# An Algebra of Pseudodifferential Operators and the Asymptotics of Quantum Mechanics 

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#### Abstract

We discuss in detail the regularity properties of a class of pseudodifferential operators on $\mathbb{R}^{b}$ introduced by Grossmann, Loupias, and Stein, which are more regular and more symmetrical than usual pseudodifferential operators. Those operators that are self-adjoint form a suitable class of smooth observables for a nonrelativistic quantum theory. If their symbols are allowed to depend smoothly upon Planck's constant $\hbar$, those operators provide the framework for regular asymptotics expansions as $\hbar \rightarrow 0$ of quantum mechanics around classical mechanics.


## Introduction

This article studies the relevance of a ring of pseudodifferential operators on $\mathbb{R}^{l}$, introduced by Grossmann et al. [1], to the usual formulation of nonrelativistic quantum mechanics using self-adjoint operators on a Hilbert space. The symbols in [1] lead to operators more regular than the usual pseudodifferential operators [2,3]; they satisfy global $L^{2}$ estimates, and some operators with real semibounded symbols admit Friedrichs extensions. When the symbols are allowed to depend explicitly and smoothly upon $\hbar$, the relationship between symbol calculus and the classical limit of quantum mechanics is recovered [1, 4,5]. Moreover, we thus obtain a framework for regular asymptotic expansions of quantum mechanics as $\hbar \rightarrow 0$; the discussion of relevant physical examples is started in [5].

We use the following notation: $Q=\mathbb{R}^{l}$ is configuration space; our pseudodifferential operators will act on various function or distribution spaces $\mathscr{P}(Q)$, $L^{2}(Q), \mathscr{S}^{\prime}(Q), \ldots$ (always complex-valued unless stated otherwise, and noted as in [6]); their symbols will be functions on the phase space $X=Q \oplus P\left(P=Q^{*}\right)$ which carries the canonical symplectic 2-form $\omega$; $\omega\left(x_{1}, x_{2}\right)=q_{1} p_{2}-q_{2} p_{1}$ (vector indices will usually be suppressed and $q p$ will denote the scalar product).

[^0]In any dimension $n$, we shall use multi-indices $\alpha, \beta \in \mathbb{N}^{n}$ with the usual notation:

$$
\begin{aligned}
&|\alpha|=\sum_{1}^{n} \alpha_{j} ; \quad \alpha \pm \beta=\left(\alpha_{1} \pm \beta_{1}, \ldots, \alpha_{n} \pm \beta_{n}\right) ; \\
& \alpha \geqslant \beta \quad \text { if } \quad \alpha-\beta \in \mathbb{N}^{n} ; \quad \alpha!=\prod_{1}^{n} \alpha_{j}!
\end{aligned}
$$

if $r \in \mathbb{R}^{n}, \boldsymbol{r}^{\alpha}=r_{1}^{\alpha_{1}} \cdots \boldsymbol{r}_{n}^{\alpha_{n}} ; \partial_{r}^{\alpha}=(\partial / \partial r)^{\alpha}=\partial^{\alpha_{1}} / \partial r_{1}^{\alpha_{1}} \cdots \partial^{\alpha_{n}} / \partial r_{n}^{\alpha_{n}}$.
The Hilbert space $L^{2}(Q, d q)$ is represented by $\mathscr{H}$ and its scalar product is represented by $\langle$,$\rangle .$

## 1. The Weyl Correspondence

The starting idea in any pseudodifferential operator theory is to establish a 1-1 correspondence between polynomials on $X$ and differential operators on $Q: a(q, p) \leftrightarrow \hat{a}(q,-i \partial / \partial q)$ and then to extend this correspondence to larger classes of functions; the function $a$ is called the (full) symbol of the operator $\hat{a}$. The purpose of this correspondence is to study the operators $\hat{a}$ via their symbols.

The symbol map is usually such that if $\hat{a}$ has constant coefficients, $\hat{a}=$ $P\left(-i \partial / \partial q_{j}\right)$, then its symbol is the same function as that of the scalar variables $p_{j}$ $\left(a-P\left(p_{j}\right)\right)$ (thus, since the symbol of $-i \partial / \partial q_{j}$ is $p_{j}$, we can write $\hat{p}$ for $\left.-i(\partial / \partial q)\right)$.

The extension of the symbol map to differential operators with variable coefficients involves some arbitrariness. For operators on an affine space $Q$, the Weyl prescription [7] has important symmetry advantages. Given a $C^{\infty}$ function $a(q, p)$ polynomial in $p$, it defines $\hat{a}$ by its action on the functions $u \in C_{0}{ }^{\infty}(Q)$ as

$$
\begin{equation*}
(\hat{a} u)(q)=(2 \pi)^{-l} \int_{p_{\times Q}} a\left(\frac{q+q^{\prime}}{2}, p\right) u\left(q^{\prime}\right) e^{i p\left(q-q^{\prime}\right)} d p d q^{\prime} \tag{1.1}
\end{equation*}
$$

This formula defines a differential operator: Let

$$
a(q, p)=\sum_{|\alpha| \leqslant m} a_{\alpha}(q) p^{\alpha} \quad \text { and } \quad q^{\prime}=q+r
$$

Then

$$
\begin{aligned}
(\hat{a} u)(q) & =(2 \pi)^{-l} \sum_{\alpha} \int a_{\alpha}(q+r / 2) p^{\alpha} u(q+r) e^{-i p r} d p d r, \\
\text { (by parts) } & =(2 \pi)^{-l} \sum_{\alpha} \int\left(-i \frac{\partial}{\partial r}\right)^{\alpha}\left[a_{\alpha}\left(q+\frac{r}{2}\right) u(q+r)\right] e^{-i p r} d p d r \\
& =\left.\sum_{\alpha}\left(-i \frac{\partial}{\partial r}\right)^{\alpha}\left[a_{\alpha}\left(q+\frac{r}{2}\right) u(q+r)\right]\right|_{r=0}
\end{aligned}
$$

(by the Fourier inversion theorem), and the result has the form $\sum_{|\alpha| \leqslant m} c_{\alpha}(q) \partial_{q}{ }^{\alpha} u(q)$ by the Leibniz formula.

Equation (1.1) defines $\hat{a}$ as an integral operator of kernel

$$
\begin{equation*}
\langle q| \hat{a}\left|q^{\prime}\right\rangle=(2 \pi)^{-l} \int_{P} a\left(\frac{q+q^{\prime}}{2}, p\right) e^{i p\left(q-q^{\prime}\right)} d p \tag{1.2}
\end{equation*}
$$

Let $q^{\prime \prime}=\left(q+q^{\prime}\right) / 2, r=q^{\prime}-q$, corresponding to the decomposition $Q \oplus Q=$ $\Delta \oplus R(\Delta=$ diagonal subspace of $Q \oplus Q ; R=$ antidiagonal subspace of $Q \oplus Q$ ). Equation (1.2) reads

$$
\left\langle q^{\prime \prime}-r / 2\right| \hat{a}\left|q^{\prime \prime}+r / 2\right\rangle=(2 \pi)^{-l} \int_{P} a\left(q^{\prime \prime}, p\right) e^{-i p r} d p ;
$$

hence

$$
\langle q| \hat{\boldsymbol{a}}\left|q^{\prime}\right\rangle=(\mathscr{F} a)\left(q, q^{\prime}\right),
$$

where $\mathscr{F}$ is a partial Fourier transformation between the conjugate spaces $P$ and $R$, not acting on the space $\Delta$.

In the form (1.2), the Weyl correspondence can be considerably extended $[5,8]: \mathscr{F}$ is a continuous isomorphism of

$$
\mathscr{S}^{\prime}(X)\left(=\mathscr{S}^{\prime}(Q) \times \mathscr{S}^{\prime}(P)\right) \quad \text { onto } \quad \mathscr{S}^{\prime}(Q \times Q)\left(=\mathscr{S}^{\prime}(\Delta) \times \mathscr{S}^{\prime}(R)\right) ;
$$

hence to any $a \in \mathscr{S}^{\prime}(X)$ it associates an "operator" $\hat{a}$ of integral kernel in $\mathscr{S}^{\prime}(\underset{\sim}{Q} \times Q): \hat{a}$ makes sense as continuous sesquilinear form on $\mathscr{P}(Q) \times \mathscr{S}(Q)$ : $(u, v) \rightarrow\langle v| \hat{a}|u\rangle$, or equivalently as a continuous operator $\mathscr{S}(Q) \rightarrow \mathscr{S}^{\prime}(Q)$.

The map $a \rightarrow \hat{a}$ was introduced by Weyl [7] to "quantize" classical observables; in this context we have here Planck's constant $\hbar=1$, until stated otherwise. The map $\hat{a} \rightarrow a$ is a Wigner transformation [9], and we call $a(x)$ the Wigner symbol of the operator $\hat{a}$.

Mathematicians commonly associate with a given symbol $a(x)$ a different operator $a^{M}$, which is such that $[1,3]$

$$
\left(a^{M} u\right)(q)=(2 \pi)^{-l} \int_{P \times Q} a(q, p) u\left(q^{\prime}\right) e^{i p\left(q-q^{\prime}\right)} d p d q^{\prime} \quad\left(u \in C_{0}^{\infty}(Q)\right)
$$

hence

$$
\left(\sum_{|\alpha| \leqslant m} a_{\alpha}(q) p^{\alpha}\right)^{M}=\sum_{|\alpha| \leqslant m} a_{\alpha}(q)\left(-i \frac{\partial}{\partial q}\right)^{\alpha} .
$$

Another most interesting operator prescription is the Wick (or "normal") ordering $a \rightarrow: a:[10,11]$. Some properties of the correspondences $a \rightarrow a^{M}$ and $a \rightarrow: a$ : are listed in the Appendix.

### 1.1. Some Properties of the Weyl Correspondence [12, 13]

It is immediately checked that in the sense of operators on $\mathscr{P}(Q)$,

$$
\begin{aligned}
& \widehat{1}=\mathbb{1} \quad \text { (the identity operator) }, \\
& \widehat{q_{j}}=q_{j} \quad \text { (multiplication by the coordinate function } q_{j} \text { ), } \\
& \widehat{p_{j}}=-i \frac{\partial}{\partial q_{j}} .
\end{aligned}
$$

If $T(p)+V(q) \in O_{M}(X)$ we can define $\widehat{T(p)+V(q)}$ by Eq. (1.2) and $T(\hat{p})+V(\hat{q})$ by the operator functional calculus [14]: These two definitions coincide (the latter is the Schrödinger quantization for Hamiltonians of the standard type $\left.p^{2} / 2 m+V(q)\right)$.

The most interesting are the various symmetry properties (not satisfied by other symbol maps).
-For all $m \in \mathbb{N}$ and $\alpha, \beta \in \mathbb{C}$,

$$
\begin{equation*}
\widehat{(\alpha q+\beta p)^{m}}=(\alpha \hat{q}+\beta \hat{p})^{m} \tag{1.3}
\end{equation*}
$$

and this property characterizes the restriction of the Weyl correspondence to polynomials (proof in [13])
-For all $\alpha, \beta \in \mathbb{R}$,

$$
\begin{equation*}
\widehat{\exp i(\alpha q+\beta p})=\exp i(\alpha \hat{q}+\beta \hat{p}) \tag{1.4}
\end{equation*}
$$

Proof. The operator $(\alpha \hat{q}+\beta \hat{p})$ admits the Hermite functions as entire vectors [14]; for any finite linear combinations $u$ and $v$ of such vectors, $\langle u| \exp i(\alpha \hat{q}+\beta \hat{p})|v\rangle$ is defined (as $\left.\sum_{0}^{\infty}\left(i^{m} / m!\right)\langle u|(\alpha \hat{q}+\beta \hat{p})^{m}|v\rangle\right)$, and

$$
\begin{aligned}
\langle u| \exp i(\alpha \hat{q}+\beta \hat{p})|v\rangle & =\langle u| e^{i \alpha \beta / 2} \exp i \alpha \hat{q} \exp i \beta \hat{p}|v\rangle \\
& =e^{i \alpha \beta / 2} \int_{O} u^{*}(q) e^{i \alpha q} v(q+\beta) d q \\
& \left.=(2 \pi)^{-l} \int_{X} \mathscr{F}-1(v \otimes)^{*}\right) e^{i(\alpha q+\beta p)} d x \\
& =\int_{O \times Q}\left(v \otimes u^{*}\right) \mathscr{F} e^{i(\alpha q+\beta p)}=\langle u| \widehat{e^{i(\alpha q+\beta p)}}|v\rangle
\end{aligned}
$$

and this extends by continuity to $u, v \in \mathscr{S}$ (the closed linear span of the Hermite functions).

Conversely, extending (1.4) by linearity to superpositions of exponentials, one can recover the Weyl correspondence [1].

- $\widehat{\boldsymbol{a}}{ }^{*}=\hat{\boldsymbol{a}}^{\dagger}$ (in the sense of sesquilinear forms), since the kernel given by (1.2) trivially satisfies $\langle q| \widehat{a^{*}}\left|q^{\prime}\right\rangle=\left\langle q^{\prime}\right| \hat{a}|q\rangle^{*}$. By contrast, the mathematicians' symbol for the adjoint of an operator is rather complicated [2,3] (see Appendix).

In particular if $a(x)$ is real-valued, the form $\hat{a}$ is symmetric (Hermitian). There is no known characterization of those real symbols for which a selfadjoint operator (densely defined) on the Hilbert space $L^{2}(Q)$ can be associated to the form $\hat{a}$ (a symmetric operator as simple as $\hat{q} \hat{p} \hat{q}$ has no self-adjoint extension [15]). We shall give a partial solution to this problem for a class of symbols.

### 1.2. Metaplectic Covariance

Let $\mathrm{Mp}(\mathscr{H}) \rightarrow^{\pi} \mathrm{Sp}(X)$ denote the metaplectic representation of the symplectic group of $X$ [16]: It is a double covering of $\operatorname{Sp}(X)$ by a subgroup of the unitary group of $\mathscr{H}$; every $\mathscr{U} \in \mathrm{Mp}(\mathscr{H})$ is not only a unitary map on $L^{2}(Q)$ but it extends to a continuous isomorphism of $\mathscr{P}^{\prime}(Q)$ [17]. The Weyl quantization has the interesting covariance property

$$
\begin{equation*}
\forall \mathscr{U} \in \operatorname{Mp}(\mathscr{H}), \quad \forall a \in \mathscr{S}^{\prime}(X), \quad \mathscr{U} \hat{a} \mathscr{U}^{+}=\widehat{a \circ \pi(\mathscr{U})^{-1}} \tag{1.5}
\end{equation*}
$$

which implies that Weyl quantizations with respect to different symplectic frames on $X$ are isomorphically equivalent.

