M}^{0}=C_{0}\left(\mathbb{R}^{n}\right)$, namely $\mathcal{M}^{s}$ is the completion of $\mathscr{S}$ for the norm $\|\varphi\|_{\mathcal{M}^{s}}:=\sup _{\|u\|_{s}=1}\|\varphi u\|_{s}$, and then we have $\mathcal{M}^{s}=\mathcal{M}^{-s}$ isometrically and $\mathcal{M}^{s} \subset \mathcal{M}^{t}$ if $s \geq t \geq 0$ (by interpolation).

Definition 2.11 A couple $(\mathscr{G}, \mathscr{H})$ consisting of a Hilbert module $\mathscr{H}$ and a Hilbert space $\mathscr{G}$ such that $\mathscr{G} \subset \mathscr{H}$ continuously and densely will be called a Friedrichs module. If $\mathscr{H}$ is a Hilbert $X$-module over a locally compact space $X$, we say that $(\mathscr{G}, \mathscr{H})$ is a Friedrichs $X$-module. If $\mathcal{M}(\mathscr{H}) \subset \mathcal{K}(\mathscr{G}, \mathscr{H})$, we say that $(\mathscr{G}, \mathscr{H})$ is $a$ compact Friedrichs module.

In the situation of this defi nition we always identify $\mathscr{H}$ with its adjoint space, which gives us a Gelfand triplet $\mathscr{G} \subset \mathscr{H} \subset \mathscr{G}^{*}$. If $(\mathscr{G}, \mathscr{H})$ is a compact Friedrichs module then each operator $M$ from $\mathcal{M}(\mathscr{H})$ extends to a compact operator $M: \mathscr{H} \rightarrow \mathscr{G}^{*}$ (this is the adjoint of the compact operator $M^{*}: \mathscr{G} \rightarrow \mathscr{H}$ ). Thus we shall have $\mathcal{M}(\mathscr{H}) \subset \mathcal{K}(\mathscr{G}, \mathscr{H}) \cap \mathcal{K}\left(\mathscr{H}, \mathscr{G}^{*}\right)$.

Example 2.12 With the notations of Example 2.10, if we set $\mathscr{H}=\mathscr{H}^{0}$ and take $s>$ 0 , then $\left(\mathscr{H}^{s}, \mathscr{H}\right)$ is a compact Friedrichs module and the associated Gelfand triplet is $\mathscr{H}^{s} \subset \mathscr{H} \subset \mathscr{H}^{-s}$. Indeed, if $\varphi \in C_{0}\left(\mathbb{R}^{n}\right)$ then the operator of multiplication by $\varphi$ is a compact operator $\mathscr{H}^{s} \rightarrow \mathscr{H}$.

### 2.2 Decay improving operators

Let $\mathscr{H}$ and $\mathscr{K}$ be Banach spaces. If $\mathscr{K}$ is a Banach module then we shall denote by $\mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ the norm closed linear subspace generated by the operators $M T$, with $T \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and $M \in \mathcal{M}(\mathscr{K})$. We say that an operator in $\mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ is decay improving, or left vanishes at infinity (with respect to $\mathcal{M}(\mathscr{K})$, if this is not obvious from the context). If $J_{\alpha}$ is an approximate unit for $\mathcal{M}(\mathscr{K})$, then for an operator $S \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ we have:

$$
\begin{align*}
S \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K}) & \Leftrightarrow \lim _{\alpha}\left\|J_{\alpha} S-S\right\|=0  \tag{2.6}\\
& \Leftrightarrow \quad S=M T \text { for some } M \in \mathcal{M}(\mathscr{K}) \text { and } T \in \mathcal{B}(\mathscr{H}, \mathscr{K}) .
\end{align*}
$$

The second equivalence follows from the Cohen-Hewitt theorem (Theorem 2.1).
If $\mathscr{H}$ is a Banach module then one can similarly defi ne $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$ as the norm closed linear subspace generated by the operators $T M$ with $T \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and $M \in$ $\mathcal{M}(\mathscr{H})$. We say that the elements of $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$ right vanish at infinity. If both $\mathscr{H}$ and $\mathscr{K}$ are Banach modules then both spaces $\mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ and $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$ make sense and we set $\mathcal{B}_{0}(\mathscr{H}, \mathscr{K})=\mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K}) \cap \mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$.

Some simple properties of these spaces are described below.
Proposition 2.13 If $\mathscr{K}$ is a reflexive Banach module and $S \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ then $S^{*}$ belongs to $\mathcal{B}_{0}^{r}\left(\mathscr{K}^{*}, \mathscr{H}^{*}\right)$.

Proof: We have $S=M T$ with $M \in \mathcal{M}(\mathscr{K})$ and $T \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ by (2.6), which implies $S^{*}=T^{*} M^{*}$ and we have $M^{*} \in \mathcal{M}\left(\mathscr{K}^{*}\right)$ by defi nition.

Proposition 2.14 If $\mathscr{H}$ is a Hilbert module then $\mathcal{B}_{0}(\mathscr{H})$ is a $C^{*}$-algebra and an operator $S \in \mathcal{B}(\mathscr{H})$ belongs to it if and only if one can write $S=M T N$ with $M, N \in \mathcal{M}(\mathscr{H})$ and $T \in \mathcal{B}(\mathscr{H})$.

Proof: $\mathcal{B}_{0}(\mathscr{H})$ is clearly a $C^{*}$-algebra, so if $S \in \mathcal{B}_{0}(\mathscr{H})$ then $S=S_{1} S_{2}$ for some operators $S_{1}, S_{2} \in \mathcal{B}_{0}(\mathscr{H})$. Thus $S_{1}=M T_{1}$ and $S_{2}=T_{2} N$ for some $M, N \in$ $\mathcal{M}(\mathscr{H})$ and $T_{1}, T_{2} \in \mathcal{B}(\mathscr{H})$, hence $S=M T_{1} T_{2} N$.

Proposition 2.15 If $\mathscr{K}$ is a Banach module then $\mathcal{K}(\mathscr{H}, \mathscr{K}) \subset \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$. If $\mathscr{H}$ is a reflexive Banach module, then $\mathcal{K}(\mathscr{H}, \mathscr{K}) \subset \mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$.

Proof: If $\left\{J_{\alpha}\right\}$ is an approximate unit for $\mathcal{M}(\mathscr{K})$ then $s-\lim _{\alpha} J_{\alpha} u=u$ uniformly in $u$ if $u$ belongs to a compact subset of $\mathscr{K}$. Hence if $S \in \mathcal{K}(\mathscr{H}, \mathscr{K})$ then $\lim _{\alpha} \| J_{\alpha} S-$ $S \|=0$ and thus $S \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ by (2.6). To prove the second part of the proposition, note that if $S \in \mathcal{K}(\mathscr{H}, \mathscr{K})$ then $S^{*} \in \mathcal{K}\left(\mathscr{K}^{*}, \mathscr{H}^{*}\right)$, so $S^{*} \in \mathcal{B}_{0}^{l}\left(\mathscr{K}^{*}, \mathscr{H}^{*}\right)$ by what we just proved, hence $S^{* *} \in \mathcal{B}_{0}^{r}\left(\mathscr{H}, \mathscr{K}^{* *}\right)$ by Proposition 2.13. Thus we get $\lim _{\alpha}\left\|S^{* *} J_{\alpha}-S^{* *}\right\|=0$ if $\left\{J_{\alpha}\right\}$ is an approximate unit for $\mathcal{M}(\mathscr{H})$. But clearly $S=S^{* *}$, hence $S \in \mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$.

Proposition 2.16 Let $\mathscr{H}$ be a Banach module and $\mathscr{G}$ a Banach space continuously embedded in $\mathscr{H}$ and such that for each $M \in \mathcal{M}(\mathscr{H})$ the restriction of $M$ to $\mathscr{G}$ is a compact operator $\mathscr{G} \rightarrow \mathscr{H}$. If $R \in \mathcal{B}_{0}^{l}(\mathscr{H})$ and $R \mathscr{H} \subset \mathscr{G}$, then $R \in \mathcal{K}(\mathscr{H})$.

Proof: According to (2.6) we have $R=\lim _{\alpha} J_{\alpha} R$, the limit being taken in norm. But $R \in \mathcal{B}(\mathscr{H}, \mathscr{G})$ by the closed graph theorem and $J_{\alpha} \in \mathcal{K}(\mathscr{G}, \mathscr{H})$ by hypothesis, so that $J_{\alpha} R \in \mathcal{K}(\mathscr{H})$.

### 2.3 Decay preserving operators

Definition 2.17 Let $\mathscr{H}$, $\mathscr{K}$ be Banach modules and let $S \in \mathcal{B}(\mathscr{H}, \mathscr{K})$. We say that $S$ is left decay preserving if for each $M \in \mathcal{M}(\mathscr{H})$ we have $S M \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$. We say that $S$ is right decay preserving if for each $M \in \mathcal{M}(\mathscr{K})$ we have $M S \in$ $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$. If $S$ is left and right decay preserving, we say that $S$ is decay preserving.

We denote $\mathcal{B}_{q}^{l}(\mathscr{H}, \mathscr{K}), \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$ and $\mathcal{B}_{q}(\mathscr{H}, \mathscr{K})$ these classes of operators (the index q comes from quasilocal, a terminology which is sometimes more convenient than "decay preserving"). These are closed subspaces of $\mathcal{B}(\mathscr{H}, \mathscr{G})$. The next result is obvious; a similar assertion holds in the right decay preserving case.

Proposition 2.18 Let $\left\{J_{\alpha}\right\}$ be an approximate unit for $\mathcal{M}(\mathscr{H})$ and let $S$ be an operator in $\mathcal{B}(\mathscr{H}, \mathscr{K})$. Then $S$ is left decay preserving if and only if one of the following conditions is satisfied:
(1) $S J_{\alpha} \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ for all $\alpha$.
(2) for each $M \in \mathcal{M}(\mathscr{H})$ there are $T \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and $N \in \mathcal{M}(\mathscr{K})$ such that $S M=N T$.

The next proposition, which says that the decay preserving property is stable under the usual algebraic operations, is an immediate consequence of Proposition 2.18. There is, of course, a similar statement with "left" and "right" interchanged. We denote by $\mathscr{G}, \mathscr{H}$ and $\mathscr{K}$ Banach modules.

Proposition 2.19 (1) $S \in \mathcal{B}_{q}^{l}(\mathscr{H}, \mathscr{K})$ and $T \in \mathcal{B}_{q}^{l}(\mathscr{G}, \mathscr{H}) \Rightarrow S T \in \mathcal{B}_{q}^{l}(\mathscr{G}, \mathscr{K})$.
(2) If $\mathscr{H}, \mathscr{K}$ are reflexive and $S \in \mathcal{B}_{q}^{l}(\mathscr{H}, \mathscr{K})$, then $S^{*} \in \mathcal{B}_{q}^{r}\left(\mathscr{K}^{*}, \mathscr{H}^{*}\right)$.
(3) If $\mathscr{H}$ is a Hilbert module then $\mathcal{B}_{q}(\mathscr{H})$ is a unital $C^{*}$-subalgebra of $\mathcal{B}(\mathscr{H})$.

Obviously $\mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K}) \subset \mathcal{B}_{q}^{l}(\mathscr{H}, \mathscr{K})$ and $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K}) \subset \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$ but this fact is of no interest. The main results of this paper depend on fi nding other, more interesting examples of decay preserving operators. We shall give in the rest of this subsection some elementary examples of such operators and in Subsections 4 and 6 more subtle ones.

From now on in this subsection $X$ will be a locally compact non-compact topological space. The support supp $u \subset X$ of an element $u$ of a Banach $X$-module $\mathscr{H}$ is defi ned as the smallest closed set such that its complement $U$ has the property $\varphi(Q) u=0$ if $\varphi \in C_{\mathrm{c}}(U)$. Clearly, the set $\mathscr{H}_{\mathrm{c}}$ of elements $u \in \mathscr{H}$ such that supp $u$ is compact is a dense subspace of $\mathscr{H}$.

Let $\mathscr{H}, \mathscr{K}$ be Banach $X$-modules, let $S \in \mathcal{B}(\mathscr{H}, \mathscr{K})$, and let $\varphi, \psi \in C(X)$, not necessarily bounded. We say that $\varphi(Q) S \psi(Q)$ is a bounded operator if there is a constant $C$ such that

$$
\|\xi(Q) \varphi(Q) S \psi(Q) \eta(Q)\| \leq C \sup |\xi| \sup |\eta|
$$

for all $\xi, \eta \in C_{\mathrm{c}}(X)$. The lower bound of the admissible constants $C$ in this estimate is denoted $\|\varphi(Q) S \psi(Q)\|$. If $\mathscr{K}$ is a reflexive Banach $X$-module, then the product $\varphi(Q) S \psi(Q)$ is well defi ned as sesquilinear form on the dense subspace $\mathscr{K}_{\mathrm{c}}^{*} \times \mathscr{H}_{\mathrm{c}}$ of $\mathscr{K}^{*} \times \mathscr{H}$ and the preceding boundedness notion is equivalent to the continuity of this form for the topology induced by $\mathscr{K}^{*} \times \mathscr{H}$. We similarly defi ne the boundedness of the commutator $[S, \varphi(Q)]$.

Proposition 2.20 Assume that $S \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and let $\theta: X \rightarrow[1, \infty[$ be a continuous function such that $\lim _{x \rightarrow \infty} \theta(x)=\infty$. If $\theta(Q) S \theta^{-1}(Q)$ is a bounded operator then $S$ is left decay preserving. If $\theta^{-1}(Q) S \theta(Q)$ is a bounded operator then $S$ is right decay preserving.

Proof: Let $K \subset X$ be compact, let $U \subset X$ be a neighbourhood of infi nity in $X$, and let $\varphi, \psi \in C_{\mathrm{b}}(X)$ such that $\operatorname{supp} \varphi \subset K$, $\operatorname{supp} \psi \subset U$ and $|\varphi| \leq 1,|\psi| \leq 1$. Then $\theta \varphi$ is a bounded function and $\psi \theta^{-1}$ is bounded and can be made as small as we wish by choosing $U$ conveniently. Thus given $\varepsilon>0$ we have

$$
\|\psi(Q) S \varphi(Q)\| \leq\left\|\psi \theta^{-1}\right\| \cdot\left\|\theta(Q) S \theta^{-1}(Q)\right\| \cdot\|\theta \varphi\| \leq \varepsilon
$$

if $U$ is a suffi ciently small neighbourhood of infi nity. Then the result follows from Proposition 2.18(1) and relation (2.6).

The boundedness of $\theta(Q) S \theta^{-1}(Q)$ can be checked by estimating the commutator $[S, \theta(Q)]$; we give an example for the case of metric spaces. Note that on metric spaces one has a natural class of regular functions, namely the Lipschitz functions, for example the functions which give the distance to subsets: $\rho_{K}(x)=\inf _{y \in K} \rho(x, y)$ for $K \subset X$.

We say that a locally compact metric space $(X, \rho)$ is proper if the metric $\rho$ has the property $\lim _{y \rightarrow \infty} \rho(x, y)=\infty$ for some (hence for all) points $x \in X$. Equivalently, if $X$ is not compact but the closed balls are compact.

Corollary 2.21 Let $(X, \rho)$ be a proper locally compact metric space. If $S$ belongs to $\mathcal{B}(\mathscr{H}, \mathscr{K})$ and if $[S, \theta(Q)]$ is bounded for each positive Lipschitz function $\theta$, then $S$ is decay preserving.

Proof: Indeed, by taking $\theta=1+\rho_{K}$ and by using the notations of the proof of Proposition 2.20, we easily get the following estimate: there is $C<\infty$ depending only on $K$ such that

$$
\|\varphi(Q) S \psi(Q)\| \leq C(1+\rho(K, U))^{-1}
$$

where $\rho(K, U)$ is the distance from $K$ to $U$. Since $S^{*}$ has the same properties as $S$, this proves that $S$ is decay preserving. Note that the boundedness of $\left[S, \rho_{x}(Q)\right]$ for some $x \in X$ suffi ces in this argument.

## 3 Compact perturbations in Banach modules

In this section $(\mathscr{G}, \mathscr{H})$ will be a compact Friedrichs module in the sense of Defi nition 2.11. As usual, we associate to it a Gelfand triplet $\mathscr{G} \subset \mathscr{H} \subset \mathscr{G}^{*}$ and we set
$\|\cdot\|=\|\cdot\|_{\mathscr{H}}$. We are interested in criteria which ensure that an operator $B$ is a compact perturbation of an operator $A$, both operators being unbounded operators in $\mathscr{H}$ obtained as restrictions of some bounded operators $\mathscr{G} \rightarrow \mathscr{G}^{*}$. More precisely, the following is a general assumption (suggested by the statement of Theorem 2.1 in [OS]) which will always be fulfi lled:
$(A B)\left\{\begin{array}{l}A, B \text { are closed densely defi ned operators in } \mathscr{H} \text { with } \rho(A) \cap \rho(B) \neq \emptyset \\ \text { and having the following properties: } \mathcal{D}(A) \subset \mathscr{G} \text { densely, } \mathcal{D}\left(A^{*}\right) \subset \mathscr{G}, \\ \mathcal{D}(B) \subset \mathscr{G} \text { and } A, B \text { extend to continuous operators } \widetilde{A}, \widetilde{B} \in \mathcal{B}\left(\mathscr{G}, \mathscr{G}^{*}\right) .\end{array}\right.$
Example 3.1 One can construct interesting classes of operators with the properties required in (AB) as follows. Let $\mathscr{G}_{a}, \mathscr{G}_{b}$ be Hilbert spaces such that $\mathscr{G} \subset \mathscr{G}_{a} \subset \mathscr{H}$ and $\mathscr{G} \subset \mathscr{G}_{b} \subset \mathscr{H}$ continuously and densely. Thus we have two scales

$$
\begin{aligned}
& \mathscr{G} \subset \mathscr{G}_{a} \subset \mathscr{H} \subset \mathscr{G}_{a}^{*} \subset \mathscr{G}^{*}, \\
& \mathscr{G} \subset \mathscr{G}_{b} \subset \mathscr{H} \subset \mathscr{G}_{b}^{*} \subset \mathscr{G}^{*} .
\end{aligned}
$$

Then let $A_{0} \in \mathcal{B}\left(\mathscr{G}_{a}, \mathscr{G}_{a}^{*}\right)$ and $B_{0} \in \mathcal{B}\left(\mathscr{G}_{b}, \mathscr{G}_{b}^{*}\right)$ such that $A_{0}-z: \mathscr{G}_{a} \rightarrow \mathscr{G}_{a}^{*}$ and $B_{0}-z: \mathscr{G}_{b} \rightarrow \mathscr{G}_{b}^{*}$ are bijective for some number $z$. According to Lemma A. 1 we can associate to $A_{0}, B_{0}$ closed densely defi ned operators $A=\widehat{A_{0}}, B=\widehat{B_{0}}$ in $\mathscr{H}$, such that the domains $\mathcal{D}(A)$ and $\mathcal{D}\left(A^{*}\right)$ are dense subspaces of $\mathscr{G}_{a}$ and the domains $\mathcal{D}(B)$ and $\mathcal{D}\left(B^{*}\right)$ are dense subspaces of $\mathscr{G}_{b}$. If we also have $\mathcal{D}(A) \subset \mathscr{G}$ densely, $\mathcal{D}\left(A^{*}\right) \subset \mathscr{G}$ and $\mathcal{D}(B) \subset \mathscr{G}$, then all the conditions of the assumption (AB) are fulfi lled with $\widetilde{A}=A_{0} \mid \mathscr{G}$ and $\widetilde{B}=B_{0} \mid \mathscr{G}$. Such a construction will be used in Corollary 3.4.

The rôle of the assumption ( AB ) is to allow us to give a rigorous meaning to the formal relation, where $z \in \rho(A) \cap \rho(B)$,

$$
\begin{equation*}
(A-z)^{-1}-(B-z)^{-1}=(A-z)^{-1}(B-A)(B-z)^{-1} . \tag{3.7}
\end{equation*}
$$

Recall that $z \in \rho(A)$ if and only if $\bar{z} \in \rho\left(A^{*}\right)$ and then $\left(A^{*}-\bar{z}\right)^{-1}=(A-z)^{-1 *}$. Thus we have $(A-z)^{-1 *} \mathscr{H} \subset \mathscr{G}$ by the assumption $(\mathrm{AB})$ and this allows one to deduce that $(A-z)^{-1}$ extends to a unique continuous operator $\mathscr{G}^{*} \rightarrow \mathscr{H}$, that we shall denote for the moment by $R_{z}$. From $R_{z}(A-z) u=u$ for $u \in \mathcal{D}(A)$ we get, by density of $\mathcal{D}(A)$ in $\mathscr{G}$ and continuity, $R_{z}(\widetilde{A}-z) u=u$ for $u \in \mathscr{G}$, in particular

$$
(B-z)^{-1}=R_{z}(\widetilde{A}-z)(B-z)^{-1} .
$$

On the other hand, the identity

$$
(A-z)^{-1}=(A-z)^{-1}(B-z)(B-z)^{-1}=R_{z}(\widetilde{B}-z)(B-z)^{-1}
$$

is trivial. Subtracting the last two relations we get

$$
(A-z)^{-1}-(B-z)^{-1}=R_{z}(\widetilde{B}-\widetilde{A})(B-z)^{-1}
$$

Since $R_{z}$ is uniquely determined as extension of $(A-z)^{-1}$ to a continuous map $\mathscr{G}^{*} \rightarrow$ $\mathscr{H}$, we shall keep the notation $(A-z)^{-1}$ for it. With this convention, the rigorous version of (3.7) that we shall use is:

$$
\begin{equation*}
(A-z)^{-1}-(B-z)^{-1}=(A-z)^{-1}(\widetilde{B}-\widetilde{A})(B-z)^{-1} \tag{3.8}
\end{equation*}
$$

Theorem 3.2 Let $A, B$ satisfy assumption $(A B)$ and let us assume that there are a Banach module $\mathscr{K}$ and operators $S \in \mathcal{B}\left(\mathscr{K}, \mathscr{G}^{*}\right)$ and $T \in \mathcal{B}_{0}^{l}(\mathscr{G}, \mathscr{K})$ such that $\widetilde{B}-\widetilde{A}=S T$ and $(A-z)^{-1} S \in \mathcal{B}_{q}^{l}(\mathscr{K}, \mathscr{H})$ for some $z \in \rho(A) \cap \rho(B)$. Then the operator $B$ is a compact perturbation of the operator $A$ and $\sigma_{\mathrm{ess}}(B)=\sigma_{\mathrm{ess}}(A)$.

Proof: It suffi ces to show that $R \equiv(A-z)^{-1}-(B-z)^{-1} \in \mathcal{B}_{0}^{l}(\mathscr{H})$, because the domains of $A$ and $B$ are included in $\mathscr{G}$, hence $R \mathscr{H} \subset \mathscr{G}$, which fi nishes the proof because of Proposition 2.16. Now due to (3.8) and to the factorization assumption, we can write $R$ as a product $R=\left[(A-z)^{-1} S\right]\left[T(B-z)^{-1}\right]$ where the first factor is in $\mathcal{B}_{q}^{l}(\mathscr{K}, \mathscr{H})$ and the second in $\mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$, so the product is in $\mathcal{B}_{0}^{l}(\mathscr{H})$.

Remarks 3.3 (1) We could have stated the assumptions of Theorem 3.2 in an apparently more general form, namely $B-A=\sum_{k=1}^{n} S_{k} T_{k}$ with operators $S_{k} \in$ $\mathcal{B}\left(\mathscr{K}_{k}, \mathscr{G}^{*}\right)$ and $T_{k} \in \mathcal{B}\left(\mathscr{G}, \mathscr{K}_{k}\right)$. But we are reduced to the stated version of the assumption by considering the Hilbert module $\mathscr{K}=\oplus \mathscr{K}_{k}$ and $S=\oplus S_{k}, T=\oplus T_{k}$.
(2) If $V \in \mathcal{K}\left(\mathscr{G}, \mathscr{G}^{*}\right)$ and if $\mathscr{K}$ is an infi nite dimensional module, then there are operators $S \in \mathcal{B}\left(\mathscr{K}, \mathscr{G}^{*}\right)$ and $T \in \mathcal{K}(\mathscr{G}, \mathscr{K})$ such that $V=S T$ (the proof is an easy exercise). This and the preceding remark show that compact contributions to $\widetilde{B}-\widetilde{A}$ are trivially covered by the factorization assumption.

