# SCATTERING THEORY FOR PSEUDO-DIFFERENTIAL OPERATORS I. THE EXISTENCE OF WAVE OPERATORS 

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## 1. Introduction

This paper and the following paper are concerned with scattering theory for pseudo-differential operators. The operators we consider are of the form

$$
\begin{equation*}
P(D)+A(X, D) \tag{1.1}
\end{equation*}
$$

in $\boldsymbol{R}^{n}$, where the unperturbed operator $P(D)$ and the perturbation $A(X, D)$ are pseudo-differential operators. We examine in the present paper the existence of wave operators, while we shall prove the completeness of wave operators in the succeeding paper.

We briefly recall the definition of wave operators $W_{ \pm}$. Let $H$ and $H_{0}$ be the self-adjoint realizations of $P(D)+A(X, D)$ and $P(D)$ in $L^{2}\left(\boldsymbol{R}^{n}\right)$, respectively. Then $W_{ \pm}$are defined by the limits

$$
\begin{equation*}
W_{ \pm}=s-\lim _{t \rightarrow \infty} e^{i t H} e^{-i t H_{o}} P_{\mathrm{ac}}\left(H_{0}\right) \tag{1.2}
\end{equation*}
$$

Here $P_{\mathrm{ac}}\left(H_{0}\right)$ denotes the orthogonal projection onto the subspace of absolutely continuity with respect to $H_{0}$. (We refer to Kato [7, Chapter $X$ ] for definitions and results from spectral theory.)

Some authors have studied scattering theory for pseudo-differential operators. Among others, recently Simon [10] has considered operators $H=H_{0}+V$ where $H_{0}$ is a pseudo-differential operator and showed that the main conclusions of scattering theory hold (namely the wave operators exist and are complete etc.). He does not necessarily require that $V$ be a differential operator, or even a pseudo-differential operator. In fact, he only needs that $V$ be a symmetric operator with some falloff at infinity. He used the methods which have been originally found by Enss [5]. The condition that Simon calls the Enss condition (see [10]) plays an important role in proving the completeness of wave operators. Enss and Simon used purely time-dependent methods.

Schechter [9] considered the operators of the form (1.1) and proved the
main theorems of scattering theory. For the perturbations he took the specific operators of the form

$$
\Sigma q_{k}(X) \sigma_{k}(D)
$$

or of the form

$$
\Sigma \tau_{k}(D) q_{k}(X) \sigma_{k}(D) .
$$

Here $q_{k}(X)$ is an operator of multiplication by a function $q_{k}(x) ; \sigma_{k}(D)$ and $\tau_{k}(D)$ are pseudo-differential operators with symbols $\sigma_{k}(\xi)$ and $\tau_{k}(\xi)$ respectively. He exploited time-independent methods.

The main purpose of the present paper is to give a sufficient condition for wave operators to exist. It is easy to find examples which are covered by the present paper but which are not included in the results of [9], [10]. Such an example is given in Section 2. However, the hypotheses of both [9] and [10] assure the completeness of wave operators, which suggests that our hypotheses are too weak to assure the completeness of wave operators. In addition, we prove the symmetry and the self-adjointness of pseudo-differential operators under suitable hypotheses.

We make assumptions directly on the symbols of $P(D)$ and $A(X, D)$. But Schechter [9] and Simon [10] did not do so.

We use Cook's method, which is the main time-dependent technique, to show the existence of wave operators. Our proof is similar to that of Hörmander [6] which is based on the method of stationary phase, though he treated differential operators only.

Finally, we sketch the contents of this paper. Section 2 contains the main theorems and some examples. The proofs of the main theorems are given in Sections 3 and 4 after we prove several lemmas except a key lemma. We prove the key lemma in Section 5. In Section 6, we give a necessary and sufficient condition for pseudo-differential operators to be symmetric. In Section 7, we give sufficient conditions for pseudo-differential operators to have self-adjoint extentions.

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## 2. Existence of wave operators

In this section, we shall mention two basic theorems and some examples. One of the theorems is a characterization of the subspace of absolute continuity with respect to $H_{0}=P(D)$. The other theorem gurantees the existence of wave operators.

Before giving the assumption of $P(D)$ and $A(X, D)$ we shall list the notations which will be employed in the sequel without further reference.
$\boldsymbol{R}^{n}$; product of $n$ copies of the real line $\boldsymbol{R}$.
$\langle x, \xi\rangle=x_{1} \xi_{1}+\cdots+x_{n} \xi_{n},\left(x, \xi \in \boldsymbol{R}^{n}\right)$.
$\boldsymbol{Z}_{+}^{n}$; product of $n$ copies of the set $\mathbf{Z}_{+}$of nonnegative integers.
$|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$; the length of the multi-index $\alpha \in \mathbf{Z}_{+}^{n} . \quad \alpha!=\alpha_{1}!\cdots \alpha_{n}!$.
supp $u$; the support of a function $u$.
$D_{j}=-i \partial / \partial x_{j}, \quad(i=\sqrt{-1})$.
$(\partial / \partial x)^{\alpha}=\left(\partial / \partial x_{1}\right)^{\alpha_{1}} \cdots\left(\partial / \partial x_{n}\right)^{\alpha_{n}}, \quad\left(\alpha \in \mathbf{Z}_{+}^{n}\right)$.
$f^{\prime}=\left(\partial f / \partial x_{1}, \cdots, \partial f / \partial x_{n}\right)$; the gradient of a function $f$.
$f^{\prime \prime}=\left(\partial^{2} f / \partial x_{j} \partial x_{k}\right)_{j, k=1, \cdots, n}$; the Hesse matrix of $f$.
$\mathscr{F} u(\xi)=\hat{u}(\xi)=\int e^{-i\langle x, \xi\rangle} u(x) d x ;$ Fourier transform of $u=u(x)$.
$\overline{\mathscr{F}} v(x)=\int e^{i\langle x, \xi\rangle} v(\xi) d \xi$; inverse Fourier transform of $v=v(\xi),\left(d \xi=(2 \pi)^{-n} d \xi\right)$.
$\Delta=\left(\partial / \partial x_{1}\right)^{2}+\cdots+\left(\partial / \partial x_{n}\right)^{2}$; the Laplace operator.
$C^{N}(\Omega)$; the space of complex-valued functions, defined and $N$ times continuously differentiable in an open set $\Omega$, equipped with the topology of uniform convergence on every compact subset of $\Omega$, of the functions and of each one of their derivatives of order $<N+1\left(N \in \mathbf{Z}_{+}\right.$or $\left.N=+\infty\right)$.
$C_{0}^{N}(\Omega)$; the subspace of $C^{N}(\Omega)$ consisting of the functions having a compact support; if $u \in C_{0}^{N}(\Omega)$ we write

$$
|u|_{N}=\max _{|\alpha| \leqslant N} \sup _{x}\left|(\partial / \partial x)^{\alpha} u(x)\right| .
$$

$L^{2}\left(\boldsymbol{R}^{n}\right)$; the Hilbert space of measurable functions $u$ square integrable over $\boldsymbol{R}^{n}$, equipped with the norm

$$
\|u\|_{L^{2}}=\left(\int|u(x)|^{2} \mathrm{~d} x\right)^{1 / 2}
$$

(if there are no risks of confusion we will omit the subscript $L^{2}$ and write as $\|u\|)$.
$\mathcal{S}\left(\boldsymbol{R}^{n}\right)$; the space of $C^{\infty}$ functions $u$ in $\boldsymbol{R}^{n}$ such that, for any nonnegative integer $N$,

$$
|u|_{N, S} \equiv \max _{|\alpha|+k \leqslant N} \sup _{x}\left(1+|x|^{2}\right)^{k / 2}\left|(\partial / \partial x)^{\infty} u(x)\right|<\infty
$$

equipped with the topology defined by the seminorms $\left|\left.\right|_{N, S}\right.$.
$\mathcal{S}^{\prime}\left(\boldsymbol{R}^{n}\right)$; the dual space of $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$, also the space of tempered distributions in $\boldsymbol{R}^{n}$.
$H_{s}\left(\boldsymbol{R}^{n}\right)$; the Sobolev space of order $s \in \boldsymbol{R}$ in $\boldsymbol{R}^{n}$, i.e., the space of tempered distributions $u$ in $\boldsymbol{R}^{n}$ whose Fourier transform $\hat{u}$ is a measurable function such that

$$
\|u\|_{s}=\left(\int\left(1+|\xi|^{2}\right)^{s}|\hat{u}(\xi)|^{2} d \xi\right)^{1 / 2}<+\infty
$$

equipped with the Hilbert space structure defined by the norm || $\|_{s}$.
meas $_{k}$; the Lebesgue measure on $\boldsymbol{R}^{k}$.
[s] ; the integral part of $s \in \boldsymbol{R}$.
$\mathcal{H}_{\mathrm{p}}(H)$; the closed linear manifold spanned by all eigenvectors of an operator $H$ in $L^{2}\left(\boldsymbol{R}^{n}\right)$.
$\mathcal{H}_{\mathrm{ac}}(H)$; the subspace of absolute continuity with respect to a self-adjoint operator $H$ in $L^{2}\left(\boldsymbol{R}^{n}\right)$.
$P_{\mathrm{ac}}(H)$; the orthogonal projection onto $\mathcal{H}_{\mathrm{ac}}(H)$.
$\mathcal{A}_{\mathrm{s}}(H)$; the subspace of singularity with respect to a self-adjoint operator $H$ in $L^{2}\left(\boldsymbol{R}^{n}\right)$.

Now let us make the following assumption on symbols $p(\xi)$ of $P(D)$ and $a(x, \xi)$ of $A(X, D)$. The set of all critical points of $p(\xi)$ will be denoted by $\Sigma$ :

$$
\Sigma=\left\{\xi \in \boldsymbol{R}^{n} \mid p^{\prime}(\xi)=0\right\}
$$

## Assumption 2.1.

(A) $\quad a(x, \xi)$ is a complex-valued $C^{\infty}$ function on $\boldsymbol{R}_{x}^{n} \times \boldsymbol{R}_{\xi}^{n}$ such that for all multi-indicies $\alpha, \beta$ the estimate

$$
\begin{aligned}
& \left|(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} a(x, \xi)\right| \\
& \leqslant C_{\alpha \beta}(1+|x|)^{l+\tau|\alpha|}(1+|\xi|)^{m+\delta|\beta|}, \quad x \in \boldsymbol{R}^{n}, \xi \in \boldsymbol{R}^{n},
\end{aligned}
$$

is valid for some constant $C_{\alpha \beta}$, where $l, m, \delta$ and $\tau$ are constants with $l, m \geqslant 0,0 \leqslant \delta, \tau<1$
(P.1) $p(\xi)$ is a real-valued $C^{\infty}$ function on $\boldsymbol{R}_{\xi}^{n}$ such that for every multiindex $\alpha$, the estimate

$$
\left|(\partial / \partial \xi)^{\omega} p(\xi)\right| \leqslant C_{\infty}(1+|\xi|)^{N_{\infty}}
$$

is valid for some constants $C_{\alpha}$ and $N_{\alpha}$.
(P.2) There is a closed set $\Xi \subset \boldsymbol{R}^{n}$ with the following properties:
(a) $\mathrm{meas}_{n}(\Xi \backslash \Sigma)=0$;
(b) every point $\xi_{0} \in \boldsymbol{R}^{n} \backslash \Xi$ has a neighborhood where $\operatorname{rank} p^{\prime \prime}(\xi)$ is constant;
(c) $\operatorname{rank} p^{\prime \prime}(\xi) \geqslant 1, \xi \in \boldsymbol{R}^{n} \backslash \Xi$.
(H) $\quad P(D)+A(X, D)$ with domain $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ has a self-adjoint extention in $L^{2}\left(\boldsymbol{R}^{n}\right)$.