With minor changes in notation, we have given a proof of Eq. (1.5) in [5]. Here we mention it only as another symmetry of the Weyl quantization, which is not shared by the other ordering procedures listed in the Appendix.

### 1.3. Hilbert-Schmidt Operators

The Fourier transformation $\mathscr{F}$ introduced above is a unitary map from $L^{2}\left(X, d x /(2 \pi)^{l}\right)$ onto $L^{2}(Q \times Q, d q \times d q)$ by Parseval's identity. But the operators with $L^{2}$ kernels are the Hilbert-Schmidt operators [14] and they form a Hilbert space $\mathscr{L}^{2}(\mathscr{H})$ for the norm

$$
\langle B, A\rangle_{H S}=\operatorname{Tr} B^{\dagger} A=\int_{Q \times Q}\langle q| B\left|q^{\prime}\right\rangle^{*}\langle q| A\left|q^{\prime}\right\rangle d q d q^{\prime}
$$

Hence the Weyl correspondence restricted to $L^{2}\left(X, d x /(2 \pi)^{l}\right)$ is a unitary map onto $\mathscr{L}^{2}(\mathscr{H})$.

## 2. An Algebra of Symbols

Grossmann, Loupias, and Stein [1] have defined a space of symbols, denoted here GLS symbols, which give rise to pseudodifferential-like operators under Weyl quantization. We review and strengthen their results here. Symbols presenting some analogies have been defined by other authors [18].

### 2.1. GLS Symbols of Order $\leqslant m$

A GLS symbol of order $\leqslant m$ is any function $a \in C^{\infty}(X)$ satisfying (for a real number $m$ )

$$
\begin{equation*}
\forall \alpha \in \mathbb{N}^{2 l} \exists C_{\alpha}>0, \quad\left|\partial^{\alpha} a(x)\right| \leqslant C_{\alpha}\left(1 \mid \cdot\|x\|^{2}\right)^{(m-|\alpha|) / 2} \tag{2.1}
\end{equation*}
$$

where we let $\|x\|^{2}=q^{2}+p^{2}$ (but any other quadratic norm on $X$ would do).
The space $S_{m}$ of these symbols can be topologized by the scminorms $\|a\|_{\alpha}=$ $\sup _{x \in X}\left(\left|\partial^{\alpha} a(x)\right| /\left(1+\|x\|^{2}\right)^{(m-|\alpha|) / 2}\right)$; it is then a Fréchet space.
$S_{m}$ is stable under complex conjugation; any operator $\partial_{x}{ }^{\alpha}$ maps $S_{m}$ into $S_{m-|\alpha|}$ continuously; multiplication of functions sends $S_{m} \times S_{m^{\prime}}$ into $S_{m+m^{\prime}}$. We have $S_{m^{\prime}} \subset S_{m}$ (a continuous inclusion) if $m^{\prime} \leqslant m$. The space $S_{-\infty}=\bigcap_{m \in \mathbb{R}} S_{m}$ (with the coarsest topology making all inclusions $S_{-\infty} \subset S_{m}$ continuous) is the Schwartz space $\mathscr{S}(X)$. We denote $S=\bigcup_{m \in \mathbb{R}} S_{m}$ the space of all GLS symbols.

The main differences from the usual definitions of symbols [2,3] are: (1) The latter involve a growth condition in the $p$ coordinates only (while we take stronger growth conditions, isotropic in phase space and $\mathrm{Sp}(X)$-invariant); (2) That growth condition is expressed [3] by an asymptotic series

$$
a(q, p) \sim \sum_{j=0}^{\infty} a_{m-j}(q, p)(\|p\| \rightarrow+\infty)
$$

where each $a_{m-j}$ is positively homogeneous in $p$ of degree $=m-j$, while our condition (2.1) is only an inequality. Here we can also define subspaces of symbols $\mathbf{S}_{m} \subset S_{m}$ by a condition stronger than (2.1):
$\left.\begin{array}{rl}a(x) \sim \sum_{j=0}^{\infty} a_{m-j}(x) \text { in the sense that } \\ \forall \alpha \in \mathbb{N}^{2 l}, \forall n \in \mathbb{N}, \quad\left|\partial^{\alpha}\left(a-\sum_{j=0}^{n-1} a_{m-j}\right)\right| \leqslant \mathcal{O}\left(\|\left. x\right|^{m-n-|\alpha|}\right)\end{array}\right\} \quad(\|x\| \rightarrow+\infty)$,
where each $a_{m-j}(x)$ is positively homogeneous in $x$ of degree ( $m-j$ ); but this requirement is too strong in view of realistic physical applications (see Section 4). (Note: By convention, any sum of the form $\sum_{j=0}^{-1}(\ldots)$ has value 0 .)

### 2.2. Symbols Defined by Asymptotic Series

We say that the symbol $a \in S_{m}$ admits an asymptotic expansion $\sum_{j=0}^{\infty} a_{m-j}$ with $a_{m-j} \in S_{m-j}$, if for all $n \in \mathbb{N},\left(a-\sum_{j=0}^{n-1} a_{m-j}\right) \in S_{m-n}$; this is denoted $a \sim$ $\sum_{j} a_{m-j}$ (more general descending asymptotic expansions could be allowed). A given symbol admits an infinity of expansions of this form, but the converse is more interesting: to find a symbol having a given expansion.

Theorem 2.2.1. (i) Given any sequence $\left\{a_{m-j} \in S_{m-j}\right\}_{j \in \mathbb{N}}$ there exists $\mathbf{a} \in S_{m}$ such that $\mathbf{a} \sim \sum_{j} a_{m-j}$.
(ii) The set of symbols a such that $a \sim \sum_{j} a_{m-j}$ is the class of a modulo $S_{-\infty}$.

Proof. Let $\left\{\lambda_{j}\right\}_{j \in \mathbb{N}}$ be a sequence of positive numbers increasing to $+\infty$, and let $\chi \in C^{\infty}(\mathbb{R})$ be a monotonic function such that $\chi(r)=0$ for $r \leqslant 1$ and $\chi(r)=1$ for $r \geqslant 2$. We define

$$
\mathbf{a}(x)=\sum_{j=0}^{\infty} \chi\left(\frac{\|x\|}{\lambda_{j}}\right) a_{m-j}(x)=\sum_{j=0}^{\infty} \mathbf{a}_{m-j}(x) .
$$

In any compact subset of $X$ this sum is actually finite, so it defines a $C^{\infty}$ function $\mathbf{a}(x)$. Now we must control the behavior of $\mathbf{a}(x)$ when $\|x\| \rightarrow+\infty$ for a suitable choice of the sequence $\left\{\lambda_{j}\right\}$. For any $\alpha \in \mathbb{N}^{2 l}$ and $j \in \mathbb{N}$,

$$
\partial^{\alpha}\left[x\left(\frac{\|x\|}{\lambda}\right) a_{m-j}(x)\right]=\sum_{\beta \leqslant \alpha} \frac{c_{j \alpha \beta}}{\lambda^{\beta \mid}}\left(\partial^{\beta} \chi\right)\left(\frac{\|x\|}{\lambda}\right) \partial^{\alpha-\beta} a_{m-j}(x),
$$

by the Leibniz formula; estimate (2.1) then yields

$$
\left|\partial^{\alpha}\left[\chi\left(\frac{\|x\|}{\lambda}\right) a_{m-j}(x)\right]\right| \leqslant \sum_{\beta \leqslant \alpha} \frac{c_{j \alpha \beta}^{\prime}}{\lambda^{|\beta|}}\left[\left(\partial^{\beta} \chi\right)\left(\frac{\|x\|}{\lambda}\right)\right]\|x\|^{m-j-|\alpha|+\mid \beta_{i}} .
$$

In the RHS we can replace each $\|x\|^{|\beta|} / \lambda^{|\beta|}$ by $2^{|\beta|}$ (obvious if $|\beta|=0$; if $|\beta|>0$, it is due to $\left(\partial^{\beta} \chi\right)(\|x\| / \lambda) \equiv 0$ if $\left.\|x\|>2 \lambda\right)$. We now choose

$$
\lambda_{j} \geqslant \sup _{|\alpha| \leqslant j}\left(\sum_{|\beta| \leqslant \alpha} 2^{|\beta|} c_{j \alpha \beta}^{\prime} \cdot \sup _{r \in \mathbb{R}}\left|\partial^{\beta} \chi(r)\right|\right) ;
$$

let $\chi\left(\|x\| / \lambda_{j}\right) a_{m-j}(x)=\mathbf{a}_{m-j}(x)$. Whenever $|\alpha| \leqslant j$,

$$
\left|\partial^{\alpha} \mathbf{a}_{m-j}(x)\right| \leqslant \lambda_{j}\|x\|^{m-j-|x|} \leqslant\|x\|^{m-j-|\alpha|+1}
$$

since the left-hand side vanishes when $\|x\| \leqslant \lambda_{j}$. The series $\sum_{j} \partial^{\alpha} \mathbf{a}_{m-j}(x)$ is thus uniformly bounded by a geometric series if $\|x\|>1+\epsilon$, and we have

$$
\begin{aligned}
\left|\partial^{\alpha} \mathbf{a}(x)\right| & =\left|\partial^{\alpha} \mathbf{a}_{m}(x)+\sum_{j=1}^{\infty} \partial^{\alpha} \mathbf{a}_{m-j}(x)\right| \\
& \leqslant C_{\alpha}\|x\|^{m-|\alpha|}+\frac{\|x\|^{m-|\alpha|+1}}{\|x\|-1}=\mathcal{O}\left(\|x\|^{m-|\alpha|}\right) \text { as }\|x\| \rightarrow \infty
\end{aligned}
$$

where we have used (2.1) for $\mathbf{a}_{m}(x)\left(=a_{m}(x)\right.$ for $\left.\|x\|>2 \lambda_{0}\right)$. So we have $\mathbf{a}(x) \in S_{m}$. The same argument shows that

$$
\mathbf{a}(x)-\sum_{j=0}^{n-1} \mathbf{a}_{m-j}(x) \in S_{m-n}
$$

Now:

$$
\mathbf{a}(x)-\sum_{j=0}^{n-1} a_{m-j}(x)=\left(\mathbf{a}(x)-\sum_{0}^{n-1} \mathbf{a}_{m-j}\right)+\sum_{0}^{n-1}\left(\mathbf{a}_{m-j}-a_{m-j}\right)
$$

and the second term vanishes identically for $\|x\| \geqslant 2 \lambda_{n 1}$, owing to the form of the cutoff function $\chi$. Hence $\mathbf{a}-\sum_{j=0}^{n-1} a_{m-j} \in S_{m-n}$, proving (i).
(ii) $a \sim \sum_{j} a_{m-j}$ means that for all $n, a-\sum_{j=0}^{n-1} a_{m-j} \in S_{m-n}$, which is equivalent to $\left(a-\sum_{j=0}^{n-1} a_{m-j}\right)-\left(\mathbf{a}-\sum_{j=0}^{n-1} a_{m-j}\right) \in S_{m-n}$, or $(a-\mathbf{a}) \in S_{m-n}$ for all $n$, i.e., $a-\mathbf{a} \in S_{-\infty}$.

This motivates the following:
Definition 2.2.2. We say that $a_{1} \sim a_{2}$ if $a_{1}-a_{2} \in S_{-\infty}$; the class of $a \in S_{m}$ in the quotient space $\tilde{S}_{m}=S_{m} / S_{-\infty}$, or asymptotic class of $a$, is denoted $\tilde{a}$; we let $\tilde{S}=S / S_{-\infty}=\bigcup_{m} \tilde{S}_{m}$. (In this language, Theorem 2.2.1 just states that any formal series $\sum_{j=0}^{\infty} a_{m-j}$ with $a_{m-j} \in S_{m-j}$ defines a (unique) element of $\tilde{S}_{m}$.)

### 2.3. Quantization of Symbols

Every GLS symbol $a(x)$ is an element of $O_{M}(X) \subset \mathscr{S}^{\prime}(X)$; hence by Weyl quantization it defines a continuous map $\hat{a}: \mathscr{S}(X) \rightarrow \mathscr{S}^{\prime}(X)$; we call $\hat{S}_{m}$ (resp. $\hat{S}$ ) the image of $S_{m}$ (resp. $S$ ) under ${ }^{\wedge}$. The properties of the operators $\hat{a} \in \hat{S}$ are discussed in Section 3, but we give a preliminary result here.

Thforfm 2.3.1. (i) à maps $\mathscr{P}(Q)$ into $\mathscr{P}(Q)$ continuously; hence by duality it extends to a continuous operator $\mathscr{P}^{\prime}(Q) \rightarrow \mathscr{P}^{\prime}(Q)$ (also denoted â).
(ii) If $a \in S_{m}$ with $m<-l$, a maps $\mathscr{H}=L^{2}(Q)$ into itself, and $\hat{a} \in \mathscr{L}^{2}(\mathscr{H})$.
(iii) If $a \in S_{-\infty}$, $\hat{a}$ maps $\mathscr{P}^{\prime}(Q)$ into $\mathscr{S}(Q)$ continuously.

Proof. (i) Let $a \in S_{m}$ and $u \in \mathscr{P}(Q)$. In the sense of distributions,

$$
\begin{align*}
(\hat{a} u)(q) & =(2 \pi)^{-i} \int_{P \times Q} d p d r a(q+r / 2, p) u(q+r) e^{-i p r} \\
& =\int_{P \times Q} d p d r I(q, p, r) e^{-i p r} \tag{2.3}
\end{align*}
$$

and $\hat{a} u \in \mathscr{S}(Q)$ if every distribution $q^{\alpha} \partial^{a}(\hat{a} u)$ is actually a bounded function. Since $\partial^{B}(\hat{\alpha} u)=(2 \pi)^{-l} \int d p d r \partial_{q}{ }^{\beta}(a(q+r / 2, p) \cdot u(q+r)) e^{-i p r}$ is a sum of integrals of the same type as (2.3), it suffices to prove that $q^{\alpha}(\hat{a} u)(q)$ is a bounded function of $q$.