If $A$ is self-adjoint then the conditions on $A$ in assumption (AB) are satisfi ed if $\mathcal{D}(A) \subset \mathscr{G} \subset \mathcal{D}\left(|A|^{1 / 2}\right)$ densely (see the Appendix). Moreover, if $A$ is semibounded, then this condition is also necessary. In particular, we have:
Corollary 3.4 Let $A, B$ be self-adjoint operators on $\mathscr{H}$ such that

$$
\mathcal{D}(A) \subset \mathscr{G} \subset \mathcal{D}\left(|A|^{1 / 2}\right) \text { and } \mathcal{D}(B) \subset \mathscr{G} \subset \mathcal{D}\left(|B|^{1 / 2}\right) \text { densely }
$$

Let $\widetilde{A}, \widetilde{B}$ be the unique extensions of $A, B$ to operators in $\mathscr{B}\left(\mathscr{G}, \mathscr{G}^{*}\right)$. Assume that there is a Hilbert module $\mathscr{K}$ and that $\widetilde{B}-\widetilde{A}=S^{*} T$ for some $S \in \mathcal{B}(\mathscr{G}, \mathscr{K})$ and $T \in \mathcal{B}_{0}^{l}(\mathscr{G}, \mathscr{K})$ such that $S(A-z)^{-1} \in \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$ for some $z \in \rho(A) \cap \rho(B)$. Then $B$ is a compact perturbation of $A$ and $\sigma_{\text {ess }}(B)=\sigma_{\text {ess }}(A)$.

The next theorem is convenient for applications to differential operators in divergence form. Observe that if $(\mathscr{E}, \mathscr{K})$ is a Friedrichs module then $\mathcal{B}(\mathscr{K}) \subset \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ hence we can defi ne

$$
\begin{equation*}
\mathcal{B}_{00}^{l}\left(\mathscr{E}, \mathscr{E}^{*}\right)=\text { norm closure of } \mathcal{B}_{0}^{l}(\mathscr{K}) \text { in } \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right) \tag{3.9}
\end{equation*}
$$

We shall use the terminology and the facts established in the Appendix, in particular Lemma A.1: the operators $D^{*} a D$ and $D^{*} b D$ considered below belong to $\mathcal{B}\left(\mathscr{G}^{G}, \mathscr{G}^{*}\right)$ and we denote by $\Delta_{a}$ and $\Delta_{b}$ the operators on $\mathscr{H}$ associated to them.

Theorem 3.5 Let $(\mathscr{E}, \mathscr{K})$ be an arbitrary Friedrichs module and let $D \in \mathcal{B}(\mathscr{G}, \mathscr{E})$, $a, b \in \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ and $z \in \mathbb{C}$ such that:
(1) The operators $D^{*} a D-z$ and $D^{*} b D-z$ are bijective maps $\mathscr{G} \rightarrow \mathscr{G}^{*}$,
(2) $a-b \in \mathcal{B}_{00}^{l}\left(\mathscr{E}, \mathscr{E}^{*}\right)$,
(3) $D\left(\Delta_{a}^{*}-\bar{z}\right)^{-1} \in \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$.

Then $\Delta_{b}$ is a compact perturbation of $\Delta_{a}$.

Proof: We give a proof independent of Theorem 3.2, although we could apply this theorem. Clearly $\Delta_{a}-z$ and $\Delta_{b}-z$ extend to bijections $\mathscr{G} \rightarrow \mathscr{G}^{*}$ and the identity

$$
R:=\left(\Delta_{a}-z\right)^{-1}-\left(\Delta_{b}-z\right)^{-1}=\left(\Delta_{a}-z\right)^{-1} D^{*}(b-a) D\left(\Delta_{b}-z\right)^{-1}
$$

holds in $\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)$, hence in $\mathcal{B}(\mathscr{H})$. Since the domains of $\Delta_{a}$ and $\Delta_{b}$ are included in $\mathscr{G}$, we have $R \mathscr{H} \subset \mathscr{G}$. Thus, according to Proposition 2.16, it suffi ces to show that $R \in \mathcal{B}_{0}^{l}(\mathscr{H})$. Since the space $\mathcal{B}_{0}^{l}(\mathscr{H})$ is norm closed and since by hypothesis we can approach $b-a$ in norm in $\mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ by operators in $\mathcal{B}_{0}^{l}(\mathscr{K})$, it suffi ces to show that

$$
\left(D\left(\Delta_{a}^{*}-\bar{z}\right)^{-1}\right)^{*} c D\left(\Delta_{b}-z\right)^{-1} \in \mathcal{B}_{0}^{l}(\mathscr{H})
$$

if $c \in \mathcal{B}_{0}^{l}(\mathscr{K})$. But this is clear because $c D\left(\Delta_{b}-z\right)^{-1} \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$ and $\left(D\left(\Delta_{a}^{*}-\right.\right.$ $\left.\bar{z})^{-1}\right)^{*} \in \mathcal{B}_{q}^{l}(\mathscr{K}, \mathscr{H})$ by Proposition 2.19.

The spaces $\mathcal{B}_{00}^{r}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ and $\mathcal{B}_{00}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ are defi ned in an obvious way and we have

$$
\begin{equation*}
\mathcal{K}\left(\mathscr{E}, \mathscr{E}^{*}\right) \subset \mathcal{B}_{00}\left(\mathscr{E}, \mathscr{E}^{*}\right) \tag{3.10}
\end{equation*}
$$

because $\mathcal{K}(\mathscr{K})$ is a dense subset of $\mathcal{K}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ and $\mathcal{K}(\mathscr{K}) \subset \mathcal{B}_{0}(\mathscr{K})$. So we could assume $a-b \in \mathcal{K}\left(\mathscr{E}, \mathscr{E}^{*}\right)$, but this case is trivial from the point of view of this paper. Although the space $\mathcal{B}_{00}^{l}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ is much larger than $\mathcal{K}\left(\mathscr{E}, \mathscr{E}^{*}\right)$, it is not satisfactory in some applications, cf. Example 3.6 below and Remark 5.2. However, we can allow still more general perturbations and obtain more explicit results if we impose more structure on the modules. In Subsection 4.3 we describe such improvements for a class of Banach modules over abelian groups.

Example 3.6 In the context of Example 2.12 we may consider the two classes of operators $\mathcal{B}_{00}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{-s}\right)$ and $\mathcal{B}_{0}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{-s}\right)$. The first space is included in the second one and the inclusion is strict, for example $\mathcal{B}_{00}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{-s}\right)$ does not contain operators of order $2 s$, while $\mathcal{B}_{0}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{-s}\right)$ contains such operators.

The only condition of Theorem 3.5 which, in some concrete situations, is not easy to check is condition (3). We now give a perturbative method for checking it.

For the rest of this section we fix two Friedrichs modules $(\mathscr{G}, \mathscr{H})$ and $(\mathscr{E}, \mathscr{K})$ and a continuous operator $D: \mathscr{G} \rightarrow \mathscr{E}$. Let $a \in \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ such that the operator $D^{*} a D$ is coercive (see the Appendix), more precisely we have

$$
\begin{equation*}
\operatorname{Re}\langle D u, a D u\rangle \geq \mu\|u\|_{\mathscr{G}}^{2}-\nu\|u\|_{\mathscr{H}}^{2} \tag{3.11}
\end{equation*}
$$

for some strictly positive constants $\mu, \nu$ and all $u \in \mathscr{G}$ Then, as explained in the Appendix, if $\operatorname{Re} z \leq-\nu$ the operator $D^{*} a D-z$ is a bijective map $\mathscr{G} \rightarrow \mathscr{G}^{*}$ and

$$
\begin{equation*}
\left\|\left(D^{*} a D-z\right)^{-1}\right\|_{\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)} \leq \mu^{-1} \tag{3.12}
\end{equation*}
$$

Note that $a^{*}$ has all these properties too so the closed densely defi ned operators $\Delta_{a}$ and $\Delta_{a^{*}}$ in $\mathscr{H}$ are well defi ned, their domains are dense subsets of $\mathscr{G}$, and we have $\Delta_{a}^{*}=\Delta_{a^{*}}$. It is easy to check that $\left\|\left(\Delta_{a}-z\right)^{-1}\right\|_{\mathcal{B}(\mathscr{H})} \leq|\operatorname{Re} z+\nu|^{-1}$ if $\operatorname{Re} z+\nu<0$. Since $a$ and $a^{*}$ play a symmetric role, it will suffi ce to consider $\Delta_{a}-z$ in place of $\Delta_{a}^{*}-\bar{z}$ in condition (3) of Theorem 3.5.

Now let $c$ be a second operator with the same properties as $a$. We assume, without loss of generality, that it satisfi es an estimate like (3.11) with the same constants $\mu, \nu$.

Proposition 3.7 Assume that

$$
D\left(\Delta_{c}-z\right)^{-1} \in \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K}) \text { and } D\left(D^{*} c D-z\right)^{-1} D^{*} \in \mathcal{B}_{q}^{r}(\mathscr{K})
$$

for some $z$ with $\operatorname{Re} z \leq-\nu$. If $a-c \in \mathcal{B}_{q}^{r}(\mathscr{K})$ then

$$
D\left(\Delta_{a}-z\right)^{-1} \in \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K}) \text { and } D\left(D^{*} a D-z\right)^{-1} D^{*} \in \mathcal{B}_{q}^{r}(\mathscr{K})
$$

A similar assertion holds for the spaces $\mathcal{B}_{q}^{l}$.
Proof: Let $V=D^{*}(a-c) D$ and $L_{t}=(1-t) D^{*} c D+t D^{*} a D=D^{*} c D+t V$. For $z$ as in the statement of the proposition we have $\operatorname{Re}\left\langle u,\left(L_{t}-z\right) u\right\rangle \geq \mu\|u\|_{\mathscr{G}}^{2}$ if $0 \leq t \leq 1$. Hence there is $\varepsilon>0$ such that $\operatorname{Re}\left\langle u,\left(L_{t}-z\right) u\right\rangle \geq \mu / 2\|u\|_{\mathscr{G}}^{2}$ if $-\varepsilon \leq t \leq 1+\varepsilon$, in particular $\left\|\left(L_{t}-z\right)^{-1}\right\|_{\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)} \leq 2 / \mu$ for all such $t$. If $-\varepsilon \leq s \leq 1+\varepsilon$ and $|t-s|\left\|V\left(L_{t}-z\right)^{-1}\right\|_{\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)}<1$ we get a norm convergent expansion in $\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)$

$$
\left(L_{t}-z\right)^{-1}=\left(L_{s}-z-(s-t) V\right)^{-1}=\sum_{k \geq 0}(s-t)^{k}\left(L_{s}-z\right)^{-1}\left[V\left(L_{s}-z\right)^{-1}\right]^{k}
$$

so the map $t \mapsto\left(L_{t}-z\right)^{-1} \in \mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)$ is real analytic on the interval $]-\varepsilon, 1+\varepsilon[$. Let us denote $\Delta_{t}$ the operator in $\mathscr{H}$ associated to $L_{t}$ then we see that the maps $t \mapsto$ $D\left(\Delta_{t}-z\right)^{-1} \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and $t \mapsto D\left(L_{t}-z\right)^{-1} D^{*} \in \mathcal{B}(\mathscr{K})$ are real analytic on the same interval. The set of decay preserving operators is a closed subspace of the Banach space $\mathcal{B}(\mathscr{H}, \mathscr{K})$ and an analytic function which on a nonempty open set takes values in a closed subspace remains in that subspace for ever. Thus it suffi ces to show that $D\left(\Delta_{t}-z\right)^{-1} \in \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$ for small positive values of $t$. Similarly, we need to prove $D\left(L_{t}-z\right)^{-1} D^{*} \in \mathcal{B}(\mathscr{K})$ only for small $t$. To prove the first assertion for example, we take $s=0$ above and get a norm convergent series in $\mathcal{B}(\mathscr{H}, \mathscr{K})$ :

$$
D\left(L_{t}-z\right)^{-1}=\sum_{k \geq 0}(-t)^{k} D\left(D^{*} c D-z\right)^{-1}\left[D^{*}(a-c) D\left(D^{*} c D-z\right)^{-1}\right]^{k}
$$

It is clear that each term belongs to $\mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$.

## 4 Banach modules over abelian groups

## 4.1 $X$-modules over locally compact abelian groups

Since a locally compact abelian group $X$ is a locally compact space, we can consider $X$-modules in the sense of Example 2.4. However, the group structure of $X$ allows us to associate to it more interesting classes of Banach modules that we shall also call $X$-modules. Whenever necessary in order to avoid ambiguities we shall speak of $X$ module over the topological space $X$ if we have in mind the context of Example 2.4 and of $X$-module over the group $X$ when we refer to the structure introduced in the next Defi nition 4.1.

In this section we fix a locally compact non-compact abelian group $X$ with the group operation denoted additively. For example, $X$ could be $\mathbb{R}^{n}, \mathbb{Z}^{n}$, or a fi nite dimensional vector space over a local fi eld, e.g. over the fi eld of p -adic numbers. Let $X^{*}$ be the abelian locally compact group dual to $X$.

Definition 4.1 $A$ Banach $X$-module over the group $X$ is a Banach space $\mathscr{H}$ equipped with a strongly continuous representation $\left\{V_{k}\right\}$ of $X^{*}$ on $\mathscr{H}$.

Note that we shall use the same notation $V_{k}$ for the representations of $X^{*}$ in different spaces $\mathscr{H}$ whenever this does not lead to ambiguities.

Such a Banach X-module has a canonical structure of Banach module that we now defi ne. We choose Haar measures $d x$ and $d k$ on $X$ and $X^{*}$ normalized by the following condition: if the Fourier transform of a function $\varphi$ on $X$ is given by $(\mathcal{F} \varphi)(k) \equiv \widehat{\varphi}(k)=\int_{X} \overline{k(x)} \varphi(x) d x$ then $\varphi(x)=\int_{X^{*}} k(x) \widehat{\varphi}(k) d k$. Recall that $X^{* *}=X$. Let $C^{(a)}(X):=\mathcal{F} L_{\mathrm{c}}^{1}\left(X^{*}\right)$ be the set of Fourier transforms of integrable functions with compact support on $X^{*}$. It is easy to see that $C^{(a)}(X)$ is a $*$-algebra for the usual algebraic operations; more precisely, it is a dense subalgebra of $C_{0}(X)$ stable under conjugation. For $\varphi \in C^{(a)}(X)$ we set

$$
\begin{equation*}
\varphi(Q)=\int_{X^{*}} V_{k} \widehat{\varphi}(k) d k \tag{4.13}
\end{equation*}
$$

This defi nition is determined by the formal requirement $k(Q)=V_{k}$. Then

$$
\begin{equation*}
\mathcal{M}:=\text { norm closure of }\left\{\varphi(Q) \mid \varphi \in C^{(a)}(X)\right\} \text { in } \mathcal{B}(\mathscr{H}) \tag{4.14}
\end{equation*}
$$

is a Banach subalgebra of $\mathcal{B}(\mathscr{H})$. By using the next lemma we see that the couple $(\mathscr{H}, \mathcal{M})$ satisfi es the conditions of Defi nition 2.2, which gives us the canonical Banach module structure on $\mathscr{H}$.

Lemma 4.2 The algebra $\mathcal{M}$ has an approximate unit consisting of elements of the form $e_{\alpha}(Q)$ with $e_{\alpha} \in C^{(a)}(X)$.

Proof: Let us fix a compact neighborhood $K$ of the identity in $X^{*}$. The set of compact neighborhoods of the identity $\alpha$ such that $\alpha \subset K$ is ordered by $\alpha_{1} \geq \alpha_{2} \Leftrightarrow$ $\alpha_{1} \subset \alpha_{2}$. For each such $\alpha$ defi ne $e_{\alpha}$ by $\widehat{e}_{\alpha}=\chi_{\alpha} /|\alpha|$, where $|\alpha|$ is the Haar measure of $\alpha$. Then $\left\|e_{\alpha}(Q)\right\| \leq \sup _{k \in K}\left\|V_{k}\right\|<\infty$, from which it is easy to infer that $\lim _{\alpha}\left\|e_{\alpha}(Q) \varphi(Q)-\varphi(Q)\right\|=0$ for all $\varphi \in C^{(a)}(X)$.

Example 4.3 Let $X=\mathbb{R}^{n}$ with the additive group structure and let $\mathscr{H}$ be the Sobolev space $\mathscr{H}^{s}(X)$ for some real number $s$. We identify as usual $X^{*}$ with $X$ by setting $k(x)=\exp (i\langle x, k\rangle)$, where $\langle x, k\rangle$ is the scalar product in $X$. Then we get a Banach $X$-module structure on $\mathscr{H}$ by setting $\left(V_{k} u\right)(x)=\exp (i\langle x, k\rangle) u(x)$, where $\langle x, k\rangle$ is the scalar product. Note that $V_{k} \mathscr{H}^{s} \subset \mathscr{H}^{s}$ and $\left\|V_{k}\right\| \leq C(1+|k|)^{s}$. It is easy to see that the Banach module structure associated to this $X$-module structure coincides with that defi ned in Example 2.10.

Remark 4.4 Algebras $\mathcal{A}$ as in Lemma 2.3 can be easily constructed in this context. Indeed, if $\omega$ is a sub-multiplicative function on $X^{*}$, i.e. a Borel map $X^{*} \rightarrow[1, \infty[$ satisfying $\omega\left(k^{\prime} k^{\prime \prime}\right) \leq \omega\left(k^{\prime}\right) \omega\left(k^{\prime \prime}\right)$ (hence $\omega$ is locally bounded), let $C^{(\omega)}(X)$ be the set of functions $\varphi$ whose Fourier transform $\widehat{\varphi}$ satisfi es

$$
\begin{equation*}
\|\varphi\|_{C^{(\omega)}}:=\int_{X^{*}}|\widehat{\varphi}(k)| \omega(k) d k<\infty . \tag{4.15}
\end{equation*}
$$

Then $C^{(\omega)}(X)$ is a subalgebra of $C_{0}(X)$ and is a Banach algebra for the norm (4.15). Moreover, $C^{(a)}(X) \subset C^{(\omega)}(X)$ densely and the net $\left\{e_{\alpha}\right\}$ defi ned in the proof of Lemma 4.2 is an approximate unit of $C^{(\omega)}(X)$. If $\left\|V_{k}\right\|_{\mathcal{B}(\mathscr{H})} \leq c \omega(k)$ for some number $c>0$ then $\varphi(Q)$ is well defi ned for each $\varphi \in C^{(\omega)}(X)$ by the relation (4.13) and $\Phi(\varphi)=\varphi(Q)$ is a continuous morphism with dense range of $C^{(\omega)}(X)$ into $\mathcal{M}(\mathscr{H})$. We could take $\omega(k)=\sup \left(1,\left\|V_{k}\right\|_{\mathcal{B}(\mathscr{H})}\right)$ but if a second Banach $X$-module $\mathscr{K}$ is given then it is more convenient to take $\omega(k)=\sup \left\{1,\left\|V_{k}\right\|_{\mathcal{B}(\mathscr{H})},\left\|V_{k}\right\|_{\mathcal{B}(\mathscr{K})}\right\}$.

The adjoint of a reflexive Banach $X$-module has a natural structure of Banach $X$ module. Indeed, a weakly continuous representation is strongly continuous, so we can equip the adjoint space $\mathscr{H}^{*}$ with the Banach $X$-module structure defi ned by the representation $k \mapsto\left(V_{\bar{k}}\right)^{*}$, where $\bar{k}=k^{-1}$ is the complex conjugate of $k$.

The group $X$ is, in particular, a locally compact topological space, hence the notion of Banach $X$-module in the sense of Example 2.4 makes sense. But this is in fact a particular case of that of Banach $X$-module in the sense of Defi nition 4.1. Indeed, according to Remark 2.5 , we get a strongly continuous representation of $X^{*}$ on $\mathscr{H}$ by setting $V_{k}=k(Q)$. In the case of Hilbert $X$-modules we have a more precise fact.

Lemma 4.5 Let $\mathscr{H}$ be a Hilbert space. Then giving a Hilbert $X$-module structure on $\mathscr{H}$ is equivalent with giving on $\mathscr{H}$ a Banach $X$-module structure over the group $X$ such that the representation $\left\{V_{k}\right\}_{k \in X^{*}}$ is unitary. The relation between the two structures is determined by the condition $V_{k}=k(Q)$.

Proof: If $\mathscr{H}$ is a Hilbert $X$-module then we can defi ne $V_{k}=k(Q) \in \mathcal{B}(\mathscr{H})$ and check that $\left\{V_{k}\right\}_{k \in X^{*}}$ is a strongly continuous unitary representation of $X^{*}$ on $\mathscr{H}$ with the help of Remark 2.6. Reciprocally, it is well known that such a representation allows one to equip $\mathscr{H}$ with a Hilbert $X$-module structure. The main point is that the estimate $\|\varphi(Q)\| \leq \sup |\varphi|$ holds, see [Lo].

Banach $X$-modules over the group $X$ which are not Hilbert $X$-modules often appear in the following context (cf. Example 4.3 in the case $s>0$ ).

Definition 4.6 A stable Friedrichs $X$-module over the group $X$ is a Friedrichs $X$ module $(\mathscr{G}, \mathscr{H})$ satisfying $V_{k} \mathscr{G} \subset \mathscr{G}$ for all $k \in X^{*}$ and such that if $u \in \mathscr{G}$ and if $K \subset X^{*}$ is compact then $\sup _{k \in K}\left\|V_{k} u\right\|_{\mathscr{G}}<\infty$.

Here $V_{k}=k(Q)$. It is clear that $V_{k} \mathscr{G} \subset \mathscr{G}$ implies $V_{k} \in \mathcal{B}(\mathscr{G})$ and that the local boundedness condition implies that the map $k \mapsto V_{k} \in \mathcal{B}(\mathscr{G})$ is a weakly, hence strongly, continuous representation of $X^{*}$ on $\mathscr{G}$ (not unitary in general). The local boundedness condition is automatically satisfi ed if $X^{*}$ is second countable.

Thus, if $(\mathscr{G}, \mathscr{H})$ is a stable Friedrichs $X$-module, then $\mathscr{G}$ is equipped with a canonical Banach $X$-module structure. Then, by taking adjoints, we get a natural Banach $X$-module structure on $\mathscr{G}^{*}$ too. Our defi nitions are such that after the identifi cations $\mathscr{G} \subset \mathscr{H} \subset \mathscr{G}^{*}$ the restriction to $\mathscr{H}$ of the operator $V_{k}$ acting in $\mathscr{G}^{*}$ is just the initial $V_{k}$. Indeed, we have $V_{k}^{*}=V_{k}^{-1}=V_{\bar{k}}$ in $\mathscr{H}$. Thus there is no ambiguity in using the same notation $V_{k}$ for the representation of $X^{*}$ in the spaces $\mathscr{G}, \mathscr{H}$ and $\mathscr{G}^{*}$.

Proposition 4.7 If $\mathscr{K}$ is a Banach space then $\mathcal{B}_{0}^{l}(\mathscr{K}, \mathscr{G}) \subset \mathcal{B}_{0}^{l}(\mathscr{K}, \mathscr{H})$, and if $\mathscr{K}$ is a Banach module then $\mathcal{B}_{q}^{l}(\mathscr{K}, \mathscr{G}) \subset \mathcal{B}_{q}^{l}(\mathscr{K}, \mathscr{H})$.

Proof: If $S \in \mathcal{B}_{0}^{l}(\mathscr{K}, \mathscr{G})$ then $S=\varphi(Q) T$ for some $\varphi \in C^{(\omega)}(X)$ with $\omega(k)=$ $\sup \left(1,\left\|V_{k}\right\|_{\mathcal{B}(\mathscr{G})}\right)$ and some $T \in \mathcal{B}(\mathscr{K}, \mathscr{G})$ (see Remark 4.4). But clearly such a $\varphi(Q)$ belongs to the multiplier algebra of $\mathscr{H}$ and $T \in \mathcal{B}(\mathscr{K}, \mathscr{H})$.

### 4.2 Regular operators are decay preserving

We show now that, in the case of Banach $X$-modules over groups, the decay preserving property is related to regularity in the sense of the next defi nition.

Definition 4.8 Let $\mathscr{H}$ and $\mathscr{K}$ be Banach $X$-modules. We say that a continuous operator $S: \mathscr{H} \rightarrow \mathscr{K}$ is of class $C^{\mathrm{u}}(Q)$, and we write $S \in C^{\mathrm{u}}(Q ; \mathscr{H}, \mathscr{K})$, if the map $k \mapsto V_{k}^{-1} S V_{k} \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ is norm continuous.