Let $H$ be a self-adjoint extention in $L^{2}\left(\boldsymbol{R}^{n}\right)$ of $P(D)+A(X, D)$ with domain $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ and let $H_{0}$ be the closure in $L^{2}\left(\boldsymbol{R}^{n}\right)$ of $P(D)$ with domain $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$. Our main results are:

Theorem 1. Let hypothesis (P.1) be fulfilled. Then we have

$$
\begin{aligned}
& \mathscr{A}_{\mathrm{ac}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \Sigma\right\} \\
& \mathcal{H}_{\mathrm{s}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \boldsymbol{R}^{n} \backslash \Sigma\right\} .
\end{aligned}
$$

Theorem 2. Let Assumption 2.1 be fulfilled and put

$$
\Omega_{k}=\left\{\xi \in \boldsymbol{R}^{n} \backslash \Xi \mid p^{\prime}(\xi) \neq 0, \operatorname{rank} p^{\prime \prime}(\xi)=k\right\} \quad(k=1, \cdots, n) .
$$

For every $k$ such that $\Omega_{k} \neq \phi$ assume that,
( $c_{k}$ ) for any $r>0$ and any compact set $K \subset \Omega_{k}$ there is an integer $N_{k} \geqslant[k / 2]+2$ such that

$$
\begin{aligned}
& \int_{1}^{\infty} t^{(n-k) / 2}\left[\int_{r<|y|<2 r}\left(1+t d_{K}(y)\right)^{-N_{k}}\right. \\
& \left.\times\left\{\max _{|\alpha|<2 N_{k}}\left(\sup _{\xi \in K}\left|(\partial / \partial \xi)^{\alpha} a( \pm t y, \xi)\right|\right)^{2}\right\} d y\right]^{1 / 2} d t<\infty .
\end{aligned}
$$

Then the wave operators (1.2) exist. Here $[k / 2]$ is the integral part of $k / 2$ and $d_{K}(y)$ is the distance from $y$ to the set

$$
p^{\prime}(K)=\left\{p^{\prime}(\xi) \mid \xi \in K\right\}
$$

Theorems 1 and 2 will be proven in Sections 3 and 4, respectively.
Remark 2.1. Since $p(\xi)$ is a real-valued $C^{\infty}$ function, $P(D)$ with domain $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ is essentially self-adjoint. Its unique self-adjoint extention $H_{0}$ is given by its closure. The domain $\mathscr{D}\left(H_{0}\right)$ consists of those $u$ in $L^{2}\left(\boldsymbol{R}^{n}\right)$ such that $p \hat{u}$ is also in $L^{2}\left(\boldsymbol{R}^{n}\right)$.

In the rest of this section we illustrate some applications of Theorems 1 and 2.

Example 2.2. Let $n=1$ and let $p(\xi)=\varphi(\xi) \xi^{2}(\xi \in \boldsymbol{R})$ where $\varphi \in C^{\infty}(\boldsymbol{R})$ with $\varphi(\xi)=0$ (resp. 1) for $|\xi| \leqslant 1$, (resp. $\geqslant 2$ ). Suppose that

$$
\xi \varphi^{\prime}(\xi)+2 \varphi(\xi) \neq 0, \quad|\xi|>1
$$

Then it follows immediatly from Theorem 1 that

$$
\mathcal{H}_{\mathrm{ac}}\left(H_{0}\right)=\left\{u \in L^{2}(\boldsymbol{R}) \mid \hat{u}(\xi)=0 \text { for almost every }|\xi| \leqslant 1\right\}
$$

It is easy to see that $H_{0}$ has a single eigenvalue which is equal to zero (see the proof of Proposition 3.5).

Let $A$ be a multiplication operator by a real-valued $C^{\infty}$ function $a(x)$ such that

$$
\begin{equation*}
|a(x)| \leqslant C(1+|x|)^{-1-\varepsilon}, \quad x \in \boldsymbol{R}, \varepsilon>0 . \tag{2.1}
\end{equation*}
$$

Put $H=H_{0}+A$. Then $H$ is self-adjoint since $A$ is bounded. Applying Theorem 2, we will show the existence of wave operators (1.2). Suppose that $p^{\prime \prime}(\xi) \neq 0$ for $|\xi|>1$ except finitely many points $\xi_{1}, \cdots, \xi_{N}$. Put

$$
\Xi=\{\xi \in \boldsymbol{R} \mid-1 \leqslant \xi \leqslant 1\} \cup\left\{\xi_{1}, \cdots, \xi_{N}\right\} .
$$

Then Assumption 2.1 is fulfilled. Since

$$
\int_{1}^{\infty}\left(\int_{r<|y|<2 r}|a( \pm t y)|^{2} d y\right)^{1 / 2} d t \leqslant C \int_{1}^{\infty} t^{-1-\varepsilon} d t
$$

for every $r>0$, the condition $\left(c_{1}\right)$ holds. Hence, from Theorem 2, it follows that the wave operators (1.2) exist.

On the other hand, the proper wave operators $W_{ \pm}=s-\lim _{t \rightarrow \pm \infty} e^{i t H} e^{-i t H_{0}}$ do not necessarily exist. For example, let us assume, in addition to (2.1), that $a(x)$ does not vanish almost everywhere and suppose that the proper wave operators exist. Let $u$ be an eigenvector of $H_{0}$ corresponding to the eigenvalue zero. Then $e^{-i t H_{0}} u=u$ for all $t \in \boldsymbol{R}$. Since the proper wave operators exist,

$$
\left\|e^{i t H} u-u\right\|=\left\|e^{i(t+s) H} e^{-i(t+s) H_{0}} u-e^{i s H} e^{-i s H_{0}} u\right\|
$$

converges to zero as $s$ tends to $\infty$. Therefore

$$
\begin{equation*}
e^{i t H} u=u \tag{2.2}
\end{equation*}
$$

It follows immediately from (2.2) that $u \in \mathscr{D}(H)$ and $H u=0$. Since $H=H_{0}+A$ and since $H_{0} u=0$, we have $a u=0$. Thus $u=0$ because $a \neq 0$ a.e.. This contradicts the fact that $u$ is an eigenvector. Hence the proper wave operators do not exist.

Next, we shall give an example which satisfies the conditions of Theorem 2 but which does not satisfy the Enss condition [10]. Let us first recall the definition of the Enss condition. Let $A$ be a symmetric operator in $L^{2}\left(\boldsymbol{R}^{n}\right)$ such that $\mathscr{D}(A) \supset H_{2 N}$ for some $N$. Put

$$
\begin{equation*}
h(R)=\left\|A\left((-\Delta)^{N}+1\right)^{-1} F(|x| \geqslant R)\right\| \tag{2.3}
\end{equation*}
$$

where

$$
F(|x| \geqslant R) u(x)= \begin{cases}0 & (|x|<R) \\ u(x) & (|x| \geqslant R) .\end{cases}
$$

The norm in (2.3) is the operator norm. Simon [10] calls the condition

$$
\begin{equation*}
h(0)<\infty, \int_{0}^{\infty} h(R) d R<\infty \tag{2.4}
\end{equation*}
$$

the Enss condition. Roughly speaking, the Enss condition implies the existence and the completeness of the wave operators.

Example 2.3. Let $n=2$ and let $A$ be a multiplication operator by a realvalued $C^{\infty}$ function $a(x)$ on $\boldsymbol{R}^{2}$ such that

$$
\left\{\begin{array}{l}
|a(x)| \leqslant C, x=\left(x_{1}, x_{2}\right) \in \boldsymbol{R}^{2}  \tag{2.5}\\
|a(x)| \geqslant C(1+|x|)^{-\varepsilon}, x_{2}^{4}>x_{1}^{2}+1,0 \leqslant \varepsilon \leqslant 1 \\
a(x)=0, x_{2}^{4}<x_{1}^{2}
\end{array}\right.
$$

Then $A$ is a bounded operator on $L^{2}\left(\boldsymbol{R}^{2}\right)$. But the integral in (2.4) diverges. We show this by using the fact that when $h(0)<\infty$,

$$
\int_{0}^{\infty} h(R) d R<\infty
$$

if and only if

$$
\begin{equation*}
\int_{0}^{\infty}\left\|A J_{R}\left((-\Delta)^{N}+1\right)^{-1}\right\| d R<\infty \tag{2.6}
\end{equation*}
$$

Here $J_{R}$ is the multiplication operator by the function $J_{R}(x)=\varphi(x / R)$ with $\varphi \in C^{\infty}\left(\boldsymbol{R}^{2}\right)$ and $\varphi(x)=0$ (resp. 1) for $|x| \leqslant 1$, (resp. $\geqslant 2$ ). (The details can be found on p. 124 of Simon [10].)

Now, let $N$ be an integer. Choose $v \in C_{0}^{\infty}$ such that

$$
\text { supp } v \subset\left\{x \in \boldsymbol{R}^{2}| | x \mid<1\right\},\left\|\left((-\Delta)^{N}+1\right) v\right\|_{L^{2}}=1
$$

and put

$$
\begin{aligned}
& v_{R}(x)=v\left(x-3 R e_{2}\right), \quad e_{2}=(0,1) \\
& u_{R}(x)=\left((-\Delta)^{N}+1\right) v_{R}(x)
\end{aligned}
$$

Then it follows that

$$
\left\|u_{R}\right\|^{2}=\int\left|\left(|\xi|^{2 N}+1\right) \hat{v}(\xi)\right|^{2} d \xi=1
$$

Noting that

$$
\operatorname{supp} v_{R} \subset\left\{x \in \boldsymbol{R}^{2}| | x \mid \geqslant 2 R\right\}
$$

we have

$$
\begin{aligned}
\left\|A J_{R}\left((-\Delta)^{N}+1\right)^{-1}\right\| & \geqslant\left\|A J_{R}\left((-\Delta)^{N}+1\right)^{-1} u_{R}\right\| \\
& =\left\|A J_{R} v_{R}\right\| \\
& =\left\|A v_{R}\right\| \\
& \geqslant C(1+R)^{-\varepsilon}\|v\| .
\end{aligned}
$$

Since $\varepsilon \leqslant 1$, the integral in (2.6) diverges. Thus $\int_{0}^{\infty} h(R) d R$ diverges and the Enss condition is not satisfied.

On the other hand, if we regard $A$ as a perturbation of $P(D)=\left(1+D_{1}^{2}\right)^{1 / 2}$, then the existence of wave operators follows from Theorem 2. In fact, Assumption 2.1 is clearly fulfilled (with $\Xi=\left\{\xi \in \boldsymbol{R}^{2} \mid \xi_{1}=0\right\}$ ). Since the symbol of $P(D)$ is $\left(1+\xi_{1}^{2}\right)^{1 / 2}$, we have

$$
\begin{aligned}
& p^{\prime}(\xi)=\left(\xi_{1}\left(1+\xi_{1}^{2}\right)^{-1 / 2}, 0\right) \\
& p^{\prime \prime}(\xi)=\left[\begin{array}{cc}
\left(1+\xi_{1}^{2}\right)^{-3 / 2} & 0 \\
0 & 0
\end{array}\right]
\end{aligned}
$$

and

$$
\begin{aligned}
& \Omega_{1}=\left\{\xi \in \boldsymbol{R}^{2} \mid \xi_{1} \neq 0\right\} \\
& \Omega_{2}=\phi
\end{aligned}
$$

(see Theorem 2 for the definition of $\Omega_{k}$ ). We examine that the integral in ( $c_{1}$ ) converges for any $r>0$ and any compact set $K \subset \Omega_{1}$. Set

$$
\mathcal{O}_{t}=\left\{y \in \boldsymbol{R}^{2}\left|r \leqslant|y| \leqslant 2 r, t^{4} y_{2}^{4} \geqslant t^{2} y_{1}^{2}\right\}\right.
$$

Noting that the set $p^{\prime}(K)$ is in the $y_{1}$-axis, we have

$$
\begin{equation*}
d_{K}(y) \geqslant C t^{-1 / 2}, \quad y \in \mathcal{O}_{t} \tag{2.7}
\end{equation*}
$$

By (2.5) and (2.7)

$$
\begin{aligned}
& \int_{1}^{\infty} t^{1 / 2}\left(\int_{r \leqslant|y| \leqslant 2 r}\left(1+t d_{K}(y)\right)^{-N_{1}}|a( \pm t y)|^{2} d y\right)^{1 / 2} d t \\
& \quad \leqslant C \int_{1}^{\infty} t^{1 / 2}\left(\int_{O_{t}}\left(1+t^{1 / 2}\right)^{-N_{1}} d y\right)^{1 / 2} d t \\
& \quad \leqslant C \int_{1}^{\infty}(1+t)^{\left(2-N_{1}\right) / 4} d t
\end{aligned}
$$

Hence the condition $\left(c_{1}\right)$ holds if we take $N_{1}=7$. Thus the wave operators exist. Incidentally it follows from Theorem 1 that $\mathcal{H}_{\mathrm{ac}}\left(H_{0}\right)=L^{2}\left(\boldsymbol{R}^{2}\right)$. We should note that if we take $P(D)=\Delta$ then $\left(c_{1}\right)$ does not hold.