To do this, we shall exploit the oscillatory character of the integrand (due to $e^{-i p r}$ ) as $\|p\|$ and $\|r\| \rightarrow \infty$, by the well-known method of successive integrations by parts.

First, for any integer $k$,

$$
\begin{aligned}
\hat{a} u(q) & =\int_{Q} d r I(q, p, r) e^{-i p r}=\int_{Q} d r \frac{I(q, p, r)}{\left(1+\|p\|^{2}\right)^{k}}\left(1+\|p\|^{2}\right)^{k} e^{-i p r} \\
& =\int_{Q} d r \frac{I(q, p, r)}{\left(1+\|p\|^{2}\right)^{k}}\left[\left(1-\Delta_{q}\right)^{k} e^{-i p r}\right]=\int_{Q} d r\left(1-\Delta_{Q}\right)^{k} \frac{I(q, p, r)}{\left(1+\|p\|^{2}\right)^{k}} e^{-i p r}
\end{aligned}
$$

(by parts, using the rapid decrease of all $\partial_{r}{ }^{\alpha} I(q, p, r)$ as $\left.\|r\|>\infty\right)$. We takc $k$ large enough to let $\left(1-\Delta_{r}\right)^{k}\left(I(q, p, r) /\left(1+\|p\|^{2}\right)^{k}\right)=J(q, p, r)$ be absolutely integrable in $p$; this is possible because $a(x)$ satisfies condition (2.1).

Then, for any integer $k^{\prime}$ we get, similarly,

$$
\begin{equation*}
\int_{P \times Q} d p d r J(q, p, r) e^{-i p r}=\int_{P \times O} d p d r\left(1-\Delta_{p}\right)^{k^{\prime}} \frac{J(q, p, r)}{\left(1+\|r\|^{2}\right)^{k^{\prime}}} e^{-i p r} \tag{2.4}
\end{equation*}
$$

Thus we have reduced ourselves to (definite) integrals of the form

$$
\int K(q, p, r) d p d r=\int d p d r \frac{a^{\prime}(q+r / 2, p)\left(\partial^{\beta} u\right)(q+r) e^{-i p r}}{\left(1+\|p\|^{2}\right)^{k}\left(1+\|r\|^{2}\right)^{k^{\prime}}} \quad\left(a^{\prime} \in S_{m}\right)
$$

since (2.4) is a sum of such terms.
We now bound the integrand $K(q, p, r)$ from above. If $a^{\prime} \in S_{m}$, let $m_{+}=$ $\max \{m / 2,0\}$; by $(2.1):\left|a^{\prime}(q, p)\right| \leqslant \mathcal{O}\left[\left(1+\|q\|^{2}+\|p\|^{2}\right)^{m_{+}}\right] \leqslant \mathcal{O}\left[\left(1+\|q\|^{2}\right)^{m_{+}}\right.$ $\left.\left(1+\|p\|^{2}\right)^{m_{+}}\right]$. Also, for every $N \in \mathbb{N}:\left|\partial^{\beta} u(q)\right| \leqslant C_{N \beta}^{\prime}\left(1+\|q\|^{2}\right)^{-N}$. Lemma 2.3.2 (below) then yields, for $\|q\|$ or $\|r\| \rightarrow \infty$,

$$
\begin{aligned}
a^{\prime}(q+r / 2, p) \mid & \leqslant \mathcal{O}\left[\left(1+\|q\|^{2}\right)^{m_{+}}\left(1+\|r\|^{2}\right)^{m_{+}}\left(1+\|p\|^{2}\right)^{m_{+}}\right] \\
\left|\partial^{\beta} u(q+r)\right| & \leqslant \mathcal{O}\left[C_{N \beta}^{\prime}\left(1+\|q\|^{2}\right)^{-N}\left(1+\|r\|^{2}\right)^{N}\right]
\end{aligned}
$$

hence,

$$
\begin{aligned}
& |K(q, p, r)| \leqslant \mathcal{O}\left[C_{N \beta}^{\prime}\left(1+\|q\|^{2}\right)^{m_{+}-N}\left(1+\|p\|^{2}\right)^{m_{+}-k}\left(1+\|r\|^{2}\right)^{m_{+}+N-k^{\prime}}\right. \\
\Rightarrow & \int K d p d r \leqslant \mathcal{O}\left[C_{N \beta}^{\prime}\left(1+\|q\|^{2}\right)^{-M}\right] \quad \text { for all } \quad M \in \mathbb{N},
\end{aligned}
$$

by a suitable choice of $N\left(\geqslant M+m_{+}\right)$and of $k$ and $k^{\prime}$ (depending on $N$ ). More generally we would obtain, for all $M, \alpha$,

$$
\left|\left(1+\|q\|^{2}\right)^{M} \partial^{\alpha}(\hat{a} u)(q)\right| \leqslant c C_{N \beta}^{\prime}=c \cdot \sup _{q \in Q}\left\{\left|\left(1+\|q\|^{2}\right)^{N} \partial^{\beta} u(q)\right|\right\}
$$

for some $N, \beta$; hence $\mathscr{S} \rightarrow{ }^{\hat{a}} \mathscr{S}$ is continuous.
The adjoint operator of $\mathscr{S} \rightarrow^{\hat{a}} \mathscr{S}$ is then an extension of $\hat{a}$ to a continuous operator $\mathscr{S}^{\prime} \rightarrow{ }^{a} \mathscr{S}^{\prime}$.

Proof of (ii). Trivial, since $m<-l$ implies $S_{m} \subset L^{2}(X)$, and by Section 1.3 Weyl quantization sends $L^{2}(X)$ onto $\mathscr{L}^{2}(\mathscr{H})$.

Proof of (iii). If $a \in S_{-\infty}=\mathscr{S}(X)$, the kernel of $\hat{a}$ is in $\mathscr{S}(Q \times Q)$ by Fourier transformation, it then defines a continuous sesquilinear functional on $\mathscr{S}^{\prime}(Q)$, i.e., a continuous map $\mathscr{S}^{\prime}(Q) \rightarrow \mathscr{S}(Q)$. Operators of this type are thus "infinitcly regularizing."

The proof of (i) required the following:
Lemma 2.3.2. For all $q, r \in Q$,

$$
\begin{equation*}
\left(1+\|q+r\|^{2}\right)^{ \pm 1}<2\left(1+\|q\|^{2}\right)^{ \pm 1}\left(1+\|r\|^{2}\right) . \tag{2.5}
\end{equation*}
$$

Proof. For $(+1)$ we have

$$
\left(1+\|q+r\|^{2}\right)<2\left(1+\|q\|^{2}+\|r\|^{2}\right) \leqslant 2\left(1+\|q\|^{2}\right)\left(1\| \| r \|^{2}\right) .
$$

For $(-1)$ the result follows from

$$
\left(1+\|q+r\|^{2}\right) \geqslant 1+(\|q\|-\|r\|)^{2}>\frac{1}{2}\left(\frac{1+\|q\|^{2}}{1+\|r\|^{2}}\right),
$$

since

$$
\begin{aligned}
& 2\left[1+\left(\|q\|-\|r\|^{2}\right)\right]\left(1+\|r\|^{2}\right)-\left(1+\|q\|^{2}\right) \\
& \quad=1+(\|q\|-2\|r\|)^{2}+2\|r\|^{2}\|q\|-\|r\|^{2}>0 .
\end{aligned}
$$

### 2.4. The Twisted Product of Symbols

By Theorem 2.3.1(i), the operator product $\hat{a} \hat{b}$ of two GLS operators makes sense as a continuous mapping $\mathscr{S}(Q) \rightarrow \mathscr{S}(Q)$. âb has a kernel in $\mathscr{S}^{\prime}(Q \times Q)$, and by isomorphism (1.2) there is a unique $c \in \mathscr{S}^{\prime}(X)$ such that $\hat{c}=\hat{a} \hat{b} . c$ is called the twisted product of $a$ and $b$; it can be computed by expressing $(\hat{a} \hat{b}) u=\hat{c} u$ (for $u \in \mathscr{S}(Q)$ ) via the transformation $\mathscr{F}$ (Eq. (1.2)). By a straightforward computation,

$$
\begin{align*}
c(x) & =\mathscr{F}^{-1}\left(\langle q| \hat{a} b\left|q^{\prime}\right\rangle\right)=\mathscr{F}-1\left[\int_{Q} d q^{\prime \prime}\langle q| \hat{a}\left|q^{\prime \prime}\right\rangle\left\langle q^{\prime \prime}\right| b\left|q^{\prime}\right\rangle\right] \\
& =\mathscr{F}^{1}\left(\int_{O} d q^{\prime \prime}(\mathscr{F} a)\left(q, q^{\prime \prime}\right)(\mathscr{F} b)\left(q^{\prime \prime}, q^{\prime}\right)\right) \\
& =\pi^{-2 l} \int_{X^{b}} d q_{1} d p_{1} d q_{2} d p_{2} a\left(q+q_{1}, p+p_{1}\right) b\left(q+q_{2}, p+p_{2}\right) e^{2 i\left(q_{1} p_{2}-q_{2} p_{1}\right)} \\
& =\pi^{-2 l} \int_{X^{l}} d x_{1} d x_{2} a\left(x+x_{1}\right) b\left(x+x_{2}\right) e^{2 i \omega\left(x_{1}, x_{2}\right)} \tag{2.6}
\end{align*}
$$

which illustrates the special role played by the symplectic form $\omega$ of $X$ in the Weyl correspondence $[1,8,9]$. The integral (2.6) is defined in the sense of distributions.

Theorem 2.4.1. (i) $S$ is an associative *-algebra for twisted multiplication, which is a continuous map from $S_{m_{1}} \times S_{m_{2}}$ into $S_{m_{1}+m_{2}}$ (for any $m_{1}, m_{2} \geqslant-\infty$ ). The twisted product $c$ of $a$ and $b \in S$ is given by Eq. (2.6).
(ii) $c(x)$ admits the asymptotic expansion (Groenewold [19])

$$
\begin{equation*}
c(x) \sim \sum_{n=0}^{\infty} \frac{1}{n!}\left(\frac{i}{2}\right)^{n}\left[\left(\partial_{q_{1}} \partial_{p_{2}}-\partial_{a_{2}} \partial_{p_{1}}\right)^{n} a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right]_{x_{1}=x_{2}=0} \tag{2.7}
\end{equation*}
$$

denoted $c(x) \sim a(x)\left[\exp (i / 2)\left(\stackrel{\partial}{\partial}_{q} \partial_{p}-\hat{\partial}_{p} \partial_{q}\right)\right] b(x)$.
(iii) If $a$ or $b$ is a polynomial in $x$, expansion (2.7) is a finite sum which is exactly equal to $c(x)$.

Corollary 2.4.2. (i) $S$ is a Lie algebra for the twisted commutator $(a, b) \rightarrow$ $\{\{a, b\}\}=c$ defined $b y$

$$
\begin{equation*}
\hat{c}=[\hat{a}, b]_{-} i, \quad \text { and } \quad\left\{\left\{S_{m_{1}}, S_{m_{2}}\right\}\right\} \subset S_{m_{1}+m_{2}-2} \tag{2.8}
\end{equation*}
$$

(ii) $\{\{a, b\}\}(x) \sim a(x)\left[2 \sin \left(\left(\dot{\partial}_{q} \dot{\partial}_{p}-\dot{\partial}_{p} \dot{\partial}_{q}\right) / 2\right)\right] b(x)$.
(iii) If $a$ or $b$ is a polynomial in $x$ of degree $\leqslant 2$,

$$
\{\{a, b\}\} \equiv\{a, b\} \quad \text { (the ordinary Poisson bracket) }
$$

(The operation $\{\{\}$,$\} is the Moyal bracket [20].)$
Lemma 2.4.3. For any polynomial $Q\left(x_{1}, x_{2}\right)$ there exists a differential operator $P\left(\partial_{x_{1}}, \partial_{x_{2}}\right)$ (with constant coefficients) such that

$$
Q\left(x_{1}, x_{2}\right) e^{2 i \omega\left(x_{1}, x_{2}\right)} \equiv P\left(\partial_{x_{1}}, \partial_{x_{2}}\right) e^{2 i \omega\left(x_{1}, x_{2}\right)}
$$

Proof. For convenience we denote $Y=X^{2}$; let $y=\left(x_{1}, x_{2}\right) \in Y$. The statement to be proved is invariant by linear transformations of $Y$; but there exists such a transformation under which $e^{2 i \omega\left(x_{1}, x_{2}\right)}$ becomes a nondegenerate gaussian $\exp \left(i \sum_{j=1}^{4 l} \epsilon_{j} y_{j}^{2}\right)\left(\forall j: \epsilon_{j}= \pm 1\right)$ (this follows solely from the nondegeneracy of the 2 -form $\omega$ ). Then it is known that the Hermite-like polynomials

$$
\left\{\exp \left(-i \sum \epsilon_{j} y_{j}^{2}\right) \partial_{y^{\alpha}}{ }^{\alpha} \exp \left(i \sum \epsilon_{j} y_{j}^{2}\right)\right\} \quad \alpha \in \mathbb{N}^{4 l}
$$

generate all polynomials in $y$; hence the result.
Proof of Theorem 2.4.1. (i) We need only prove the continuity of the mapping (2.6): $(a, b) \rightarrow c$ from $S_{m_{1}} \times S_{m_{2}}$ into $S_{m_{1}+m_{2}}$; the other algebraic properties are those of the corresponding operator multiplication. We thus want

$$
\begin{equation*}
\partial^{\alpha} c(x)=\pi^{-2 l} \int_{X^{2}} d x_{1} d x_{2} \partial_{x}^{\alpha}\left[a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right] e^{2 i \omega\left(x, x_{2}\right)} \tag{2.9}
\end{equation*}
$$

and

$$
\sup _{x \in X} \frac{\left\|\partial^{\alpha} c(x)\right\|}{\left(1+\|x\|^{2}\right)^{\left(m_{1}+m_{2}-|\alpha| / / 2\right.}}
$$

bounded by some seminorms of $a$ and $b$.
We shall prove both statements by defining the RHS of (2.9) as an oscillatory integral (the derivation under $\int$ in (2.9) is then legitimate). The proof is essentially the same as for Theorem 2.3.1. By the Leibniz formula, the RHS of (2.9) is a sum of terms of the form