Note that norm continuity at the origin implies norm continuity everywhere. The class of regular operators is stable under algebraic operations:

Proposition 4.9 Let $\mathscr{G}, \mathscr{H}, \mathscr{K}$ be Banach $X$-modules.
(i) If $S \in C^{\mathrm{u}}(Q ; \mathscr{H}, \mathscr{K})$ and $T \in C^{\mathrm{u}}(Q ; \mathscr{G}, \mathscr{H})$ then $S T \in C^{\mathrm{u}}(Q ; \mathscr{G}, \mathscr{K})$.
(ii) If $S \in C^{\mathrm{u}}(Q ; \mathscr{H}, \mathscr{K})$ is bijective, then $S^{-1} \in C^{\mathrm{u}}(Q ; \mathscr{K}, \mathscr{H})$.
(iii) If $S \in C^{\mathrm{u}}(Q ; \mathscr{H}, \mathscr{K})$ and $\mathscr{H}, \mathscr{G}$ are reflexive, then $S^{*} \in C^{\mathrm{u}}\left(Q ; \mathscr{K}^{*}, \mathscr{H}^{*}\right)$.

Proof: We prove only (ii). If we set $S_{k}=V_{k}^{-1} S V_{k}$ then $V_{k}^{-1} S^{-1} V_{k}=S_{k}^{-1}$, hence

$$
\left\|V_{k}^{-1} S^{-1} V_{k}-S^{-1}\right\|=\left\|S_{k}^{-1}-S^{-1}\right\|=\left\|S_{k}^{-1}\left(S-S_{k}\right) S^{-1}\right\| \leq C\left\|S-S_{k}\right\|
$$

if $k$ is in a compact set, and this tends to zero as $k \rightarrow 0$.

Proposition 4.10 If $T \in C^{\mathrm{u}}(Q ; \mathscr{H}, \mathscr{K})$ then $T$ is decay preserving.
Proof: We show that $\varphi(Q) T \in \mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$ if $\varphi \in C^{(a)}(X)$. A similar argument gives $T \varphi(Q) \in \mathcal{B}_{0}^{l}(\mathscr{H}, \mathscr{K})$. Set $T_{k}=V_{k} T V_{k}^{-1}$, then

$$
\varphi(Q) T=\int_{X^{*}} \widehat{\varphi}(k) V_{k} T d k=\int_{X^{*}} T_{k} \widehat{\varphi}(k) V_{k} d k
$$

Since $k \mapsto T_{k}$ is norm continuous on the compact support of $\widehat{\varphi}$, for each $\varepsilon>0$ we can construct, with the help of a partition of unity, functions $\theta_{i} \in C_{\mathrm{c}}\left(X^{*}\right)$ and operators $S_{i} \in \mathcal{B}(\mathscr{H}, \mathscr{K})$, such that $\left\|T_{k}-\sum_{i=1}^{n} \theta_{i}(k) S_{i}\right\|<\varepsilon$ if $\widehat{\varphi}(k) \neq 0$. Thus

$$
\left\|\varphi(Q) T-\sum_{i=1}^{n} \int_{X^{*}} \theta_{i}(k) S_{i} \widehat{\varphi}(k) V_{k} d k\right\| \leq \varepsilon \sum_{i=1}^{n} \int_{X^{*}}|\widehat{\varphi}(k)|\left\|V_{k}\right\|_{\mathcal{B}(\mathscr{H})} d k
$$

Now, since $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$ is a norm closed subspace, it suffi ces to show that the operator $\int_{X^{*}} \theta_{i}(k) S_{i} \widehat{\varphi}(k) V_{k} d k$ belongs to $\mathcal{B}_{0}^{r}(\mathscr{H}, \mathscr{K})$ for each $i$. But if $\psi_{i}$ is the inverse Fourier transform of $\theta_{i} \widehat{\varphi}$ then this is $S_{i} \psi_{i}(Q)$ and $\psi_{i} \in C^{(a)}(X)$.

Let $\mathscr{H}, \mathscr{K}$ be Hilbert $X$-modules over the group $X$. We say that an operator $S \in$ $\mathcal{B}(\mathscr{H}, \mathscr{K})$ is of finite range ${ }^{3}$ if there is a compact neighborhood $\Lambda$ of the origin such that for any compact sets $H, K \subset X$ with $(H-K) \cap \Lambda=\emptyset$ we have $\chi_{H}(Q) S \chi_{K}(Q)=$ 0 . From Remark 2.6 we get that this is equivalent to $S \chi_{K}(Q)=\chi_{K+\Lambda}(Q) S \chi_{K}(Q)$ for any Borel set $K$. A fi nite range operator is clearly decay preserving. Moreover, the set of fi nite range operators is stable under sums and products, and the adjoint of such an operator is also of fi nite range.
Proposition 4.11 If $\mathscr{H}, \mathscr{K}$ are Hilbert $X$-modules over the group $X$, then each operator of class $C^{\mathrm{u}}(Q)$ is a norm limit of a sequence of finite range operators.
Proof: We fix a Haar measure $d k$ on $X^{*}$ and if $S \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and $\theta \in L^{1}\left(X^{*}\right)$ we defi ne

$$
\begin{equation*}
S_{\theta}=\int_{X^{*}} V_{k}^{*} S V_{k} \theta(k) d k \tag{4.16}
\end{equation*}
$$

The integral is well defi ned because $k \mapsto V_{k}^{*} S V_{k} \in \mathcal{B}(\mathscr{H})$ is a bounded strongly continuous map. In order to explain the main idea of the proof we shall make a formal computation involving the spectral measure $E(A)=\chi_{A}(Q)$, see Remark 2.6 and Lemma 4.5 (we shall use the same notation for the spectral measures in $\mathscr{H}$ and $\mathscr{K}$ ). We have for $k \in X^{*}$ and $\varphi(Q) \in B(X)$

$$
\varphi(Q) V_{k}^{*}=\varphi(Q) k(Q)^{*}=(\varphi \bar{k})(Q)=\int \varphi(x) \bar{k}(x) E(d x)
$$

Note also that for $x, y \in X$ we have $\bar{k}(x) k(y)=k(-x) k(y)=k(y-x)$. Let $\widehat{\theta}(x)=\int \overline{k(x)} \theta(k) d k$ be the Fourier transform of $\theta$. Then if $\varphi, \psi \in B(X)$ :

$$
\begin{align*}
\varphi(Q) S_{\theta} \psi(Q) & =\int_{X^{*}} \theta(k) d k \int_{X} \int_{X} \varphi(x) \bar{k}(x) k(y) \psi(y) E(d x) S E(d y) \\
& =\int_{X} \int_{X} \widehat{\theta}(x-y) \varphi(x) \psi(y) E(d x) S E(d y) \tag{4.17}
\end{align*}
$$

This clearly implies the following:
$(*)\left\{\begin{array}{l}\text { If the support of } \hat{\theta} \text { is a compact set } \Lambda \text { and if } \operatorname{supp} \varphi \cap(\Lambda+\operatorname{supp} \psi)=\varnothing \\ \text { then } \varphi(Q) S_{\theta} \psi(Q)=0 .\end{array}\right.$
We shall note give a rigorous justifi cation of (4.17) but we shall prove the preceding assertion, which suffi ces for our purposes. Observe that if $(*)$ holds for a certain set of operators $S$ then it also holds for the strongly closed linear subspace of $\mathcal{B}(\mathscr{H}, \mathscr{K})$ generated by it. So it suffi ces to prove ( $*$ ) for $S$ an operator of rank one $S f=v\langle u, f\rangle$ with some fi xed $u \in \mathscr{H}$ and $v \in \mathscr{K}$. Now the computation giving (4.17) obviously makes sense in the weak topology and gives for $f \in \mathscr{H}$ and $g \in \mathscr{K}$ :

$$
\left\langle g, \varphi(Q) S_{\theta} \psi(Q) f\right\rangle=\int_{X} \int_{X} \widehat{\theta}(x-y) \varphi(x) \psi(y)\langle g, E(d x) u\rangle\langle u, E(d y) f\rangle
$$

[^2]hence ( $*$ ) holds for such $S$.
Finally, note that if $S \in C^{\mathrm{u}}(Q)$ then $S$ is norm limit of operators of the form $S_{\theta}$. For this it suffi ces to take $\theta=|K|^{-1} \chi_{K}$ where $K$ runs over the set of open relatively compact neighbourhoods of the neutral element of $X^{*},|K|$ being the Haar measure of $K$. Then, by approximating conveniently $\theta$ in $L^{1}$ norm, one shows that $S$ is norm limit of operators $S_{\theta}$ such that $\widehat{\theta}$ has compact support.

Proposition 4.12 Assume that $X$ is a disjoint union $X=\cup_{a \in A} X_{a}$ of Borel sets $X_{a}$ such that: 1) there is a compact set $K$ such that each $X_{a}$ is a translate of a subset of $K$, and 2) for each compact neighborhood $\Lambda$ of the origin, the number of sets $X_{b}+\Lambda$ which intersects a given $X_{a}$ is bounded by a constant independent of a. Then, if $\mathscr{H}, \mathscr{K}$ are Hilbert $X$-modules over the group $X$, a finite range operator is of class $C^{\mathrm{u}}(Q)$.
Proof: Let $S$ be a fi nite range operator and let $\Lambda$ be such that $\chi_{H}(Q) S \chi_{K}(Q)=0$ if $H, K$ are compact sets with $(H-K) \cap \Lambda=\emptyset$. Let $\chi_{a}$ be the characteristic function of $X_{a}$ and $\varphi_{a}$ that of $Y_{a}=X_{a}+\Lambda$. We can assume that $A \subset X$ and that $X_{a}=a+K_{a}$ for some $K_{a} \subset \Lambda$. We shall abbreviate $\chi_{a}=\chi_{a}(Q)$ and $\varphi_{a}=\varphi_{a}(Q)$. We have $\sum_{a} \chi_{a}=1$ strongly on $\mathscr{H}$, cf. Remark 2.6, and $\left[V_{k}, S\right] \chi_{a}=\varphi_{a}\left[V_{k}, S\right] \chi_{a}$ because $V_{k}=k(Q)$. Thus there is a constant $C$, depending only on an upper bound for the number of $Y_{b}$ which intersects a fi xed $X_{a}$, such that for $u \in \mathscr{H}$ with compact support:

$$
\begin{aligned}
\left\|\left[V_{k}, S\right] u\right\|^{2} & \leq C \sum\left\|\varphi_{a}\left[V_{k}, S\right] \chi_{a} u\right\|^{2} \\
& =C \sum\left\|\varphi_{a}\left[V_{k}-k(a), S\right] \chi_{a} \cdot \chi_{a} u\right\|^{2} \\
& \leq 2 C L_{k} \sum\left\|\chi_{a} u\right\|^{2}=2 C L_{k}\|u\|^{2}
\end{aligned}
$$

where $L_{k}=\sup _{a}\left\|\left(V_{k}-k(a)\right) \varphi_{a}\right\|$. But

$$
\left\|\left(V_{k}-k(a)\right) \varphi_{a}\right\| \leq \sup _{y \in Y_{a}}|k(y)-k(a)|=\sup _{y \in Y_{a}}|k(y-a)-1| \leq \sup _{x \in L}|k(x)-1|
$$

where $L=\Lambda+\Lambda$ is a compact set. Thus $L_{k} \rightarrow 0$ if $k \rightarrow 0$ in $X^{*}$.
If $X$ is an abelian locally compact group then there is enough structure in order to develop a rich pseudo-differential calculus in $L^{2}(X)$ and Proposition 4.10 shows that many pseudo-differential operators are decay preserving. We give a simple example below. If $\varphi$ and $\psi$ are Borel functions on $X$ and $X^{*}$ respectively then, following standard quantum mechanical conventions, we denote by $\varphi(Q)$ the operator of multiplication by $\varphi$ in $L^{2}(X)$ and we set $\psi(P)=\mathcal{F}^{-1} M_{\psi} \mathcal{F}$, where $M_{\psi}$ is the operator of multiplication by $\psi$ in $L^{2}\left(X^{*}\right)$.

Let $C_{\mathrm{b}}^{\mathrm{u}}(X)$ and $C_{\mathrm{b}}^{\mathrm{u}}\left(X^{*}\right)$ be the algebras of bounded uniformly continuous functions on $X$ and $X^{*}$ respectively. Below the space $L^{2}(X)$ is equipped with its natural Hilbert $X$-module structure.

Proposition 4.13 The $C^{*}$-algebra generated by the operators $\varphi(Q)$ and $\psi(P)$, with $\varphi \in C_{\mathrm{b}}^{\mathrm{u}}(X)$ and $\psi \in C_{\mathrm{b}}^{\mathrm{u}}\left(X^{*}\right)$, consists of decay preserving operators.
Proof: By Proposition 2.19, $\mathcal{B}_{q}\left(L^{2}(X)\right)$ is a $C^{*}$-algebra, hence it suffi ces to show that each $\varphi(Q)$ and $\psi(P)$ is decay preserving. For $\varphi(Q)$ the assertion is trivial while for $\psi(P)$ we apply Proposition 4.10.

### 4.3 Compact perturbations in modules over abelian groups

In the present context it is possible to improve the results of Section 3.
Lemma 4.14 Let $(\mathscr{G}, \mathscr{H})$ and $(\mathscr{E}, \mathscr{K})$ be stable Friedrichs $X$-modules over the group X. Let $D \in \mathcal{B}(\mathscr{G}, \mathscr{E})$ and $a \in \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ be operators of class $C^{\mathrm{u}}(Q)$ such that $D^{*} a D-z: \mathscr{G} \rightarrow \mathscr{G}^{*}$ is bijective for some complex number $z$ and let $\Delta_{a}$ the operator on $\mathscr{H}$ associated to $D^{*} a D$. Then the operator $D\left(\Delta_{a}-z\right)^{-1} \in \mathcal{B}(\mathscr{H}, \mathscr{E})$ is decay preserving.

Proof: The lemma is an easy consequence of Propositions 4.9 and 4.10. Indeed, due to Proposition 4.10, it suffi ces to show that the operator $D\left(\Delta_{a}-z\right)^{-1}$ is of class $C^{\mathrm{u}}(Q ; \mathscr{H}, \mathscr{E})$. We shall prove more, namely that $D\left(D^{*} a D-z\right)^{-1}$ is of class $C^{\mathrm{u}}\left(Q ; \mathscr{G}^{*}, \mathscr{E}\right)$. Since $D$ is of class $C^{\mathrm{u}}(Q ; \mathscr{G}, \mathscr{E})$, and due to (i) of Proposition 4.9, it suffi ces to show that $\left(D^{*} a D-z\right)^{-1}$ is of class $C^{\mathrm{u}}\left(Q ; \mathscr{G}^{*}, \mathscr{G}\right)$. But $D^{*} a D-z$ is of class $C^{\mathrm{u}}\left(Q ; \mathscr{G}, \mathscr{G}^{*}\right)$ by (i) and (iii) of Proposition 4.9 and is a bijective map $\mathscr{G} \rightarrow \mathscr{G}^{*}$, so the result follows from (ii) of Proposition 4.9.

Theorem 4.15 Let $X$ be an abelian locally compact group and let $(\mathscr{G}, \mathscr{H})$ be a compact stable Friedrichs $X$-module and $(\mathscr{E}, \mathscr{K})$ a stable Friedrichs $X$-module. Assume that $D \in \mathcal{B}(\mathscr{G}, \mathscr{E})$ and $a, b \in \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ are operators of class $C^{\mathrm{u}}(Q)$ such that the operators $D^{*} a D-z$ and $D^{*} b D-z$ are bijective maps $\mathscr{G} \rightarrow \mathscr{G}^{*}$ for some complex number $z$. If $a-b \in \mathcal{B}_{0}^{l}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ then $\Delta_{b}$ is a compact perturbation of $\Delta_{a}$.

Proof: The proof is a repetition of that of Theorem 3.5. The only difference is that we write directly

$$
R=\left(D\left(\Delta_{a}^{*}-\bar{z}\right)^{-1}\right)^{*}(b-a) D\left(\Delta_{b}-z\right)^{-1}
$$

and observe that $(b-a) D\left(\Delta_{b}-z\right)^{-1} \in \mathcal{B}_{0}^{l}\left(\mathscr{H}, \mathscr{E}^{*}\right)$ and that $\left(D\left(\Delta_{a}^{*}-\bar{z}\right)^{-1}\right)^{*}$ as an operator $\mathscr{E}^{*} \rightarrow \mathscr{H}$ is decay preserving by (2) of Proposition 2.19 and because the operator $D\left(\Delta_{a}^{*}-\bar{z}\right)^{-1}: \mathscr{H} \rightarrow \mathscr{E}$ is decay preserving by Lemma 4.14.

We fi nish with a simple corollary of Theorem 3.2 which nevertheless covers interesting examples of differential operators of any order.

Theorem 4.16 Assume that $(\mathscr{G}, \mathscr{H})$ is a compact stable Friedrichs $X$-module over the group $X$ and that condition $(A B)$ from page 13 is satisfied. Let us also assume that $\widetilde{A}-z: \mathscr{G} \rightarrow \mathscr{G}^{*}$ is bijective for some $z \in \rho(A) \cap \rho(B)$ and that $\widetilde{A} \in C^{\mathrm{u}}\left(Q ; \mathscr{G}, \mathscr{G}^{*}\right)$. If $\widetilde{B}-\widetilde{A} \in \mathcal{B}_{0}^{l}\left(\mathscr{G}, \mathscr{G}^{*}\right)$, then $B$ is a compact perturbation of $A$.

Proof: We apply Theorem 3.2 with $\mathscr{K}=\mathscr{G}^{*}, S$ the identity operator and $T=\widetilde{B}-\widetilde{A}$. Then $(\widetilde{A}-z)^{-1}$ is of class $C^{\mathrm{u}}\left(Q ; \mathscr{G}^{*}, \mathscr{G}\right)$ by (ii) of Proposition 4.9, hence $(\widetilde{A}-z)^{-1} \in$ $\mathcal{B}_{q}\left(\mathscr{G}^{*}, \mathscr{G}\right)$ by Proposition 4.10. But this is stronger than $(\widetilde{A}-z)^{-1} \in \mathcal{B}_{q}^{l}\left(\mathscr{G}^{*}, \mathscr{H}\right)$, as follows from Proposition 4.7.

### 4.4 A class of hypoelliptic operators on abelian groups

In this subsection we assume that $X$ is non-discrete, so $X^{*}$ is non-compact. We also fi x a fi nite dimensional complex Hilbert space $E$ and take $\mathscr{H}=E(X ; E)$ equipped with its natural Hilbert $X$-module structure. Note that, according to our conventions, the unitary representation of $X^{*}$ is given by the multiplication operators $V_{k}=k(Q)$.

Let $w: X^{*} \rightarrow[1, \infty[$ be a continuous function satisfying $w(k) \rightarrow \infty$ as $k \rightarrow \infty$ and such that $w\left(k^{\prime} k\right) \leq \omega\left(k^{\prime}\right) w(k)$ holds for some function $\omega$ and all $k^{\prime}, k$. We shall assume that $\omega$ is the smallest function satisfying the preceding estimate. It is clear then that $\omega$ is sub-multiplicative in the sense defi ned in Remark 4.4 (see [Ho, Section 10.1] for this construction).

Then $w(P)$ is a self-adjoint operator on $\mathscr{H}$ with $w(P) \geq 1$ (see page 21 for this notation). We denote $\mathscr{H}^{w}=\mathcal{D}(w(P))$ and equip it with the Banach $X$-module structure given by the norm $\|u\|_{w}=\|w(P) u\|$ and the representation $V_{k} \mid \mathscr{H}^{w}$. Obviously, this space is a generalization of the usual notion of Sobolev spaces.

Lemma $4.17\left(\mathscr{H}^{w}, \mathscr{H}\right)$ is a compact stable Friedrichs $X$-module.
Proof: If $\varphi \in C_{0}(X)$ then $\varphi(Q) w(P)^{-1}$ is a compact operator because $w^{-1}$ belongs to $C_{0}(X)$, hence $\varphi(Q) \in \mathcal{K}\left(\mathscr{H}^{w}, \mathscr{H}\right)$. Then observe that $V_{k}^{-1} w(P) V_{k}=w(k P)$ and $w(k P) \leq \omega(k) w(P)$. Thus $V_{k}$ leaves stable $\mathscr{H}^{w}$ and we have the estimate $\left\|V_{k}\right\|_{\mathcal{B}\left(\mathscr{H}^{w}\right)} \leq \omega(k)$.

We call uniformly hypoelliptic an operator $A$ on $\mathscr{\sim} \mathscr{H}$ such that there are $w$ as above and an operator $\widetilde{A} \in \mathcal{B}\left(\mathscr{H}^{w}, \mathscr{H}^{w *}\right)$ such that $\widetilde{A}-z: \mathscr{H}^{w} \rightarrow \mathscr{H}^{w *}$ is bijective for some complex $z$ and such that $A$ is the operator induced by $\widetilde{A}$ in $\mathscr{H}$ (see the Appendix). For example, the constant coeffi cients case with $E=\mathbb{C}$ corresponds to the choice $A=h(P)$ with $h: X^{*} \rightarrow \mathbb{C}$ a Borel function such that $c^{\prime} w^{2} \leq 1+|h| \leq c^{\prime \prime} w^{2}$ and such that the range of $h$ is not dense in $\mathbb{C}$. We shall justify our terminology in the remark at the end of this subsection.

Theorem 4.16 is quite well adapted to show the stability of the essential spectrum of such operators under perturbations which are small at infi nity. We stress that the differential operators covered by these results can be of any order and that in the usual case when the coeffi cients are complex measurable functions a condition of the type $\widetilde{A} \in C^{\mathrm{u}}\left(Q ; \mathscr{H}^{w}, \mathscr{H}^{w *}\right)$ is very general, if not automatically satisfi ed (see the remark at the end of this subsection). Hence the only condition really relevant in this context is $\widetilde{B}-\widetilde{A} \in \mathcal{B}_{0}^{l}\left(\mathscr{H}^{w}, \mathscr{H}^{w *}\right)$ and the main point is that it allows perturbations of the higher order coeffi cients even in the non-smooth case.

It is clear that these results can be used to establish the stability of the essential spectrum of pseudo-differential operators on fi nite dimensional vector spaces over local fi elds (see [Sa, Ta]) under perturbations of the same order.

We shall give an application of physical interest to Dirac operators. Let $X=\mathbb{R}^{n}$ and let $\alpha_{0} \equiv \beta, \alpha_{1}, \ldots, \alpha_{n}$ be symmetric operators on $E$ such that $\alpha_{j} \alpha_{k}+\alpha_{k} \alpha_{j}=\delta_{j k}$. Then the free Dirac operator is $D=\sum_{k=1}^{n} \alpha_{k} P_{k}+m \beta$ for some real number $m$. The natural compact stable Friedrichs $X$-module in this context is $\left(\mathscr{H}^{1 / 2}, \mathscr{H}\right)$. Note that we use the same notation $\mathscr{H}^{s}$ for Sobolev spaces of $E$-valued functions.

Proposition 4.18 Let $V$, $W$ be measurable functions on $X$ with values symmetric operators on $E$ and such that the operators of multiplication by $V$ and $W$ define continuous maps $\mathscr{H}^{1 / 2} \rightarrow \mathscr{H}^{-1 / 2}$ and $V-W \in \mathcal{B}_{0}\left(\mathscr{H}^{1 / 2}, \mathscr{H}^{-1 / 2}\right)$. Assume that $D+V+i$ and $D+W+i$ are bijective maps $\mathscr{H}^{1 / 2} \rightarrow \mathscr{H}^{-1 / 2}$. Then $D+V$ and $D+W$ induce self-adjoint operators $A$ and $B$ in $\mathscr{H}, B$ is a compact perturbation of $A$, and $\sigma_{\text {ess }}(B)=\sigma_{\text {ess }}(A)$.

This follows immediately from Theorem 4.16. We stress that the main new feature of this result is that the "unperturbed" operator $A$ is locally as singular as the "perturbed" one $B$. The assumptions imposed on $V, W$ are quite general, compare with [Ar, AY, K1, N1, N2].

Remark: In order to clarify the relation between the notion of uniform hypoellipticity introduced above and the original notion of hypoellipticity due to Hörmander, we shall consider the case of differential operators on $\mathbb{R}^{n}$ (which is identifi ed with its dual group in the standard way). Assume first that $h$ is a polynomial on $\mathbb{R}^{n}$ and that $A=h(P)$. Then the function defi ned by $w(k)^{4}=\sum_{\alpha}\left|h^{(\alpha)}(k)\right|^{2}$ satisfi es $w\left(k^{\prime}+k\right) \leq\left(1+c\left|k^{\prime}\right|\right)^{m / 2} w(k)$, where $c$ is a number and $m$ is the order of $h$, see [Ho, Example 10.1.3]. Now the "form domain" of the operator $h(P)$ in $L^{2}\left(\mathbb{R}^{n}\right)$ is the space $\mathscr{G}=\mathcal{D}\left(|h(P)|^{1 / 2}\right)$ and this domain is stable under $V_{k}=\exp i\langle k, Q\rangle$ if and only the function $w$ satisfi es $w^{2} \leq c(1+|h|)$, see Lemma 7.6.7 in [ABG]. On the other hand, Defi nition 11.1.2 and Theorem 11.1.3 from [Ho] show that $A$ is hypoelliptic if and only if $h^{(\alpha)}(k) / h(k) \rightarrow 0$ when $k \rightarrow \infty$, for all $\alpha \neq 0$. So in this case we have $c^{\prime} w^{2} \leq 1+|h| \leq c^{\prime \prime} w^{2}$ and the operator $h(P)$ is uniformly hypoelliptic in our sense if $h\left(\mathbb{R}^{n}\right)$ is not dense in $\mathbb{C}$. If $n=2$ then $h(k)=k_{1}^{4}+k_{2}^{2}$ is a simple example of polynomial which satisfi es all these conditions but is not elliptic. See [GM, Subsections 2.7-2.10] for the case of matrix valued functions $h$.