Example 2.4 (Higher order perturbations). Let us consider two selfadjoint operators

$$
\begin{gather*}
H_{0}=(1-\Delta)^{1 / 2}  \tag{2.8}\\
H=(1-\Delta)^{1 / 2}+\left(1+|x|^{2}\right)^{-8 / 4}(1-\Delta)^{3 / 2}\left(1+|x|^{2}\right)^{-8 / 4} \tag{2.9}
\end{gather*}
$$

in $L^{2}\left(\boldsymbol{R}^{3}\right)$, where $\varepsilon>1$. Since the right side of (2.9) with domain $\mathcal{S}\left(\boldsymbol{R}^{3}\right)$ is a real operator, it admits a self-adjoint extention. We want to check the hypotheses of Theorem 2. The symbol of the perturbation

$$
\left(1+|x|^{2}\right)^{-8 / 4}(1-\Delta)^{3 / 2}\left(1+|x|^{2}\right)^{-8 / 4}
$$

is given by the following oscillatory integral:

$$
a(x, \xi)=\mathrm{O}_{\mathrm{s}}-\iint e^{-i\langle y, \eta\rangle}\left(1+|x|^{2}\right)^{-8 / 4}\left(1+|\xi+\eta|^{2}\right)^{3 / 2}\left(1+|x+y|^{2}\right)^{-\varepsilon / 4} d y d \eta
$$

(see Kumano-go [8]). By repeated integration by parts, we have

$$
\begin{align*}
& a(x, \xi)  \tag{2.10}\\
&=\left(1+|x|^{2}\right)^{-\varepsilon / 4} \iint e^{-i\langle y, \eta\rangle}\left(1+|y|^{2}\right)^{-2}\left(1-\Delta_{\eta}\right)^{2}\left(1+|\eta|^{2}\right)^{-4}\left(1+|\xi+\eta|^{2}\right)^{3 / 2} \\
& \quad \times\left(1-\Delta_{y}\right)^{4}\left(1+|x+y|^{2}\right)^{-\varepsilon / 4} d y d \eta . \\
& \equiv\left(1+\left|x^{2}\right|\right)^{-\varepsilon / 4} \iint e^{-i\langle y, \eta\rangle} I d y d \eta .
\end{align*}
$$

Since, as can be easily verified,

$$
\begin{aligned}
& \left|(\partial / \partial \xi)^{\infty}(\partial / \partial x)^{\beta} I\right| \\
\leqslant & C_{\alpha \beta}^{1}(1+|y|)^{-4+\varepsilon / 2}(1+|\eta|)^{-5}(1+|x|)^{-8 / 2}(1+|\xi|)^{3-|\alpha|}
\end{aligned}
$$

with a constant $C_{\alpha \beta}^{1}$ for all multi-indices $\alpha, \beta$, it follows by differentiation under the integral sign that

$$
a(x, \xi) \in C^{\infty}\left(\boldsymbol{R}^{3} \times \boldsymbol{R}^{3}\right)
$$

and

$$
\begin{equation*}
\left|(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} a(x, \xi)\right| \leqslant C_{\alpha \beta}^{2}(1+|x|)^{-\varepsilon}(1+|\xi|)^{3-|\alpha|} \tag{2.11}
\end{equation*}
$$

with a constant $C_{\alpha \beta}^{2}$ for all $\alpha, \beta$. Hence the hypothesis $(A)$ is satisfied. The symbol $p(\xi)$ of $H_{0}$ is $\left(1+|\xi|^{2}\right)^{1 / 2}$, so it follows that $\Sigma=\{0\}$ and $\operatorname{rank} p^{\prime \prime}(\xi)=3$ at every $\xi \in \boldsymbol{R}^{3}$. Thus Assumption 2.1 is fulfilled (with $\Xi=\{0\}$ ). Finally, we check the condition $\left(c_{k}\right)$. In this case, we have $\Omega_{1}=\Omega_{2}=\phi$ and $\Omega_{3}=\boldsymbol{R}^{3} \backslash\{0\}$. Notice that, by (2.11),

$$
\int_{r<|y|<2 r}\left(\sup _{\xi \in \mathbb{K}}\left|(\partial / \partial \xi)^{\infty} a( \pm t y, \xi)\right|\right)^{2} d y \leqslant C(1+|t|)^{-2 \varepsilon}
$$

for any $r>0$, any compact set $K \subset \Omega_{3}$ and any multi-index $\alpha$. Then it is obvious that ( $c_{3}$ ) holds. Thus the wave operators exist. By Theorem 1, we have also $\mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)=L^{2}\left(\boldsymbol{R}^{3}\right)$.

## 3. Proof of Theorem 1

In the present section we shall prove Theorem 1 and some propositions on the spectral structure of $H_{0}$. First we should note that the hypothesis (P.1) can be further relaxed. We only need that $p(\xi)$ be $n$ times continuously differentiable.

As a preliminary to the proof, we note that since the Fourier transformation is unitary, it is sufficient to consider the self-adjoint operator $\hat{H}_{0}$ defined by

$$
\hat{H}_{0}=\mathscr{F} H_{0} \overline{\mathscr{F}}
$$

instead of $H_{0}$. It is easy to see that

$$
\begin{equation*}
u \in \mathscr{H}_{\mathrm{ac}}\left(H_{0}\right) \Leftrightarrow \hat{u} \in \mathscr{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right) . \tag{3.1}
\end{equation*}
$$

This implies that

$$
\begin{equation*}
u \in \mathscr{H}_{\mathrm{s}}\left(H_{0}\right) \Leftrightarrow \hat{u} \in \mathcal{H}_{\mathrm{s}}\left(\hat{H}_{0}\right) . \tag{3.2}
\end{equation*}
$$

The operator $\hat{H}_{0}$ coincides with the multiplication operator defined by

$$
\left\{\begin{array}{l}
\mathscr{D}\left(\hat{H}_{0}\right)=\left\{\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid p \hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right)\right\} \\
\hat{H}_{0} \hat{u}(\xi)=p(\xi) \hat{u}(\xi)
\end{array}\right.
$$

Let $\hat{E}(B)$ be the spectral measure associated with $\hat{H}_{0}$, where $B$ varies over all Borel sets of real line. Then

$$
\begin{equation*}
\|\hat{E}(B) \hat{u}\|^{2}=\int_{D^{-1}(B)}|\hat{u}(\xi)|^{2} d \xi, \quad \hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right) \tag{3.3}
\end{equation*}
$$

(see Kato [7], p. 520).
We prepare a lemma which will be used later.
Lemma 3.1. Let $\varphi \in C^{1}(\Omega)$ be a real-valued function in an open set $\Omega \subset \boldsymbol{R}^{n}$ and assume that $\varphi^{\prime}(\xi) \neq 0$ for every $\xi \in \Omega$. If $B$ is a Borel set of $\boldsymbol{R}$ with meas $_{1}(B)=$ 0 , then $\varphi^{-1}(B)$ is also a Borel set of $\boldsymbol{R}^{n}$, and meas $_{n}\left(\varphi^{-1}(B)\right)=0$.

Proof. It is easily seen that $\varphi^{-1}(B)$ is a Borel set of $\boldsymbol{R}^{n}$. We have only to show that $\operatorname{meas}_{n}\left(\varphi^{-1}(B)\right)=0$. Note that if

$$
\operatorname{meas}_{n}(\{\xi \in K \mid \varphi(\xi) \in B\})=0
$$

for every compact set $K \subset \Omega$ then

$$
\operatorname{meas}_{n}\left(\varphi^{-1}(B)\right)=0 .
$$

So it suffices to show that to every $\xi_{0} \in \Omega$ there corresponds a neighborhood $U$ of $\xi_{0}$ such that

$$
\operatorname{meas}_{n}(\{\xi \in U \mid \varphi(\xi) \in B\})=0
$$

Suppose $\xi_{0} \in \Omega$. Since $\varphi^{\prime}\left(\xi_{0}\right) \neq 0$, we may assume, without loss of generality, that $\partial \varphi\left(\xi_{0}\right) / \partial \xi_{1} \neq 0$. The Jacobian matrix of the map

$$
f: \xi \rightarrow\left(\varphi(\xi), \xi_{2}, \cdots, \xi_{n}\right)
$$

is non-singular at $\xi_{0}$. By the inverse function theorem, there is a neighborhood $U$ of $\xi_{0}$ such that $f$ is a diffeomorphism of class $C^{1}$ between $U$ and $f(U)$. Since

$$
f(\{\xi \in U \mid \varphi(\xi) \in B\}) \subset B \times \boldsymbol{R}^{n-1}
$$

we have

$$
\operatorname{meas}_{n}(f(\{\xi \in U \mid \varphi(\xi) \in B\})) \leqslant \operatorname{meas}_{n}\left(B \times \boldsymbol{R}^{n-1}\right)=0
$$

Thus

$$
\operatorname{meas}_{n}(\{\xi \in U \mid \varphi(\xi) \in B\})=\operatorname{meas}_{n}\left(f^{-1}(f(\{\xi \in U \mid \varphi(\xi) \in B\}))\right)=0
$$

where we use the fact that if $g$ is a $C^{1}$ map of a neighborhood of a null set $N \subset \boldsymbol{R}^{n}$ into $\boldsymbol{R}^{n}$, then $\operatorname{meas}_{n}(g(N))=0$.
Q.E.D.

Throughout this section, we assume that $p$ is real-valued.
Proposition 3.2. If $p \in C^{n}\left(\boldsymbol{R}^{n}\right)$, then

$$
\begin{align*}
& \mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \Sigma\right\}  \tag{3.4}\\
& \mathscr{H}_{\mathrm{s}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \boldsymbol{R}^{n} \backslash \Sigma\right\} \tag{3.5}
\end{align*}
$$

Remark 3.3. To prove Proposition 3.2, we apply Sard's theorem (see Sternberg [11], p. 47):

Let $M_{1}$ and $M_{2}$ be $C^{k}$ manifolds of dimension $n_{1}$ and $n_{2}$ respectively. Let $f$ be a map of class $C^{k}$ of $M_{1} \rightarrow M_{2}$. The critical values of $f$ form a set of measure zero if $k-1 \geqslant \max \left(n_{1}-n_{2}, 0\right)$.