$$
\begin{align*}
& \int d x_{1} d x_{2} \partial_{x}^{\alpha_{1}} a\left(x+x_{1}\right) \partial_{x}^{\alpha_{2}} b\left(x+x_{2}\right) e^{2 i \omega\left(x_{1} . x_{2}\right)} \quad\left(\alpha_{1}+\alpha_{2}=\alpha\right) \\
& \quad=\int d x_{1} d x_{2} \frac{\partial_{x}^{\alpha_{1}} a\left(x+x_{1}\right)}{\left(1+\left\|x_{1}\right\|^{2}\right)^{k_{1}}} \cdot \frac{\partial_{x}^{\alpha_{2}} b\left(x+x_{2}\right)}{\left(1+\left\|x_{2}\right\|^{2}\right)^{k_{2}}}\left[P_{k_{1}, k_{2}}\left(\partial_{x_{1}}, \partial_{x_{2}}\right) e^{2 i \omega\left(x_{1}, x_{2}\right)}\right] \tag{2.10}
\end{align*}
$$

for all $k_{1}, k_{2} \in \mathbb{N}$, where $P_{k_{1}, k_{2}}$ is the polynomial in $\partial_{x_{1}}$ and $\partial_{x_{2}}$ such that $P_{k_{1}, k_{2}}\left(\partial_{x_{1}}, \partial_{x_{2}}\right) e^{2 i \omega\left(x_{1}, x_{2}\right)}=\left(1+\left\|x_{1}\right\|^{2}\right)^{k_{1}}\left(1+\left\|x_{2}\right\|^{2}\right)^{k_{2}} e^{2 i \omega\left(x_{1}, x_{2}\right)}$ (by Lemma 2.4.3). In the sense of distributions, (2.10) means

$$
\int_{X^{1}} d x_{1} d x_{2} P_{k_{1}, k_{2}}\left(-\partial_{x_{1}},-\partial_{x_{2}}\right)\left[\frac{\partial_{x}^{\alpha_{1}} a\left(x+x_{1}\right)}{\left(1+\left\|x_{1}\right\|^{2}\right)^{k_{1}}} \cdot \frac{\partial_{x}^{\alpha_{2}} b\left(x+x_{2}\right)}{\left(1+\left\|x_{2}\right\|^{2}\right)^{k_{2}}}\right]^{2 i \omega\left(x_{1}, x_{2}\right)}
$$

and this integral converges for $k_{1}$ and $k_{2}$ large enough, since $a$ and $b$ are GLS symbols. Since the derivations $\partial_{x_{1}}$ and $\partial_{x_{2}}$ can only improve the behavior of the integrand $L\left(x, x_{1}, x_{2}\right)$, the latter can be bounded by

$$
\begin{aligned}
& \mathcal{O}\left[\frac{\partial_{x}^{\alpha_{1}} a\left(x+x_{1}\right)}{\left(1+\left\|x_{1}\right\|^{2}\right)^{k_{1}}} \cdot \frac{\partial_{x}^{\alpha_{2}} b\left(x+x_{2}\right)}{\left(1+\left\|x_{2}\right\|^{2}\right)^{k_{2}}}\right] \\
& \quad \leqslant \mathcal{O}\left[\frac{\|a\|_{\alpha_{1}}\|b\|_{\alpha_{2}}\left(1+\left\|x+x_{1}\right\|^{2}\right)^{\left(m_{1}-\left|\alpha_{1}\right| \mid / 2\right.}}{\left(1+\left\|x_{1}\right\|^{2}\right)^{k_{1}}} \times \frac{\left(1+\left\|x+x_{2}\right\|^{\left(m_{2}-\left|\alpha_{2}\right|\right\rangle / 2}\right)}{\left(1+\left\|x_{2}\right\|^{2}\right)^{k_{2}}}\right]
\end{aligned}
$$

and using Lemma 2.3.2,

$$
\begin{aligned}
\leqslant \mathcal{O}\left[\|a\|_{\alpha_{1}}\|b\|_{\alpha_{2}}(1+\right. & \left.\|x\|^{2}\right)^{\left(m_{1}+m_{2}-|x|\right) / 2} \\
& \left.\times\left(1 \quad \mid \cdot\left\|x_{1}\right\|^{2}\right)^{\left|\left(m_{1}-\left|\alpha_{1}\right|\right) / 2\right|-k_{1}}\left(1+\left\|x_{2}\right\|^{2}\right)^{\left|\left(m_{2}-\left|\alpha_{2}\right|\right) / 2\right|-k_{2}}\right]
\end{aligned}
$$

This upper bound is also ( $x_{1}, x_{2}$ )-integrable for $k_{1}$ and $k_{2}$ large enough, so that we get (with $m: m_{1}+m_{2}$ )

$$
\|\boldsymbol{c}\|_{\alpha}=\sup _{x} \frac{\left|\partial^{\alpha} c(x)\right|}{\left(1+\|x\|^{2}\right)^{(m-|\alpha|) / 2}} \leqslant \text { const } \times \max _{\alpha_{1}+\alpha_{2}=\alpha}\left\{\|a\|_{\alpha_{1}} \cdot\|b\|_{\alpha_{2}}\right\} \text {. Q.E.D. }
$$

The proof of (ii) essentially amounts to a stationary phase expansion of integral (2.6) around the critical point $x_{1}=x_{2}=0$ (we recall (cf. Lemma 2.4.3) that $e^{2 i \omega\left(x_{1}, x_{2}\right)}$ is a Gaussian phase function around $\left.y=\left(x_{1}, x_{2}\right)=0\right)$.

We first replace $a\left(x+x_{1}\right) b\left(x+x_{2}\right)$ by its Taylor polynomial to some order $n$ :

$$
\left[a\left(x+x_{1}\right) b\left(x+x_{1}\right)\right]_{n}=\sum_{|\alpha| \leqslant n} \frac{y^{\alpha}}{\alpha!} \partial_{y}^{\alpha}\left[a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right]_{y=0}
$$

The remainder $R_{n}(x, y)$ has the form

$$
\sum_{|\alpha|=n+1} \frac{y^{\alpha}}{\alpha!} \partial_{v}{ }^{\alpha}\left[a\left(x+\theta x_{1}\right) b\left(x+\theta x_{2}\right)\right] \quad \text { for some } \quad \theta \in[0,1]
$$

hence by Lemma 2.3.2, $\left|R_{n}(x, y)\right| \leqslant \mathcal{O}\left(\|x\|^{m-n-1}\|y\|^{k}\right)$ for some $k>0$; similarly, $\left|\partial_{x}{ }^{\beta} R_{n}(x, y)\right| \leqslant \mathcal{O}\left(\|x\|^{m-n-|\beta|-1}\|y\|^{k(\beta)}\right)$ for some $k(\beta)$, because $\partial_{x}{ }^{\beta} R_{n}$ is the remainder of the Taylor expansion in $y$ of $\partial_{x}{ }^{\beta}\left(a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right)$. By the same argument as that in the proof of (i), all this implies

$$
\int_{X^{2}} R_{n}(x, y) e^{2 i \omega\left(x_{1} x_{2}\right)} d x_{1} d x_{2} \in S_{m-n-1}
$$

It remains to compute $\pi^{-2 l} \int f(x, y) e^{2 i \omega\left(x_{1}, x_{2}\right)} d y$ for $f$ polynomial in $y$ of degree $\leqslant n$. For each monomial $y^{9}\left(\beta \in \mathbb{N}^{4 l}\right), \int\left(y^{\beta} / \beta!\right) e^{2 i \omega\left(x_{1}, x_{2}\right)} d y$ is the coefficient of $(i \xi)^{\beta}$ in

$$
\begin{aligned}
\int e^{i \xi y} e^{2 i \omega\left(x_{1} x_{2}\right)} d y & =\int e^{i\left(\xi_{1} x_{1}+\xi_{2} x_{2}+2 \omega\left(x_{1} x_{2}\right)\right)} d y \\
& =e^{-(i / 2) \omega\left(\# \xi_{1} \# \xi_{2}\right)} \int e^{2 i \omega\left(x_{1}+\# \xi_{2} / 2 . x_{2}-\# \xi \xi_{1} / 2\right)} d y \\
& =\pi^{22} e^{-(i / 2)\left(\xi_{2} n_{2}-\zeta_{2} n_{1}\right)}
\end{aligned}
$$

(where $\xi=\left(\xi_{1}, \xi_{2}\right)-\left(\zeta_{1}, \eta_{1}, \zeta_{2}, \eta_{2}\right) \in \mathbb{R}^{4 l}$ is dual to $y=\left(q_{1}, p_{1}, q_{2}, p_{2}\right)$, and $\#: X^{*} \rightarrow X$ is the isomorphism such that $\left.\xi \cdot x \equiv \omega(\# \xi, x) \forall x \in X \forall \xi \in X^{*}\right)$. Hence,

$$
\pi^{-2 l} \int \frac{y^{\beta}}{\beta!} e^{2 i \omega\left(x_{1} x_{2}\right)} d y=\frac{1}{\beta!}\left[-i \frac{\partial}{\partial \xi}\right]^{\beta}\left[e^{-(i / 2)\left(\xi_{1} n_{2}-\zeta_{2} n_{1}\right)}\right]_{\xi=0}
$$

By linearity we get for any $f(x, y)$ polynomial in $y$ of degree $n$ :

$$
\pi^{-2 l} \int f(x, y) e^{2 i \omega\left(x_{1}, x_{2}\right)}=\sum_{|\beta| \leqslant n} \frac{1}{\beta!}\left[-i \frac{\partial}{\partial \xi}\right]^{\beta}\left[e^{-(i / 2)\left(\left(_{1} \eta_{2}-\epsilon_{2} \eta_{1}\right)\right.}\right] \cdot \partial_{y}^{\beta} f(x, y)_{y \sim 0}
$$

which is formally equal (by Taylor's formula) to

$$
\begin{aligned}
{\left[\left.e^{-(i / 2)\left(\tau_{1} \eta_{2}-\tau_{2} \eta_{1}\right)}\right|_{\xi=-i \partial_{y}}\right] \cdot f_{y=0} } & =\exp (i / 2)\left(\partial_{\alpha_{1}} \partial_{p_{2}}-\partial_{q_{2}} \partial_{p_{1}}\right) f(x, y)_{y=0} \\
& =\sum_{\substack{k \in \mathbb{N} \\
k \leqslant n / 2}} \frac{1}{k!}\left(\frac{i}{2}\right)^{k}\left[\left(\partial_{q_{1}} \partial_{p_{2}}-\partial_{\alpha_{2}} \partial_{p_{1}}\right)^{k} f(x, y)\right]_{y=0}
\end{aligned}
$$

Summing up: For each $n \in \mathbb{N}$ we have

$$
\begin{aligned}
\pi^{-2 l} & \int d x_{1} d x_{2} a\left(x+x_{1}\right) b\left(x+x_{2}\right) e^{2 i \omega\left(x_{1}, x_{2}\right)} \\
& =\pi^{-2 l} \int d x_{1} d x_{2}\left[a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right]_{n} e^{2 i \omega\left(x_{1}, x_{2}\right)} \quad\left(\bmod S_{m-n-1}\right] \\
& =\sum_{k \leqslant n / 2} \frac{1}{k!}\left(\frac{i}{2}\right)^{k}\left(\partial_{q_{1}} \partial_{p_{2}}-\partial_{a_{2}} \partial_{p_{1}}\right)^{k}\left[a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right]_{x_{1}=x_{2}=0} \quad\left[\bmod S_{m-n-1}\right] \\
& =\sum_{k \leqslant n / 2} c_{k}(x) \quad\left[\bmod S_{m-n-1}\right] .
\end{aligned}
$$

Hence the twisted product satisfies $c(x) \sim \sum_{k=0}^{\infty} c_{k}(x)$ with $c_{k} \in S_{m-2 k}$, in the sense of Section 2.2.

This is precisely what is meant by Eq. (2.7).
Q.E.D.
(iii) If $a(x)$ is a polynomial, $a\left(x+x_{1}\right)$ is a polynomial in $x_{1}$, and $c(x)$ is the sum of the finite series

$$
\begin{aligned}
& \pi^{-2 l} \sum_{\gamma, \gamma^{\prime}} \int d x_{1} d x_{2} \partial_{q}^{\nu} \partial_{p}^{\gamma^{\prime}} a(x) \frac{q_{1}^{\nu} p_{1}^{\gamma^{\prime}}}{\gamma!\gamma^{\prime}!} b\left(x+x_{2}\right) e^{2 i\left(a_{1} p_{2}-q_{2} p_{1}\right)} \\
&=(\text { by parts }) \pi^{-2 l} \sum_{\gamma \gamma^{\prime}} \frac{\partial_{q}^{\gamma} \partial_{p}^{\gamma^{\prime}} a(x)}{\gamma!\gamma^{\prime}!} \int d x_{1} d x_{2} \\
& \times\left[\left(\frac{i}{2} \frac{\partial}{\partial p_{2}}\right)^{\gamma}\left(-\frac{i}{2} \frac{\partial}{\partial q_{2}}\right)^{\gamma^{\prime}} b\left(x+x_{2}\right)\right] e^{2 i\left(a_{1} p_{2}-q_{2} p_{1}\right)} \\
&= \sum_{\gamma \gamma^{\prime}} \frac{\partial_{q}^{\gamma} \partial_{p}^{\gamma^{\prime}} a(x)}{\gamma!\gamma!} \int d x_{2}\left[\left(\frac{i}{2} \frac{\partial}{\partial p_{2}}\right)^{\gamma}\left(-\frac{i}{2} \frac{\partial}{\partial q_{2}}\right)^{\gamma^{\prime}} b\left(x+x_{2}\right)\right] \delta\left(p_{2}\right) \delta\left(q_{2}\right) \\
&=\left(\sum_{\gamma} \frac{1}{\gamma!}\left(\frac{i}{2} \partial_{a_{1}} \partial_{p_{2}}\right)^{\gamma}\right)\left(\sum_{\gamma^{\prime}} \frac{1}{\gamma^{\prime}!}\left(-\frac{i}{2} \partial_{p_{1}} \partial_{q_{2}}\right)^{\gamma^{\prime}}\right)\left[a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right]_{x_{1}=x_{2}=0} \\
& c(x)=\sum_{n=0}^{\infty} \frac{1}{n!}\left(\frac{i}{2}\right)^{n}\left(\partial_{a_{1}} \partial_{p_{2}}-\partial_{a_{2}} \partial_{p_{1}}\right)^{n}\left[a\left(x+x_{1}\right) b\left(x+x_{2}\right)\right]_{x_{1}=x_{2}=0} \\
& \text { (a finite sum) Q.E.D. }
\end{aligned}
$$

Remarks. A formal proof of (ii) (which also works for any stationary phase expansion) is to expand $e^{2 i \omega\left(x_{1}, x_{2}\right)} /(2 \pi)^{l}$ in a series of distributions

$$
(2 \pi)^{-l} e^{2 i \omega\left(x_{1}, x_{2}\right)} \sim \sum_{n=0}^{\infty} \frac{1}{n!}\left(\frac{i}{2}\right)^{n}\left(\partial_{a_{1}} \partial_{p_{2}}-\partial_{q_{2}} \partial_{p_{1}}\right)^{n} \delta\left(x_{1}\right) \delta\left(x_{2}\right)
$$

understood as the formal Fourier transform of the Taylor series

$$
e^{-(i / 2)\left(\zeta_{1} n_{2}-\zeta_{2} \eta_{1}\right)} \sim \sum_{n=0}^{\infty} \frac{1}{n!}\left(-\frac{i}{2}\right)^{n}\left(\zeta_{1} \eta_{2}-\zeta_{2} \eta_{1}\right)^{n}
$$

that series is then substituted into (2.6) to yield expansion (2.7).
-The use of a different symbol map (see Appendix) leads to a different (less symmetrical) operation on symbols corresponding to operator multiplication.