In the variable coeffi cient case the notion of hypoellipticity defi ned in [Ho, Defi nition 13.4.3] is a local one and one may consider different global versions. For instance, [Ho, Theorem 13.4.4] suggests that the notion we introduced above is natural for operators of uniform constant strength. But the uniform constant strength condition is not satisfi ed by the operators with polynomial coeffi cients, for example, hence such operators are not uniformly hypoelliptic in our sense in general.

## 5 Operators in divergence form on Euclidean spaces

The results of this section are corollaries of Theorem 4.15. We shall take $X=\mathbb{R}^{n}$, we fi x a fi nite dimensional Hilbert space $E$, and choose $\mathscr{H}=E(X ; E)$ with the obvious Hilbert $X$-module structure. If $s \in \mathbb{R}$ then $\mathscr{H}^{s}$ is the usual Sobolev space of $E$ valued functions. Then for each $s>0$ the couple $\left(\mathscr{H}^{s}, \mathscr{H}\right)$ is a compact stable Friedrichs $X$-module, cf. Examples 2.10, 2.12 and 4.3.

Let us describe the objects which appear in Theorem 4.15 in the present context. We fix an integer $m \geq 1$ and take $\mathscr{G}=\mathscr{H}^{m}$. Let $\mathscr{K}=\bigoplus_{|\alpha| \leq m} \mathscr{H}_{\alpha}$, where $\mathscr{H}_{\alpha} \equiv \mathscr{H}$, with the natural direct sum Hilbert $X$-module structure. Here $\alpha$ are multi-
indices $\alpha \in \mathbb{N}^{n}$ and $|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$. Then we defi ne

$$
\mathscr{E}=\bigoplus_{|\alpha| \leq m} \mathscr{H}^{m-|\alpha|}=\left\{\left(u_{\alpha}\right)_{|\alpha| \leq m} \in \mathscr{K} \mid u_{\alpha} \in \mathscr{H}^{m-|\alpha|}\right\}
$$

equipped with the Hilbert direct sum structure. It is obvious that $(\mathscr{E}, \mathscr{K})$ is a stable Friedrichs $X$-module (but not compact).

We set $P_{k}=-i \partial_{k}$, where $\partial_{k}$ is the derivative with respect to the $k$-th variable, and $P^{\alpha}=P_{1}^{\alpha_{1}} \ldots P_{n}^{\alpha_{n}}$ if $\alpha \in \mathbb{N}^{n}$. Then for $u \in \mathscr{G}$ let $D u=\left(P^{\alpha} u\right)_{|\alpha| \leq m} \in \mathscr{K}$. Since

$$
\|D u\|^{2}=\sum_{|\alpha| \leq m}\left\|P^{\alpha} u\right\|^{2}=\|u\|_{\mathscr{H}^{m}}^{2}
$$

we see that $D: \mathscr{G} \rightarrow \mathscr{K}$ is a linear isometry. Moreover, we have defi ned $\mathscr{E}$ such as to have $D \mathscr{G} \subset \mathscr{E}$, hence $D \in \mathcal{B}(\mathscr{G}, \mathscr{E})$. We have $D \in C^{\mathrm{u}}(Q ; \mathscr{G}, \mathscr{E})$ because

$$
V_{k}^{-1} D V_{k}=\left(V_{k}^{-1} P^{\alpha} V_{k}\right)_{|\alpha| \leq m}=\left((P+k)^{\alpha}\right)_{|\alpha| \leq m}
$$

and this a polynomial in $k$ with coeffi cients in $\mathcal{B}(\mathscr{G}, \mathscr{E})$.
We shall identify $\mathscr{H}^{*}=\mathscr{H}$ and $\mathscr{K}^{*}=\mathscr{K}$, which implies $\mathscr{G}^{*}=\mathscr{H}^{-m}$ and

$$
\mathscr{E}^{*}=\oplus_{|\alpha| \leq m} \mathscr{H}^{|\alpha|-m}
$$

The operator $D^{*} \in \mathcal{B}\left(\mathscr{E}^{*}, \mathscr{G}^{*}\right)$ acts as follows:

$$
D^{*}\left(u_{\alpha}\right)_{|\alpha| \leq m}=\sum_{|\alpha| \leq m} P^{\alpha} u_{\alpha} \in \mathscr{H}^{-m},
$$

because $u_{\alpha} \in \mathscr{H}{ }^{|\alpha|-m}$.
By taking into account the given expressions for $\mathscr{E}$ and $\mathscr{E}^{*}$ we see that we can identify an operator $a \in \mathcal{B}\left(\mathscr{E}, \mathscr{E}^{*}\right)$ with a matrix of operators $a=\left(a_{\alpha \beta}\right)_{|\alpha|,|\beta| \leq m}$, where $a_{\alpha \beta} \in \mathcal{B}\left(\mathscr{H}^{m-|\beta|}, \mathscr{H}^{|\alpha|-m}\right)$ and

$$
a\left(u_{\beta}\right)_{|\beta| \leq m}=\left(\sum_{|\beta| \leq m} a_{\alpha \beta} u_{\beta}\right)_{|\alpha| \leq m} .
$$

Then we clearly have

$$
\begin{equation*}
D^{*} a D=\sum_{|\alpha|,|\beta| \leq m} P^{\alpha} a_{\alpha \beta} P^{\beta} . \tag{5.18}
\end{equation*}
$$

which is a general version of a differential operator in divergence form. We must, however, emphasize that our $a_{\alpha \beta}$ are not necessarily ( $B(E)$ valued) functions, they could be pseudo-differential or more general operators.

In view of the statement of the next theorem, we note that, since the Sobolev spaces are Banach $X$-modules over the group $X$, the class of regularity $C^{\mathrm{u}}\left(Q ; \mathscr{H}^{s}, \mathscr{H}^{t}\right)$ is well defi ned for all real $s, t$. A bounded operator $S: \mathscr{H}^{s} \rightarrow \mathscr{H}^{t}$ belongs to this class if and only if the map $k \mapsto V_{-k} S V_{k} \in \mathcal{B}\left(\mathscr{H}^{s}, \mathscr{H}^{t}\right)$ is norm continuous. In particular, this condition is trivially satisfi ed if $S$ is the operator of multiplication
by a function, because then $V_{k}$ commutes with $S$. Since the coeffi cients $a_{\alpha \beta}$ of the differential expression (5.18) are usually assumed to be functions, this is a quite weak restriction in the setting of the next theorem. The condition $S \in \mathcal{B}_{0}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{t}\right)$ is also well defi ned and it is easily seen that it is equivalent to

$$
\begin{equation*}
\lim _{r \rightarrow \infty}\|\theta(Q / r) S\|_{\mathscr{H}_{s}^{s} \rightarrow \mathscr{H}^{t}}=0 \tag{5.19}
\end{equation*}
$$

where $\theta$ is a $C^{\infty}$ function on $X$ equal to zero on a neighborhood of the origin and equal to one on a neighborhood of infi nity. Now we can state the following immediate consequence of Theorem 4.15 .

Proposition 5.1 Let $a_{\alpha \beta}$ and $b_{\alpha \beta}$ be operators of class $C^{\mathrm{u}}\left(\mathscr{H}^{m-|\beta|}, \mathscr{H}^{|\alpha|-m}\right)$ and such that the operators $D^{*} a D-z$ and $D^{*} b D-z$ are bijective maps $\mathscr{H}^{m} \rightarrow \mathscr{H}^{-m}$ for some complex $z$. Let $\Delta_{a}$ and $\Delta_{b}$ be the operators in $\mathscr{H}$ associated to $D^{*} a D$ and $D^{*} b D$ respectively. Assume that

$$
\begin{equation*}
\lim _{r \rightarrow \infty}\left\|\theta(Q / r)\left(a_{\alpha \beta}-b_{\alpha \beta}\right)\right\|_{\mathscr{H}^{m-|\beta|} \rightarrow \mathscr{H}^{|\alpha|-m}}=0 \tag{5.20}
\end{equation*}
$$

for each $\alpha$, $\beta$, where $\theta$ is a function as above. Then $\Delta_{b}$ is a compact perturbation of $\Delta_{a}$ and the operators $\Delta_{a}$ and $\Delta_{b}$ have the same essential spectrum.

Example: In the simplest case the coeffi cients $a_{\alpha \beta}$ and $b_{\alpha \beta}$ of the principal parts (i.e. $|\alpha|=|\beta|=m$ ) are functions. Then the conditions become: $a_{\alpha \beta}$ and $b_{\alpha \beta}$ belong to $L^{\infty}(X)$ and $\left|a_{\alpha \beta}(x)-b_{\alpha \beta}(x)\right| \rightarrow 0$ as $|x| \rightarrow \infty$. Of course, the assumptions on the lowest order coeffi cients are much more general.

Example: We show here that "highly oscillating potentials" do not modify the essential spectrum. If $m=1$ then the terms of order one of $D^{*} a D$ are of the form $S=$ $\sum_{k=1}^{n}\left(P_{k} v_{k}^{\prime}+v_{k}^{\prime \prime} P_{k}\right)$, where $v_{k}^{\prime} \in \mathcal{B}\left(\mathscr{H}^{1}, \mathscr{H}\right)$ and $v_{k}^{\prime \prime} \in \mathcal{B}\left(\mathscr{H}, \mathscr{H}^{-1}\right)$. Choose $v_{k} \in \mathcal{B}\left(\mathscr{H}^{1}, \mathscr{H}\right)$ symmetric in $\mathscr{H}$ and let $v_{k}^{\prime}=i v_{k}, v_{k}^{\prime \prime}=-i v_{k}$. Then $S=[i P, v] \equiv$ $\operatorname{div} v$, with natural notations, can also be thought as a term of order zero. Now assume that $v_{k}$ are bounded Borel functions and consider a similar term $T=[i P, w]$ for $D^{*} b D$. Then the condition $\left|v_{k}(x)-w_{k}(x)\right| \rightarrow 0$ as $|x| \rightarrow \infty$ suffi ces to ensure the stability of the essential spectrum. However, the difference $S-T$ could be a function which does not tend to zero at infi nity in a simple sense, being only "highly oscillating". An explicit example in the case $n=1$ is the following: a perturbation of the form $\exp (x)(1+$ $|x|)^{-1} \cos (\exp (x))$ is allowed because it is the derivative of $(1+|x|)^{-1} \sin (\exp (x))$ plus a function which tends to zero at infi nity.

In order to apply Proposition 5.1 we need that $D^{*} a D-z: \mathscr{H}^{m} \rightarrow \mathscr{H}^{-m}$ be bijective for some $z \in \mathbb{C}$, and similarly for $b$. A standard way of checking this is to require the following coercivity condition:
(C) $\left\{\begin{array}{l}\text { there are } \mu, \nu>0 \text { such that for all } u \in \mathscr{H}^{m}: \\ \sum_{|\alpha|,|\beta| \leq m} \operatorname{Re}\left\langle P^{\alpha} u, a_{\alpha \beta} P^{\beta} u\right\rangle \geq \mu\|u\|_{\mathscr{H}^{m}}^{2}-\nu\|u\|_{\mathscr{H}}^{2}\end{array}\right.$

Example: One often imposes a stronger ellipticity condition that we describe below. Observe that the coeffi cients of the highest order part of $D^{*} a D$ defi ned by
$A_{0}=\sum_{|\alpha|=|\beta|=m} P^{\alpha} a_{\alpha \beta} P^{\beta}$ are operators $a_{\alpha \beta} \in \mathcal{B}(\mathscr{H})$. Then ellipticity means:
(Ell) $\left\{\begin{array}{l}\text { there is } \mu>0 \text { such that if } u_{\alpha} \in \mathscr{H} \text { for }|\alpha|=m \text { then } \\ \sum_{|\alpha|=|\beta|=m} \operatorname{Re}\left\langle u_{\alpha}, a_{\alpha \beta} u_{\beta}\right\rangle \geq \mu \sum_{|\alpha|=m}\left\|u_{\alpha}\right\|_{\mathscr{H}}^{2} .\end{array}\right.$
But we emphasize that, our conditions on the lower order terms being quite general, e.g. the $a_{\alpha \beta}$ could be differential operators, so the terms of formally lower order could be of order $2 m$ in fact, we have to supplement the ellipticity condition (Ell) with a condition saying that the rest of the terms $A_{1}=\sum_{|\alpha|+|\beta|<2 m} P^{\alpha} a_{\alpha \beta} P^{\beta}$ is small with respect to $A_{0}$. For example, we may require the existence of some $\delta<\mu$ and $\gamma>0$ such that

$$
\begin{equation*}
\left|\sum_{|\alpha|+|\beta|<2 m} \operatorname{Re}\left\langle P^{\alpha} u, a_{\alpha \beta} P^{\beta} u\right\rangle\right| \leq \delta\|u\|_{\mathscr{H}^{m}}^{2}+\gamma\|u\|_{\mathscr{H}}^{2} \tag{5.21}
\end{equation*}
$$

This is satisfi ed if $A_{1} \mathscr{H}^{m} \subset \mathscr{H}^{-m+\theta}$ for some $\theta>0$, because for each $\varepsilon>0$ there is $c(\varepsilon)<\infty$ such that $\|u\|_{\mathscr{H}^{m-\theta}} \leq \varepsilon\|u\|_{\mathscr{H}^{m}}+c(\varepsilon)\|u\|_{\mathscr{H}}$.

Remark 5.2 If we use Theorem 3.5 in the context of this section then we get the same conditions on the coeffi cients $a_{\alpha \beta}-b_{\alpha \beta}$ of the principal part (i.e. such that $|\alpha|=|\beta|=$ $m$ ) of the operator $a-b$ but those on the lower order coeffi cients are less general. Indeed, if $s+t>0$ the space $\mathcal{B}_{00}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{-t}\right)$ defi ned as the closure of $\mathcal{B}_{0}^{l}(\mathscr{H})$ in $\mathcal{B}\left(\mathscr{H}^{s}, \mathscr{H}^{-t}\right)$ does not contain operators of order $s+t$, while $\mathcal{B}_{0}^{l}\left(\mathscr{H}^{s}, \mathscr{H}^{-t}\right)$ contains such operators.

## 6 Weak decay preserving operators

The purpose of the next two sections is to reconsider the examples treated in Section 5 and to prove some stability results for perturbations which decay in a generalized sense, as described in Examples 2.7-2.9. This will be done in the next section, this one contains some preparatory material concerning weak decay preserving operators.

We first consider the setting of Example 2.9: $(X, \mu)$ is a positive measure such that $\mu(X)=\infty, \mathcal{F}_{\mu}$ is the fi lter of sets of co-fi nite measure, and $B_{\mu}(X)$ is the algebra of bounded measurable $\mathcal{F}_{\mu}$-vanishing functions. We recall that any direct integral of Hilbert spaces over $X$ has a canonical Hilbert module structure with $B_{\mu}(X)$ as multiplier algebra. To avoid ambiguities, we shall speak of $\mathcal{F}_{\mu}$-decay preserving operators when we refer to this algebra. Let $\{\mathscr{H}(x)\}_{x \in X}$ and $\{\mathscr{K}(x)\}_{x \in X}$ be measurable families of Hilbert spaces with dimensions $\leq N$ for some fi nite $N$. We shall use the notations introduced before Corollary 9.2.

Theorem 6.1 Let $S \in \mathcal{B}(\mathscr{H}, \mathscr{K}) \cap \mathcal{B}\left(\mathscr{H}_{p}, \mathscr{K}_{p}\right)$ for some $p \neq 2$. If $p<2$ then $S$ is left $\mathcal{F}_{\mu}$-decay preserving and if $p>2$ then $S$ is right $\mathcal{F}_{\mu}$-decay preserving

[^3]Proof: We shall consider only the case $p<2$, the assertion for $p>2$ follows by observing that $S^{*} \in \mathcal{B}(\mathscr{K}, \mathscr{H}) \cap \mathcal{B}\left(\mathscr{K}_{p^{\prime}}, \mathscr{H}_{p^{\prime}}\right)$ and then using Proposition 2.19. We prove that for each measurable set $N$ of fi nite measure the operator $T=S \chi_{N}(Q)$ has the property: if $\varepsilon>0$ then there is a Borel set $F \in \mathcal{F}_{\mu}$ such that $\left\|\chi_{F}(Q) T\right\| \leq \varepsilon$ (then Proposition 2.18 implies that $S$ is left $\mathcal{F}_{\mu}$-decay preserving). Since $N$ is of fi nite measure, $\chi_{N}(Q)$ is a bounded operator $\mathscr{H} \rightarrow \mathscr{H}_{p}$, hence $T \in \mathcal{B}\left(\mathscr{H}, \mathscr{K}_{p}\right)$. The rest of the proof is a straightforward application of Corollary 9.2. Let $a>0$ real and let $F$ be the set of points $x$ such that $|g(x)| \leq a$. Since $g \in L^{q}$ with $q<\infty$, we have $F \in \mathcal{F}_{\mu}$ and

$$
\left\|\chi_{F}(Q) T\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})}=\left\|\chi_{F}(Q) g(Q) R\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} \leq a\|R\|_{\mathscr{H}, \mathscr{K})} .
$$

Thus it suffi ces to choose $a$ such that $a\|R\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})}=\varepsilon$.
Let $X$ be a locally compact non-compact topological space and let $\mathscr{H}$ be a Hilbert $X$-module. Then, due to Remark 2.6, the operator $\varphi(Q) \in \mathcal{B}(\mathscr{H})$ is well defi ned for all $\varphi \in B(X)$. If $\mathcal{F}$ is a fi lter fi ner than the $\mathrm{Fr}^{\prime}$ echet fi lter on $X$ then

$$
\begin{equation*}
B_{\mathcal{F}}(X):=\left\{\varphi \in B(X) \mid \lim _{\mathcal{F}} \varphi=0\right\} \tag{6.22}
\end{equation*}
$$

is a $C^{*}$-algebra and we can consider on $\mathscr{H}$ the Hilbert module structure defi ned by the multiplier algebra $\mathcal{M}_{\mathcal{F}}:=\left\{\varphi(Q) \mid \varphi \in B_{\mathcal{F}}(X)\right\}$. We are interested in the corresponding classes of decay improving or decay preserving operators. To be precise, we shall speak in this context of (left or right) $\mathcal{F}$-vanishing at infinity or of (left or right) $\mathcal{F}$-decay preserving operators. Below and later on we use the notation $N^{\mathrm{c}}=X \backslash N$.
Lemma 6.2 Let $\mathscr{H}, \mathscr{K}$ be Hilbert $X$-modules. Then an operator $S \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ is left $\mathcal{F}$-decay preserving if and only if for each Borel set $N$ with $N^{\mathrm{c}} \in \mathcal{F}$ and for each $\varepsilon>0$ there is a Borel set $F \in \mathcal{F}$ such that $\left\|\chi_{F}(Q) S \chi_{N}(Q)\right\| \leq \varepsilon$.

Proof: We note first that the family of operators $\chi_{N}(Q)$, where $N$ runs over the family of Borel sets with complement in $\mathcal{F}$, is an approximate unit for $B_{\mathcal{F}}(X)$. Indeed, if $\varepsilon>0$ and $\varphi \in B_{\mathcal{F}}(X)$ then the set $N=\{x| | \varphi(x) \mid>\varepsilon\}$ is Borel, its complement is in $\mathcal{F}$, and $\sup _{x}\left|\varphi(x)\left(1-\chi_{N}(x)\right)\right| \leq \varepsilon$. Thus, according to Proposition 2.18, $S$ is left $\mathcal{F}$-decay preserving if and only if $S \chi_{N}(Q)$ is left $\mathcal{F}$-vanishing at infi nity for each $N$. Now the result follows from (2.6).

The main restriction we have to impose on $\mathcal{F}$ comes from the fact that the Friedrichs couple $(\mathscr{G}, \mathscr{H})$ which is involved in our abstract compactness criteria must be such that $\varphi(Q) \in \mathcal{K}(\mathscr{G}, \mathscr{H})$ if $\varphi \in B_{\mathcal{F}}(X)$. Sometimes this can be stated quite explicitly:

Lemma 6.3 Let $X$ be an Euclidean space, $\mathscr{H}=L^{2}(X)$, and let $\mathscr{G}=\mathscr{H}^{s}$ be a Sobolev space of order $s>0$. If $\varphi \in B(X)$ then $\varphi(Q) \in \mathcal{K}(\mathscr{G}, \mathscr{H})$ if and only if

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \int_{|x-a| \leq 1}|\varphi(x)| d x=0 \tag{6.23}
\end{equation*}
$$

The importance of such a condition in questions of stability of the essential spectrum has been noticed in $[\mathrm{He}, \mathrm{LV}, \mathrm{OS}, \mathrm{We}]$. That it is a natural condition follows also from the characterizations that we shall give below in a more general context.

Let $X$ be a locally compact non-compact abelian group. We shall say that a function $\varphi \in B(X)$ is weakly vanishing (at infinity) if

$$
\begin{equation*}
\lim _{a \rightarrow \infty} \int_{a+K}|\varphi(x)| d x=0 \text { for each compact set } K \tag{6.24}
\end{equation*}
$$

We shall denote by $B_{\mathrm{w}}(X)$ the set of functions $\varphi$ satisfying (6.24). This is clearly a $C^{*}$-algebra. Note that it suffi ces that the convergence condition in (6.24) be satisfi ed for only one compact set $K$ with non-empty interior.

Let us now express the condition (6.24) in terms of convergence to zero along a fi lter. We denote $|K|$ the exterior (Haar) measure of a set $K \subset X$ and we set $K_{a}=$ $a+K$ if $a \in X$. A subset $N$ is called w-small (at infinity) if there is a compact neighborhood $K$ of the origin such that $\lim _{a \rightarrow \infty}\left|N \cap K_{a}\right|=0$. The complement of a w-small set will be called w-large (at infinity). The family $\mathcal{F}_{\mathrm{w}}$ of all w-large sets is clearly a fi lter on $X$ fi ner than the Fr'echet fi lter.

We give now a characterization of weakly vanishing functions in terms of compactness properties. This characterization implies that of Lemma 6.3 if $X=\mathbb{R}^{n}$. Observe that a Borel set is w -small if and only if its characteristic function weakly vanishes at infi nity. Denote $f * g$ the convolution of two functions on $X$.

Lemma 6.4 For a function $\varphi \in B(X)$ the following conditions are equivalent: (1) $\varphi$ is weakly vanishing; (2) $\theta *|\varphi| \in C_{0}(X)$ if $\theta \in C_{\mathrm{c}}(X)$; (3) $\lim _{\mathcal{F}_{\mathrm{w}}} \varphi=0$; (4) $\varphi(Q) \psi(P)$ is a compact operator on $L^{2}(X)$ for all $\psi \in C_{0}(X)$.

Proof: The equivalence of (1) and (2) is clear because $\int_{K_{a}}|\varphi| d x=\left(\chi_{K} *|\varphi|\right)(a)$. Then (3) means that for each $\varepsilon>0$ the Borel set $N$ where $|\varphi(x)|>\varepsilon$ is w-small. Since $\chi_{N} \leq \varphi / \varepsilon$, the implication (2) $\Rightarrow(3)$ is clear, while the reciprocal implication follows from $\chi_{K} *|\varphi| \leq \sup |\varphi| \chi_{K} * \chi_{N}+\varepsilon|K|$. If (4) holds, let us choose $\psi$ such that its Fourier transform $\widehat{\psi}$ be a positive function in $C_{\mathrm{c}}(X)$ and let $f \in C_{\mathrm{c}}(X)$ be positive and not zero. Since $\psi(P) f$ is essentially the convolution of $\widehat{\psi}$ with $f$, there is a compact set $K$ with non-empty interior such that $\psi(P) f \geq c \chi_{K}$ with a number $c>0$. Let $U_{a}$ be the unitary operator of translation by $a$ in $L^{2}(X)$, then $U_{a} f \rightarrow 0$ weakly when $a \rightarrow \infty$, hence $\left\|\varphi(Q) U_{a} \psi(P) f\right\|=\left\|\varphi(Q) \psi(P) U_{a} f\right\| \rightarrow 0$. Since $U_{a}^{*} \varphi(Q) U_{a}=\varphi(Q-a)$ we get $\left\|\varphi(Q-a) \chi_{K}\right\| \rightarrow 0$, hence (1) holds.