Proof. As remarked before, we consider $\hat{H}_{0}$ instead of $H_{0}$. Define

$$
\begin{aligned}
& \mathcal{L}=\left\{\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \Sigma\right\} \\
& \mathscr{M}=\left\{\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \boldsymbol{R}^{n} \backslash \Sigma\right\} .
\end{aligned}
$$

Let $\hat{u} \in \mathcal{L}$, and let $B$ be a Borel set of the real line with $\operatorname{meas}_{1}(B)=0$. Then by (3.3) we have

$$
\|\hat{E}(B) \hat{u}\|^{2}=\int_{p^{-1}(B) \cap\left(R^{n} \backslash \Sigma\right)}|\hat{u}(\xi)|^{2} d \xi
$$

Since $p^{\prime}(\xi) \neq 0$ for every $\xi \in \boldsymbol{R}^{n} \backslash \Sigma$, we have

$$
\operatorname{meas}_{n}\left(p^{-1}(B) \cap\left(\boldsymbol{R}^{n} \backslash \Sigma\right)\right)=0
$$

by Lemma 3.1. Thus

$$
\|\hat{E}(B) \hat{u}\|^{2}=0
$$

which means that $\hat{u} \in \mathscr{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right)$. Thus we have $\mathcal{L} \subset \mathscr{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right)$.
Similarly, we show that $\mathscr{M} \subset \mathcal{H}_{\mathrm{s}}\left(\hat{H}_{0}\right)$. Let $\hat{u} \in \mathcal{M}$, and let $B_{0}$ be the set of critical values of $p$. Then, by Sard's theorem

$$
\operatorname{meas}_{1}\left(B_{0}\right)=0
$$

Moreover, $B_{0}$ is a Borel set. In fact, with $B_{j}=\left\{\xi \in \boldsymbol{R}^{n}| | \xi \mid \leqslant j\right\}$

$$
B_{0}=p(\Sigma)=\bigcup_{j=1}^{\infty} p\left(\Sigma \cap B_{j}\right)
$$

decomposes $B_{0}$ as the union of a countable collection of closed sets. If $B \subset$ $\boldsymbol{R} \backslash B_{0}$, then $p^{-1}(B) \cap \Sigma=\phi$, and

$$
\|\hat{E}(B) \hat{u}\|^{2}=\int_{p^{-1}(B) \cap \Sigma}|\hat{u}(\xi)|^{2} d \xi=0 .
$$

Thus the measure $\|\hat{E}(B) \hat{u}\|^{2}$ is singular with respect to the Lebesgue measure, so $\hat{u} \in \mathscr{H}_{\mathrm{s}}\left(\hat{H}_{0}\right)$. Therefore we have $\mathscr{M} \subset \mathcal{H}_{\mathrm{s}}\left(\hat{H}_{0}\right)$.

It is easy to see that $\mathcal{L}$ and $\mathscr{M}$ are closed linear subspaces of $L^{2}\left(\boldsymbol{R}^{n}\right)$, are orthogonal complements to each other. This implies that $\mathcal{L}=\mathcal{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right)$ and $\mathscr{M}=\mathscr{H}_{\mathrm{s}}\left(\hat{H}_{0}\right)$. Thus, by (3.1) and (3.2), we obtain the conclusions. Q.E.D.

Proof of Theorem 1. It is an immediate consequence of Proposition 3.2.
Q.E.D.

Remark 3.4. As mentioned before, the hypothesis ( $P$. 1) can be further relaxed. Simon [10] allows the possibility of singular points to include an example like $p(\xi)=|\xi|$. In such cases we are also able to show a result similar to Theorem 1:

$$
\begin{aligned}
& \mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \overline{C_{p} \cup S_{p}}\right\} \\
& \mathscr{H}_{\mathrm{s}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \notin \overline{\left.C_{p} \cup S_{p}\right\}} .\right.
\end{aligned}
$$

Here $C_{p}\left(S_{p}\right)$ is the set of critical (singular) points of $p$. (See [10] for the definition.) The proof exactly follows from that of Proposition 3.2 with a minor change.

If meas $_{n}(\Sigma)>0$ and the boundary $\partial \Sigma=\Sigma \backslash$ Int $\Sigma$ is of measure zero, then we can improve Proposition 3.2. It should be noted that that $p \in C^{1}\left(\boldsymbol{R}^{n}\right)$ does not imply that $\partial \Sigma$ is of measure zero.

Proposition 3.5. Let meas $_{n}(\Sigma)>0$. If $p \in C^{1}\left(\boldsymbol{R}^{n}\right)$ and meas $_{n}(\partial \Sigma)=0$, then we have

$$
\begin{aligned}
& \mathcal{H}_{\mathrm{ac}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \Sigma\right\} \\
& \mathcal{H}_{\mathrm{s}}\left(H_{0}\right)=\left\{u \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \boldsymbol{R}^{n} \backslash \Sigma\right\} .
\end{aligned}
$$

Moreover,

$$
\mathcal{H}_{p}\left(H_{0}\right)=\mathscr{H}_{\mathbf{s}}\left(H_{0}\right) .
$$

Proof. We also consider the operator $\hat{H}_{0}$ and define $\mathcal{L}$ and $\mathscr{M}$ as in the proof of Proposition 3.2. It is obvious that $\mathcal{L} \subset \mathcal{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right)$. We show that $\mathscr{M}=\mathscr{H}_{p}\left(\hat{H}_{0}\right)$ which implies that $\mathscr{M} \subset \mathscr{H}_{\mathrm{s}}\left(\hat{H}_{0}\right)$.

Let $\left\{O_{j}\right\}$ be the collection of connected components of Int $\Sigma$ (note that Int $\Sigma \neq \phi)$. Then, as can be easily verified, there exists $\lambda_{j} \in \boldsymbol{R}$ corresponding to $O_{j}$ such that

$$
p(\xi)=\lambda_{j}, \quad \xi \in O_{j} .
$$

Since the characteristic function $1_{o_{j}}(\xi)$ satisfies that

$$
p(\xi) 1_{o_{j}}(\xi)=\lambda_{j} 1_{o_{j}}(\xi), \quad \xi \in \boldsymbol{R}^{n}
$$

we have $\left\{\lambda_{j}\right\} \subset \sigma_{\mathrm{p}}\left(\hat{H}_{0}\right)$, where $\sigma_{\mathrm{p}}\left(\hat{H}_{0}\right)$ denotes the set of eigenvalues of $\hat{H}_{0}$. Furthermore, if $\hat{H}_{0} \hat{u}=\lambda \hat{u}$ with $\lambda \in \boldsymbol{R} \backslash\left\{\lambda_{j}\right\}$, then $\hat{u}=0$ by

$$
\begin{equation*}
\operatorname{meas}_{n}\left(\left\{\xi \in \boldsymbol{R}^{n} \mid p(\xi)=\lambda\right\}\right)=0 . \tag{3.6}
\end{equation*}
$$

In (3.6) we used the hypothesis

$$
\operatorname{meas}_{n}(\partial \Sigma)=0
$$

and the fact

$$
\operatorname{meas}_{n}\left(\left\{\xi \in \boldsymbol{R}^{n} \backslash \Sigma \mid p(\xi)=\lambda\right\}\right)=0
$$

which follows from Lemma 3.1. Thus we have $\sigma_{\mathrm{p}}\left(H_{0}\right)=\left\{\lambda_{j}\right\}$. Finally, writing

$$
\mathscr{M}_{j}=\left\{\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \boldsymbol{R}^{n} \backslash O_{j}\right\},
$$

we see that $\mathscr{M}_{j}$ is the eigenspace of $\hat{H}_{0}$ corresponding to the eigenvalue $\lambda_{j}$ and that the closed linear manifold spanned by all $\mathscr{M}_{j}$ coincides with $\mathscr{M}$. Hence, by the definition of $\mathscr{H}_{\mathrm{p}}\left(\hat{H}_{0}\right), \mathscr{M}$ coincides with $\mathscr{H}_{\mathrm{p}}\left(\hat{H}_{0}\right)$.

Thus we have shown that $\mathcal{L} \subset \mathcal{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right)$ and $\mathscr{M}=\mathcal{H}_{\mathrm{p}}\left(\hat{H}_{0}\right)$. By the arguments in the last step of the proof of Proposition 3.2, the result follows.
Q.E.D.

It is well known that when $p \in C^{1}\left(\boldsymbol{R}^{n}\right)$, that meas $_{n}(\Sigma)=0$ implies that $\mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)=L^{2}\left(\boldsymbol{R}^{n}\right)$. The converse is also true when $p \in C^{n}\left(\boldsymbol{R}^{n}\right)$.

Proposition 3.6. Suppose $p \in C^{n}\left(\boldsymbol{R}^{n}\right)$. If $\mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)=L^{2}\left(\boldsymbol{R}^{n}\right)$, then meas ${ }_{n}(\Sigma)$ $=0$.

Proof. We consider $\hat{H}_{0}$ instead of $H_{0}$. Let $B_{0}$ be the set of critical values of $p$. Then, by Sard's theorem, meas $\left(B_{0}\right)=0$. Recall (3.1). Since $\mathscr{H}_{\mathrm{ac}}\left(\hat{H}_{0}\right)=$ $L^{2}\left(\boldsymbol{R}^{n}\right)$, the measure $\|\hat{E}(B) \hat{u}\|^{2}$ is absolutely continuous with respect to the Lebesgue measure for all $\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right)$, and so $\left\|\hat{E}\left(B_{0}\right) \hat{u}\right\|^{2}=0$ for all $\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right)$. Noting (3.3), we have $\operatorname{meas}_{n}\left(p^{-1}\left(B_{0}\right)\right)=0$. Thus we obtain

$$
\operatorname{meas}_{n}(\Sigma)=0
$$

since $\Sigma \subset p^{-1}\left(B_{0}\right)$.
Q.E.D.

## 4. Proof of Theorem 2

In the present section we shall prove Theorem 2. For the proof we need several lemmas.

Lemma 4.1. Under hypothesis $(A)$ the psuedo-differential operator $A=$ $A(X, D)$ is continuous from $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ to $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$.

Proof. Suppose $u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$. Then

$$
A u(x)=\int e^{i\langle x, \xi\rangle} a(x, \xi) \hat{u}(\xi) d \xi
$$

Since $\hat{u} \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$, we can differentiate with respect to $x$ under the integral sign as often as we like. Therefore, it follows that $A u \in C^{\infty}\left(\boldsymbol{R}^{n}\right)$. Moreover, for every multi-index $\alpha$ and every integer $k \geqslant 0$ we obtain

$$
\begin{equation*}
(1+|x|)^{-l+2(1-\tau) k}\left|(\partial / \partial x)^{\infty} A u(x)\right| \leqslant C_{k, \infty}|u|_{N, S} \tag{4.1}
\end{equation*}
$$

where $C_{k, \infty}$ is a constant and $N=|\alpha|+[m]+2(n+k)+3$. Indeed, by differentiation under the integral sign and integration by parts, we see that

$$
\begin{aligned}
& \left(1+|x|^{2}\right)^{k}(\partial / \partial x)^{\alpha} A u(x) \\
= & \int e^{i\langle x, \xi\rangle}\left(1-\Delta_{\xi}\right)^{k}\left\{\sum_{\beta \leqslant \infty} \frac{\alpha!}{\beta!(\alpha-\beta)!}(\partial / \partial x)^{\alpha-\beta} a(x, \xi)(i \xi)^{\beta} \hat{u}(\xi)\right\} d \xi,
\end{aligned}
$$

where the integrand can be estimated by

$$
(1+|x|)^{l+2 \tau k}(1+|\xi|)^{-n-1}|\hat{u}|_{|\alpha|+[m]+n+2 k+2, S}
$$

Noting that for every integer $j \geqslant 0$

$$
|\hat{u}|_{j, S} \leqslant C_{j}|u|_{j+n+1, \mathcal{S}}
$$

with a constant $C_{j}$, we obtain (4.1). Since $\tau<1,-l+2(1-\tau) k \rightarrow \infty$ as $k \rightarrow \infty$. Hence (4.1) means that $A u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$ and $A$ is a continuous map of $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ into $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$.
Q.E.D.