The proof of Corollary 2.4.2 is straightforward.
Proposition 2.4.4. The following subspaces of $S$ are subalgebras for twisted multiplication:
-the subspace $\mathbf{S}$ defined by condition (2.2),
-the subspace of polynomials in $x$,
-the subspace of symbols $a(q, p)$ polynomial in $p$, i.e., symbols of standard partial differential operators on $\mathscr{S}(Q)($ resp. polynomial in $q)$.

Proposition 2.4.5. $\quad S_{-\infty}$ is a tro-sided ideal of $S$, hence $\tilde{S}=S / S_{-\infty}$ is an algebra for "twisted multiplication modulo $S_{-\infty}$ " (defined by (2.7) [1].

We omit the proofs (elementary).

### 2.5. The Inversion of Symbols

We seek conditions under which the inverse $\hat{a}^{-1}$ of a GLS operator $\hat{a} \in \hat{S}$ exists in $S$ and can be computed. We shall actually solve this problem modulo $S_{-\infty}$. If $\hat{a}^{-1}=\hat{b}$, the asymptotic class $\hat{b}$ of the symbol $b$ is the inverse for twisted multiplication $\bmod \widehat{S_{-\infty}}$ of $\tilde{a}$ in the sense that

$$
\begin{equation*}
\tilde{a} \exp \frac{i}{2}\left(\grave{\partial}_{q} \vec{\partial}_{p}-\hat{\partial}_{p} \ddot{\partial}_{q}\right) \tilde{b}=1 \tag{2.11}
\end{equation*}
$$

We look for $\tilde{b}$ in the form of a descending expansion $\tilde{b}(x)=\sum_{j=0}^{\infty} b_{-2 j}(x)$ with $b_{-2 j}(x) \in S_{m^{\prime}-2 j}$ (for some $m^{\prime}$ ) computed recursively so as to satisfy (the parametrix method):

$$
\begin{gathered}
a b_{0}=1 \\
a b_{-2}+(i / 2) u\left(\hat{\partial}_{q} \vec{\partial}_{p}-\hat{\partial}_{p} \vec{\partial}_{q}\right) b_{0}=0 \\
a b_{-2 j}+\sum_{k=0}^{j \ldots-1} \mathscr{A}_{j-k}\left(x, \partial_{x}\right) b_{-2 k}=0
\end{gathered}
$$

where $\mathscr{A}_{k}\left(x, \partial_{x}\right)$ is the differential operator on $X$ defined by

$$
\mathscr{A}_{k}\left(x, \partial_{x}\right) u=a\left[\frac{1}{k!}\left(\frac{i}{2}\right)^{k}\left(\partial_{q} \vec{\partial}_{p}-\ddot{\partial}_{p} \vec{\partial}_{q}\right)^{k}\right] u \quad(\forall u \in \mathscr{S}) ;
$$

i.e.,

$$
\begin{equation*}
\mathscr{A}_{k}\left(x, \partial_{x}\right)=\sum_{|\gamma|+\left|\gamma^{\prime}\right|=k} \frac{1}{\gamma!\gamma^{\prime}!}\left(\frac{i}{2}\right)^{k}(-1)^{\left|\nu^{\prime}\right|}\left(\partial_{q}{ }^{\gamma} \partial_{p}^{\gamma^{\prime}} a\right) \partial_{q}^{\gamma^{\prime}} \partial_{p}{ }^{\gamma} . \tag{2.12}
\end{equation*}
$$

We can thus choose

$$
\begin{equation*}
b_{0}(x)=\frac{1}{a(x)}, \ldots, b_{-2 j}(x)=-\frac{1}{a(x)}\left[\sum_{k=0}^{j-1} \mathscr{A}_{j-k}\left(x, \partial_{x}\right) b_{-2 k}\right] . \tag{2.13}
\end{equation*}
$$

Theorem 2.5.1. A sufficient condition for the existence in $\tilde{S}$ of the asymptotic inverse $\tilde{b}$ is that $b_{-2 j} \in S_{m^{\prime}-2 j}$ for some $m^{\prime} \in \mathbb{R}\left(\right.$ then $\left.\tilde{b}=\sum_{j=0}^{\infty} b_{m^{\prime}-2 j} \in \tilde{S}_{m^{\prime}}\right)$. This is realized under the following condition: $a(x) \neq 0$ for all $x$, and

$$
\begin{equation*}
\forall \alpha \in \mathbb{N}^{2 l}, \quad\left|\hat{\partial}^{\alpha} a(x)\right| \leqslant c_{\alpha}|a(x)|\|x\|^{-|x|} \quad \text { for } \quad\|x\| \rightarrow \infty \tag{2.14}
\end{equation*}
$$

Remark. This condition states not only that $a(x)$ should not vanish for finite $x$, but that it should not tend to 0 too rapidly as $\|x\| \rightarrow \infty$. This condition is not necessary: Given any expansion of $\tilde{a}$ such as $\tilde{a}=\sum_{0}^{\infty} a_{m-j}\left(a_{m-j} \in S_{m-j}\right)$, the previous method can be adapted in an obvious way provided $a_{m}(x) \neq 0(\forall x)$ and $\left|\partial^{\alpha} a_{m-j}(x)\right| \leqslant c_{j \alpha}\left|a_{m}(x)\right|\|x\|^{-|\alpha|-j}\left(\forall j \in \mathbb{N}, \forall \alpha \in \mathbb{N}^{2 l}\right)$, the latter condition being satisfied, for instance, if $a_{m}$ is homogeneous of degree $m$ (then, moreover, $\left.\tilde{b} \in S_{-m}\right)$.

Proof of the theorem. Condition (2.14) with $|\alpha|=1$ implies $\left\|\vec{\nabla}_{x}(\log |a(x)|)\right\| \leqslant$ $m^{\prime} /\|x\|$ for $m^{\prime}=\max _{|\alpha|=1}\left\{c_{\alpha}\right\}$; hence $|1 / a(x)| \leqslant \mathcal{O}\left(\|x\|^{m^{\prime}}\right)$ as $\|x\| \rightarrow \infty$. A repeated use of Eqs. (2.12)-(2.13) allows us to show recursively on $j$ and $|\alpha|$ that

$$
\forall j \in N, \quad \forall \alpha \in \mathbb{N}^{2 l}, \quad\left|\partial^{\alpha} b_{-2 j}(x)\right| \leqslant c_{j \alpha}\left|b_{0}(x)\right|\|x\|^{-|\alpha|-2 j} ;
$$

hence

$$
\left|\partial^{\alpha} b_{-2 j}(x)\right| \leqslant \mathcal{O}\left(\|x\|^{m^{\prime}-|\alpha|-2 j}\right) \quad \text { and } \quad b_{-2 j} \in S_{m^{\prime}-2 j} . \quad \text { Q.E.D. }
$$

This almost settles the question of the existence of the inverse operator $\hat{a}^{-1}$. For if $b(x)$ is any GLS symbol admitting the previously found $\tilde{b}$ as an asymptotic class, we have

$$
\begin{align*}
& \hat{a} \hat{b} \quad 1+\hat{c} \quad \text { for some } \quad \hat{c} \in S_{-\infty} \\
& \rightarrow \hat{a}^{-1}=b(\mathbb{1}+\hat{c})^{-1} \tag{2.15}
\end{align*}
$$

Operator $\hat{c}$ is compact on $\mathscr{H}$; except if by accident (-1) happens to be an eigenvalue of $\hat{c},(\mathbb{1}+\hat{c})^{-1}$ exists (by the Fredholm alternative). Writing
$(\hat{1}+\hat{c})^{-1}=\hat{1}+\hat{d}: \mathscr{S} \rightarrow^{\hat{a}} \mathscr{S}$ is continuous; hence by the relation $\hat{d}=-\hat{c}-\hat{d} \hat{c}$, $\mathscr{S}^{\prime} \rightarrow{ }^{\hat{d}} \mathscr{S}$ is continuous, so that $\hat{d} \in S_{-\infty} \Rightarrow \hat{a}^{-1}=\hat{b}(1+\hat{d}) \in \hat{S}_{m^{\prime}}$ and its symbol class is $\tilde{b}$.

Remark. Procedure (2.13) can be simplified by noting that $\hat{a} \hat{b}=1 \Rightarrow \hat{b} \hat{a}=\mathbb{1}$; hence Eq. (2.11) implies

$$
\begin{equation*}
\tilde{a}\left[\exp -(i / 2)\left(\hat{\partial}_{q} \vec{\partial}_{p}-\widehat{\partial}_{p} \vec{\partial}_{q}\right)\right] \tilde{b}=1 \tag{2.16}
\end{equation*}
$$

System (2.11)-(2.16) splits as

$$
\begin{align*}
& \tilde{a}\left[\cos \left(\left(\stackrel{\partial}{\partial}_{q} \vec{\partial}_{p}-\stackrel{\partial}{\partial}_{p} \vec{\partial}_{q}\right) / 2\right)\right] \tilde{b}=1  \tag{+}\\
& \tilde{a}\left[\sin \left(\left(\stackrel{\rightharpoonup}{\partial}_{q} \stackrel{\rightharpoonup}{\partial}_{p}-\stackrel{\rightharpoonup}{\partial}_{p} \vec{\partial}_{q}\right) / 2\right)\right] \tilde{b}=0 . \tag{-}
\end{align*}
$$

Equation $\left(2.17^{+}\right)$alone determines the $b_{-2 j}$ by separate recursions on even and odd $j$ :

$$
\begin{aligned}
& b_{0}=\frac{1}{a}, \ldots, b_{-4 j}=-\frac{1}{a}\left[\sum_{k=0}^{j-1} \mathscr{A}_{2(j-k)} b_{-4 k}\right], \\
& b_{-2}=0, \ldots, b_{-2(2 j+1)}=-\frac{1}{a}\left[\sum_{k=0}^{j-1} \mathscr{A}_{2(j-k)} b_{-2(2 k+1)}\right]=0 ;
\end{aligned}
$$

hence all $b_{-2 j}$ with $j$ odd are actually zero, and Eq. (2.17-) then yields an infinite set of identities satisfied by the $b_{-4 j}$.

### 2.6. The Twisted Square Root of a Symbol

In Section 2.7, we shall seek conditions under which a GLS operator $\hat{a} \in \mathcal{S}$ is positive. For this it is natural to check whether $\hat{a}$ admits a symmetric operator square root $\hat{a}=\hat{b} \hat{b}$, with $b \in S$ real-valued. We can solve this problem $\bmod S_{-\infty}$ along the same lines as those in Section 2.5:

Theorem 2.6.1. If $a \in S_{m}$ satisfies (2.14) and $a(x)>0$ for all $x$, there exists a real-valued $\tilde{b} \in \tilde{S}_{m / 2}$ such that

$$
\tilde{a}=\tilde{b}\left[\exp (i / 2)\left(\grave{\partial}_{q} \vec{\partial}_{p}-\dot{\partial}_{p} \vec{\partial}_{q}\right)\right] \tilde{b} .
$$

Proof. Since all $\tilde{b}\left(\tilde{\partial}_{q} \vec{\partial}_{p}-\bar{\partial}_{p} \vec{\partial}_{q}\right)^{2 k+1} \tilde{b}=0$, by symmetry,

$$
\left.\tilde{a}=\tilde{b}\left[\cos \left(\dot{\partial}_{q} \vec{\partial}_{p}-\dot{\partial}_{p} \vec{\partial}_{q}\right) / 2\right)\right] \tilde{b} .
$$

We try $\tilde{b}=\sum_{j=0}^{\infty} b_{-4 j}$ with $b_{-4 j} \in S_{m-4 j}$. We can choose $b_{0}(x)=(a(x))^{1 / 2}$, and $b_{-4 j}$ recursively as the solution of the equation

$$
\sum_{j^{\prime}+j^{\prime}+k=j} b_{-4 j^{j}}\left[\frac{(-1)^{k}}{(2 k)!}\left(\frac{\stackrel{\partial}{\partial}_{p} \dot{\partial}_{q}-\stackrel{\partial}{\partial}_{p} \ddot{\partial}_{q}}{2}\right)^{2 k}\right] b_{-4 j^{*}}=0
$$

i.e.,

$$
b_{-4 j}=\frac{-1}{2 a^{1 / 2}}\left(\sum_{\substack{j^{\prime}+j^{\prime \prime}+k=j \\ j^{\prime}, j^{\prime \prime} \leqslant j}} b_{-4 j^{\prime}}\left[\frac{(-1)^{k}}{(2 k)!}\left(\frac{\tilde{\partial}_{q} \vec{\partial}_{n}-\dot{\partial}_{n} \vec{\partial}_{a}}{2}\right)^{2 k}\right] b_{-4 j^{\prime \prime}}\right) .
$$

$b_{-4 j}$ is real-valued, and condition (2.14) implies recursively on $j$ and $|\alpha|$ :

$$
\left|\partial^{\alpha} b_{-4 j}(x)\right| \leqslant c_{j \alpha}\left|b_{0}(x)\right| \cdot\|x\|^{-|\alpha|-4 j}
$$

hence

$$
\left|\partial^{\alpha} b_{-4 j}(x)\right| \leqslant \mathcal{O}\left(\|x\|^{m / 2-|x|-4 j}\right) \quad \text { and } \quad \tilde{b} \in \tilde{S}_{m / 2}
$$

If $b \in S_{m / 2}$ is any real symbol of asymptotic class $\tilde{b}$, then

$$
\begin{equation*}
\hat{a}=\hat{b} \hat{b}+\hat{c} \quad \text { for some real-valued } c \in S_{-\infty} \tag{2.18}
\end{equation*}
$$

Remark. As in Section 2.5, we can relax the hypothesis of Theorem 2.6.1: If $\tilde{a}=\sum_{0}^{\infty} a_{m-j}$ with $a_{m-j} \in S_{m-j}$, the previous method of finding $b$ can be adapted if we only require: $a_{m}(x)>0, a$ and all $a_{m-j}$ real-valued, and $\left|\partial^{\alpha} a_{m-j}(x)\right| \leqslant c_{j \alpha}\left|a_{m}(x)\right|| | x| |-|\alpha|-j(\forall j, \alpha)$.