Finally, let us prove that (1) $\Rightarrow$ (4). It suffi ces to prove that $\varphi(Q) \psi(P)$ is compact if $\widehat{\psi} \in C_{\mathrm{c}}(X)$ and for this it suffi ces that $\bar{\psi}(P)|\varphi|^{2}(Q) \psi(P)$ be compact. Since $\xi:=$ $|\varphi|^{2} \in B_{\mathrm{w}}(X)$ and since $\psi(P)$ is the operator of convolution by a function $\theta \in C_{\mathrm{c}}(X)$, we are reduced to proving that the integral operator $S$ with kernel $S(x, y)=\int \bar{\theta}(z-$ $x) \xi(z) \theta(z-y) d z$ is compact. If $K=\operatorname{supp} \theta$ and $\Lambda$ is the compact set $K-K$, then clearly there is a number $C$ such that

$$
|S(x, y)| \leq C \int_{K_{x}} \xi(z) d z \chi_{\Lambda}(x-y) \equiv \phi(x) \chi_{\Lambda}(x-y)
$$

where $\phi \in C_{0}(X)$. The last term here is a kernel which defi nes a compact operator $T$. Thus $\eta(Q) S$ is a Hilbert-Schmidt operator for each $\eta \in C_{\mathrm{c}}(X)$ and from the preceding estimate we get $\|(S-\eta(Q) S) u\| \leq\|(1-\eta(Q)) T|u|\|$ for each $u \in L^{2}(X)$. Thus
$\|S-\eta(Q) S\| \leq\|(1-\eta(Q)) T\|$ and the right hand side tends to zero if $\eta \equiv \eta_{\alpha}$ is an approximate unit for $C_{0}(X)$.

We shall consider now a general class of filters defi ned in terms of the metric and measure space structure. We consider only the case of an Euclidean space $X$, the extension to the case of locally compact groups or metric spaces being obvious. We set $B_{a}(r)=\{x \in X| | x-a \mid<r\}, B_{a}=B_{a}(1)$ and $B(r)=B_{0}(r)$. To each function $\nu: X \rightarrow] 0, \infty\left[\right.$ such that $\liminf _{a \rightarrow \infty} \nu(a)=0$ we associate a set of subsets of $X$ as follows:

$$
\begin{equation*}
\mathscr{N}_{\nu}=\left\{N \subset X\left|\limsup _{a \rightarrow \infty} \nu(a)^{-1}\right| N \cap B_{a} \mid<\infty\right\} \tag{6.25}
\end{equation*}
$$

Clearly $\mathscr{F}_{\nu}=\left\{F \subset X \mid F^{c} \in \mathscr{N}_{\nu}\right\}$ is a fi lter on $X$ fi ner than the Fr'echet filter.
Theorem 6.5 Let $X=\mathbb{R}^{n}$ and let $\left.\nu: X \rightarrow\right] 0, \infty\left[\right.$ such that $\liminf _{a \rightarrow \infty} \nu(a)=0$ and $\sup _{|b-a| \leq r} \nu(b) / \nu(a)<\infty$ for each real $r$. If $S \in \mathcal{B}\left(L^{2}(X)\right)$ is of class $C^{\mathrm{u}}(Q)$ and if $S \in \mathcal{B}\left(L^{p}(X)\right)$ for some $p<2$, then $S$ is left $\mathscr{F}_{\nu}$-decay preserving.

Proof: We can approximate in norm in $\mathcal{B}\left(L^{2}(X)\right)$ the operator $S$ by operators which are in $\mathcal{B}\left(L^{2}(X)\right) \cap \mathcal{B}\left(L^{p}(X)\right)$ and have fi nite range. Indeed, the approximation procedure (4.16) used in the proof of Proposition 4.11 is such that it leaves $\mathcal{B}\left(L^{2}(X)\right) \cap$ $\mathcal{B}\left(L^{p}(X)\right)$ invariant (because $V_{k}$ are isometries in $L^{p}$ too). Since the set of left $\mathscr{F}_{\nu^{-}}$ decay preserving operators is norm closed in $\mathcal{B}\left(L^{2}(X)\right)$, we may assume in the rest of the proof that $S$ is of fi nite range. According to Lemma 6.2, it suffi ces to show that, for a given Borel set $N \in \mathscr{N}_{\nu}$ and for any number $\varepsilon>0$, there is a Borel set $M \in \mathscr{N}_{\nu}$ such that $\left\|\chi_{M^{c}}(Q) S \chi_{N}(Q)\right\|<\varepsilon$.

In the rest of the proof we shall freely use the notations introduced in Section 9 (see also the proof of Proposition 4.11). In particular, $q$ is defi ned by $\frac{1}{p}=\frac{1}{2}+\frac{1}{q}$. If $f \in L^{2}(X)$ we have

$$
\left\|\chi_{N} f\right\|_{L^{p}\left(K_{a}\right)} \leq\left\|\chi_{N}\right\|_{L^{q}\left(K_{a}\right)}\|f\|_{L^{2}\left(K_{a}\right)} \leq\left|N \cap K_{a}\right|^{1 / q}\|f\|_{L^{2}\left(K_{a}\right)}
$$

Since $N \in \mathscr{N}_{\nu}$ we can find a constant $c$ such that $\left|N \cap K_{a}\right| \leq c \nu(a)$ (note that the defi nition (6.25) does not involve the restriction of $\nu$ to bounded sets). Thus, if we take $\lambda_{a}=\nu(a)^{-1 / q}$ for $a \in Z \equiv \mathbb{Z}^{n}$, we get $\chi_{N} f \in \mathscr{L}$ with the notations of Section 9. In other terms, we see that we have $\chi_{N}(Q) \in \mathcal{B}\left(L^{2}(X), \mathscr{L}\right)$. Let $T=S \chi_{N}(Q)$ and let us assume that we also have $S \in \mathcal{B}(\mathscr{L})$. Then $T \in \mathcal{B}\left(L^{2}(X), \mathscr{L}\right)$ and we can apply the Maurey type factorization theorem Theorem 9.7, where $\mathscr{H}=L^{2}(X)$. Thus we can write $T=g(Q) R$ for some $R \in \mathcal{B}\left(L^{2}(X)\right)$ and some function $g \in \mathscr{M}$, which means that $G:=\sup _{a \in Z} \nu(a)^{-1 / q}\|g\|_{L^{q}\left(K_{a}\right)}$ is a fi nite number. If $t>0$ and $M=\{x \mid g(x)>t\}$ then we get for all $a \in Z$ :

$$
\left|M \cap K_{a}\right|=\left\|\chi_{M}\right\|_{L^{q}\left(K_{a}\right)}^{q} \leq\|g / t\|_{L^{q}\left(K_{a}\right)}^{q} \leq(G / t)^{q} \nu(a) .
$$

Note that the second condition imposed on $\nu$ in Theorem 6.5 can be stated as follows: there is an increasing strictly positive function $\delta$ on $[0, \infty[$ such that $\nu(b) \leq$ $\delta(|b-a|) \nu(a)$ for all $a, b$. Indeed, we may take $\delta(r)=\sup _{|b-a| \leq r} \nu(b) / \nu(a)$. Now
let $a \in X$ and let $D(a)$ be the set of $b \in Z$ such that $K_{b}$ intersects $B_{a}$. Clearly $D(a)$ contains at most $2^{n}$ points $b$ all of them satisfying $|b-a| \leq \sqrt{n}+1$. Hence:

$$
\left|M \cap K_{a}\right| \leq \sum_{b \in D(a)}\left|M \cap K_{b}\right| \leq 2^{n} \sup _{b \in D(a)}(G / t)^{q} \nu(b) \leq 2^{n}(G / t)^{q} \delta(\sqrt{n}+1) \nu(a),
$$

which proves that $M$ belongs to $\mathscr{N}_{\nu}$. On the other hand, we have:

$$
\left\|\chi_{M^{\mathrm{c}}}(Q) T\right\|=\left\|\chi_{M^{\mathrm{c}}}(Q) g(Q) R\right\| \leq\left\|\chi_{M^{\mathrm{c}}} g\right\|_{L^{\infty}}\|R\| \leq t\|R\|
$$

To fi nish the proof of the theorem it suffi ces to take $t=\varepsilon /\|R\|$.
We still have to prove that $S \in \mathcal{B}(\mathscr{L})$. Since $S$ is of fi nite range, there is a number $r$ such that $\chi_{a}(Q) \chi_{b}(Q)=0$ if $|a-b| \geq r$. Then for any $f \in \mathscr{L}$ :

$$
\sum_{a} \lambda_{a}^{2}\left\|\chi_{a} S f\right\|_{L^{p}}^{2}=\sum_{a} \lambda_{a}^{2}\left\|\sum_{|b-a|<r} \chi_{a} S \chi_{b} f\right\|_{L^{p}}^{2} \leq C \sum_{|b-a|<r} \lambda_{a}^{2}\left\|\chi_{a} S \chi_{b} f\right\|_{L^{p}}^{2}
$$

where $C$ is a number depending only on $r$ and $n$. Since $S$ is bounded in $L^{p}$ the last term is less than $C C^{\prime} \sum_{|b-a|<r} \lambda_{a}^{2}\left\|\chi_{b} f\right\|_{L^{p}}^{2}$ for some constant $C^{\prime}$. Finally, from $\nu(b) \leq$ $\delta(|b-a|) \nu(a) \leq \delta(r) \nu(a)$ we get

$$
\sum_{|a-b|<r} \lambda_{a}^{2}=\sum_{|a-b|<r} \nu(a)^{-2 / q} \leq L(r) \delta(r)^{2 / q} \lambda_{b}^{2}
$$

where $L(r)$ is the maximum number of points from $Z$ inside a ball of radius $r$. Thus we have $\|S\|_{\mathcal{B}(\mathscr{L})}^{2} \leq C C^{\prime} L(r) \delta(r)^{2 / q}$.

Theorem 6.6 Let $X=\mathbb{R}^{n}$ and let $S$ be a pseudo-differential operator of class $S^{0}$. Then $S$ is $\mathcal{F}_{\mathrm{w}}$-decay preserving in $L^{2}(X)$, i.e. if $\varphi \in B_{\mathrm{w}}(X)$ then $\varphi(Q) S=T_{1} \psi_{1}(Q)$ and $S \varphi(Q)=\psi_{2}(Q) T_{2}$ for some $\psi_{1}, \psi_{2} \in B_{\mathrm{w}}(X)$ and $T_{1}, T_{2} \in \mathcal{B}\left(L^{2}(X)\right)$.

Proof: Since the adjoint of $S$ is also a pseudo-differential operator of class $S^{0}$, it suffi ces to show that $S$ is left $\mathcal{F}_{\mathrm{w}}$-decay preserving. We have $S \in \mathcal{B}\left(L^{p}(X)\right)$ for all $1<p<\infty$ and $S$ is of class $C^{\mathrm{u}}(Q)$ because the commutators $\left[Q_{j}, S\right]$ are bounded operators for all $1 \leq j \leq n$. Thus we can apply Theorem 6.5 and deduce that for any function $\nu$ as in the statement of the theorem, for any $\varepsilon>0$, and for any $N \in \mathscr{N}_{\nu}$ there is $M \in \mathscr{N}_{\nu}$ such that $\left\|\chi_{M^{c}}(Q) S \chi_{N}(Q)\right\| \leq \varepsilon$. Now let $N$ be a Borel w-small set, i.e. such that $\left|N \cap B_{a}\right| \rightarrow 0$ if $a \rightarrow \infty$. We shall prove that there is a function $\nu$ with the properties required in Theorem 6.5 and with $\lim _{a \rightarrow \infty} \nu(a)=0$ such that $N \in \mathscr{N}_{\nu}$. This fi nishes the proof of the corollary because the relation $M \in \mathscr{N}_{\nu}$ implies now that $M$ is w-small.

We construct $\nu$ as follows. The relation $\theta(r)=\sup _{|a| \geq r}\left|N \cap B_{a}\right|$ defi nes a positive decreasing function on $\left[0, \infty\left[\right.\right.$ which tends to zero at infi nity and such that $\left|N \cap B_{a}\right| \leq$ $\theta(|a|)$ for all $a \in X$. We set $\xi(t)=\theta(0)$ if $0 \leq t<1$ and for $k \geq 0$ integer and $2^{k} \leq t<2^{k+1}$ we defi ne $\xi(t)=\max \left\{\xi\left(2^{k-1}\right) / 2, \theta\left(2^{k}\right)\right\}$. So $\xi$ is a strictly positive decreasing function on $[0, \infty[$ which tends to zero at infi nity and such that $\theta \leq \xi$. Moreover, if $2^{k} \leq s<2^{k+1}$ and $2^{k+p} \leq t<2^{k+p+1}$ then

$$
\xi(t)=\xi\left(2^{k+p}\right) \geq \xi\left(2^{k+p-1}\right) / 2 \geq \ldots \geq 2^{-p} \xi\left(2^{k}\right)=2^{-p} \xi(s)
$$

hence $\xi(s) \geq \xi(t) \geq \frac{s}{2 t} \xi(s)$ if $1 \leq s \leq t$. We take $\nu(a)=\xi(|a|)$, so $\nu$ is a bounded strictly positive function on $X$ with $\lim _{a \rightarrow \infty} \nu(a)=0$ and $\left|N \cap B_{a}\right| \leq \nu(a)$ for all $a$. If $a, b$ are points with $|a|,|b| \geq 1$ and $|a-b| \leq r$ then $\nu(b) / \nu(a) \leq 1$ if $|a| \leq|b|$ and if $|a|>|b|$ then

$$
\frac{\nu(b)}{\nu(a)}=\frac{\xi(|b|)}{\xi(|a|)} \leq \frac{2|a|}{|b|} \leq 2(1+r)
$$

Thus the second condition imposed on $\nu$ in Theorem 6.5 is also satisfi ed.
As a fi nal example, we introduce now classes of vanishing at infi nity functions of a more topological nature. Let us fi x a uniformly discrete set $L \subset X$, i.e. a set such that inf $|a-b|>0$ where the infi mum is taken over couples of distinct points $a, b \in L$. Let $L_{\varepsilon}=L+B(\varepsilon)$ be the set of points at distance $<\varepsilon$ from $L$. We say that a subset $N \subset X$ is $L$-thin if for each $\varepsilon>0$ there is $r<\infty$ such that $N \backslash B(r) \subset L_{\varepsilon}$. In other terms, $N$ is $L$-thin if there is a family $\left\{\delta_{a}\right\}_{a \in L}$ of positive real numbers with $\delta_{a} \rightarrow 0$ as $a \rightarrow \infty$ such that $N \subset \bigcup B_{a}\left(\delta_{a}\right)$. The complement of such a set will be called $L$-fat. We denote $\mathcal{F}_{L}$ the family of $L$-fat sets, we note that $\mathcal{F}_{L}$ is a filter on $X$ contained in $\mathcal{F}_{\mathrm{w}}$ and fi ner than the $\mathrm{Fr}{ }^{\prime}$ echet filter, and we denote $B(X)$ the set of bounded Borel functions such that $\lim _{\mathcal{F}_{L}} \varphi=0$. So $\varphi \in B(X)$ belongs to $B_{L}(X)$ if and only if the set $\{|\varphi| \geq \lambda\}$ is $L$-thin for each $\lambda>0$.

Proposition 6.7 Let $X=\mathbb{R}^{n}$ and let $S$ be a bounded operator on $L^{2}(X)$ such that on the region $x \neq y$ its distribution kernel is a function satisfying the estimate $|S(x, y)| \leq$ $c|x-y|^{-m}$ for some $m>n$. Then $S$ is $\mathcal{F}_{L}$-decay preserving.

Proof: Let $\theta \in C_{\mathrm{b}}(X)$ such that $\theta(x)=0$ on a neighborhood of the origin and $S_{\theta}(x, y)=\theta(x-y) S(x, y)$. If $\xi(x)=\theta(x)|x|^{-m}$ then for the operator $S_{\theta}$ of kernel $S_{\theta}(x, y)$ we have $\left\|S_{\theta} u\right\| \leq c\|\xi *|u|\|$ hence $\left\|S_{\theta}\right\| \leq c\|\xi\|_{L^{1}}$ By choosing a convenient sequence of functions $\theta$ we see that $S$ is the norm limit of a sequence of operators which besides the properties from the statement of the proposition are such that $S(x, y)=0$ if $|x-y|>R(S)$. Since the set of $\mathcal{F}_{L}$-decay preserving operators is closed in norm (see Subsection 2.3), we may assume in the rest of the proof that the kernel of $S$ has this property. In fact, in order to simplify the notations and without loss of generality, we shall assume $S(x, y)=0$ if $|x-y|>1$.

Let $N$ be an $L$-thin Borel set and let $\varepsilon>0$. We shall construct an $L$-fat Borel set with $F \subset N^{\mathrm{c}}$ such that $\left\|\chi_{N}(Q) S \chi_{F}(Q)\right\| \leq \varepsilon$. Since the adjoint operator $S^{*}$ has the same properties as $S$, this suffi ces to prove that it is decay preserving.

We shall only need two simple estimates. First, if $\rho_{x}(G)$ is the distance from a Borel set $G$ to a point $x$, then

$$
\begin{equation*}
\int_{G} \frac{d y}{|x-y|^{2 m}} \leq C(m, n) \rho_{x}(G)^{n-2 m} \tag{6.26}
\end{equation*}
$$

Then, if $B_{0}, B$ are two balls with the same center and radiuses $\delta$ and $\delta+\varepsilon$, then

$$
\begin{equation*}
\int_{B_{0}} \rho_{x}\left(B^{\mathrm{c}}\right)^{n-2 m} d x \leq C(m, n) \varepsilon^{n-2 m} \delta^{n} \tag{6.27}
\end{equation*}
$$

We shall choose $\varepsilon=\delta^{n / 2 m}$. Then $\chi_{B_{0}}(Q) S \chi_{B^{c}}(Q)$ is an operator with integral kernel and we can estimate its Hilbert-Schmidt norm as follows:

$$
\begin{align*}
\left\|\chi_{B_{0}}(Q) S \chi_{B^{\mathrm{c}}}(Q)\right\|_{H S}^{2} & =\int_{X \times X} \chi_{B_{0}}(x)|S(x, y)|^{2} \chi_{B^{\mathrm{c}}}(y) d x d y \\
& \leq c \int_{B_{0}} d x \int_{B^{\mathrm{c}}} \frac{d y}{|x-y|^{2 m}} \leq C \int_{B_{0}} \rho_{x}\left(B^{\mathrm{c}}\right)^{n-2 m} d x \\
& \leq C^{\prime} \varepsilon^{n-2 m} \delta^{n}=C^{\prime} \delta^{\lambda} \tag{6.28}
\end{align*}
$$

where $\lambda=n^{2} / 2 m>0$.
We can assume that $N=\bigcup_{a} B_{a}\left(\delta_{a}\right)$, where the sequence of numbers $\delta_{a}$ satisfi es $\delta_{a} \rightarrow 0$ as $a \rightarrow \infty$. Denote $N_{a}=B_{a}\left(\delta_{a}\right)$ and $M_{a}=B_{a}\left(\delta_{a}+\varepsilon_{a}\right)$, where we choose $\varepsilon_{a}=\delta_{a}^{n / 2 m}$ as above. Choose $r$ such that the balls $N_{a}$ are pairwise disjoint and $\delta_{a}+\varepsilon_{a}<1$ if $|a|>r$ and let $R$ such that $\chi_{N_{a}}(Q) S \chi_{B(R)^{\mathrm{c}}}(Q)=0$ if $|a| \leq r$. Let $M=\bigcup M_{a}$ and $F=M^{\text {c }} \backslash B(R)$, so that $F$ is a closed $L$-fat set. Then for any $u \in L^{2}(X)$ we have:

$$
\left\|\chi_{N}(Q) S \chi_{F}(Q) u\right\|^{2}=\sum_{|a|>r}\left\|\chi_{N_{a}}(Q) S \chi_{F}(Q) u\right\|^{2}
$$

Since $S$ is of range 1 we have $\chi_{N_{a}}(Q) S \chi_{B_{a}(2)^{c}}(Q)=0$ if $\delta_{a}<1$. Thus

$$
\left\|\chi_{N}(Q) S \chi_{F}(Q) u\right\|^{2} \leq \sum_{|a|>r}\left\|\chi_{N_{a}}(Q) S \chi_{F \cap B_{a}(2)}(Q)\right\|^{2}\left\|\chi_{B_{a}(2)}(Q) u\right\|^{2}
$$

The number of $b \in L$ such that $B_{b}(2)$ meets $B_{a}(2)$ is a bounded function of $a$, hence there is a constant $C$ depending only on $L$ such that

$$
\left\|\chi_{N}(Q) S \chi_{F}(Q) u\right\| \leq C \sup _{|a|>r}\left\|\chi_{N_{a}}(Q) S \chi_{F \cap B_{a}(2)}(Q)\right\|\|u\|
$$

We have $F \subset M^{\mathrm{c}} \subset M_{a}^{\mathrm{c}}$ hence

$$
\left\|\chi_{N_{a}}(Q) S \chi_{F \cap B_{a}(2)}(Q)\right\| \leq\left\|\chi_{N_{a}}(Q) S \chi_{M_{a}^{c}}(Q)\right\|_{H S} \leq C^{\prime} \delta_{a}^{\lambda / 2}
$$

because of (6.28). So the norm $\left\|\chi_{N}(Q) S \chi_{F}(Q)\right\|$ can be made as small as we wish by choosing $r$ large enough.

Corollary 6.8 Let $X=\mathbb{R}^{n}$, $\mu$ the Lebesgue measure, and $L$ a uniformly discrete subset of $\mathbb{R}^{n}$. Then a pseudo-differential operator of class $S^{0}$ on $L^{2}(X)$ is decay preserving with respect to $\mathcal{F}_{\mu}, \mathcal{F}_{\mathrm{w}}$ and $\mathcal{F}_{L}$.

Proof: In the first case we use Theorem 6.1 by taking into account that a pseudodifferential operator of class $S^{0}$ belongs to $\mathcal{B}\left(L^{p}(X)\right)$ for all $1<p<\infty$ and that the adjoint of such an operator is also pseudo-differential of class $S^{0}$. The second case has already been considered in Theorem 6.6. For the third case, note that the distribution kernel of such an operator verifi es the estimates $|S(x, y)| \leq C_{k}|x-y|^{-n}(1+|x-y|)^{-k}$ for any $k>0$, see [Ho].

## 7 Weakly vanishing perturbations

In this subsection we reconsider the framework of Subsection 5 and improve, but with a stronger assumption $a \in \mathcal{B}(\mathscr{K})$, the decay condition (5.20). We shall consider on $\mathscr{H}$ the class of "vanishing at infi nity" functions corresponding to the algebra $B_{\mathrm{w}}(X)$, in other terms we equip $\mathscr{H}$ with the Hilbert module structure associated to the multiplier algebra $\left\{\varphi(Q) \mid \varphi \in B_{\mathrm{w}}(X)\right\}$. By Lemma 6.3, $(\mathscr{G}, \mathscr{H})$ remains a compact Friedrichs module. The space $\mathscr{K}$ inherits a natural direct sum Hilbert module structure.

We keep the notations and terminology of Sections 5 and 6 . We recall that an operator $D^{*} a D: \mathscr{H}^{m} \rightarrow \mathscr{H}^{-m}$ is coercive if there are numbers $\mu, \nu>0$ such that

$$
\begin{equation*}
\operatorname{Re}\langle D u, a D u\rangle \geq \mu\|u\|_{\mathscr{H}^{m}}^{2}-\nu\|u\|_{\mathscr{H}}^{2} \quad \forall u \in \mathscr{H}^{m} . \tag{7.29}
\end{equation*}
$$

Clearly the next lemma remains true if the filter $\mathcal{F}_{\mathrm{w}}$ is replaced by $\mathcal{F}_{\mu}$ or $\mathcal{F}_{L}$.
Lemma 7.1 Assume that $a \in \mathcal{B}(\mathscr{K})$ is $\mathcal{F}_{\mathrm{w}}$-decay preserving and that the operator $D^{*} a D: \mathscr{H}^{m} \rightarrow \mathscr{H}^{-m}$ is coercive. Then $D\left(\Delta_{a}-z\right)^{-1}$ is $\mathcal{F}_{\mathrm{w}}$-decay preserving if $\operatorname{Re} z \leq-\nu$, where $\nu$ is as in (7.29).