Lemma 4.2. Under hypothesis ( $P$. 1) we have

$$
\begin{equation*}
\left(e^{-i t H_{0} u}\right)(x)=\int e^{i(\langle x, \xi\rangle-t p(\xi))} \hat{u}(\xi) d \xi \tag{4.2}
\end{equation*}
$$

for all $u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$.
Proof. Put $U_{t}=\overline{\mathscr{F}} e^{-i t p(\xi)} \mathcal{F}$. Then $U_{t},-\infty<t<\infty$, is a group of unitary operators in $L^{2}\left(\boldsymbol{R}^{n}\right)$. Let $K$ be the infinitesimal generator of $U_{t}$. For $u \in \mathscr{D}\left(H_{0}\right)$

$$
\left\|t^{-1}\left(U_{t}-I\right) u-\left(-i H_{0}\right) u\right\|^{2}=\int\left|t^{-1}\left(e^{-i t p(\xi)}-1\right)-(-i p(\xi))\right|^{2}|\hat{u}(\xi)|^{2} d \xi \rightarrow 0
$$

as $t \rightarrow{ }^{\prime} 0$, where we use the Lebesgue dominated convergence theorem. Thus we have

$$
\begin{equation*}
-i H_{0} \subset K \tag{4.3}
\end{equation*}
$$

According to the theory of semi-groups, $K$ has the resolvent $(K-\lambda I)^{-1}$ for every $\lambda>0$. Since $H_{0}$ is self-adjoint, $-i H_{0}$ also has the resolvent $\left(-i H_{0}-\lambda I\right)^{-1}$ for every $\lambda>0$. Therefore, both

$$
K-\lambda I: \mathscr{D}(K) \rightarrow L^{2}\left(\boldsymbol{R}^{n}\right)
$$

and

$$
-i H_{0}-\lambda I: \mathscr{D}\left(H_{0}\right) \rightarrow L^{2}\left(\boldsymbol{R}^{n}\right)
$$

are bijective, provided that $\lambda>0$. Combining this fact and (4.3), we see that $\mathscr{D}(K)=\mathscr{D}\left(H_{0}\right)$. Thus

$$
-i H_{0}=K
$$

which implies that $U_{t}=e^{-i t H_{0}}$.
Q.E.D.

Remark 4.3. The hypothesis ( $P .1$ ) in Lemma 4.2 is not essential; by the arguments in the proof of Lemma 4.2, one can show that $e^{-i t H_{0}}=\overline{\mathscr{H}} e^{-i t p(\xi)} \mathscr{F}$ even if $p$ is a real-valued continuous function.

Lemma 4.4. Let hypothesis (P.1) be fulfilled. If $u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$ then

$$
t \rightarrow e^{-i t H_{0} u}
$$

is a continuous function with values in $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$.
Proof. By Lemma 4.2, it is easy to see that $\left(e^{-i t H_{0}} u\right)(x)$ is a complex-valued $C^{\infty}$ function. Integrating by parts after differentiating under the integral sign, we have

$$
\begin{equation*}
\left(1+|x|^{2}\right)^{k}(\partial / \partial x)^{\alpha}\left(e^{-i t H_{0}} u\right)(x)=\int e^{i\langle x, \xi\rangle}\left(1-\Delta_{\xi}\right)^{k}\left\{e^{-i t p(\xi)}(i \xi)^{\alpha} \hat{u}(\xi)\right\} d \xi \tag{4.4}
\end{equation*}
$$

for every multi-index $\alpha$ and every integer $k \geqslant 0$. The integrand is integrable by ( $P .1$ ), thus

$$
\sup _{x}\left(1+|x|^{2}\right)^{k}\left|(\partial / \partial x)^{\infty}\left(e^{-i t H_{0}} u\right)(x)\right|<\infty .
$$

Thus $e^{-i t H_{0}} u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$. It follows from (4.4) that for every integer $N \geqslant 0$

$$
\begin{align*}
& \left|e^{-i t H_{0}} u-e^{-i s H_{0}} u\right|_{N, S}  \tag{4.5}\\
\leqslant & C_{N_{k}} \max _{k+|\alpha| \leqslant N} \int\left|\left(1-\Delta_{\xi}\right)^{k}\left\{\left(e^{-i t p(\xi)}-e^{-i s p(\xi)}\right)(i \xi)^{\alpha} \hat{u}(\xi)\right\}\right| d \xi
\end{align*}
$$

with a constant $C_{N}$. Letting $|t-s| \rightarrow 0$, we see that the right side of (4.5) tends to zero, by the Lebesgue dominated convergence theorem, so that the
left side of (4.5) tends to zero, which proves the lemma.
Q.E.D.

Lemma 4.5. Let hypotheses $(A),(P .1)$ and $(H)$ be fulfilled. If $u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$, then

$$
t \rightarrow e^{i t H} e^{-i t H_{0}} u
$$

is a $C^{1}$ function with values in $L^{2}\left(\boldsymbol{R}^{n}\right)$ and

$$
\begin{equation*}
\frac{d}{d t}\left(e^{i t H} e^{-i t H_{0}} u\right)=i e^{i t H} A e^{-i t H_{0}} u \tag{4.6}
\end{equation*}
$$

Proof. First note that $\mathcal{S}\left(\boldsymbol{R}^{n}\right) \subset \mathscr{D}\left(H_{0}\right)$ and $\mathcal{S}\left(\boldsymbol{R}^{n}\right) \subset \mathscr{D}(H)$. Since, by Lemma 4.4, $e^{-i t H_{0}} u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$, we have

$$
\begin{equation*}
H e^{-i t H_{0}} u-H_{0} e^{-i t H_{0}} u=A e^{-i t H_{0}} u \tag{4.7}
\end{equation*}
$$

Using (4.7), we write

$$
\begin{align*}
& h^{-1}\left(e^{i(t+h) H} e^{-i(t+h) H_{0}} u-e^{i t H_{H}} e^{-i t H_{0}} u\right)-i e^{i t H} A e^{-i t H_{0}} u  \tag{4.8}\\
= & e^{i(t+h) H}\left\{h^{-1}\left(e^{-i(t+h) H_{0}}-e^{-i t H_{0}}\right) u-\left(-i H_{0}\right) e^{-i t H_{0}} u\right\} \\
& +\left(e^{i(t+h) H}-e^{i t H}\right)\left(-i H_{0}\right) e^{-i t H_{0}} u \\
& +\left\{h^{-1}\left(e^{i(t+h) H}-e^{i t H}\right) e^{-i t H_{0}} u-(i H) e^{i t H} e^{-i t H_{0}} u\right\} .
\end{align*}
$$

By the unitarity of $e^{i t H}$ and the fact that $u \in \mathscr{D}\left(H_{0}\right)$, the first term in the right side of (4.8) tends to zero in the $L^{2}$ norm as $h \rightarrow 0$. Note that for every $v \in L^{2}\left(\boldsymbol{R}^{n}\right)$

$$
t \rightarrow e^{i t H} v
$$

is a continuous function with values in $L^{2}\left(\boldsymbol{R}^{n}\right)$. Then it is clear that the second term tends to zero in the $L^{2}$ norm as $h \rightarrow 0$. The third term goes to zero in the $L^{2}$ norm as $h \rightarrow 0$, since $e^{-i t H_{0}} u \in \mathscr{D}(H)$. Thus we obtain (4.6). It follows immediately from Lemmas 4.1 and 4.4 that

$$
t \rightarrow A e^{-i t H_{0}} u
$$

is a continuous function with values in $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$, so that it is a continuous function with values in $L^{2}\left(\boldsymbol{R}^{n}\right)$. Thus,

$$
t \rightarrow i e^{i t H} A e^{-i t H_{0}} u
$$

is an $L^{2}$-valued continuous function.
Q.E.D.

The following is a key to the proof of Theorem 2. The proof will be given in Section 5, as it is fairly long.

Lemma 4.6. Let all of the hypotheses of Theorem 2 be fulfilled. Then for every non-empty $\Omega_{k}$, we have

$$
\begin{equation*}
\int_{-\infty}^{\infty}\left\|A e^{-i t H_{0}} u\right\| d t<\infty, \quad \hat{u} \in C_{0}^{\infty}\left(\Omega_{k}\right) \tag{4.9}
\end{equation*}
$$

Proof of Theorem 2. First let $\mathcal{O}$ be the subspace of $L^{2}\left(\boldsymbol{R}_{\xi}^{n}\right)$ defined by

$$
\mathcal{q}=C_{1}^{\infty}\left(\Omega_{1}\right) \oplus \cdots \oplus C_{0}^{\infty}\left(\Omega_{n}\right) .
$$

If $\Omega_{k}$ is empty, we interpret $C_{0}^{\infty}\left(\Omega_{k}\right)$ as $\{0\}$, i.e., the space consisting of the only function which is identically zero in $\Omega_{k}$. Since

$$
\bigcup_{k=1}^{n} \Omega_{k}=\boldsymbol{R}^{n} \backslash(\Xi \cup \Sigma)
$$

and since

$$
\operatorname{meas}_{n}\left(\left(\boldsymbol{R}^{n} \backslash \Sigma\right) \backslash \bigcup_{k=1}^{n} \Omega_{k}\right)=0
$$

$q \mathcal{i s}$ dense in the subspace

$$
\left\{\hat{u} \in L^{2}\left(\boldsymbol{R}^{n}\right) \mid \hat{u}(\xi)=0 \text { for almost every } \xi \in \Sigma\right\}
$$

of $L^{2}\left(\boldsymbol{R}_{\xi}^{n}\right)$. Writing $\left.\mathscr{X}=\overline{\mathscr{F}} Q\right\}$, it follows from Theorem 1 that $\mathscr{X}$ is dense in $\mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)$.

Now, let $u \in \mathscr{X}$. Then, by Lemmas 4.5 and 4.6,

$$
\int_{-\infty}^{\infty}\left\|\frac{d}{d t}\left(e^{i t H} e^{-i t H_{0}} u\right)\right\| d t<\infty .
$$

This implies that $e^{i t H} e^{-i t H_{0}} u$ converges in the $L^{2}$ norm (i.e., converges strongly) as $t \rightarrow \pm \infty$. Since $\mathscr{X}$ is dense in $\mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)$ and since $\left\|e^{i t H} e^{-i t H_{0}}\right\|=1$, we see that the limits

$$
s-\lim _{t \rightarrow \pm \infty} e^{i t H} e^{-i t H_{0}} u
$$

exist for every $u \in \mathscr{H}_{\mathrm{ac}}\left(H_{0}\right)$. Thus, we have proved the existence of the wave operators (1.2).
Q.E.D.

## 5. Proof of the key lemma 4.6

We now turn to the proof of Lemma 4.6. We shall divide the proof into two cases: (i) $k=n$ (ii) $k<n$.