### 2.7. Some $L^{2}$ Properties of GLS Operators

We begin by introducing a class of comparison operators: the powers of the quantized harmonic oscillator. Let $H=\hat{h}$ with $h(x)=\|x\|^{2} / 2=\frac{1}{2}\left(q^{2}+p^{2}\right)$. Then for all $n \in \mathbb{Z}:(1+H)^{n} \in S_{2 n}$ and its symbol behaves like $\left(1+\|x\|^{2} / 2\right)^{n} \times$ $(1+o(1))$ as $\|x\| \rightarrow \infty$. The proof is tedious but simple, using Theorems 2.4.1 (for $n>0$ ) and 2.5.1 (for $n<0$ ).

We now consider $A \in S_{m}$ as a (possibly unbounded) operator on $\mathscr{H}=L^{2}(Q)$.
Theorem 2.7.1. (i) If $m<0, A$ is compact.
(ii) If $m<-l$, $A$ is of Hilbert-Schmidt class.
(iii) If $m<-2 l, A$ is of trace class.

Pronf. (ii) see Theorem 2.3.1(ii). (i) If $m<0$, some power $A^{k} \in \hat{S}_{k m c}$ (with $k \in \mathbb{N}, k>l /|m|$ ) is Hilbert-Schmidt by (ii); hence $A$ is compact. (iii) If $m<-2 l, A$ is the product $B\left(B^{-1} A\right)$ with $B=\widehat{(1+h)^{m / 4}} \in S_{m / 2}$ and $B^{-1} \in$ $\hat{S}_{-m / 2}$, hence $B^{-1} A \in S_{m / 2}$; both $B$ and $B^{-1} A$ are Hilbert-Schmidt, so that $A$ is of trace class.

Theorem 2.7.2. (i) If $a \in S_{m}$ is a semibounded function: $a(x) \geqslant M>-\infty$ and satisfies (2.14), then $\hat{a}$ is a semibounded symmetric operator for which the Friedrichs (self-adjoint) extension exists.
(ii) The operators in $\hat{S}_{0}$ are bounded.
(iii) For all $m \in \mathbb{R}$ and $n \in \mathbb{N}$ with $n \geqslant m / 2$ : If $A \in S_{m}, A(1+H)^{-n}$ is a bounded operator; hence $A$ is defined on the dense domain of $H^{n}$ in $\mathscr{H}$.

Proof. (i) We let $M>0$ (the general case follows trivially). Then, hypothesis (2.14) is automatically satisfied if $a \in S_{0}$ (this is needed for proving (ii)). Using Theorem 2.6.1, we can write Eq. (2.18): $\hat{a}=\hat{b}^{2}+\hat{c}$ in the sense of continuous operators $\mathscr{P} \rightarrow \mathscr{P}$, with $b$ and $c$ real-valued, $c \in S_{-\infty}$. Then $\hat{a}$ is a densely defined symmetric operator on $\mathscr{H}$, whose associated quadratic form is semibounded ( $\hat{b}^{2}$ defines a positive form and $c$ is a symmetric bounded operator by Theorem 2.7.1(i)). The Friedrichs extension theorem [21] then holds and defines a semibounded self-adjoint extension $A$ of $\hat{a}: A \geqslant-\|\hat{c}\| 1$.

Note that $a(x)>0$ does not imply that $A$ is a positive operator; as a onedimensional counterexample, let $a(x)=p^{2} q^{2}+\epsilon$ with $0<\epsilon<\frac{1}{4}$; then $\hat{a}=(\widehat{p q})^{2}+\left(\epsilon-\frac{1}{4}\right) \mathbb{1}$ admits $\left(\epsilon-\frac{1}{4},+\infty\right)$ as spectrum and is not a positive operator. Our proof provides no estimate for the lower bound of $\hat{a}$ in the general case. The result is mainly useful as a criterion of self-adjointness. (Note: Berezin and Shubin [31] prove that if $a(x)$ is real-valued and if $a(x) \pm i$ satisfies (2.14), then $\hat{a}$ has deficiency indices ( 0,0 ); hence its closure is self-adjoint.)
(ii) If $a \in S_{0}$, then $\operatorname{Re} a(x)$ is bounded above and below; (i) then implies that $\widehat{\operatorname{Re} a}$ is a bounded operator, and similarly for $\widehat{\operatorname{Im} a}$; hence $\hat{a}$ is bounded.
(iii) If $a \in S_{m}$ and $n \geqslant m / 2, a(1+H)^{-n} \in S_{0}$ and is a bounded operator by (ii).

We can somewhat refine the result of (iii). For each $n \in \mathbb{Z}$, let $\mathscr{H}_{n}$ be the Hilbert space obtained by completion of the domain of $(1+H)^{n}$ in $\mathscr{H}$ for the inner product $(u, v) \rightarrow\langle u, v\rangle_{n}=\left\langle(1+H)^{n} u,(1+H)^{n} v\right\rangle$. We obtain a sequence of spaces

$$
\cdots \subset \mathscr{H}_{+k} \subset \cdots \subset \mathscr{H}_{+1} \subset \mathscr{H}_{0}=\mathscr{H} \subset \mathscr{H}_{-1} \subset \cdots \subset \mathscr{H}_{-k} \subset \cdots
$$

( $k>0$ ), where $\mathscr{H}_{k}$ and $\mathscr{H}_{-k}$ are dual of each other for the inner product of $\mathscr{H}$. Moreover, $\mathscr{S}(Q)=\bigcap_{k \in \mathbb{Z}} \mathscr{H}_{k}$ and its topology is given by the directed family of seminorms $\left\|\|_{n} ;\right.$ hence $\mathscr{S}^{\prime}(Q)=\bigcup_{k \in \mathbb{Z}} \mathscr{H}_{k}[14,22]$.

Proposition 2.7.3. Let $A \in S_{m}$. For every $n \in \mathbb{Z}$ with $n \geqslant m / 2$, and for all $k \in \mathbb{Z}, A$ is a continuous map $\mathscr{H}_{k} \rightarrow \mathscr{H}_{k-n}$.

Proof. It suffices to show that $A$ is bounded on the domain of $(1+H)^{k}$ in $\mathscr{H}$, which is dense in $\mathscr{H}_{k}$ for the topology of $\mathscr{H}_{k}$. For any vector $u$ in the domain, $A u=(1+H)^{n-k}\left[(1+H)^{k-n} A(1+H)^{-k}\right](1+H)^{k} u$. But $(1+H)^{k}$ is an isometry of $\mathscr{H}_{k}$ into $\mathscr{H},\left[(1+H)^{k-n} A(1+H)^{-k}\right] \in S_{0}$; hence it is a bounded operator on $\mathscr{H}$, and $(1+H)^{n-k}$ is an isometry of $\mathscr{H}$ into $\mathscr{H}_{k-n}$.

Letting $k \rightarrow+\infty$ (resp. $-\infty$ ) we would recover the continuity of $A$ from $\mathscr{S}$ to $\mathscr{S}$ (resp. $\mathscr{S}^{\prime}$ to $\mathscr{S}^{\prime}$ ).

## 3. Asymptotic Quantum Mechanics and Admissible Operators

We refer to "asymptotic quantum mechanics" when dealing with power series expansions of quantum mechanical quantities in Planck's constant $\hbar$. The mathematical structure of quantum mechanics [23] depends explicitly on $\hbar \in \mathbb{R}+\backslash\{0\}$, which can be considered as a variable parameter of the theory. Although $\hbar$ is a constant in nature, this viewpoint can be useful because certain quantum mechanical quantities tend as $\hbar \rightarrow 0^{+}$to their "analogs" of classical Hamiltonian mechanics [24], and in many physical instances the actual value of $\hbar$ is relatively so small as to make a perturbative computation "around" classical mechanics a sensible operation. The works of Maslov [25], Leray [26], and Duistermaat [27] have brought out the analogy between these small- $\hbar$ expansions and the ordinary asymptotic expansions in pseudodifferential operator theory. A specific pseudodifferential algebra can be introduced [4,5] to deal directly with $\hbar$-expansions. It is analogous to the GLS algebra, but with additional requirements on the explicit behavior in $\hbar$ of the symbols. We thus need to rewrite Sections 1 and 2 showing the explicit dependence in $\hbar$ everywhere $(\hbar>0)$.

### 3.1. The Weyl Quantization

It associates to a classical function $a(q, p)$ the integral operator on $\mathscr{H}=L^{2}(Q)$ of kernel (cf. Eq. (1.2)):

$$
\begin{equation*}
\langle q| \hat{a}_{\hbar}\left|q^{\prime}\right\rangle=\left(\mathscr{F}_{\hbar} a\right)\left(q, q^{\prime}\right)=(2 \pi \hbar)^{-t} \int_{P} a\left(\left(q+q^{\prime}\right) / 2, p\right) e^{i p\left(q-q^{\prime}\right) / \hbar} d p \tag{3.1}
\end{equation*}
$$

The conditions of validity and the properties of this Weyl quantization are the same as those in Section 1 (up to some normalization constants); in particular, $\mathscr{F}_{\hbar}$ is a continuous isomorphism $\mathscr{S}^{\prime}(X) \rightarrow \mathscr{S}^{\prime}(Q \times Q)$ and also an isometry $L^{2}\left(X, d x /(2 \pi \hbar)^{l}\right) \rightarrow L^{2}(Q \times Q) \approx \mathscr{L}^{2}(\mathscr{H})$. The coordinate functions become quantized as

$$
\begin{aligned}
& \hat{q}_{j}=q_{j} \quad\left(\text { multiplication by } q_{j}\right), \\
& \vec{\partial}_{j}=-i \hbar \frac{\partial}{\partial q_{j}}
\end{aligned}
$$

hence they generate the well-known Schrodinger representation of the canonical commutation rules: $\left[\hat{q}_{j}, \hat{p}_{k}\right]_{-}=i \hbar \delta_{j k} \downarrow\left(=i \hbar \delta_{j k} \hat{l}\right)$. The other properties listed in Section 1 stay unchanged.

### 3.2. Admissible Symbols

An admissible symbol of order $\leqslant m(m \in \mathbb{R})$ is any $C^{\infty}$ function $a \in C^{\infty}$ $\left(\left[0, \hbar_{0}\right) \times X\right)$ (for some unspecified $0<\hbar_{0} \leqslant \infty$ ) such that for all $j \in \mathbb{N}$ the mapping $\hbar \rightarrow\left(\partial^{j} / \partial \hbar^{j}\right) a(\hbar, \cdot)$ is $C^{\infty}$ from [0, $\hbar_{0}$ ) into $S_{m-2 j}$ (note that regularity in $\hbar$ is required up to $\hbar=0$ ).

The space of such symbols is denoted ${\pi_{m}}_{m}$; we put on it the topology of uniform convergence on compacts of $\left[0, \hbar_{0}\right.$ ) of every derivative $\partial^{j} a / \partial \hbar^{j}$ in $S_{m-2 j}$. We have $O \eta_{m} \subset C l_{m}$ if $m^{\prime}<m$; we let $O \eta_{-\infty}=\bigcap_{m} C l_{m}$ and $C l=\bigcup_{m} C \eta_{m}$.

The Taylor series at $\hbar=0$ of a symbol $a \in O Z_{m}$,

$$
\begin{equation*}
a(\hbar, \cdot) \sim \sum_{0}^{\infty} a_{j} \hbar^{j}, \quad a_{j} \in S_{m-2 j} \tag{3.2}
\end{equation*}
$$

defines an asymptotic expansion both in the sense of Section 2.2 and in the sense of functions of $\hbar$ as $\hbar \rightarrow 0^{+}$; for all $n \in \mathbb{N}$,

$$
\left(a(\hbar, \cdot)-\sum_{k=0}^{n-1} a_{k} \hbar^{k}\right) / \hbar^{n}
$$

stays in a bounded set of $S_{m-2 n}$ as $\hbar \rightarrow 0^{+}$.
The leading term $a_{0}(x)$ is the principal symbol.
The asymptotic study for $\hbar \rightarrow 0$ will identify as zero any symbol admitting the expansion $\sum_{0}^{\infty} 0 \cdot \hbar^{n}$ in the sense of (3.2). We are thus led to define:
-The subspace $\mathcal{N}$ of negligible symbols as the space of symbols $a \in C$ such that $a \sim 0$ in the sense of (3.2); i.e., $a(\hbar, \cdot) \in S_{-\infty}$ and for all $n \in \mathbb{N}: a(\hbar, \cdot) / \hbar^{n}$ stays in a bounded set of $S_{-\infty}$ as $\hbar \rightarrow 0$.
-The quotient spaces $\widetilde{\mathscr{a}}_{m}=\mathscr{a}_{m} / \mathcal{N}, \widetilde{\mathscr{a}}_{-\infty}=\mathscr{a}_{-\infty} / \mathfrak{N}, \widetilde{\mathscr{G}}=O / \mathcal{N}$ of asymptotic classes of symbols modulo $\mathscr{N}$; the class of $a$ is noted $\tilde{a}$. Note that here we let the space $\mathscr{N}$ (and not ${\theta_{-\infty}}^{\text {}}$ ) play the same role as the space $S_{-\infty}$ in the GLS case.