Proof: We shall use Proposition 3.7 with $c$ the the identity operator in $\mathscr{K}$, so $\Delta \equiv \Delta_{c}$ is the operator in $\mathscr{H}$ associated to $D^{*} D=\sum_{|\alpha| \leq m} P^{2 \alpha}$, which is the canonical (Riesz) positive isomorphism of $\mathscr{G}$ onto $\mathscr{G}^{*}$ and (7.29) means $\operatorname{Re} D^{*} a D \geq \mu D^{*} D-\nu$. We have $D(\Delta-z)^{-1} \in \mathcal{B}_{q}(\mathscr{H}, \mathscr{K})$ and $D\left(D^{*} D-z\right)^{-1} D^{*} \in \mathcal{B}_{q}(\mathscr{K})$ if $\operatorname{Re} z<0$ because these operators consist of matrices of pseudo-differential operators with constant coeffi cients of class $S^{0}$, so we can use Theorem 6.6.

We now consider two operators $\mathscr{H}^{m} \rightarrow \mathscr{H}^{-m}$ of the form

$$
D^{*} a D=\sum_{|\alpha|,|\beta| \leq m} P^{\alpha} a_{\alpha \beta} P^{\beta} \quad \text { and } \quad D^{*} b D=\sum_{|\alpha|,|\beta| \leq m} P^{\alpha} b_{\alpha \beta} P^{\beta}
$$

where the coeffi cients are continuous operators $a_{\alpha \beta}, b_{\alpha \beta}: \mathscr{H}^{m-|\beta|} \rightarrow \mathscr{H}^{|\alpha|-m}$ satisfying some other conditions stated below and denote as usual $\Delta_{a}$ and $\Delta_{b}$ the operators in $\mathscr{H}$ associated to them.

Theorem 7.2 Assume that the operators $D^{*} a D$ and $D^{*} b D$ are coercive and that their coefficients satisfy the following conditions: (1) $a_{\alpha \beta} \in \mathcal{B}(\mathscr{H})$ and are $\mathcal{F}_{\mathrm{w}}$-decay preserving operators; (2) if $|\alpha|+|\beta|=2 m$ then $a_{\alpha \beta}-b_{\alpha \beta}$ is left $\mathcal{F}_{\mathrm{w}}$-vanishing at infinity; (3) if $|\alpha|+|\beta|<2 m$ then $a_{\alpha \beta}-b_{\alpha \beta} \in \mathcal{K}\left(\mathscr{H}^{m-|\beta|}\right.$, $\left.\mathscr{H}^{|\alpha|-m}\right)$. Then the operator $\Delta_{b}$ is a compact perturbation of $\Delta_{a}$, in particular $\sigma_{\mathrm{ess}}\left(\Delta_{a}\right)=\sigma_{\mathrm{ess}}\left(\Delta_{b}\right)$.

Proof: We check the conditions of Theorem 3.5. Because of the coercivity assumptions, condition (1) is fulfi lled, and (3) is satisfi ed by Lemma 7.1. The part of condition (2) involving the coeffi cients such that $|\alpha|+|\beta|=2 m$ is satisfi ed by defi nition, for the lower order coeffi cients it suffi ces to use (3.10).

Remark 7.3 If $a_{\alpha \beta}$ and $b_{\alpha \beta}$ are bounded Borel functions and $a_{\alpha \beta}-b_{\alpha \beta} \in B_{\mathrm{w}}(X)$ for all $\alpha, \beta$, then the conditions (1)-(3) of the theorem are satisfi ed. Indeed, in order to check the compactness conditions on the lower order coeffi cients note that, by Lemma 6.3, if $\varphi \in B_{\mathrm{w}}(X)$ then the operator $\varphi(Q): \mathscr{H}^{s} \rightarrow \mathscr{H}^{-t}$ is compact if $s, t \geq 0$ and one of them is not zero.

The next result is a more general but less explicit version of Theorem 7.2. This is an improvement of [OS, Theorem 2.1], thus it covers some subelliptic operators.

Theorem 7.4 Assume that $D^{*} a D$ satisfies (7.29) and that $\Delta_{b}$ is a closed densely defined operator such that there is $z \in \rho\left(\Delta_{b}\right)$ with $\operatorname{Re} z \leq-\nu$. Moreover, assume that $a, b$ satisfy the conditions (1)-(3) of Theorem 7.2. Then the operator $\Delta_{b}$ is a compact perturbation of $\Delta_{a}$.

Proof: We shall apply Theorem 3.2 with $A=\Delta_{a}$ and $B=\Delta_{b}$. The assumption (AB) is clearly satisfi ed and we take $\widetilde{A}=D^{*} a D$ and $\widetilde{B}=D^{*} b D$, hence $\widetilde{B}-\widetilde{A}=$ $D^{*}(b-a) D$. Then let $S=D^{*}$ and $T=(b-a) D$.

Finally, let us note that one should be able to use Theorem 3.2 to treat situations when the coeffi cients $a_{\alpha \beta}$ and $b_{\alpha \beta}$ are unbounded operators even if $|\alpha|=\beta \mid=m$ (as in [OS, Theorem 3.1] and [Ba1, Ba2]), see the framework of Example 3.1 and Corollary 3.4, but we shall not pursue this idea here.

## 8 Riemannian manifolds

Let $\mathscr{H}, \mathscr{K}$ be two Hilbert spaces identifi ed with their adjoints and d a closed densely defi ned operator mapping $\mathscr{H}$ into $\mathscr{K}$. Let $\mathscr{G}=\mathcal{D}(\mathrm{d})$ equipped with the graph norm, so $\mathscr{G} \subset \mathscr{H}$ continuously and densely and $\mathrm{d} \in \mathcal{B}(\mathscr{G}, \mathscr{K})$.

Then the quadratic form $\|\mathrm{d} u\|_{\mathscr{K}}^{2}$ on $\mathscr{H}$ with domain $\mathscr{G}$ is positive densely defi ned and closed. Let $\Delta$ be the positive self-adjoint operator on $\mathscr{H}$ associated to it. In fact $\Delta=\mathrm{d}^{*} \mathrm{~d}$, where the adjoint $\mathrm{d}^{*}$ of d is a closed densely defi ned operator mapping $\mathscr{K}$ into $\mathscr{H}$.

Now let $\lambda \in \mathcal{B}(\mathscr{H})$ and $\Lambda \in \mathcal{B}(\mathscr{K})$ be self-adjoint and such that $\lambda \geq c$ and $\Lambda \geq c$ for some real $c>0$. Then we can defi ne new Hilbert spaces $\widetilde{\mathscr{H}}$ and $\widetilde{\mathscr{K}}$ as follows:
$(*)\left\{\begin{array}{l}\widetilde{\mathscr{H}}=\mathscr{H} \text { as vector space and }\langle u \mid v\rangle_{\widetilde{\mathscr{H}}}=\langle u \mid \lambda v\rangle_{\mathscr{H}}, \\ \widetilde{\mathscr{K}}=\mathscr{K} \text { as vector space and }\langle u \mid v\rangle_{\widetilde{K}}=\langle u \mid \Lambda v\rangle_{\mathscr{K}} .\end{array}\right.$
Since $\mathscr{H}=\widetilde{\mathscr{H}}$ and $\mathscr{K}=\widetilde{\mathscr{K}}$ as topological vector spaces, the operator $\mathrm{d}: \mathscr{G} \subset$ $\widetilde{\mathscr{H}} \rightarrow \widetilde{\mathscr{K}}$ is still a closed densely defi ned operator, hence the quadratic form $\|\mathrm{d} u\|_{\widetilde{K}}^{2}$ on $\widetilde{\mathscr{H}}$ with domain $\mathscr{G}$ is positive, densely defi ned and closed. We shall denote by $\widetilde{\Delta}$ the positive self-adjoint operator on $\widetilde{\mathscr{H}}$ associated to it.

We can express $\widetilde{\Delta}$ in more explicit terms as follows. Denote by $\widetilde{\mathrm{d}}$ the operator d when viewed as acting from $\widetilde{\mathscr{H}}$ to $\widetilde{\mathscr{K}}$. Then $\widetilde{\Delta}=\widetilde{\mathrm{d}^{*}} \widetilde{\mathrm{~d}}$, where $\widetilde{\mathrm{d}}^{*}: \mathcal{D}\left(\widetilde{\mathrm{d}^{*}}\right) \subset \widetilde{\mathscr{K}} \rightarrow$
$\widetilde{\mathscr{H}}$ is the adjoint of $\widetilde{\mathrm{d}}=\mathrm{d}$ with respect to the new Hilbert space structures (the spaces $\widetilde{\mathscr{H}}, \widetilde{\mathscr{K}}$ being also identifi ed with their adjoints). It is easy to check that $\widetilde{\mathrm{d}}{ }^{*}=\lambda^{-1} \mathrm{~d}^{*} \Lambda$. Thus $\widetilde{\Delta}=\lambda^{-1} \mathrm{~d}^{*} \Lambda \mathrm{~d}$.

Now let $(X, \rho)$ be a proper locally compact metric space (see the defi nition before Corollary 2.21) and let us assume that $\mathscr{H}$ and $\mathscr{K}$ are Hilbert $X$-modules.

Definition 8.1 A closed densely defined map $\mathrm{d}: \mathcal{D}(\mathrm{d}) \subset \mathscr{H} \rightarrow \mathscr{K}$ is a first order operator if there is $C \in \mathbb{R}$ such that for each bounded Lipschitz function $\varphi$ on $X$ the form $[\mathrm{d}, \varphi(Q)]$ is a bounded operator and $\|[\mathrm{d}, \varphi(Q)]\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} \leq C \operatorname{Lip} \varphi$.

Here

$$
\operatorname{Lip} \varphi=\inf _{x \neq y}|\varphi(x)-\varphi(y)| \rho(x, y)^{-1}
$$

In more explicit terms, we require

$$
\left|\left\langle\mathrm{d}^{*} u, \varphi(Q) v\right\rangle_{\mathscr{H}}-\langle u, \varphi(Q) \mathrm{d} v\rangle_{\mathscr{K}}\right| \leq C \operatorname{Lip} \varphi\|u\|_{\mathscr{K}}\|v\|_{\mathscr{H}}
$$

for all $u \in \mathcal{D}\left(\mathrm{~d}^{*}\right)$ and $v \in \mathcal{D}(d)$. Thus $\left\langle\mathrm{d}^{*} u, \varphi(Q) v\right\rangle-\langle u, \varphi(Q) \mathrm{d} v\rangle$ is a sesquilinear form on the dense subspace $\mathcal{D}\left(\mathrm{d}^{*}\right) \times \mathcal{D}(\mathrm{d})$ of $\mathscr{K} \times \mathscr{H}$ which is continuous for the topology induced by $\mathscr{H} \times \mathscr{K}$. Hence there is a unique continuous operator $[\mathrm{d}, \varphi(Q)]$ : $\mathscr{H} \rightarrow \mathscr{K}$ such that

$$
\left\langle\mathrm{d}^{*} u, \varphi(Q) v\right\rangle_{\mathscr{H}}-\langle u, \varphi(Q) \mathrm{d} v\rangle_{\mathscr{K}}=\langle u,[\mathrm{~d}, \varphi(Q)] v\rangle_{\mathscr{K}}
$$

for all $u \in \mathcal{D}\left(\mathrm{~d}^{*}\right), v \in \mathcal{D}(\mathrm{~d})$ and $\|[\mathrm{d}, \varphi(Q)]\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} \leq C \operatorname{Lip} \varphi$.
Lemma 8.2 The operator $\mathrm{d}(\Delta+1)^{-1}$ is decay preserving.
Proof: We shall prove that $S:=\mathrm{d}(\Delta+1)^{-1}$ is a decay preserving operator with the help of Corollary 2.21, more precisely we show that $[S, \varphi(Q)]$ is a bounded operator if $\varphi$ is a positive Lipschitz function. Let $\varepsilon>0$ and $\varphi_{\epsilon}=\varphi(1+\varepsilon \varphi)^{-1}$. Then $\varphi_{\varepsilon}$ is a bounded function with $\left|\varphi_{\varepsilon}\right| \leq \varepsilon^{-1}$ and

$$
\left|\varphi_{\varepsilon}(x)-\varphi_{\varepsilon}(y)\right|=\frac{|\varphi(x)-\varphi(y)|}{(1+\varepsilon \varphi(x))(1+\varepsilon \varphi(y))} \leq|\varphi(x)-\varphi(y)|
$$

hence $\operatorname{Lip} \varphi_{\varepsilon} \leq \operatorname{Lip} \varphi$. Let $v \in \mathcal{D}(\mathrm{~d})$ we have for all $u \in \mathcal{D}\left(\mathrm{~d}^{*}\right)$ :

$$
\begin{aligned}
\left|\left\langle\mathrm{d}^{*} u, \varphi_{\epsilon}(Q) v\right\rangle_{\mathscr{H}}\right| & =\left|\left\langle u, \varphi_{\epsilon}(Q) \mathrm{d} v\right\rangle_{\mathscr{K}}+\left\langle u,\left[\mathrm{~d}, \varphi_{\varepsilon}(Q)\right] v\right\rangle_{\mathscr{K}}\right| \\
& \leq\|u\|_{\mathscr{K}}\left(\varepsilon^{-1}\|\mathrm{~d} v\|_{\mathscr{K}}+C \operatorname{Lip} \varphi_{\varepsilon}\|u\|_{\mathscr{H}}\right) .
\end{aligned}
$$

Hence $\varphi_{\varepsilon}(Q) v \in \mathcal{D}\left(\mathrm{~d}^{* *}\right)=\mathcal{D}(\mathrm{d})$ because d is closed. Thus $\varphi_{\varepsilon}(Q) \mathcal{D}(\mathrm{d}) \subset \mathcal{D}(\mathrm{d})$ and by the closed graph theorem we get $\varphi_{\varepsilon}(Q) \in \mathcal{B}(\mathscr{G})$, where $\mathscr{G}$ is the domain of d equipped with the graph topology. This also implies that $\varphi_{\varepsilon}(Q)$ extends to an operator in $\mathcal{B}\left(\mathscr{G}^{*}\right)$ (note that $\varphi_{\varepsilon}(Q)$ is symmetric in $\mathscr{H}$ ).

Now, if we think of d as a continuous operator $\mathscr{G} \rightarrow \mathscr{K}$, then it has an adjoint $\mathrm{d}^{*}: \mathscr{K} \rightarrow \mathscr{G}^{*}$ which is the unique continuous extension of the operator $\mathrm{d}^{*}: \mathcal{D}\left(\mathrm{d}^{*}\right) \subset$ $\mathscr{K} \rightarrow \mathscr{H} \subset \mathscr{G}^{*}$. Thus the canonical extension of $\Delta$ to an element of $\mathcal{B}\left(\mathscr{G}^{( } \mathscr{G}^{*}\right)$ is
the product of d: $\mathscr{G} \rightarrow \mathscr{K}$ with $\mathrm{d}^{*}: \mathscr{K} \rightarrow \mathscr{G}^{*}$ (note $\mathcal{D}(\mathrm{d})$ is the form domain of $\Delta$ ). Then it is trivial to justify that we have in $\mathcal{B}\left(\mathscr{G}, \mathscr{G}^{*}\right)$ :

$$
\left[\Delta, \varphi_{\varepsilon}(Q)\right]=\left[\mathrm{d}^{*}, \varphi_{\varepsilon}(Q)\right] \mathrm{d}+\mathrm{d}^{*}\left[\mathrm{~d}, \varphi_{\varepsilon}(Q)\right]
$$

Here $\left[\mathrm{d}^{*}, \varphi_{\varepsilon}(Q)\right]=\left[\varphi_{\varepsilon}(Q), \mathrm{d}\right]^{*} \in \mathcal{B}(\mathscr{K}, \mathscr{H})$. Since $\Delta+1: \mathscr{G} \rightarrow \mathscr{G}^{*}$ is a linear homeomorphism, we then have in $\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)$ :

$$
\begin{aligned}
{\left[\varphi_{\varepsilon}(Q),(\Delta+1)^{-1}\right] } & =(\Delta+1)^{-1}\left[\Delta, \varphi_{\varepsilon}(Q)\right](\Delta+1)^{-1} \\
& =(\Delta+1)^{-1}\left[\varphi_{\varepsilon}(Q), \mathrm{d}\right]^{*} \mathrm{~d}(\Delta+1)^{-1} \\
& +(\Delta+1)^{-1} \mathrm{~d}^{*}\left[\mathrm{~d}, \varphi_{\varepsilon}(Q)\right](\Delta+1)^{-1}
\end{aligned}
$$

Finally, taking once again into account the fact that $\varphi_{\varepsilon}(Q)$ leaves $\mathscr{G}$ invariant, we have:

$$
\begin{aligned}
{\left[\varphi_{\varepsilon}(Q), \mathrm{d}(\Delta+1)^{-1}\right]=} & {\left[\varphi_{\varepsilon}(Q), \mathrm{d}\right](\Delta+1)^{-1} } \\
& +\mathrm{d}(\Delta+1)^{-1}\left[\varphi_{\varepsilon}(Q), \mathrm{d}\right]^{*} \mathrm{~d}(\Delta+1)^{-1} \\
& +\mathrm{d}(\Delta+1)^{-1} \mathrm{~d}^{*}\left[\mathrm{~d}, \varphi_{\varepsilon}(Q)\right](\Delta+1)^{-1}
\end{aligned}
$$

Hence:

$$
\begin{aligned}
& \left\|\left[\varphi_{\varepsilon}(Q), \mathrm{d}(\Delta+1)^{-1}\right]\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} \leq\left\|\left[\varphi_{\varepsilon}(Q), \mathrm{d}\right]\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})}\left\|(\Delta+1)^{-1}\right\|_{\mathcal{B}(\mathscr{H})} \\
& \quad+\left\|\mathrm{d}(\Delta+1)^{-1}\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})}\left\|\left[\varphi_{\varepsilon}(Q), \mathrm{d}\right]^{*}\right\|_{\mathcal{B}(\mathscr{K}, \mathscr{H})}\left\|\mathrm{d}(\Delta+1)^{-1}\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} \\
& \quad+\left\|\mathrm{d}(\Delta+1)^{-1} \mathrm{~d}^{*}\right\|_{\mathcal{B}(\mathscr{K}, \mathscr{K})}\left\|\left[\mathrm{d}, \varphi_{\varepsilon}(Q)\right]\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})}\left\|(\Delta+1)^{-1}\right\|_{\mathcal{B}(\mathscr{H})} .
\end{aligned}
$$

The most singular factor here is

$$
\left\|\mathrm{d}(\Delta+1)^{-1} \mathrm{~d}^{*}\right\|_{\mathcal{B}(\mathscr{K}, \mathscr{K})} \leq\|\mathrm{d}\|_{\mathcal{B}(\mathscr{G}, \mathscr{K})}\left\|(\Delta+1)^{-1}\right\|_{\mathcal{B}\left(\mathscr{G}^{*}, \mathscr{G}\right)}\left\|\mathrm{d}^{*}\right\|_{\mathcal{B}\left(\mathscr{K}, \mathscr{G}^{*}\right)}
$$

and this is fi nite. Thus we get for a fi nite constant $G$ :

$$
\begin{aligned}
\left\|\left[\varphi_{\varepsilon}(Q), \mathrm{d}(\Delta+1)^{-1}\right]\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} & \leq C_{1}\left\|\left[\mathrm{~d}, \varphi_{\varepsilon}(Q)\right]\right\|_{\mathcal{B}(\mathscr{H}, \mathscr{K})} \\
& \leq C_{1} C \operatorname{Lip} \varphi_{\varepsilon} \leq C_{1} C \operatorname{Lip} \varphi
\end{aligned}
$$

Now let $u \in \mathscr{K}_{\mathrm{c}}$ and $v \in \mathscr{H}_{\mathrm{c}}$. We get:

$$
\begin{aligned}
& \left|\left\langle\varphi(Q) u, \mathrm{~d}(\Delta+1)^{-1} v\right\rangle-\left\langle u, \mathrm{~d}(\Delta+1)^{-1} \varphi(Q) v\right\rangle\right|= \\
& \quad=\lim _{\varepsilon \rightarrow 0}\left|\left\langle\varphi_{\varepsilon}(Q) u, \mathrm{~d}(\Delta+1)^{-1} v\right\rangle-\left\langle u, \mathrm{~d}(\Delta+1)^{-1} \varphi_{\varepsilon}(Q) v\right\rangle\right| \\
& \quad \leq C_{1} C \operatorname{Lip} \varphi
\end{aligned}
$$

Thus $\left[\varphi(Q), \mathrm{d}(\Delta+1)^{-1}\right]$ is a bounded operator.
Theorem 8.3 Let $(X, \rho)$ be a proper locally compact metric space. Assume that $(\mathscr{G}, \mathscr{H})$ is a compact Friedrichs $X$-module and that $\mathscr{K}$ is a Hilbert $X$-module. Let $\mathrm{d}, \lambda, \Lambda$ be operators satisfying the following conditions:
(i) d is a closed first order operator from $\mathscr{H}$ to $\mathscr{K}$ with $\mathcal{D}(\mathrm{d})=\mathscr{G}$;
(ii) $\lambda$ is a bounded self-adjoint operator on $\mathscr{H}$ with $\inf \lambda>0$ and such that $\lambda-1 \in$ $\mathcal{K}(\mathscr{G}, \mathscr{H})\left(\right.$ e.g. $\left.\lambda-1 \in \mathcal{B}_{0}(\mathscr{H})\right)$;
(iii) $\Lambda$ is a bounded self-adjoint operator on $\mathscr{K}$ with $\inf \Lambda>0$ and such that $\Lambda-1 \in$ $\mathcal{B}_{0}(\mathscr{K})$.
Then the self-adjoint operators $\Delta$ and $\widetilde{\Delta}$ have the same essential spectrum.

Proof: In this proof, we shall consider $\widetilde{\Delta}$ as an operator acting on $\mathscr{H}$. Since $\widetilde{\mathscr{H}}=\mathscr{\sim}$ as topological vector spaces and the notion of spectrum is purely topological, $\widetilde{\Delta}$ is a closed densely defi ned operator on $\mathscr{H}$ and it has the same spectrum as the self-adjoint $\widetilde{\Delta}$ on $\widetilde{\mathscr{H}}$. Moreover, if we defi ne the essential spectrum $\sigma_{\text {ess }}(A)$ as the set of $z \in \mathbb{C}$ such that either $\operatorname{ker}(A-z)$ is infi nite dimensional or the range of $A-z$ is not closed, we see that the essential spectrum is a topological notion, so $\sigma_{\text {ess }}(\widetilde{\Delta})$ is the same, whether we think of $\widetilde{\Delta}$ as operator on $\mathscr{H}$ or on $\widetilde{\mathscr{H}}$. Finally, with this defi nition of $\sigma_{\text {ess }}$ we have $\sigma_{\text {ess }}(A)=\sigma_{\text {ess }}(B)$ if $(A-z)^{-1}-(B-z)^{-1}$ is a compact operator for some $z \in \rho(A) \cap \rho(B)$.

Thus it suffi ces to prove that $(\Delta+1)^{-1}-(\widetilde{\Delta}+1)^{-1} \in \mathcal{K}(\mathscr{H})$. Now we observe that

$$
\widetilde{\Delta}+1=\lambda^{-1} d^{*} \Lambda d+1=\lambda^{-1}\left(\mathrm{~d}^{*} \Lambda \mathrm{~d}+\lambda\right)
$$

and $\Delta_{\Lambda}=\mathrm{d}^{*} \Lambda \mathrm{~d}$ is the positive self-adjoint operator on $\mathscr{H}$ associated to the closed quadratic form $\|\mathrm{d} u\|_{\widetilde{K}}^{2}$ on $\mathscr{H}$ with domain $\mathscr{G}$. Thus $(\widetilde{\Delta}+1)^{-1}=\left(\Delta_{\Lambda}+\lambda\right)^{-1} \lambda$ and

$$
(\widetilde{\Delta}+1)^{-1}-\left(\Delta_{\Lambda}+\lambda\right)^{-1}=\left(\Delta_{\Lambda}+\lambda\right)^{-1}(\lambda-1)=\left[(\lambda-1)\left(\Delta_{\Lambda}+\lambda\right)^{-1}\right]^{*}
$$

The range of $\left(\Delta_{\Lambda}+\lambda\right)^{-1}$ is included in the form domain of $\Delta_{\Lambda}+\lambda$, which is $\mathscr{G}$. The map $\left(\Delta_{\Lambda}+\lambda\right)^{-1}: \mathscr{H} \rightarrow \mathscr{G}$ is continuous, by the closed graph theorem, and $\lambda-1: \mathscr{G} \rightarrow \mathscr{H}$ is compact. Hence $(\widetilde{\Delta}+1)^{-1}-\left(\Delta_{\Lambda}+\lambda\right)^{-1}$ is compact. Similarly:

$$
(\Delta+1)^{-1}-\left(\Delta_{\Lambda}+\lambda\right)^{-1}=\left(\mathrm{d}^{*} \mathrm{~d}+1\right)^{-1}-\left(\mathrm{d}^{*} \Lambda \mathrm{~d}+1\right)^{-1} \in \mathcal{K}(\mathscr{H})
$$

For this we use Theorem 3.5 with: $\mathscr{E}=\mathscr{K}, D=\mathrm{d}, a=1, b=\Lambda$ and $z=-1$. Since $\mathrm{d}^{*} \mathrm{~d}$ and $\mathrm{d}^{*} \Lambda d$ are positive self-adjoint operators on $\mathscr{H}$ with the same form domain $\mathscr{G}$, the fi rst condition of Theorem 3.5 is satisfi ed. Then the second condition holds because $\Lambda-1 \in \mathcal{B}_{0}^{l}(\mathscr{K})$. Thus it remains to observe that the operator $\mathrm{d}(\Delta+1)^{-1}$ is decay preserving by Lemma 8.2.