Proof of Lemma 4.6. It suffices to show that to each $\xi_{0} \in \Omega_{k}$ there corresponds a neighborhood $U$ of $\xi_{0}$ such that

$$
\begin{equation*}
\int_{1 \leq|t|<\infty}\left\|A e^{-i t H_{0}} u\right\| d t<\infty, \quad \hat{u} \in C_{0}^{\infty}(U) \tag{5.1}
\end{equation*}
$$

In fact, by a partition of unity we then obtain (4.9).
(i) Assume that $\xi_{0} \in \Omega_{n}$. Then $\operatorname{det} p^{\prime \prime}\left(\xi_{0}\right) \neq 0$ which means that the Jacobian of the map

$$
p^{\prime}: \xi \rightarrow p^{\prime}(\xi)
$$

is non-zero at $\xi_{0}$. By the inverse function theorem there exists a neighborhood $U$ of $\xi_{0}$ such that $p^{\prime}$ is a diffeomorphism between $U$ and $p^{\prime}(U)$. Here we may assume that $U \subset \Omega_{n}$ and

$$
p^{\prime}(U) \subset\{r \leqslant|y| \leqslant 2 r\}
$$

for some positive number $r$.
Now, let $\hat{u} \in C_{0}^{\infty}(U)$. Then, by Lemma 4.2,

$$
\begin{equation*}
\left(A e^{-i t H_{0}} u\right)(t y)=\int e^{i t(\langle y, \xi\rangle-p(\xi))} a(t y, \xi) \hat{u}(\xi) d \xi \tag{5.2}
\end{equation*}
$$

Let $y_{*} \in p^{\prime}(\operatorname{supp} \hat{u})$. Then we can find a unique critical point $\eta_{*} \in U$ of the phase function

$$
\xi \rightarrow\left\langle y_{*}, \xi\right\rangle-p(\xi)
$$

with $\operatorname{det} p^{\prime \prime}\left(\eta_{*}\right) \neq 0$. By Lemma II of Appendix and the remarks following it, there is a sufficiently small number $\varepsilon_{*}>0$ such that for $y$ with $\left|y-y_{*}\right|<\varepsilon_{*}$

$$
\begin{align*}
& \mid \int e^{i t(\langle y, \xi\rangle-p(\xi))} a(t y, \xi) \hat{u}(\xi) d \xi  \tag{5.3}\\
& \quad-\left|\operatorname{det} p^{\prime \prime}(\eta) / 2 \pi\right|^{-1 / 2} e^{\pi i \sigma / 4} e^{i t p(\eta)} L_{y} a(t y, \eta) u(\eta)|t|^{-n / 2} \mid \\
& \quad \leqslant \\
& \quad C|a(t y, \cdot) \hat{u}|_{[n / 2]+3}|t|^{-n / 2-1}, \quad t \geqslant 1
\end{align*}
$$

where $L_{y}$ is a constant depending on $y$ and $\sigma=\operatorname{sgn} p^{\prime \prime}\left(\eta_{*}\right)$ (i.e., the signature of the symmetric matrix $\left.p^{\prime \prime}\left(\eta_{*}\right)\right), y=p^{\prime}(\eta)(\eta \in U)$. From (5.2) and (5.3) we have
when $\left|y-y_{*}\right|<\varepsilon_{*},|t| \geqslant 1$. Since $p^{\prime}(\operatorname{supp} \hat{u})$ is compact, we can choose a neighborhood $V$ of $p^{\prime}(\operatorname{supp} \hat{u})$ so that $\bar{V} \subset p^{\prime}(U)$ and so that (5.4) remains valid for $y \in V$. Furthermore, we choose a compact set $K \subset U$ so that $p^{\prime}(K)=\bar{V}$. Thus we have

$$
\begin{align*}
& \int_{1<|t|<\infty}\left(\int_{y \in V}\left|A e^{-i t H_{0}} u(t y)\right|^{2}|t|^{n} d y\right)^{1 / 2} d t  \tag{5.5}\\
& \quad \leqslant C \int_{1<|t|<\infty}\left[\int_{r \leqslant|y|<2 r}\left(1+|t| d_{K}(y)\right)^{-N_{n}}\right. \\
& \left.\quad \times\left\{\max _{|\alpha|<2 N_{n}}\left(\sup _{\xi \in K}\left|(\partial / \partial \xi)^{\infty} a(t y, \xi)\right|\right)^{2}\right\} d y\right]^{1 / 2} d t .
\end{align*}
$$

Here we have used the facts that $d_{K}(y)=0$ when $y \in V$ and

$$
V \subset\{r \leqslant|y| \leqslant 2 r\} .
$$

By condition $\left(c_{n}\right)$ the right side of (5.5) is convergent. Thus we have shown that

$$
\begin{equation*}
\int_{1<|t|<\infty}\left(\int_{x / t \in V}\left|A e^{-i t H_{0}} u(x)\right|^{2} d x\right)^{1 / 2} d t<\infty \tag{5.6}
\end{equation*}
$$

Now, we shall show that

$$
\begin{equation*}
\int_{1<|t|<\infty}\left(\int_{x / t \pm V}\left|A e^{-i t H_{0}} u(x)\right|^{2} d x\right)^{1 / 2} d t<\infty \tag{5.7}
\end{equation*}
$$

To do so we shall apply Lemma A. 1 of Hörmander [6] to the integral

$$
\left(A e^{-i t H_{0}} u\right)(x)=\int e^{i(|x|+|t|)\langle\langle x, \xi\rangle-t p(\xi)) /(|x|+|t|)} a(x, \xi) \hat{u}(\xi) d \xi
$$

Let $\Phi$ be the set of phase functions given by

$$
\Phi=\overline{\left\{\xi \rightarrow(\langle x, \xi\rangle-t p(\xi)) /(|x|+|t|) \mid x / t \in \boldsymbol{R}^{n} \backslash V\right\}} \quad \text { in } C^{N+1}\left(\boldsymbol{R}^{n}, \boldsymbol{R}\right)
$$

Then $\Phi$ is a compact subset of $C^{N+1}\left(\boldsymbol{R}^{n}, \boldsymbol{R}\right)$ such that for every $f \in \Phi$

$$
f^{\prime}(\xi) \neq 0 \quad \text { if } \xi \in \operatorname{supp} \hat{u} .
$$

In fact, writing

$$
Z=\overline{\left\{(z, s) \in \boldsymbol{R}^{n+1} \mid(z, s)=(x, t) /(|x|+|t|) \text { for some } x / t \in \boldsymbol{R}^{n} \backslash V\right\}}
$$

we have

$$
\Phi=\{\xi \rightarrow\langle z, \xi\rangle-s p(\xi) \mid(z, s) \in Z\}
$$

It is easily seen that

$$
(z, s) \rightarrow\langle z, \xi\rangle-s p(\xi) \in C^{N+1}\left(\boldsymbol{R}^{n}, \boldsymbol{R}\right)
$$

is continuous. Since $Z$ is compact, it follows immediately that $\Phi$ is a compact subset of $C^{N+1}\left(\boldsymbol{R}^{n}, \boldsymbol{R}\right)$. Thus Lemma A. 1 of [6] shows that for every $N$

$$
\begin{equation*}
\left|\left(A e^{-i t H_{0}} u\right)(x)\right| \leqslant C_{N}(1+|x|+|t|)^{-N}|a(x, \cdot) \hat{u}|_{N} \tag{5.8}
\end{equation*}
$$

if $x / t \in \boldsymbol{R}^{n} \backslash V$. By Assumption (2.1) ( $A$ ) the right side of (5.8) is bounded by

$$
C_{N}^{\prime}(1+|x|+|t|)^{l+(\tau-1) N} .
$$

Thus we see that

$$
\int_{x / t \in V}\left|A e^{-i t H_{0}} u(x)\right|^{2} d x \leqslant C_{N}^{\prime}\left(\int(1+|x|)^{-n-1} d x\right)(1+|t|)^{n+1+2 l-2(1-\tau) N}
$$

for every integer $N$ with $(n+1) / 2+l-(1-\tau) N<-1$. This proves (5.7). Thus we have proved (5.1) in the case $k=n$.
(ii) Assume that $\xi_{0} \in \Omega_{k}$. Since the rank of the symmetric matrix $p^{\prime \prime}\left(\xi_{0}\right)$ is $k$, we can find a principal minor of order $k$ which does not vanish. So we may assume that

$$
\operatorname{det}\left(\partial^{2} p\left(\xi_{0}\right) / \partial \xi_{i} \partial \xi_{j}\right)_{i, j=1, \cdots, k} \neq 0
$$

Let

$$
\begin{aligned}
y_{0}^{\prime} & =p_{\xi^{\prime}}\left(\xi_{0}\right) \\
& \equiv\left(\partial p\left(\xi_{0}\right) / \partial \xi_{1}, \cdots, \partial p\left(\xi_{0}\right) / \partial \xi_{k}\right) .
\end{aligned}
$$

Then, by the implicit function theorem, there is a neighborhood $W$ of $y_{0}^{\prime}$, a neighborhood $U^{\prime}$ of $\xi_{0}^{\prime}$, a neighborhood $U^{\prime \prime}$ of $\xi_{0}^{\prime \prime}$ and a unique function

$$
\varphi=\left(\varphi_{1}, \cdots, \varphi_{k}\right) \in C^{\infty}\left(W \times U^{\prime \prime}\right)
$$

such that

$$
\begin{aligned}
& \left\{\left(y^{\prime}, \xi\right) \in W \times U^{\prime} \times U^{\prime \prime} \mid y^{\prime}=p_{\xi^{\prime}}(\xi)\right\} \\
= & \left\{\left(y^{\prime}, \xi\right) \in \boldsymbol{R}^{k} \times \boldsymbol{R}^{n} \mid \xi^{\prime}=\varphi\left(y^{\prime}, \xi^{\prime \prime}\right) \text { for some }\left(y^{\prime}, \xi^{\prime \prime}\right) \in W \times U^{\prime \prime}\right\} .
\end{aligned}
$$

Here $\xi^{\prime}=\left(\xi_{1}, \cdots, \xi_{k}\right), \xi^{\prime \prime}=\left(\xi_{k+1}, \cdots, \xi_{n}\right)$ and $y^{\prime}=\left(y_{1}, \cdots, y_{k}\right)$. We now set

$$
U=\left\{\xi \in \boldsymbol{R}^{n} \mid \xi^{\prime}=\varphi\left(y^{\prime}, \xi^{\prime \prime}\right) \text { for some }\left(y^{\prime}, \xi^{\prime \prime}\right) \in W \times U^{\prime \prime}\right\}
$$

Obviously $U$ is a neighborhood of $\xi_{0}$. We may assume that $U \subset \Omega_{k}$ and that

$$
\begin{equation*}
p^{\prime}(U) \subset\left\{y \in \boldsymbol{R}^{n}|5 r / 4 \leqslant|y| \leqslant 7 r / 4\}\right. \tag{5.9}
\end{equation*}
$$

for some $r>0$. Using Assumption (2.1) (P.2) (b) we shall show that $p_{\xi^{\prime \prime}}\left(\varphi\left(y^{\prime}\right.\right.$, $\left.\xi^{\prime \prime}\right), \xi^{\prime \prime}$ ) is independent of $\xi^{\prime \prime}$. Put

$$
g_{j}\left(y^{\prime}, \xi^{\prime \prime}\right)=\partial p\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right) / \partial \xi_{j} \quad(j=k+1, \cdots, n)
$$

Then we have the differential

$$
\begin{equation*}
d g_{j}=\sum_{i=1}^{k}\left(\partial^{2} p / \partial \xi_{i} \partial \xi_{j}\right) d \varphi_{i}+\sum_{i=k+1}^{n}\left(\partial^{2} p / \partial \xi_{i} \partial \xi_{j}\right) d \xi_{i} \tag{5.10}
\end{equation*}
$$

On the other hand, since $y_{j}=\partial p\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right) / \partial \xi_{j}(j=1, \cdots, k)$ we have

$$
\begin{equation*}
d y_{j}=\sum_{i=1}^{k}\left(\partial^{2} p / \partial \xi_{i} \partial \xi_{j}\right) d \varphi_{i}+\sum_{i=k+1}^{n}\left(\partial^{2} p / \partial \xi_{i} \partial \xi_{j}\right) d \xi_{i} \tag{5.11}
\end{equation*}
$$

Since the rank of $p^{\prime \prime}\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right)$ is $k$ and since

$$
\operatorname{det} p_{\xi^{\prime} \xi^{\prime}}\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right) \neq 0
$$