Theorem 3.2.1. Every formal series $\sum_{0}^{\infty} a_{j} \hbar^{j}$ with $a_{j} \in S_{m-2 j}$ defines an element of $\widetilde{\mathscr{T}}_{m}$ uniquely.

The proof can be done similarly to the proof of the analogous theorem, Theorem 2.2.1.

### 3.3. Quantization of Admissible Symbols

For every $0<\hbar<\hbar_{0}, \widehat{a(\hbar, \cdot)_{\hbar}}$ is an operator in $\hat{S}_{m}$ with all the properties listed in Section 2.3; we call it an admissible operator.

Many important operators of quantum mechanics are indeed admissible operators: for instance ( $l=3$ ), the components of the angular momentum $\hat{\mathbf{q}} \times \hat{\mathbf{p}}$, the squared angular momentum $(\hat{\mathbf{q}} \times \hat{\mathbf{p}})^{2}$ (it has the symbol $\left.(\mathbf{q} \times \mathbf{p})^{2}-3 \hbar^{2} / 4\right)$, and all Hamiltonians $\hat{p}^{2} / 2+\widehat{V(q)}$ with polymonial interaction potentials. Nonpolynomial potentials $\left(V \in C^{\infty}(Q)\right.$ are never admissible $(V \notin S)$ [33]; ${ }^{1}$ but this restriction disappears if we replace condition (2.1) by: $\exists m, n \in \mathbb{R} \forall \alpha, \beta \in \mathbb{N}^{l},\left|\partial_{q}{ }^{\alpha} \partial_{p}{ }^{\beta} a(x)\right| \leqslant C_{\alpha \beta}(1+\|q\|)^{m-|\alpha|}(1+\|p\|)^{n-|\beta|}$
${ }^{1}$ We thank J. Chazarain for calling our attention to this fact.

### 3.4. The Twisted Multiplication of Admissible Symbols

As in Section 2.4, it is defined so as to correspond to the operator multiplication under the now $\hbar$-dependent Weyl quantization.

Theorem 3.4.1. (i) $O$ is an associative *-algebra for twisted multiplication, which is a continuous map from $O \tau_{m_{1}} \times O{I_{m_{2}}}$ into $O l_{m_{1}+m_{2}}\left(\right.$ for all $\left.m_{1}, m_{2} \geqslant-\infty\right)$ : $(a, b) \rightarrow c$, given by

$$
\begin{equation*}
c(\hbar, x)=(\pi \hbar)^{-2 l} \int_{X^{2}} d x_{1} d x_{2} a\left(\hbar, x+x_{1}\right) b\left(\hbar, x+x_{2}\right) e^{2 i \omega\left(x_{1}, x_{2}\right) / \hbar} \tag{3.3}
\end{equation*}
$$

(ii) c admits the asymptotic expansion in 0t:

$$
\begin{equation*}
c(\hbar, x) \sim \tilde{a}(\hbar, x)\left[\exp (i \hbar / 2)\left(\dot{\partial}_{q_{1}} \vec{\partial}_{p_{2}}-\vec{\partial}_{p_{1}} \vec{\partial}_{a_{2}}\right)\right] \tilde{b}(\hbar, x) \tag{3.4}
\end{equation*}
$$

(Its principal symbol is $c_{0}(x)=a_{0}(x) b_{0}(x)$. )
(iii) If $a$ or $b$ is polynomial in $x$, expansion (2.7) is a finite sum which equals $c(x)$ exactly.

The proof is identical to that of Theorem 2.4.1. All other statements in Section 2.4 have easy to find analogs; in particular, $\widetilde{\mathscr{\ell}}$ is a *-algebra for the "twisted multiplication modulo $\mathscr{N}$ " defined in Eq. (3.4).

### 3.5. The "Twisted Inverse" of an Admissible Symbol

Here we can make a stronger statement than the analog of Theorem 2.5.1:
Theorem 3.5.1. Let $a \in C t, a \sim \sum_{0}^{\infty} a_{j} \hbar^{i}$ in the sense of Section 3.2, with $a_{0}(x) \neq 0$ for all $x$ and $a_{0}(x)$ satisfying, as $\|x\| \rightarrow \infty$,

$$
\begin{equation*}
\forall \alpha \in \mathbb{N}^{2 l}, \quad\left|\partial^{\alpha} a_{j}(x)\right| \leqslant c_{j_{\alpha}}\left|a_{0}(x)\right|\|x\|^{-|\alpha|-2 j} . \tag{3.5}
\end{equation*}
$$

Then for all $\hbar$ small enough, $\hat{a}_{\hbar}$ admits an operator inverse of the form $\ddot{b}_{\hbar}$, with $b \in C l$ of principal symbol $1 / a_{0}(x)$.

Proof. We begin, as in Theorem 2.5.1, by computing the asymptotic class $\tilde{b}=\sum_{0}^{\infty} b_{j} \hbar^{j}$ using the parametrix method; in particular, $b_{0}(x)=1 / a_{0}(x)$. Condition (3.5) is needed to keep $a_{0}(x)$ from tending to zero too fast as $\|x\| \rightarrow \infty$ and it guarantees that for some $m^{\prime} \in \mathbb{R}$ and all $j \in \mathbb{N}: b_{j} \in S_{m^{\prime}-2 j}$. For any $b \in C t$ having this asymptotic class $\tilde{b}$, we then have

$$
\hat{a}_{\hbar} \hat{b}_{\hbar}=1+\hat{c}_{\hbar} \quad \text { for some } \quad c \in \mathcal{N},
$$

hence:

$$
\begin{equation*}
\hat{a}_{n}^{-1}=\hat{b}_{n}\left(\mathbb{1}+\hat{c}_{n}\right)^{-1} . \tag{3.6}
\end{equation*}
$$

The new fact here is that not only is $\hat{c}_{\boldsymbol{n}}\left(\in \hat{S}_{-\infty}\right)$ a Hilbert-Schmidt operator for all $\hbar$, but its operator norm is a rapidly decreasing function of $\hbar$ as $\hbar \rightarrow 0$, since for all $n \in \mathbb{N}$,

$$
\begin{aligned}
\left\|\hat{c}_{\hbar}\right\| & \leqslant \operatorname{Tr} \hat{c}_{\hbar}^{\dagger} \hat{c}_{\hbar}=(2 \pi \hbar)^{-l}\|c(\hbar, \cdot)\|_{L^{2}(X)} \\
& \leqslant(2 \pi \hbar)^{-l} C_{n}\left\|\frac{\hbar^{n}}{\left(1+\|x\|^{2}\right)^{n}}\right\|_{L^{2}}=\mathcal{O}\left(\hbar^{n-l}\right) .
\end{aligned}
$$

This guarantees that for $\hbar$ smaller than some (undetermined) upper bound, the operator inverse ( $\left.1 \mid \hat{c}_{\hbar}\right)^{-1}$ exists (and can be computed by its Neumann series or by the Fredholm method); also, $\left(\mathbb{1}+\hat{c}_{n}\right)^{-1}=\mathbb{1}+\hat{\vec{d}}_{\hbar}$ with $d \in \mathscr{N}$; hence $\hat{a}^{-1}$ given by (3.6) is an admissible operator with $\tilde{b}$ as asymptotic class.
Q.E.D.

Remark. The reader can check that condition (3.5) is preserved under twisted multiplication and twisted inversion. (The role of this condition is suggested by its implication that, in the sense of operator theory, $\left.\widehat{\left(a(\hbar, \cdot)_{\hbar}\right.}-\widehat{a_{0 \hbar}}\right)$ is a regular perturbation of $\widehat{a_{0 \hbar}}$, namely, that $\left.\widehat{\left(a(\hbar, \cdot)_{\hbar}\right.}-\widehat{a_{0 \hbar}}\right)\left(\widehat{a_{\mathrm{E} \hbar}}\right)^{-1} \in \hat{S}_{\mathrm{u}}$ and its operator norm is $\mathcal{O}(h))$.

### 3.6. Some L $L^{2}$ Properties of Admissible Operators

We consider for every $\hbar \in\left(0, \hbar_{0}\right)$ fixed, and for $a \in C_{m}$, the operator $\hat{a}_{\hbar}$ on $\mathscr{H}=L^{2}(Q)$. Since $a(\hbar, \cdot) \in S_{m}$, Theorems 2.7 .1 and 2.7 .2 hold without modification, but we can also strengthen point (i) of Theorem 2.7.2.

Theorem 3.6.1. Let $a \in \mathbb{A}, a \sim \sum_{0}^{\infty} a_{j} h^{j}$ in the sense of Section 3.2, with $a_{0}(x)>0, a(\hbar, x)$ and $a_{j}(x) \in \mathbb{R}$ for all $x \in X, j \in \mathbb{N}$. If a satisfies (3.5), then for all $\hbar$ small enough, $\hat{a}_{\hbar}$ is a positive operator.

Proof. We begin by assuming that $a_{0}(x) \geqslant k>0$. 'I'hen $(a-k / 2) \in \pi$ also satisfies (3.5) (use $\left|a_{0}(x) \leqslant 2\right| a_{0}(x)-k / 2 \mid$ ) and has a positive principal symbol. As in Theorem 2.6.1, we can compute a real-valued symbol $b \sim \sum_{0}^{\infty} b_{j} h^{j}$ with $b_{j} \in S_{m / 2-2 j}$ (if $a \in\left(Z_{m n}\right)$ such that $\widehat{(a-k / 2)_{\hbar}}=\widehat{b_{\hbar} b_{\hbar}}+\hat{c}_{\hbar}$ with $c \in \mathscr{N}$ real-valued. Then $\hat{a}_{\hbar}=\hat{b}_{\hbar}^{2}+(k / 2) \downarrow+\hat{c}_{\hbar}$, with $\hat{b}_{n}$ symmetric and $\hat{c}_{n}$ selfadjoint with operator norm $-o\left(\hbar^{\infty}\right)$ : for $\hbar$ small enough this norm is $<k / 2$, and then $\hat{a}_{\hbar}$ is a positive operator.

If $\inf _{x \in X}\left\{a_{0}(x)\right\}=0$, we know by condition (3.5) that $a_{0}(x) \geqslant K\|x\|^{-m^{\prime}}$ still holds for some $m^{\prime} \in R$. Fix $n \in \mathbb{N}, n \geqslant m^{\prime} / 4$; let $h(x)=\|x\|^{2}$ and define the operator $A^{\prime}=\left(1+\hat{h}_{n}\right)^{n} \hat{a}_{n}\left(1+\hat{h}_{\hbar}\right)^{n}$. By Theorem 3.4.1, $A^{\prime}=\widehat{a_{\hbar^{\prime}}}$ with $a^{\prime} \in \mathcal{O}, a^{\prime} \sim \sum_{0}^{\infty} a_{j}^{\prime} \hbar^{j}$, and $a_{0}^{\prime}=\left(1+\|x\|^{2}\right)^{2 n} a_{0}(x)$. We can check that $a^{\prime}$ satisfies the hypothesis of Theorem 3.6.1 (using the fact that (3.5) is satisfied
by $a$ and by $(1+h)$, and is preserved under twisted multiplication), and also $a_{0}^{\prime}(x) \geqslant k>0$. By the previous argument $\widehat{a_{h^{\prime}}}$ is a positive operator for small $\hbar$, but then so is $\hat{a}_{n}=\left(1+\hat{k}_{n}\right)^{-n} \hat{a}_{n}^{\prime}\left(1+\hat{h}_{n}\right)^{-n}$.

## 4. Admissible Functionals and Physical Applications

Assume that we are interested in the spectrum and in the eigenfunctions of a given self-adjoint admissible operator $A=\hat{a}_{\tilde{n}}$. It is convenient, in order to have the discrete and continuous spectra on the same footing, to work in the "Dirac formalism," i.e., with Gelfand triplets [28]. Since the algebra of admissible operators admits $\mathscr{S}(Q)$ as a common dense invariant domain such that $\mathscr{S} \rightarrow{ }^{A} \mathscr{S}$ is continuous (Theorem 2.3.1), we can use the Gelfand triplet $\mathscr{S}(Q) \subset L^{2}(Q) \subset$ $\mathscr{S}^{\prime}(Q)$ : The "eigenfunctions" of the continuous spectrum are tempered distributions (those of the discrete spectrum lie in $L^{2}(Q)$ ).

Is it then reasonable to devise a regular perturbation scheme to compute the eigenfunctions of $A$ in power series of $\hbar$ ? Even if we ignore the difficulties linked with the discrete spectrum (which becomes continuous in the classical limit $\hbar=0$ ), the answer is certainly no; even in the simplest case (cf. the WKB method), the eigenfunctions exhibit very bad caustic singularities. An alternative approach eliminates this difficulty: We associate to every $u \in \mathscr{S}^{\prime}(Q)$ the linear functional on $a_{-\infty}$ (for fixed $\hbar>0$ ),

$$
c \in O Z_{-\infty} \rightarrow\left\langle u(q) \otimes u^{*}\left(q^{\prime}\right),\left(\mathscr{F}_{n} c\right)\left(q, q^{\prime}\right)\right\rangle,
$$

in the sense of the inner product between $\mathscr{S}(Q \times Q)$ and $\mathscr{S}^{\prime}(Q \times Q)$. Let $\rho_{u}(\hbar)=(2 \pi \hbar)^{-l} \mathscr{F}_{\hbar}^{-1}\left(u \otimes u^{*}\right) \in \mathscr{S}^{\prime}(X)$. By the Parseval identity for distributions,

$$
\begin{equation*}
\left\langle u \otimes u^{*},\left(\mathscr{F}_{\hbar} c\right)\right\rangle_{\mathscr{S}^{\prime}(O \times O) \times \mathscr{P}(O \times O)}=\left\langle\rho_{u}(\hbar), c(\hbar, \cdot)\right\rangle_{\mathscr{S}^{\prime}(X) \times \mathscr{S}(X)} \tag{4.1}
\end{equation*}
$$

( $\rho_{u}$ is the "Wigner function" of the distribution $u$ [9]).
In many solvable cases, $\left\langle\rho_{u}(\hbar), c(\hbar, \cdot)\right\rangle$ (of Eq. (4.1)) is a continuous linear functional of $c(\hbar, \cdot) \in S_{-\infty}($ for $\hbar \neq 0)$ whose value is a $C^{\infty}$ function of $\hbar$ (up to $\hbar=0$ ). Since $O_{-\infty}$ consists of the $C^{\infty}$ maps $\hbar \rightarrow c(\hbar, \cdot) \in \mathscr{S}(X)$, an equivalent statement is that $\rho_{u}$ must be a $C^{\infty}$ map: $\hbar \in\left[0, \hbar_{0}\right) \rightarrow \rho_{u}(\hbar) \in \mathscr{S}^{\prime}(X)$; we call all such $C^{\infty}$ maps $\hbar \in\left[0, h_{0}\right) \rightarrow \rho(\hbar) \in \mathscr{S}^{\prime}(X)$ admissible functionals (acting on $O_{-\infty}$ ); they furm the space $\widehat{\mathscr{S}}^{\prime}(X)$.