We shall consider now an application of Theorem 8.3 to concrete Riemannian manifolds. It will be clear from what follows that we could treat Lipschitz manifolds with measurable metrics (see [ $\mathrm{DP}, \mathrm{Hi}, \mathrm{Te}, \mathrm{We}$ ] for example), but the case of $C^{1}$ manifolds with locally bounded metrics suffi ces as an example. Note, however, that the arguments of the proof of Theorem 8.4 cover without any modifi cation the case when $X$ is not $C^{1}$ but is a Lipschitz manifold and a countable atlas has been specifi ed, because then the tangent space are well defi ned almost everywhere and the absolute continuity notions that we use make sense.

From now on in this section $X$ is a non-compact differentiable manifold of class $C^{1}$. Then its cotangent manifold $T^{*} X$ is a topological vector fi ber bundle over $X$ whose fi ber over $x$ will be denoted $T_{x}^{*} X$. If $u: X \rightarrow \mathbb{R}$ is differentiable then $\mathrm{d} u(x) \in$ $T_{x}^{*} X$ is its differential at the point $x$ and its differential $\mathrm{d} u$ is a section of $T^{*} X$. Thus for the moment d is a linear map defi ned on the space of real $C^{1}(X)$ functions to the space of sections of $T^{*} X$.

A measurable locally bounded Riemannian structure on $X$ will be called an $R$ structure on $X$. To be precise, an R-structure is given on $X$ if each $T_{x}^{*} X$ is equipped
with a quadratic (i.e. generated by scalar product) norm $\|\cdot\|_{x}$ such that:
$(R)\left\{\begin{array}{l}\text { if } v \text { is a continuous section of } T^{*} X \text { over a compact set } K \text { such that } \\ v(x) \neq 0 \text { for } x \in K, \text { then } x \mapsto\|v(x)\|_{x} \text { is a bounded Borel map on } \\ K \text { and }\|v(x)\|_{x} \geq c \text { for some number } c>0 \text { and all } x \in K .\end{array}\right.$
Such a structure allows one to construct a metric compatible with the topology on $X$, the distance between two points being the infi mum of the length of the Lipschitz curves connecting the points (see the references above). Since $X$ was assumed to be non-compact, the metric space $X$ is proper in the sense defi ned in Subsection 2.3 if and only if it is complete. If this is the case, we say that the $R$-structure is complete.

It will also be convenient to complexify these structures (i.e. replace $T_{x}^{*} X$ by $T_{x}^{*} X \otimes \mathbb{C}$ and extend the scalar product as usual) and to keep the same notations for the complexifi ed objects.

We shall consider positive measures $\mu$ on $X$ such that:
$(M)\left\{\begin{array}{l}\mu \text { is absolutely continuous and its density is locally bounded } \\ \text { and locally bounded from below by strictly positive constants. }\end{array}\right.$
A couple consisting of an R-structure and a measure verifying (M) on $X$ will be called an $R M$-structure on $X$. The defi nition of a complete RM-structure is obvious. To an R-structure we may canonically associate an RM-structure by taking $\mu$ equal to the Riemannian volume element.

If an RM-structure is given on $X$ then we may consider the two Hilbert spaces $\mathscr{H}=L^{2}(X, \mu)$ and $\mathscr{K}$ defi ned as the completion of the space of continuous sections with compact support of $T^{*} X$ equipped with the norm

$$
\|v\|_{\mathscr{K}}^{2}=\int_{X}\|v(x)\|_{x}^{2} d \mu(x) .
$$

In fact, $\mathscr{K}$ is the space of (suitably defi ned) square integrable sections of $T^{*} X$.
The operator of exterior differentiation d induces a linear map $C_{\mathrm{c}}^{1}(X) \rightarrow \mathscr{K}$ which is easily seen to be closable as operator from $\mathscr{H}$ to $\mathscr{K}$ (this is a purely local problem and the hypotheses we put on the metric and the measure allow us to reduce ourselves to the Euclidean case). We shall keep the notation d for its closure and we note that its domain $\mathscr{G}$ is the fir rst order Sobolev space $\mathscr{H}^{1}$ defi ned in this context as the closure of $C_{\mathrm{c}}^{1}(X)$ under the norm

$$
\|u\|_{\mathscr{H}^{1}}^{2}=\int_{X}\left(|u(x)|^{2}+\|\mathrm{d} u(x)\|_{x}^{2}\right) d \mu(x)
$$

The self-adjoint operator $\Delta=\mathrm{d}^{*} \mathrm{~d}$ in $\mathscr{H}$ associated to the quadratic form $\|\cdot\|_{\mathscr{H}^{1}}^{2}$ is the Laplace operator associated to the given RM-structure. This is a generalized form of the Laplace operator associated to the Riemannian structure of $X$ because $\mu$ is not necessarily the Riemannian volume element.

Two RM-structures $\left(\left\{\|\cdot\|_{x}\right\}_{x \in X}, \mu\right)$ and $\left(\left\{\|\cdot\|_{x}^{\prime}\right\}_{x \in X}, \mu^{\prime}\right)$ on $X$ are called equivalent if there are bounded Borel functions $\alpha, \beta, \lambda$ on $X$ with $\alpha \geq c$ and $\lambda \geq c$ for some number $c>0$ such that $\alpha(x)\|\cdot\|_{x} \leq\|\cdot\|_{x}^{\prime} \leq \beta(x)\|\cdot\|_{x}$ for all $x$ and $\mu^{\prime}=\lambda \mu$.

The distances $\rho, \rho^{\prime}$ on $X$ associated to these structures satisfy $a \rho \leq \rho^{\prime} \leq b \rho$ for some numbers $b \geq a>0$, hence if one of the RM-structures is complete, the second one is also complete. Notice that the spaces $\mathscr{H}, \mathscr{K}$ associated to equivalent RM-structures are identical as topological vector spaces.

Two RM-structures are strongly equivalent if they are equivalent and if the functions $\alpha, \beta, \lambda$ can be chosen such that $\lambda(x) \rightarrow 1, \alpha(x) \rightarrow 1$ and $\beta(x) \rightarrow 1$ as $x \rightarrow \infty$.

Theorem 8.4 The Laplace operators associated to strongly equivalent complete RMstructures on $X$ have the same essential spectrum.

Proof: We check that the assumptions of Theorem 8.3 are satisfi ed. We noted above that $X$ is a proper metric space for the metric associated to the initial Riemann structure. The spaces $\mathscr{H}, \mathscr{K}$ have obvious $X$-module structures and for each $\varphi \in C_{\mathrm{c}}(X)$ the operator $\varphi(Q): \mathscr{H}^{1} \rightarrow \mathscr{H}$ is compact. Indeed, by using partitions of unity, we may assume that the support of $\varphi$ is contained in the domain of a local chart and then we are reduced to a known fact in the Euclidean case. Thus $(\mathscr{G}, \mathscr{H})$ is a compact Friedrichs $X$-module. To see that d is a first order operator we observe that if $\varphi$ is Lipschitz then $[\mathrm{d}, \varphi]$ is the operator of multiplication by the differential $\mathrm{d} \varphi$ of $\varphi$ and the estimate ess-sup $\|\mathrm{d} \varphi(x)\|_{x} \leq \operatorname{Lip} \varphi$ is easy to obtain. The conditions on $\lambda$ in Theorem 8.3 are trivially verifi ed. So it remains to consider the operator $\Lambda$. For each $x \in X$ there is a unique operator $\Lambda_{0}(x)$ on $T_{x}^{*} X$ such that $\langle u \mid v\rangle_{x}^{\prime}=\left\langle u \mid \Lambda_{0}(x) v\right\rangle_{x}$ for all $u, v \in T_{x}^{*} X$ and we have $\alpha(x)^{2} \leq \Lambda_{0}(x) \leq \beta(x)^{2}$ by hypothesis. Here the inequalities must be interpreted with respect to the initial scalar product on $T_{x}^{*} X$. Thus the operator $\Lambda$ on $\mathscr{K}$ is just the operator of multiplication by the function $\Lambda(x)=\lambda(x) \Lambda_{0}(x)$ and the condition (iii) of Theorem 8.3 is clearly satisfi ed.

The (strong) equivalence of two R-structures is defi ned in an obvious way. Note that if $\mu, \mu^{\prime}$ are the Riemannian measures associated to two strongly equivalent R -structures then the unique function $\lambda$ such that $\mu^{\prime}=\lambda \mu$ satisfi es $\lambda(x) \rightarrow 1$ as $x \rightarrow \infty$.

Corollary 8.5 The Laplace operators associated to strongly equivalent complete $R$ structures on $X$ have the same essential spectrum.

We stress that if one of the Riemannian structures is locally Lipschitz then this result is easy to prove by using local regularity estimates for elliptic equations.

An assumption of the form $\alpha(x) \rightarrow 1$ as $x \rightarrow \infty$ imposed in the defi nition of strong equivalence means that the set where $|\alpha(x)-1|>\varepsilon$ is relatively compact for any $\varepsilon>0$. We shall consider now a weaker notion of equivalence associated to the fil ler $\mathcal{F}_{\mu}$ introduced in Example 2.9.

We first introduce two notions which clearly depend only on the equivalence class of an RM-structure. We say that an RM-structure is of infinite volume if $\mu(X)=\infty$. We say that it has the $F$-embedding property if for each Borel set $F \subset X$ of fin nite measure the operator $\chi_{F}(Q): \mathscr{H}^{1} \rightarrow \mathscr{H}$ is compact.

Remark 8.6 The F-embedding property is satisfi ed under quite general conditions. Indeed, the compactness of $\chi_{F}(Q): \mathscr{H}^{1} \rightarrow \mathscr{H}$ is equivalent to the compactness of the operator $\chi_{F}(Q)(\Delta+1)^{-1 / 2}$ in $\mathscr{H}$. Or the set of functions $\varphi \in C([0, \infty[)$
such that $\chi_{F}(Q) \varphi(\Delta)$ is compact is a closed $C^{*}$-subalgebra of $C([0, \infty[)$ so it suffi ces to find one function $\varphi$ which generates this algebra such that $\chi_{F}(Q) \varphi(\Delta)$ be compact. But $\chi_{F}(Q) \varphi(\Delta)$ is compact if and only if $\chi_{F}(Q)|\varphi(\Delta)|^{2} \chi_{F}(Q)$ is compact, so we see that it suffi ces to show that for each Borel set $F$ of fi nite measure there is $t>0$ such that $\chi_{F}(Q) \mathrm{e}^{-t \Delta} \chi_{F}(Q)$ be compact. For example, it suffi ces that this operator be Hilbert-Schmidt, i.e. that the integral kernel $P_{t}$ of $\mathrm{e}^{-t \Delta}$ be such that $\int_{F \times F}\left|P_{t}(x, y)\right|^{2} d \mu(x) d \mu(y)<\infty$, which is true if $P_{t}$ satisfi es a Gaussian upper estimate and the measure of a ball of radius $t^{1 / 2}$ is bounded below by a strictly positive constant (see [AC, ACDH] and references there).

Two infi nite volume RM-structures will be called $\mu$-strongly equivalent if they are equivalent and if the functions $\alpha, \beta, \lambda$ can be chosen such that for each $\varepsilon>0$ the set where one of the inequalities $|\alpha(x)-1|>\varepsilon,|\beta(x)-1|>\varepsilon$ or $|\alpha(x)-1|>\varepsilon$ holds is of fi nite measure.

We say that an RM-structure is regular if there is $p>2$ such that $\mathrm{d}(\Delta+1)^{-1}$ induces a bounded operator in $L^{p}$. More precisely, this means that there is a constant $C$ such that if $u \in L^{2}(X) \cap L^{p}(X)$ then $\mathrm{d}(\Delta+1)^{-1} u$, which is a section of $T^{*} X$ of fi nite $L^{2}$ norm, has an $L^{p}$ norm bounded by $C\|u\|_{L^{p}}$. If the operator $\mathrm{d}(\Delta+1)^{-1} d^{*}$ also induces a bounded operators in $L^{p}$ (in an obvious sense), we say that the RM-structure is strongly regular. From the relation $\mathrm{d}(\Delta+1)^{-1} d^{*}=\left[\mathrm{d}(\Delta+1)^{-1 / 2}\right]\left[\mathrm{d}(\Delta+1)^{-1 / 2}\right]^{*}$ we see that strong regularity follows from: there is $\varepsilon>0$ such that $\mathrm{d}(\Delta+1)^{-1 / 2}$ induces a bounded operator in $L^{p}$ for $2-\varepsilon<p<2+\varepsilon$.

Theorem 8.7 Let $\Delta$ be the Laplace operator associated to an infinite volume complete $R M$-structure on $X$ which has the $F$-embedding property and is regular. Then the Laplace operator associated to an RM-structure $\mu$-strongly equivalent to the given structure has the same essential spectrum as $\Delta$.

Proof: Let $\Lambda(x)$ be as in the proof of Theorem 8.4. Clearly there is a number $C>0$ such that $C^{-1} \leq \Lambda(x) \leq C$ for all $x$ and such that for each $\varepsilon>0$ the set where $\|\Lambda(x)-1\|>\varepsilon$ is of fi nite measure (the inequalities and the norm are computed on $T_{x}^{*} X$, which is equipped with the initial scalar product).

Now we proceed as in the proof of Theorem 8.3 but this time we equip $\mathscr{H}$ and $\mathscr{K}$ with the Hilbert module structures described in Example 2.9. To avoid confusions, we denote $\mathcal{B}_{\mu}(\mathscr{H})$ and $\mathcal{B}_{\mu}(\mathscr{K})$ the space of decay improving operators relatively to these new module structures. The F-embedding property implies that $\left(\mathscr{H}^{1}, \mathscr{H}\right)$ is a compact Friedrichs module. Moreover, the operator $\lambda(Q)-1: \mathscr{H}^{1} \rightarrow \mathscr{H}$ is compact. Then, as in the proof of Theorem 8.3, we see that it suffi ces to prove that

$$
\left(\mathrm{d}^{*} \mathrm{~d}+1\right)^{-1}-\left(\mathrm{d}^{*} \Lambda(Q) \mathrm{d}+1\right)^{-1} \in \mathcal{K}(\mathscr{H})
$$

Clearly $\Lambda(Q)-1 \in \mathcal{B}_{\mu}(\mathscr{K})$. Now we use Theorem 3.5 exactly as in the proof of Theorem 8.3 and we see that the only condition which remains to be checked is (3) of Theorem 3.5, i.e. in our case $\mathrm{d}(\Delta+1)^{-1} \in \mathcal{B}_{q}^{r}(\mathscr{H}, \mathscr{K})$, where the decay preserving property is relatively to the algebra $B_{\mu}(X)$. But this follows from Theorem 6.1.

One may check the regularity property needed in Theorem 8.7 by using the results from [AC, ACDH] concerning the boundedness in $L^{p}$ of the operator $d \Delta^{-1 / 2}$. For
example, it suffi ces that $X$ be complete, with the doubling volume property, and such that the Poincar'e inequality holds in $E^{2}$ sense. Note, however, that these results are much stronger than necessary in our context and it seems reasonable to think that the boundedness of $d(\Delta+1)^{-1 / 2}$ holds under less restrictive assumptions.

The next result does not require regularity assumptions on any of the RM-structures that we want to compare but only on a third one in their equivalence class. Observe that each equivalence class of RM-structures contains one of the same degree of local smoothness as the manifold $X$ (make local regularizations and use a partition of unity).

Theorem 8.8 Let $\Delta_{a}, \Delta_{b}$ be the Laplace operators associated to $\mu$-strongly equivalent complete RM-structures of infinite volume and having the F-embedding property. If these structures are equivalent to a strongly regular RM-structure, then $\Delta_{a}$ and $\Delta_{b}$ have the same essential spectrum.

Proof: Let $\Delta_{c}$ be the Laplace operator associated to the third structure. From Theorem 6.1 it follows that $\mathrm{d}\left(\Delta_{c}+1\right)^{-1}$ and $\mathrm{d}\left(\Delta_{c}+1\right)^{-1} D^{*}$ are right $\mathcal{F}_{\mu}$-decay preserving. Then from Proposition 3.7 we see that $\mathrm{d}\left(\Delta_{a}+1\right)^{-1}$ is right $\mathcal{F}_{\mu}$-decay preserving and we may conclude as in the proof of Theorem 8.7.

Remark 8.9 It is natural to consider an analog of the filter $\mathcal{F}_{\mathrm{w}}$ introduced in Section 6 to get an optimal weak decay condition for the stability of the essential spectrum in the present context. The techniques of Section 6 should be relevant for this question.

Remarks on Laplace operators acting on forms: We shall describe here, without going into details, an abstract framework for the study of the Laplace operator acting on forms. Let $\mathscr{H}$ be a Hilbert space and d a closed densely defi ned operator in $\mathscr{H}$ such that $\mathrm{d}^{2}=0$. For example, $\mathscr{H}$ could be the space of square integrable differential forms over a Lipschitz manifold and d the operator of exterior differentiation. We denote $\delta=\mathrm{d}^{*}$ and we assume that $\mathscr{G}:=\mathcal{D}(\mathrm{d}) \cap \mathcal{D}(\delta)$ is dense in $\mathscr{H}$ (which is a rather strong condition in the context of this paper, e.g. in the preceding example it is a differentiability condition on the metric). Then let $D=\mathrm{d}+\delta$ with domain $\mathscr{G}$, observe that $\|D u\|^{2}=\|\mathrm{d} u\|^{2}+\|\delta u\|^{2}$ so $D$ is a closed symmetric operator, assume that $D$ is self-adjoint, and defi ne $\Delta=D^{2}=\mathrm{d} \delta+\delta \mathrm{d}$ (form sum). Then

$$
\begin{equation*}
(\Delta+1)^{-1}=(D+i)^{-1}(D-i)^{-1} \tag{8.30}
\end{equation*}
$$

Now let $a \in \mathcal{B}(\mathscr{H})$ with $a \geq \varepsilon>0$ and such that $a^{ \pm 1 \mathscr{G}} \subset \mathscr{G}$ and let $\mathscr{H}_{a}$ be the Hilbert space which is equal to $\mathscr{H}$ as vector space but is equipped with the new the scalar product $\langle u, v\rangle_{a}=\langle u, a v\rangle$. Denote $\mathrm{d}_{a}$ the operator d viewed as operator acting in $\mathscr{H}_{a}$ with adjoint $\delta_{a}=a^{-1} \delta a$. We can defi ne as above operators $D_{a}$ (with domain $\left.\mathscr{G}_{a}=\mathscr{G}\right)$ and $\Delta_{a}=D_{a}^{2}$ which are self-adjoint in $\mathscr{H}_{a}$ and satisfy a relation similar to (8.30). Then $\Delta_{a}$ is a compact perturbation of $\Delta$ if the operators $\left(D_{a} \pm i\right)^{-1}-(D \pm i)^{-1}$ are compact and this last condition is equivalent to the compactness of the operator $D_{a}-D: \mathscr{G} \rightarrow \mathscr{G}^{*}$. And this holds if $(\mathscr{G}, \mathscr{H})$ is a compact Hilbert $X$-module over a metric space $X$ and $a-1 \in \mathcal{B}_{0}(\mathscr{H})$.

## 9 On Maurey's factorization theorem

The subject of this section is quite different from that of the rest of the paper: we shall prove a version of a factorization theorem due to Bernard Maurey which plays an important role in several arguments from the main part of this article. We fi rst recall Maurey's result, cf. Theorems 2 and 8 in [Ma].

Theorem 9.1 Let $1<p<2$ and let $T$ be an arbitrary continuous linear map from a Hilbert space $\mathscr{H}$ into $L^{p}$. Then there is $R \in \mathcal{B}\left(\mathscr{H}, L^{2}\right)$ and there is a function $g \in L^{q}$, where $\frac{1}{p}=\frac{1}{2}+\frac{1}{q}$, such that $T=g(Q) R$.

We have stated only the particular case we need of the theorem (the result extends easily to larger classes of Banach spaces $\mathscr{H}$ ). The $L^{p}$ spaces refer to an arbitrary positive measure space $(X, \mu)$.

Before going on to our main purpose, we shall state an easy consequence of this theorem which is needed in Sections 7 and 8 . Let $\{\mathscr{K}(x)\}_{x \in X}$ be a measurable family of Hilbert spaces (see [Di, Ch. II]) such that the dimension of $\mathscr{K}(x)$ is $\leq N$ for some fi nite $N$. Let $\mathscr{K}=\int_{X}^{\oplus} \mathscr{K}(x) d \mu(x)$ be the corresponding direct integral and for each $p \geq 1$ let $\mathscr{K}_{p}$ be the space of ( $\mu$-equivalence classes) of measurable vector fi elds $v$ such that $\int_{X}\|v(x)\|_{\mathscr{K}(x)}^{p} d \mu(x)<\infty$. Thus $\mathscr{K}_{p}$ is naturally a Banach space and $\mathscr{K}_{2}=\mathscr{K}$.

Corollary 9.2 Let $\mathscr{H}$ be a Hilbert space and let $T \in \mathcal{B}\left(\mathscr{H}, \mathscr{K}_{p}\right)$ with $1<p<2$. Then there is $R \in \mathcal{B}(\mathscr{H}, \mathscr{K})$ and there is a function $g \in L^{q}$, where $q=2 p /(2-p)$, such that $T=g(Q) R$.

Proof: For each $n=1, \ldots, N$ let $X_{n}$ be the set of $x$ such that the dimension of $\mathscr{K}(x)$ is equal to $n$. Then $X$ is the disjoint union of the measurable sets $X_{n}$. For each $x$ there is $n$ such that $x \in X_{n}$ and we can choose a unitary map $j(x): \mathscr{K}(x) \rightarrow \mathbb{C}^{n}$ such that $\left\{j_{x}\right\}$ be a measurable family of operators. Let $J$ be the operator acting on vector fi elds according to the rule $(J v)(x)=j(x) v(x)$, let $\Pi_{n}$ be the operator of multiplication by $\chi_{X_{n}}$, and let $T_{n} \equiv \Pi_{n} J T \in \mathcal{B}\left(\mathscr{H}, L^{p}\left(X_{n} ; \mathbb{C}^{n}\right)\right)$. We can write $T_{n}=\left(T_{n}^{k}\right)_{1 \leq k \leq n}$ with $T_{n}^{k} \in \mathcal{B}\left(\mathscr{H}, L^{p}\left(X_{n}\right)\right)$ and Maurey's theorem gives us a factorization $T_{n}^{k}=g_{n}^{k}(Q) S_{n}^{k}$ with $S_{n}^{k} \in \mathcal{B}\left(\mathscr{H}, L^{2}\left(X_{n}\right)\right)$ and $g_{n}^{k} \in L^{q}\left(X_{n}\right)$, and clearly we may assume $g_{n}^{k} \geq 0$. Let $g_{n}=\sup _{k} g_{n}^{k} \in L^{q}\left(X_{n}\right)$ and $S_{n} \in \mathcal{B}\left(\mathscr{H}, L^{2}\left(X_{n} ; \mathbb{C}^{n}\right)\right)$ be the operator with components $\left(g_{n}^{k} g_{n}^{-1}\right)(Q) R_{n}^{k}$. Then $T_{n}=g_{n}(Q) S_{n}$ and if we defi ne $R_{n}=J^{-1} S_{n}$ we get

$$
g_{n}(Q) R_{n}=J^{-1} g_{n}(Q) S_{n}=J^{-1} T_{n}=\Pi_{n} T
$$

Thus, if we defi ne $g=\sum_{n} \chi_{X_{n}} g_{n}$ and $R=\sum \Pi_{n} R_{n}$, we get $T=g(Q) R$.
Our purpose in the rest of this section is to extend Theorem 9.1 (in the case $X=$ $\mathbb{R}^{n}$ ) to more general classes of spaces of measurable functions, which do not seem to be covered by the results existing in the literature, cf. [ Kr$]$. Our proof follows closely that of Maurey. We fi rst recall Ky Fan's Lemma, see [DJT, 9.10].