(5.10) and (5.11) imply that

$$
d g_{j}=\sum_{i=1}^{k} c_{i i}\left(y^{\prime}, \xi^{\prime \prime}\right) d y_{i} \quad(j=k+1, \cdots, n)
$$

where $c_{j i}\left(y^{\prime}, \xi^{\prime \prime}\right)$ are suitable functions. Hence

$$
\partial g_{j}\left(y^{\prime}, \xi^{\prime \prime}\right) / \partial \xi_{i}=0 \quad(i, j=k+1, \cdots, n)
$$

which means that $p_{\xi^{\prime \prime}}\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right)$ is independent of $\xi^{\prime \prime}$. Hence there is a function $g \in C^{\infty}(W)$ such that

$$
p^{\prime}(U)=\left\{y \in \boldsymbol{R}^{n} \mid y^{\prime} \in W, y^{\prime \prime}=g\left(y^{\prime}\right)\right\}
$$

Let us now estimate the $L^{2}$ norm of $A e^{-i t H_{0}} u$. Let $u \in C_{0}^{\infty}(U)$. We write

$$
A e^{-i t H_{0}} u(t y)=\int d \xi^{\prime \prime} \int e^{i t\langle(\langle, \xi\rangle-p(\xi))} a(t y, \xi) \hat{u}(\xi) d \xi^{\prime}
$$

Set

$$
\Gamma=\left\{\xi^{\prime \prime} \in \boldsymbol{R}^{n-k} \mid\left(\xi^{\prime}, \xi^{\prime \prime}\right) \in \operatorname{supp} \hat{u} \text { for some } \xi^{\prime} \in \boldsymbol{R}^{k}\right\}
$$

Let $y_{*} \in p^{\prime}(\operatorname{supp} \hat{u})$ and let $\xi_{*}^{\prime \prime} \in \Gamma$. We consider the phase function

$$
\xi^{\prime} \rightarrow\left\langle y_{*}^{\prime}, \xi^{\prime}\right\rangle+\left\langle y_{*}^{\prime \prime}, \xi_{*}^{\prime \prime}\right\rangle-p\left(\xi^{\prime}, \xi_{*}^{\prime \prime}\right) .
$$

Since $y_{*}^{\prime}=p_{\xi^{\prime}}\left(\varphi\left(y_{*}^{\prime}, \xi_{*}^{\prime \prime}\right), \xi_{*}^{\prime \prime}\right)$, this phase function has a unique critical point $\varphi\left(y_{*}^{\prime}, \xi_{*}^{\prime \prime}\right) \in U^{\prime}$. By Lemma II and the remarks following it, for every integer $N>0$ there is a positive number $\varepsilon_{*}$ such that when $\left|y-y_{*}\right|<\varepsilon_{*}$ and $\left|\xi^{\prime \prime}-\xi_{*}^{\prime \prime}\right|<\varepsilon_{*}$,

$$
\begin{align*}
& \left|\int e^{i t(\langle y, \xi\rangle-p(\xi))} a(t y, \xi) \hat{u}(\xi) d \xi^{\prime}-\sum_{j=0}^{N-1}\right| \operatorname{det} p_{\xi^{\prime} \xi^{\prime}}\left(\eta^{\prime}, \xi^{\prime \prime}\right) /\left.2 \pi\right|^{-1 / 2}  \tag{5.12}\\
& \quad \times\left. e^{\pi i \sigma / 4} e^{i t\left(\left\langle y^{\prime}, \eta^{\prime}\right\rangle+\left\langle y^{\prime \prime}, \xi^{\prime \prime}\right\rangle-p\left(\eta^{\prime}, \xi^{\prime \prime}\right)\right)} L_{y \xi^{\prime} j}(a(t y, \xi) \hat{u}(\xi))\right|_{\xi^{\prime}=\eta^{\prime}}|t|^{-j-k / 2} \mid \\
& \leqslant
\end{align*}
$$

Here $\eta^{\prime}=\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), s=2 N+[k / 2]+1$ and $\sigma=\operatorname{sgn} p_{\xi^{\prime}} \xi^{\prime}\left(\varphi\left(y_{*}^{\prime}, \xi_{*}^{\prime \prime}\right), \xi_{*}^{\prime \prime}\right)$. Differential operators $L_{y \xi^{\prime \prime} j}$ depend on $y, \xi^{\prime \prime}$ and by virtue of Lemma A. 5 of [6] their coefficients are $C^{\infty}$ functions of $y, \xi^{\prime \prime}$. Since $p^{\prime}(\operatorname{supp} \hat{u})$ and $\Gamma$ are compact, we can find a neighborhood $V$ of $p^{\prime}(\operatorname{supp} \hat{u})$ such that (5.12) remains valid for $y \in V$ and $\xi^{\prime \prime} \in \Gamma$. Here we may assume that

$$
p^{\prime}(\operatorname{supp} \hat{u}) \subset V \cap p^{\prime}(U)
$$

and

$$
\begin{equation*}
V \subset\left\{y \in \boldsymbol{R}^{n}|r \leqslant|y| \leqslant 2 r\}\right. \tag{5.13}
\end{equation*}
$$

with the same $r$ as in (5.9). At the critical point $\eta^{\prime}=\varphi\left(y^{\prime}, \xi^{\prime \prime}\right)$ the phase

$$
\xi^{\prime} \rightarrow\left\langle y^{\prime}, \xi^{\prime}\right\rangle+\left\langle y^{\prime \prime}, \xi^{\prime \prime}\right\rangle-p\left(\xi^{\prime}, \xi^{\prime \prime}\right)
$$

is of the form

$$
f(y)+\left\langle y^{\prime \prime}-g\left(y^{\prime}\right), \xi^{\prime \prime}\right\rangle .
$$

Indeed, writing

$$
\tilde{f}\left(y, \xi^{\prime \prime}\right)=\left\langle y^{\prime}, \varphi\left(y^{\prime}, \xi^{\prime \prime}\right)\right\rangle+\left\langle y^{\prime \prime}, \xi^{\prime \prime}\right\rangle-p\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right)
$$

and using the fact that $y^{\prime}=p_{\xi^{\prime}}\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right)$, we see that

$$
d \tilde{f}=\left\langle d y^{\prime}, \varphi\right\rangle+\left\langle d y^{\prime \prime}, \xi^{\prime \prime}\right\rangle+\left\langle y^{\prime \prime}-g\left(y^{\prime}\right), d \xi^{\prime \prime}\right\rangle .
$$

Therefore $\tilde{f}_{\xi^{\prime \prime}}=y^{\prime \prime}-g\left(y^{\prime}\right)$.
Summing up, we have for every $N>0$ and every $y \in V$

$$
\begin{align*}
& \left|A e^{-i t H_{0}} u(t y)\right|  \tag{5.14}\\
\leqslant & \left.C \sum_{j=0}^{N-1}\left|\int e^{i t\left\langle y^{\prime \prime}-g\left(y^{\prime}\right), \xi^{\prime \prime}\right\rangle}\right| \operatorname{det} p_{\xi^{\prime} \xi^{\prime}}\left(\varphi\left(y^{\prime}, \xi^{\prime \prime}\right), \xi^{\prime \prime}\right)\right|^{-1 / 2} \\
& \times\left. L_{y \xi^{\prime} \xi^{\prime \prime} j}(a(t y, \xi) \hat{u}(\xi))\right|_{\xi^{\prime}=\varphi\left(y^{\prime}, \xi^{\prime \prime}\right)} d \xi^{\prime \prime}\left|\times|t|^{-j-k / 2}\right. \\
& +C_{N}|t|^{-k / 2-N} \int\left|a\left(t y, \cdot, \xi^{\prime \prime}\right) \hat{u}\left(\cdot, \xi^{\prime \prime}\right)\right|_{s} d \xi^{\prime \prime} \\
\equiv & C \sum_{j=0}^{N-1} I_{j}|t|^{-j-k / 2}+C_{N} J|t|^{-k / 2-N}
\end{align*}
$$

Applying Lemma I to $I$, we see that for every integer $L$

$$
I_{j} \leqslant C_{L}\left(1+|t|\left|y^{\prime \prime}-g\left(y^{\prime}\right)\right|\right)^{-L} \max _{|\propto|<2 j+L}\left(\sup _{\xi \in \operatorname{supp} \hat{u}}\left|(\partial / \partial \xi)^{\infty} a(t y, \xi)\right|\right)
$$

It is obvious that

$$
J \leqslant C_{N}^{\prime} \max _{|\alpha|<s}\left(\sup _{\xi \in \operatorname{supp} \hat{u}}\left|(\partial / \partial \xi)^{\infty} a(t y, \xi)\right|\right)
$$

Choose a compact set $K$ so that supp $\hat{u} \subset K$ and so that $V \cap p^{\prime}(U) \subset p^{\prime}(K)$. Then

$$
1+|t| d_{K}(y) \leqslant 1+|t|\left|y^{\prime \prime}-g\left(y^{\prime}\right)\right|, \quad y \in V
$$

and

$$
1+|t| d_{K}(y) \leqslant C|t|, \quad y \in V,|t| \geqslant 1
$$

Thus we have

$$
\begin{align*}
& \int_{1<|t|<\infty}\left(\int_{x / t \in V}\left|A e^{-i t H_{0}} u(x)\right|^{2} d x\right)^{1 / 2} d t  \tag{5.15}\\
\leqslant & C_{N, L} \int_{1<|t|<\infty}|t|^{(n-k) / 2}\left[\int_{r<|y|<2 r}\left(1+|t| d_{K}(y)\right)^{-2 L}\right. \\
& \left.\times \max _{|\infty| \leqslant 2 N+L-2}\left(\sup _{\xi \in K}\left|(\partial / \partial \xi)^{\infty} a(t y, \xi)\right|\right)^{2} d y\right]^{1 / 2} d t
\end{align*}
$$

$$
\begin{aligned}
& +C_{N} \int_{1<|t|<\infty}|t|^{(n-k) / 2}\left[\int_{r<|y|<2 r}\left(1+|t| d_{K}(y)\right)^{-2 N}\right. \\
& \left.\times \max _{|\alpha| \leqslant s}\left(\sup _{\xi \in K}\left|(\partial / \partial \xi)^{\infty} a(t y, \xi)\right|\right)^{2} d y\right]^{1 / 2} d t
\end{aligned}
$$

where we use (5.13). If we choose $N, L$ so that $N_{k} \leqslant 2 N=2 L \leqslant N_{k}+1$, then $2 N+L-2 \leqslant 2 N_{k}$ and $s \leqslant 2 N_{k}$. Thus, by condition ( $c_{k}$ ), the integrals in the right side of (5.15) are convergent.

Finally, (5.7) is shown in quite the same way as in (i). Thus we have proved (5.1) in the case $k<n$.
Q.E.D.

## 6. Symmetry of $\boldsymbol{A}(\boldsymbol{X}, \boldsymbol{D})$

In this section we shall determine completely when the pseudo-differential operator $A=A(X, D)$ is symmetric on $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$. To avoid confusion, we denote $A \upharpoonright \mathcal{S}\left(\boldsymbol{R}^{n}\right)$ by $A_{0}$.

Lemma 6.1. Let hypothesis $(A)$ be fulfilled. Then the domain $\mathscr{D}\left(A_{0}^{*}\right)$ of the adjoint $A_{0}^{*}$ contains $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ and

$$
\left(A_{0}^{*} v\right)(x)=\int e^{i\langle x, \xi\rangle}\left(\int e^{-i\left\langle x^{\prime}, \xi\right\rangle} \bar{a}\left(x^{\prime}, \xi\right) v\left(x^{\prime}\right) d x^{\prime}\right) d \xi, \quad v \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)
$$

Here $\bar{a}(x, \xi)=\overline{a(x, \xi)}$.
Proof. See [8], p. 61.
Q.E.D.