Twisted multiplication extends to a mapping $C l \times \widehat{\mathscr{S}}^{\prime} \rightarrow \widehat{\mathscr{S}^{\prime}}$ by duality [8]: If $a \in \mathscr{O}, b \in \widehat{\mathscr{S}}^{\prime}$, the twisted product of $a$ and $b$ is the functional $d \in \mathscr{C l}_{-\infty} \rightarrow$ $\langle b, e\rangle_{\mathscr{L}^{\prime}(X) \times \mathscr{S}(X)}$, where $e \in O Z_{-\infty}$ is the twisted product of $a^{*}$ by $d: \hat{e}_{n}=\hat{a}_{n}{ }^{\prime} \hat{d}_{\hbar}$. Equations (3.3) and (3.4) remain true when $a \in O t, b \in \widehat{\mathscr{S}}^{\prime}$ [5].

A reasonable perturbative scheme is then to compute the "coefficients" $\rho_{n} \in \mathscr{S}^{\prime}(X)$ such that $\rho_{u}(\hbar) \sim \sum_{0}^{\infty} \rho_{n} \hbar^{n}$ (in the weak sense) as $\hbar \rightarrow 0^{+}$.

Unfortunately we have no theorem stating that for a given admissible operator, suitable eigenfunctions $u \in \mathscr{P}^{\prime}(Q)$ define functionals $\rho_{u}(\hbar)$ that are admissible (we know this to be true only for some exactly solvable cases; see examples below). What we know is that in many other cases the eigenvalue equations for $\rho_{u}, A \rho_{u}=\rho_{u} A=E \rho_{u}$, can be formally solved to any order $\hbar^{n}$, yielding very reasonable coefficients $\rho_{n} \in \mathscr{S}^{\prime}(X)$; no singularities are ever encountered. Moreover, in completely integrable situations the equation $\rho_{u}(\hbar) \sim \sum_{0}^{\infty} \rho_{n} \hbar^{n}$ can then be solved algebraically (again formally, to any order) to yield an eigenfunction $u$ of the Maslov (WKB) type $u(q)=\alpha(\hbar, q) e^{i S(q) / \hbar}$ (away from caustics); in the case of discrete eigenvalues corresponding to this situation, they, too, can be computed.

This whole scheme uses properties of admissible functionals which follow by duality from those of admissible operators; we refer the reader to [5] for details.

Here are two explicit examples of admissible functionals in one dimension:
-if $u$ is the ground state of the harmonic oscillator

$$
\frac{\|x\|_{\hbar}^{2}}{2}=-\frac{\hbar^{2}}{2} \frac{d^{2}}{d q^{2}}+\frac{q^{2}}{2} \quad \text { (of eigenvalue } \frac{\hbar}{2} \text { ) }
$$

then

$$
\rho_{u}(\hbar)=\frac{1}{\pi \hbar} e^{-\left(q^{2}+p^{2}\right) / \hbar} \sim \sum_{0}^{\infty}\left[\frac{1}{n!}\left(\frac{\partial^{2} / \partial q^{2}+\partial^{2} / \partial p^{2}}{4}\right)^{n} \delta(q) \delta(p)\right] \hbar^{n}
$$

-If $u$ is the eigenfunction of the "Airy operator"

$$
\widehat{\left(\frac{p^{2}}{2}-q\right)_{\hbar}}=-\frac{\hbar^{2}}{2} \frac{d^{2}}{d q^{2}}-q
$$

with "eigenvalue" 0 (in the continuous spectrum),

$$
\rho_{u}(\hbar)=\frac{1}{2 \pi \hbar} \int_{-\infty}^{+\infty} d \tau \cdot e^{-i r^{3} / 24 \hbar} e^{i \tau\left(q-p^{2} / 2\right) \hbar} \sim\left[\exp \left(\frac{\hbar^{2}}{24} \frac{\partial^{3}}{\partial q^{3}}\right)\right] \delta\left(q-\frac{p^{2}}{2}\right)
$$

whereas the function $u$ itself $\left(u(q, \hbar)=2^{+1 / 3}(2 \pi)^{1 / 2} \hbar^{-1 / 6} A i\left(-2^{+1 / 3} \hbar^{-2 / 3} q\right)\right.$ ) has a caustic singularity at $q=0$ and cannot be regularly expanded in powers of $h$.

In conclusion, although the mathematical theory is not yet complete, the theory of admissible operators and functionals seems to lead to methods for solving pseudodifferential equations asymptotically that are more regular, and probably more general, than the usual methods.

## APPENDIX: Other Symbol Maps

Different correspondences between symbols and operators can be viewed as different prescriptions for ordering the noncommuting operators $q_{j}$ and $-i\left(\partial / \partial q_{j}\right)$ in a formal expression given as $a(q,-i \partial / \partial q)$.

The usual ordering in mathematics associates to a function $a(x)$ an operator $a^{M}$ such that [2,3]

$$
\begin{equation*}
\left(a^{M} u\right)(q)=(2 \pi)^{-l} \int_{P \times Q} a(q, p) u\left(q^{\prime}\right) e^{i p\left(q-q^{\prime}\right)} d p d q^{\prime} \tag{A1}
\end{equation*}
$$

If $a(x)=b(q) c(p)$, then $a^{M}=b(q) c(-i \partial / \partial q)$; this ordering puts all differentiations to the right. In particular, as compared with Eq. (1.4), $[\exp i(\alpha q+\beta p)]^{M}=$ $\exp \left(i \alpha q^{M}\right) \cdot \exp \left(i \beta p^{M}\right)$. Note that for any $a(x)=f(q)+g(p), a^{M} \equiv \hat{a}$, in particular $q_{j}{ }^{M}-q_{j}$ and $p_{j}{ }^{M}=-i \partial / \partial q_{j}\left(\right.$ on $C_{0}{ }^{\infty}(Q)$ ).

The adjoint $\left(a^{M}\right)^{\dagger}$ (in the sense of forms) arises from the symbol

$$
\begin{equation*}
b(q, p)=(2 \pi)^{-l} \int_{P \times Q} a^{*}\left(q+r, p^{\prime}\right) e^{i\left(p-p^{\prime}\right) r} d p^{\prime} d r \tag{A2}
\end{equation*}
$$

In the sense of oscillatory integrals, the mapping $a^{*}(x) \rightarrow b(x)$ defined by (A2) is a continuous isomorphism of $S_{m}$ (same style of proof as that for Theorem 2.4.1) Equation (A2) admits the formal asymptotic expansion

$$
b(q, p) \sim a^{*}\left(q-i \frac{\partial}{\partial p}, p\right)=\sum_{\alpha \in \mathbb{N}^{2}} \frac{1}{\alpha!}(-i)^{|\alpha|} \partial_{p}^{\alpha} \partial_{q}^{\alpha} a^{*}(q, p)
$$

'I'he mapping $a \rightarrow a^{M}$ is $1-1$ from $S_{m}$ onto $S_{m}$; for every $a \in S_{m}, a^{M}=\widehat{a^{\prime}}$ with the symbol $a^{\prime} \in S_{m}$ defined by

$$
\begin{equation*}
a^{\prime}(q, p)=(2 \pi)^{-l} \int_{P \times Q} a\left(q-\frac{r}{2}, p^{\prime}\right) e^{i\left(p-p^{\prime}\right) r} d p^{\prime} d r \tag{A3}
\end{equation*}
$$

and (A3) defines a continuous isomorphism $S_{m} \rightarrow S_{m}$; the inverse formula of (A3) is

$$
\begin{equation*}
a(q, p)=(2 \pi)^{-l} \int_{P \times Q} a^{\prime}\left(q+\frac{r}{2}, p^{\prime}\right) e^{i\left(p-p^{\prime}\right) r} d p^{\prime} d r \tag{A4}
\end{equation*}
$$

The corresponding asymptotic formulas are

$$
\begin{align*}
& a^{\prime}(q, p) \sim a\left(q+\frac{i}{2} \frac{\partial}{\partial p}, p\right)=\sum_{\alpha \in \mathbb{N}^{2}} \frac{1}{\alpha!}\left(\frac{i}{2}\right)^{|\alpha|} \partial_{\alpha}^{\alpha} \partial_{p}^{\alpha} a(q, p), \\
& a(q, p) \sim a^{\prime}\left(q-\frac{i}{2} \frac{\partial}{\partial p}, p\right)=\sum_{\alpha \in \mathbb{N}^{2}} \frac{1}{\alpha!}\left(\frac{-i}{2}\right)^{|\alpha|} \partial_{q}^{\alpha} \partial_{p}^{\alpha} a^{\prime}(q, p) . \tag{A5}
\end{align*}
$$

Thus the Weyl ordering and this ordering lead to the same operator spaces.

But the operator relation $c^{M}=a^{M} b^{M}$ defines a symbol multiplication $(a, b) \rightarrow c$ which is different from twisted multiplication (Eqs. (2.6 and 2.7)):

$$
\begin{align*}
& c(q, p)=(2 \pi)^{-l} \int_{X} a\left(q, p+p^{\prime}\right) b\left(q+q^{\prime}, p\right) e^{-i p^{\prime} q^{\prime}} d p^{\prime} d q^{\prime}  \tag{A6}\\
& \sim a\left(q, p-i \partial_{q^{\prime}}\right) b\left(q+q^{\prime}, p\right)_{a^{\prime}=0}=\sum_{x \in \mathbb{N}^{l}} \frac{1}{\alpha!}(-i)^{|\alpha|} \partial_{p}^{\alpha} a(q, p) \partial_{q}^{\alpha} b(q, p) \\
&=a\left[\exp \left(-i \hat{\partial}_{p} \vec{\partial}_{q}\right)\right] b . \tag{A7}
\end{align*}
$$

An inverse ordering (putting all differentiations to the left) can similarly be defined. The two orderings clearly do not share the Weyl ordering's invariance property with respect to linear symplectic automorphisms of $X$.

A Wick (or "normal") ordering can be associated to every Euclidean structure on $X$. By a symplectic automorphism, any positive definite quadratic form on $X$ can be put in the "normal" form $h(x)=\frac{1}{2} \sum_{1}^{l} \omega_{j}\left(q_{j}{ }^{2}+p_{j}{ }^{2}\right)$ (all $\omega_{j}>0$ ). We introduce the isotropic coordinates $\bar{z}_{j}=\left(q_{j}+i p_{j}\right) / 2^{1 / 2}, z_{j}=\left(q_{j}-i p_{j}\right) / 2^{1 / 2}$; the quantized operators $\hat{z}_{j}$ (resp. $\tilde{z}_{j}$ ) are the creation (resp. annihilation) operators of the harmonic oscillator $\widehat{h(x)}$. The Wick quantization of a classical function expressed as $a(z, \bar{z})$ is an operator : $a$ : such that if $a(z, \bar{z})-\sum b_{j}(z) c_{f}(\bar{z})$ then $: a:=\sum b_{j}(\hat{z}) c_{j}(\hat{z})$. This prescription determines :a: for all polynomials $a(z, \bar{z})$, for instance. To quantize general functions, in a class of symbols for instance, the most convenient way is to define : $a$ : as a pseudodifferential operator on analytic functions $u(z)$ in the Bargmann-Fock-Segal (or holomorphic) representation [29], where $: z_{j}:$ is multiplication by $z_{j}$ and $: \bar{z}_{j}:=\partial / \partial z_{j}$; then Wick ordering in this representation is analogous to the mathematical ordering (all differentiations $\partial / \partial z_{j}$ pushed to the right). But since the Bargmann function spaces have a quite different structure, this analogy is rigorous in a nontrivial way; here we simply give the formal analogs of Eqs. (A2') and (A7) wtihout their rigorous derivation:
-For the adjoint : $a(\bar{z}, \bar{z}):^{+}=: a^{*}(\bar{z}, z):$,
-for the product of "normal" symbols (in the sense that $: c:=: a:: b:$ ),

$$
c(z, \bar{z}) \sim a(z, \bar{z})\left[\exp \left(\stackrel{\rightharpoonup}{\partial}_{\bar{z}} \stackrel{\rightharpoonup}{\partial}_{z}\right)\right] b(z, \bar{z})
$$

i.e.,

$$
: a: b: \sim \sum_{\alpha \in \mathbb{N}^{b}} \frac{1}{\alpha!} \partial_{\xi_{z}}^{\alpha} a \cdot \partial_{z}^{\alpha} b:
$$

This last formula is known as Wick's theorem for normal-ordered products [10]. It is an exact relation for polynomials.

Further properties of Wick orderings are discussed in [11].
Some other quantization procedures, with their associated symbols calculus, are discussed in [30].

Geometric quantization [32], when performed in the canonical coordinates ( $q, p$ ) of $X$, also associates an operator $a^{Q}$ on $C_{0}^{\infty}(Q)$ to suitable functions $a(x)$. According to Gawedzki [32, p. 215]: $a^{Q}=\left(a-(i / 2) \partial_{q} \partial_{p} a\right)^{M}$ (with our notation), hence formally: $a^{O}=\widehat{a^{\prime}}$ with $a^{\prime}(q, p)=\sum_{|\alpha| \geqslant 2}(1 / \alpha!)(i / 2)^{|\alpha|} \partial_{q}{ }^{\alpha} \partial_{p}{ }^{\alpha} a(q, p)$, by (A5), and $a^{O} \equiv \hat{a}$ only for $a(x)$ polynomial of degree $\leqslant 2$; otherwise the algebraic structure is not very attractive.

Note added in proof. When this manuscript was completed we received an article (previously unknown to us) by Berezin and Shubin [31] which deals with the same spaces of symbols and gives several of the results proved here in Section 2. We have nevertheless kept our text unchanged with the idea that more detailed proofs than those in Ref. [31] might be useful.

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