Proposition 9.3 Let $\mathcal{K}$ be a compact convex subset of a Hausdorff topological vector space and let $\mathscr{F}$ be a convex set of functions $F: \mathcal{K} \rightarrow]-\infty,+\infty]$ such that each
$F \in \mathscr{F}$ is convex and lower semicontinuous. If for each $F \in \mathscr{F}$ there is $g \in \mathcal{K}$ such that $F(g) \leq 0$, then there is $g \in \mathcal{K}$ such that $F(g) \leq 0$ for all $F \in \mathscr{F}$.

We need a second general fact that we state below. Let $(X, \mu)$ be a $\sigma$-fi nite positive measure space and let $L^{0}(X)$ be the space of $\mu$-equivalence classes of complex valued measurable functions on $X$ with the topology of convergence in measure. Let $\mathscr{L}$ be a Banach space with $\mathscr{L} \subset L^{0}(X)$ linearly and continuously and such that if $f \in L^{0}(X)$, $g \in \mathscr{L}$ and $|f| \leq|g|$ ( $\mu$-a.e.) then $f \in \mathscr{L}$ and $\|f\|_{\mathscr{L}} \leq\|g\|_{\mathscr{L}}$. The next result is a rather straightforward consequence of Khinchin's inequality [DJT, 1.10] (see also [Pi, Section 8]).

Proposition 9.4 There is a number $C$, independent of $\mathscr{L}$, such that for any Hilbert space $\mathscr{H}$ and any $T \in \mathcal{B}(\mathscr{H}, \mathscr{L})$ the following inequality holds

$$
\begin{equation*}
\left\|\left(\sum_{j}\left|T u_{j}\right|^{2}\right)^{1 / 2}\right\|_{\mathscr{L}} \leq C\|T\|_{\mathcal{B}(\mathscr{H}, \mathscr{L})}\left(\sum_{j}\left\|u_{j}\right\|^{2}\right)^{1 / 2} \tag{9.31}
\end{equation*}
$$

for all finite families $\left\{u_{j}\right\}$ of vectors in $\mathscr{H}$.
From now on we work in a setting adapted to our needs in Section 7, although it is clear that we could treat by the same methods a general abstract situation. Let $X=\mathbb{R}^{n}$ equipped with the Lebesgue measure, denote $Z=\mathbb{Z}^{n}$, and for each $a \in Z$ let $K_{a}=a+K$, where $\left.\left.K=\right]-1 / 2,1 / 2\right]^{n}$, so that $K_{a}$ is a unit cube centered at $a$ and we have $X=\bigcup_{a \in Z} K_{a}$ disjoint union. Let $\chi_{a}$ be the characteristic function of $K_{a}$ and if $f: X \rightarrow \mathbb{C}$ let $f_{a}=f \mid K_{a}$. We fix a number $1<p<2$ and a family $\left\{\lambda_{a}\right\}_{a \in Z}$ of strictly positive numbers $\lambda_{a}>0$ and we defi ne $\mathscr{L} \equiv \mathscr{\ell}_{\lambda}^{\ell}\left(L^{p}\right)$ as the Banach space of all (equivalence classes) of complex functions $f$ on $X$ such that

$$
\begin{equation*}
\|f\|_{\mathscr{L}}:=\left(\sum_{a \in Z}\left\|\lambda_{a} \chi_{a} f\right\|_{L^{p}}^{2}\right)^{1 / 2}<\infty \tag{9.32}
\end{equation*}
$$

Here $L^{p}=L^{p}(X)$ but note that, by identifying $\chi_{a} f \equiv f_{a}$, we can also interpret $\mathscr{L}$ as a conveniently normed direct sum of the spaces $L^{p}\left(K_{a}\right)$, see [DJT, page XIV]. If $\lambda_{a}=1$ for all $a$ we set $\ell_{\lambda}^{2}\left(L^{p}\right)=\ell^{2}\left(L^{p}\right)$. Observe that $\ell^{2}\left(L^{2}\right)=L^{2}(X)$.

Let $q$ be given by $\frac{1}{p}=\frac{1}{2}+\frac{1}{q}$, so that $1<p<2<q<\infty$. We also need the space $\mathscr{M} \equiv \ell_{\lambda}^{\infty}\left(L^{q}\right)$ defi ned by the condition

$$
\begin{equation*}
\|g\|_{\mathscr{M}}:=\sup _{a \in Z}\left\|\lambda_{a} \chi_{a} g\right\|_{L^{q}}<\infty \tag{9.33}
\end{equation*}
$$

The defi nitions are chosen such that $\|g u\|_{\mathscr{L}} \leq\|g\|_{\mathscr{M}}\|u\|_{L^{2}}$ where $L^{2}=L^{2}(X)$. As explained in [DJT, page XV], the space $\mathscr{M}$ is naturally identifi ed with the dual space of the Banach space $\mathscr{M}_{*} \equiv \ell_{\lambda^{-1}}^{1}\left(L^{q^{\prime}}\right)$, where $\frac{1}{q}+\frac{1}{q^{\prime}}=1$, defi ned by the norm

$$
\|h\|_{\mathscr{M}_{*}}:=\sum_{a \in Z}\left\|\lambda_{a}^{-1} \chi_{a} h\right\|_{L^{q^{\prime}}}
$$

Below, when we speak about $w^{*}$-topology on $\mathscr{M}$ we mean the $\sigma\left(\mathscr{M}, \mathscr{M}_{*}\right)$-topology. Clearly

$$
\mathscr{M}_{1}^{+}=\left\{g \in \mathscr{M} \mid g \geq 0,\|g\|_{\mathscr{M}} \leq 1\right\}
$$

is a convex compact subset of $\mathscr{M}$ for the $w^{*}$-topology.

Lemma 9.5 For each $f \in \mathscr{L}$ there is $g \in \mathscr{M}_{1}^{+}$such that $\|f\|_{\mathscr{L}}=\left\|g^{-1} f\right\|_{L^{2}}$.
Proof: We can assume $f \geq 0$. Since $1=\frac{p}{2}+\frac{p}{q}$, we have:

$$
\left\|f_{a}\right\|_{L^{p}}=\left\|f_{a}\right\|_{L^{p}}^{p / 2}\left\|f_{a}\right\|_{L^{p}}^{p / q}=\left\|f_{a}^{p / 2}\right\|_{L^{2}}\left\|f_{a}^{p / q}\right\|_{L^{q}}=\left\|f_{a}^{-p / q} f\right\|_{L^{2}}\left\|f_{a}^{p / q}\right\|_{L^{q}}
$$

with the usual convention $0 / 0=0$. Now we defi ne $g_{a}$ on $K_{a}$ as follows. If $f_{a}=0$ then we take any $g_{a} \geq 0$ satisfying $\lambda_{a}\left\|g_{a}\right\|_{L^{q}}=1$. If $f_{a} \neq 0$ let

$$
g_{a}=\lambda_{a}^{-1}\left(f_{a} /\left\|f_{a}\right\|_{L^{p}}\right)^{p / q}=\lambda_{a}^{-1}\left\|f_{a}^{p / q}\right\|_{L^{q}}^{-1} f_{a}^{p / q}
$$

Thus we have $\lambda_{a}\left\|g_{a}\right\|_{L^{q}}=1$ for all $a$, in particular $\|g\|_{\mathscr{M}}=1$. By the preceding computations we also have $\left\|f_{a}\right\|_{L^{p}}=\left\|g_{a}^{-1} f_{a}\right\|_{L^{2}}\left\|g_{a}\right\|_{L^{q}}$ and so

$$
\|f\|_{\mathscr{L}}^{2}=\sum \lambda_{a}^{2}\left\|f_{a}\right\|_{L^{p}}^{2}=\sum \lambda_{a}^{2}\left\|g_{a}\right\|_{L^{q}}^{2}\left\|g_{a}^{-1} f_{a}\right\|_{L^{2}}^{2}=\sum\left\|g_{a}^{-1} f_{a}\right\|_{L^{2}}^{2}
$$

which is just $\left\|g^{-1} f\right\|_{L^{2}}^{2}$.
The main technical result follows.
Proposition 9.6 Let $\left(f^{u}\right)_{u \in U}$ be a family of functions in $\mathscr{L}$ such that, for each $\alpha=$ $\left(\alpha_{u}\right)_{u \in U}$ with $\alpha_{u} \in \mathbb{R}, \alpha_{u} \geq 0$ and $\alpha_{u} \neq 0$ for at most a finite number of $u$, the function $f^{\alpha}:=\left(\sum_{u}\left|\alpha_{u} f^{u}\right|^{2}\right)^{1 / 2}$ satisfies $\left\|f^{\alpha}\right\|_{\mathscr{L}} \leq\|\alpha\|_{\ell^{2}(U)}$. Then there is $g \in$ $\mathscr{M}_{1}^{+}$such that $\left\|g^{-1} f^{u}\right\|_{L^{2}} \leq 1$ for all $u \in U$.

Proof: For each $\alpha$ as in the statement of the proposition we defi ne a function $F_{\alpha}$ : $\left.\left.\mathscr{M}_{1}^{+} \rightarrow\right]-\infty,+\infty\right]$ as follows:

$$
F_{\alpha}(g)=\left\|g^{-1} f^{\alpha}\right\|_{L^{2}}^{2}-\|\alpha\|_{\ell^{2}(U)}^{2}=\sum_{u} \alpha_{u}^{2}\left(\left\|g^{-1} f^{u}\right\|_{L^{2}}^{2}-1\right)
$$

Our purpose is to apply Proposition 9.3 with $\mathscr{K}=\mathscr{M}_{1}^{+}$equipped with the $w^{*}-$ topology and $\mathscr{F}$ equal to the set of all functions $F_{\alpha}$ defi ned above. We saw before that $\mathscr{K}$ is a convex compact set. From the second representation of $F_{\alpha}$ given above it follows that $\mathscr{F}$ is a convex set. Each $F_{\alpha}$ is a convex function because $\left\|g^{-1} f^{\alpha}\right\|_{L^{2}}^{2}=$ $\int g^{-2}\left(f^{\alpha}\right)^{2} \mathrm{~d} x$ and the map $t \mapsto t^{-2}$ is convex on $[0, \infty[$. We shall prove in a moment that $F_{\alpha}$ is lower semicontinuous. From Lemma 9.5 it follows that there is $g_{\alpha} \in \mathscr{K}$ such that $\left\|f^{\alpha}\right\|_{\mathscr{L}}=\left\|g_{\alpha}^{-1} f^{\alpha}\right\|_{L^{2}}$. Then by our assumptions we have

$$
F_{\alpha}\left(g_{\alpha}\right)=\left\|f^{\alpha}\right\|_{\mathscr{L}}^{2}-\|\alpha\|_{\ell^{2}(U)}^{2} \leq 0
$$

From Ky Fan's Lemma it follows that one can choose $g \in \mathscr{K}$ such that $F_{\alpha}(g) \leq 0$ for all $\alpha$, which fi nishes the proof of the proposition.

It remains to show the lower semicontinuity of $F_{\alpha}$. For this it suffi ces to prove that $g \mapsto\left\|g^{-1} f\right\|_{L^{2}}^{2} \in[0, \infty]$ is lower semicontinuous on $\mathscr{K}$ if $f \in \mathscr{L}, f \geq 0$. But

$$
\left\|g^{-1} f\right\|_{L^{2}}^{2}=\sum_{a} \int_{K_{a}} g_{a}^{-2} f_{a}^{2} \mathrm{~d} x
$$

and the set of lower semicontinuous functions $\mathscr{K} \rightarrow[0, \infty]$ is stable under sums and upper bounds of arbitrary families. Hence it suffi ces to prove that each map $g \mapsto$ $\int_{K_{a}} g_{a}^{-2} f_{a}^{2} \mathrm{~d} x$ is lower semicontinuous. This map can be written as a composition $\phi \circ J_{a}$ where $J_{a}: \mathscr{M} \rightarrow L^{q}\left(K_{a}\right)$ is the restriction map $J_{a} g=g_{a}$ and $\phi: L^{q}\left(K_{a}\right) \rightarrow[0, \infty]$ is defi ned by $\phi(\theta)=\int_{K_{a}} \theta^{-2} f_{a}^{2} \mathrm{~d} x$. The map $J_{a}$ is continuous if we equip $L^{q}\left(K_{a}\right)$ with the weak topology and $\mathscr{M}$ with the $w^{*}$-topology because it is the adjoint of the norm continuous map $L^{q^{\prime}}\left(K_{a}\right) \rightarrow \mathscr{M}_{*}$ which sends $u$ into the function equal to $u$ on $K_{a}$ and 0 elsewhere. Thus it suffi ces to show that $\phi$ is lower semicontinuous on the positive part of $L^{q}\left(K_{a}\right)$ equipped with the weak topology and for this we can use exactly the same argument as Maurey. We must prove that the set $\left\{\theta \in L^{q}\left(K_{a}\right) \mid \theta \geq\right.$ $0, \phi(\theta) \leq r\}$ is weakly closed for each real $r$. Since $\phi$ is convex, this set is convex, so it suffi ces to show that it is norm closed. But this is clear by the Fatou Lemma.

Theorem 9.7 Let $\mathscr{H}$ be a Hilbert space and $T: \mathscr{H} \rightarrow \mathscr{L}$ a linear continuous map. Then there exist a linear continuous map $R: \mathscr{H} \rightarrow L^{2}(X)$ and a positive function $g \in \mathscr{M}$ such that $T=g(Q) R$.

Proof: Let $U$ be the unit ball of $\mathscr{H}$ and for each $u \in U$ let $f^{u}=T u$. From Proposition 9.4 we get

$$
\left\|f^{\alpha}\right\|_{\mathscr{L}}=\left\|\left(\sum_{u}\left|T\left(\alpha_{u} u\right)\right|^{2}\right)^{1 / 2}\right\|_{\mathscr{L}} \leq A\left(\sum_{u}\left\|\alpha_{u} u\right\|^{2}\right)^{1 / 2} \leq A\left(\sum_{u}\left|\alpha_{u}\right|^{2}\right)^{1 / 2}
$$

where $A=C\|T\|_{\mathcal{B}(\mathscr{H}, \mathscr{L})}$. Since there is no loss of generality in assuming $A \leq 1$, we see that the assumptions of Proposition 9.6 are satisfi ed. So there is $g \in \mathscr{M}_{1}^{+}$ such that $\left\|g^{-1} T u\right\|_{L^{2}(X)} \leq 1$ for all $u \in U$. Thus it suffi ces to defi ne $R$ by the rule $R u=g^{-1} T u$ for all $u \in \mathscr{H}$.

## A Appendix

Let $(\mathscr{G}, \mathscr{H})$ be a Friedrichs couple and $\mathscr{G} \subset \mathscr{H} \subset \mathscr{G}^{*}$ the Gelfand triplet associated to it. To an operator $S \in \mathcal{B}\left(\mathscr{G}, \mathscr{G}^{*}\right)$ (which is the same as a continuous sesquilinear form on $\mathscr{G}$ ) we associate an operator $\widehat{S}$ acting in $\mathscr{H}$ according to the rules: $\mathcal{D}(\widehat{S})=$ $S^{-1}(\mathscr{H}), \widehat{S}=S \mid \mathcal{D}(\widehat{S})$. Due to the identifi cation $\mathscr{G}^{* *}=\mathscr{G}$, the operator $S^{*}$ is an element of $\mathcal{B}\left(\mathscr{G}, \mathscr{G}^{*}\right)$, so $\widehat{S^{*}}$ makes sense. On the other hand, if $\widehat{S}$ is densely defi ned in $\mathscr{H}$ then the adjoint $\widehat{S}^{*}$ of $\widehat{S}$ with respect to $\mathscr{H}$ is also well defi ned and we clearly have $\widehat{S^{*}} \subset \widehat{S}^{*}$.

Lemma A. 1 If $S-z: \mathscr{G} \rightarrow \mathscr{G}^{*}$ is bijective for some $z \in \mathbb{C}$, then $\widehat{S}$ is a closed densely defined operator, we have $\widehat{S}^{*}=\widehat{S^{*}}$ and $z \in \rho(\widehat{S})$. Moreover, the domains $\mathcal{D}(\widehat{S})$ and $\mathcal{D}\left(\widehat{S}^{*}\right)$ are dense subspaces of $\mathscr{G}$.

Proof: Clearly we can assume $z=0$. From the bijectivity of $S: \mathscr{G} \rightarrow \mathscr{G}^{*}$ and the inverse mapping theorem it follows that $S$ and $S^{*}$ are homeomorphisms of $\mathscr{G}$ onto $\mathscr{G}^{*}$. Since $\mathscr{H}$ is dense in $\mathscr{G}^{*}$, we see that $\mathcal{D}(\widehat{S})$ and $\mathcal{D}\left(\widehat{S^{*}}\right)$ are dense in $\mathscr{G}$, hence in $\mathscr{H}$.

Since $\widehat{S^{*}} \subset \widehat{S}^{*}$, the operator $\widehat{S}^{*}$ is also densely defi ned in $\mathscr{H}$. Thus $\widehat{S}$ is densely defi ned and closable. We now show that it is closed. Consider a sequence of elements $u_{n} \in \mathcal{D}(\widehat{S})$ such that $u_{n} \rightarrow u$ and $\widehat{S} u_{n} \rightarrow v$ in $\mathscr{H}$. Then $S u_{n} \rightarrow v$ in $\mathscr{G}^{*}$ hence, $S^{-1}$ being continuous, $u_{n} \rightarrow S^{-1} v$ in $\mathscr{G}$, so in $\mathscr{H}$. Hence $u=S^{-1} v \in \mathcal{D}(\widehat{S})$ and $\widehat{S} u=v$.

We have proved that $\widehat{S}$ is densely defi ned and closed and clearly $0 \in \rho(\widehat{S})$. Then we also have $0 \in \rho\left(\widehat{S}^{*}\right)$, so $\widehat{S}^{*}: \mathcal{D}\left(\widehat{S}^{*}\right) \rightarrow \mathscr{H}$ is bijective. Since $\widehat{S^{*}}: \mathcal{D}\left(\widehat{S^{*}}\right) \rightarrow \mathscr{H}$ is also bijective and $\widehat{S}^{*}$ is an extension of $\widehat{S^{*}}$, we get $\widehat{S^{*}}=\widehat{S}^{*}$.

A standard example of operator satisfying the condition required in Lemma A. 1 is a coercive operator, i.e. such that $\operatorname{Re}\langle u, S u\rangle \geq \mu\|u\|_{\mathscr{G}}^{2}-\nu\|u\|_{\mathscr{H}}^{2}$ for some strictly positive constants $\mu, \nu$ and all $u \in \mathscr{G}$. Indeed, replacing $S$ by $S+\nu$, we may assume $\operatorname{Re}\langle u, S u\rangle \geq \mu\|u\|_{\mathscr{G}}^{2}$. Since $S^{*}$ verifi es the same estimate, this clearly gives $\|S u\|_{\mathscr{G}_{*}} \geq \mu\|u\|_{\mathscr{G}_{\mathcal{G}}}$ and $\left\|S^{*} u\right\|_{\mathscr{G}^{*}} \geq \mu\|u\|_{\mathscr{G}}$ for all $u \in \mathscr{G}$. Thus $S$ and $S^{*}$ are injective operators with closed range, which implies that they are bijective.

If $A$ is a self-adjoint operator on $\mathscr{H}$ then there is a natural Gelfand triplet associated to it, namely $\mathcal{D}\left(|A|^{1 / 2}\right) \subset \mathscr{H} \subset \mathcal{D}\left(|A|^{1 / 2}\right)^{*}$. Then $A$ extends to a continuous operator $A_{0}: \mathcal{D}\left(|A|^{1 / 2}\right) \rightarrow \mathcal{D}\left(|A|^{1 / 2}\right)^{*}$ which fulfi lls the conditions of Lemma A. 1 and one has $\widehat{A}_{0}=A$. In our applications it is interesting to know whether there are other Gelfand triplets $\mathscr{G} \subset \mathscr{H} \subset \mathscr{G}^{*}$ with $\mathcal{D}(A) \subset \mathscr{G}$ and such that $A$ extends to a continuous operator $\mathscr{G} \rightarrow \mathscr{G}^{*}$. For not semibounded operators, e.g. for Dirac operators, many other possibilities exist such that $\mathscr{G}$ is not comparable to $\mathcal{D}\left(|A|^{1 / 2}\right)$. But if $A$ is semibounded, then the class of spaces $\mathscr{G}$ is rather restricted, as the next lemma shows.

Lemma A. 2 Assume that $A$ is a bounded from below self-adjoint operator on $\mathscr{H}$ and such that $\mathcal{D}(A) \subset \mathscr{G}$ densely. Then $A$ extends to a continuous operator $\widetilde{A}: \mathscr{G} \rightarrow \mathscr{G}^{*}$ if and only if $\mathscr{G} \subset \mathcal{D}\left(|A|^{1 / 2}\right)$ and in this case $\widetilde{A}=\left.A_{0}\right|_{\mathscr{G}}$.

Proof: We prove only the nontrivial implication of the lemma. So let us assume that $A$ extends to some $\widetilde{A} \in \mathcal{B}\left(\mathscr{G}, \mathscr{G}^{*}\right)$. Replacing $A$ by $A+\lambda$ with $\lambda$ a large enough number, we can assume that $A \geq 1$. For $u \in \mathcal{D}(A)$ we have

$$
\left\|A^{1 / 2} u\right\|_{\mathscr{H}}=\sqrt{\langle u, A u\rangle}=\sqrt{\langle u, \widetilde{A} u\rangle} \leq C\|u\|_{\mathscr{G}},
$$

where $C^{2}=\|\widetilde{A}\|_{\mathscr{G} \rightarrow \mathscr{G}_{*} .}$. Since $\mathcal{D}(A)$ is dense in $\mathscr{G}$, it follows that the inclusion map $\mathcal{D}(A) \rightarrow \mathcal{D}\left(A^{1 / 2}\right)$ extends to a continuous linear map $J: \mathscr{G} \rightarrow \mathcal{D}\left(A^{1 / 2}\right)$. If $u \in \mathscr{G}$ then there is a sequence $\left\{u_{n}\right\}$ in $\mathcal{D}(A)$ such that $u_{n} \rightarrow u$ in $\mathscr{G}$. Then $J\left(u_{n}\right) \rightarrow J(u)$ in $\mathcal{D}\left(A^{1 / 2}\right)$. Since $\mathscr{G}$ and $\mathcal{D}\left(A^{1 / 2}\right)$ are continuously embedded in $\mathscr{H}$ we shall have $u_{n} \rightarrow u$ in $\mathscr{H}$ and $u_{n}=J\left(u_{n}\right) \rightarrow J(u)$ in $\mathscr{H}$, hence $J(u)=u$ for all $u \in \mathscr{G}$. In other terms, $\mathscr{G} \subset \mathcal{D}\left(A^{1 / 2}\right)$.

We note that, under the conditions of the lemma, the inclusions $\mathcal{D}(A) \subset \mathscr{G}$ and $\mathscr{G} \subset \mathcal{D}\left(|A|^{1 / 2}\right)$ are continuous (by the closed graph theorem), so we have a scale

$$
\mathcal{D}(A) \subset \mathscr{G} \subset \mathcal{D}\left(|A|^{1 / 2}\right) \subset \mathscr{H} \subset \mathcal{D}\left(|A|^{1 / 2}\right)^{*} \subset \mathscr{G}^{*} \subset \mathcal{D}(A)^{*}
$$

with continuous and dense embeddings (because $\mathcal{D}(A)$ is dense in $\mathcal{D}\left(|A|^{1 / 2}\right)$ ).

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[^0]:    ${ }^{1}$ See [Sa, Ta] for the corresponding pseudo-differential calculus.

[^1]:    ${ }^{2}$ If $X$ is second countable then this property determines uniquely the extension. In general, uniqueness is assured by the property: if $U \subset X$ is open then $\chi_{U}(Q)=\sup _{\varphi} \varphi(Q)$, where $\varphi$ runs over the set of continuous functions with compact support such that $0 \leq \varphi \leq \chi_{U}$.

[^2]:    ${ }^{3}$ If $X$ is an euclidean space and $\mathscr{H}=\mathscr{K}=L^{2}(X)$, the next condition means that there is $r<\infty$ such that the distribution kernel of $S$ satisfi es $S(x, y)=0$ for $|x-y|>r$.

[^3]:    ${ }^{4}$ Note that if $X$ is a locally compact space and $\mu$ a Radon measure then $\mathcal{F}_{\mu}$ is fi ner than the $\mathrm{Fr}^{\prime}$ echet fi lter. Moreover, if $X$ is an abelian locally compact non-compact group then $\mathcal{F}_{\mu}$ is strictly included in the filter $\mathcal{F}_{\mathrm{w}}$ which will be defi ned below.