We wish to get the representation of the form

$$
\left(A_{0}^{*} v\right)(x)=\int e^{i\langle x, \xi\rangle} b(x, \xi) \hat{v}(\xi) d \xi
$$

Following ideas in [8], we shall consider the symbol $b(x, \xi)$ defined by the oscillatory integral

$$
\begin{equation*}
b(x, \xi)=\mathrm{O}_{\mathrm{s}}-\iint e^{-i\langle y, \eta\rangle} \bar{a}(x+y, \xi+\eta) d y d \eta \tag{6.1}
\end{equation*}
$$

For oscillatory integrals we refer to [8, Chapter I].
Lemma 6.2. Let hypothesis $(A)$ be fulfilled and let $b(x, \xi)$ be as in (6.1). Then $b \in C^{\infty}\left(\boldsymbol{R}^{n} \times \boldsymbol{R}^{n}\right)$ and for all multi-indices $\alpha, \beta$ the estimate

$$
\begin{align*}
& \left|(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} b(x, \xi)\right|  \tag{6.2}\\
\leqslant & C_{\alpha \beta}(1+|x|)^{l+(|\alpha|+2 l(\alpha))}(1+|\xi|)^{m+\delta(|\beta|+2 m(\beta))}
\end{align*}
$$

is valid for some constant $C_{\alpha \beta}$, where $l(\alpha)=[(l+n+\tau|\alpha|) / 2(1-\tau)]+1, m(\beta)=$ $[(m+n+\delta|\beta|) / 2(1-\delta)]+1$.

Proof. First put

$$
\begin{equation*}
b_{\varepsilon}(x, \xi)=\iint e^{-i\langle y, \eta\rangle} \chi\left(\varepsilon y, \varepsilon_{\eta}\right) \bar{a}(x+y, \xi+\eta) d y d \eta, \quad \varepsilon>0 \tag{6.3}
\end{equation*}
$$

where $\chi(y, \eta) \in \mathcal{S}\left(\boldsymbol{R}^{n} \times \boldsymbol{R}^{n}\right)$ and $\chi(0,0)=1$. Since we can differentiate under the sign of integration as often as we like, we conclude that $b_{\mathrm{z}} \in C^{\infty}$. Furthermore, by repeated integration by parts, we have

$$
\begin{align*}
& (\partial / \partial \xi)^{\omega}(\partial / \partial x)^{\beta} b_{\varepsilon}(x, \xi)  \tag{6.4}\\
= & \iint e^{-i\langle y, \eta\rangle}\left(1+|y|^{2}\right)^{-l(\alpha)}\left(1-\Delta_{\eta}\right)^{l(\alpha)}\left\{\left(1+|\eta|^{2}\right)^{-m(\beta)}\right. \\
& \left.\times\left(1-\Delta_{y}\right)^{m(\beta)}\left(\chi\left(\varepsilon y, \varepsilon_{\eta}\right) a_{(\beta)}\right)^{(\alpha)}(x+y, \xi+\eta)\right) d y d \eta \\
\equiv & \iint e^{-i\langle y, \eta\rangle} I_{\varepsilon}(x, \xi, y, \eta) d y d_{\eta}
\end{align*}
$$

with the notation $a_{(\beta)^{(\alpha)}}=(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} a$. By using hypothesis $(A)$ we have

$$
\begin{align*}
& \left|I_{\varepsilon}\right| \leqslant C_{\alpha \beta}^{\prime}(1+|x|)^{l+\tau|\alpha|+2 \tau l(\alpha)}(1+|\xi|)^{m+\delta|\beta|+2 \delta m(\beta)}  \tag{6.5}\\
& \quad \times(1+|y|)^{l+\tau|\alpha|+2(\tau-1) l(\alpha)}(1+|\eta|)^{m+\delta|\beta|+2(\delta-1) m(\beta)}
\end{align*}
$$

where $C_{\alpha \beta}^{\prime}$ is independent of $\varepsilon$. Since $l+\tau|\alpha|+2(\tau-1) l(\alpha)<-n$ and since $m+\delta|\beta|+2(\delta-1) m(\beta)<-n$, we can derive from (6.4) and (6.5)

$$
\begin{align*}
& \left|(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} b_{\mathrm{z}}(x, \xi)\right|  \tag{6.6}\\
\leqslant & C_{\alpha \beta}^{\prime \prime \prime}(1+|x|)^{l+\tau|\alpha|+2 \tau l(\omega)}(1+|\xi|)^{m+\delta|\beta|+2 \delta m(\beta)}
\end{align*}
$$

with a constant $C_{\alpha \beta}^{\prime \prime}$ independent of $\varepsilon$.
Now, let

$$
b_{\alpha \beta}(x, \xi)=\mathrm{O}_{\mathrm{s}}-\iint e^{-i\langle y, \eta\rangle} \bar{a}_{(\beta)}{ }^{(\alpha)}(x+y, \xi+\eta) d y d \eta
$$

Then, by repeated integration by parts,

$$
\begin{align*}
b_{\alpha \beta}(x, \xi)=\iint & e^{-i\langle y, \eta\rangle}\left(1+|y|^{2}\right)^{-l(\alpha)}\left(1-\Delta_{\eta}\right)^{l(\alpha)}\left\{\left(1+|\eta|^{2}\right)^{-m(\beta)}\right.  \tag{6.7}\\
& \left.\times\left(1-\Delta_{y}\right)^{m(\beta)} \bar{a}_{(\beta)^{(\alpha)}}(x+y, \xi+\eta)\right\} d y d \eta .
\end{align*}
$$

From (6.4) and (6.7) we see that for every compact set $K \subset \boldsymbol{R}^{n} \times \boldsymbol{R}^{n}$ and all multi-indices $\alpha, \beta$

$$
\begin{align*}
& \sup _{(x, \xi) \in K}\left|b_{\alpha \beta}(x, \xi)-(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} b_{\mathrm{e}}(x, \xi)\right|  \tag{6.8}\\
\leqslant & C_{K \alpha \beta} \int\left(\max _{|\mu|+|v| \leqslant N_{\alpha \beta}}\left|(\partial / \partial \eta)^{\mu}(\partial / \partial y)^{\nu}(1-\chi(\varepsilon y, \varepsilon \eta))\right|\right) \\
& \times(1+|y|)^{l+\tau|\alpha|+2(\tau-1) l(\alpha)}(1+|\eta|)^{m+\delta|\beta|+2(\delta-1) m(\beta)} d y d \eta
\end{align*}
$$

is valid for some constant $C_{K \alpha \beta}$, where $N_{\alpha \beta}=2(l(\alpha)+m(\beta))$. Since $(\partial / \partial \eta)^{\mu} \times$ $(\partial / \partial y)^{\nu}\left(1-\chi\left(\varepsilon y, \varepsilon_{\eta}\right)\right) \rightarrow 0$ as $\varepsilon \rightarrow 0,(6.8)$ implies that for all multi-indices $\alpha, \beta$ $(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} b_{z}$ converges uniformly on every compact set to $b_{a \beta}$ as $\varepsilon \rightarrow 0$. In particular, $b_{\varepsilon}$ converges uniformly on every compact set to $b$ as $\varepsilon \rightarrow 0$. Hence $b \in C^{\infty}$ and

$$
\begin{equation*}
(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} b(x, \xi)=b_{\alpha \beta}(x, \xi) \tag{6.9}
\end{equation*}
$$

Thus (6.2) follows from (6.6) and (6.9).
Q.E.D.

The proof of the following lemma is virtually identical to some arguments of [8, chapter 2].

Lemma 6.3. Let hypothesis $(A)$ be fulfilled and let $b(x, \xi)$ be as in (6.1). Then

$$
\begin{equation*}
\left(A_{0}^{*} v\right)(x)=\int e^{i\langle x, \xi\rangle} b(x, \xi) \hat{v}(\xi) d \xi, \quad v \in \mathcal{S}\left(\boldsymbol{R}^{n}\right) \tag{6.10}
\end{equation*}
$$

Proof. Let $b_{\mathrm{e}}(x, \xi)$ be as in the proof of Lemma 6.2. Then the Fubini theorem and a change of variables give

$$
\begin{align*}
& \int e^{i\langle x, \xi\rangle} b(x, \xi) \hat{v}(\xi) d \xi  \tag{6.11}\\
= & \iiint e^{i\langle x, \xi\rangle-i\langle y, \eta\rangle} \chi(\varepsilon y, \varepsilon \eta) \bar{a}(x+y, \xi+\eta) \hat{v}(\xi) d y d \eta d \xi \\
= & \int e^{i\left\langle x, \xi^{\prime}\right\rangle}\left\{\int e^{-i\left\langle x^{\prime} \xi^{\prime}\right\rangle} a\left(x^{\prime}, \xi^{\prime}\right)\right. \\
& \left.\times\left(\int e^{i\left\langle x^{\prime}, \xi\right\rangle} \chi\left(\varepsilon\left(x-x^{\prime}\right), \varepsilon\left(\xi^{\prime}-\xi\right)\right) \hat{v}(\xi) d \xi\right) d x^{\prime}\right\} d \xi^{\prime} \\
\equiv & \int e^{i\left\langle x, \xi^{\prime}\right\rangle} I_{2}\left(x, \xi^{\prime}\right) d \xi^{\prime} .
\end{align*}
$$

By (6.6), (6.11) and the Lebesgue dominated convergence theorem

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \int e^{i\left\langle x, \xi^{\prime}\right\rangle} I_{\mathrm{e}}\left(x, \xi^{\prime}\right) d \xi^{\prime}=\int e^{i\langle x, \xi\rangle} b(x, \xi) \hat{v}(\xi) d \xi \tag{6.12}
\end{equation*}
$$

By repeated integration by parts it follows that for every integer $k>0$

$$
\begin{equation*}
\left|\int e^{i\left\langle x^{\prime}, \xi\right\rangle} \chi\left(\varepsilon\left(x^{\prime}-x\right), \varepsilon\left(\xi^{\prime}-\xi\right)\right) \hat{v}(\xi) d \xi\right| \leqslant C_{k}\left(1+\left|x^{\prime}\right|\right)^{-k} \tag{6.13}
\end{equation*}
$$

with $C_{k}$ independent of $\varepsilon$. Therefore, by the Lebesgue dominated convergence theorem,

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} I_{\varepsilon}\left(x, \xi^{\prime}\right)=\int e^{-i\left\langle x^{\prime}, \xi^{\prime}\right\rangle} \bar{a}\left(x^{\prime}, \xi^{\prime}\right) v\left(x^{\prime}\right) d x^{\prime} \tag{6.14}
\end{equation*}
$$

Here we have used Fourier's inversion formula

$$
v\left(x^{\prime}\right)=\int e^{i\left\langle x^{\prime}, \xi\right\rangle \hat{v}(\xi) d \xi}
$$

Note that for every integer $j \geqslant 0$

$$
\begin{align*}
& \quad I_{\mathrm{e}}\left(x, \xi^{\prime}\right)  \tag{6.15}\\
& =\left(1+\left|\xi^{\prime}\right|^{2}\right)^{-j} \int e^{-i\left\langle x^{\prime}, \xi^{\prime}\right\rangle}\left(1-\Delta_{x^{\prime}}\right)^{j}\left\{a\left(x^{\prime}, \xi^{\prime}\right)\right. \\
& \left.\quad \times\left(\int e^{i\left\langle x^{\prime}, \xi\right\rangle} \chi\left(\varepsilon\left(x^{\prime}-x\right), \varepsilon\left(\xi^{\prime}-\xi\right)\right) \hat{v}(\xi) d \xi\right)\right\} d x^{\prime} .
\end{align*}
$$

Since, by (6.13) and hypothesis $(A)$, the integrand in (6.15) can be estimated by

$$
\left(1+\left|\xi^{\prime}\right|\right)^{m+2 j 8}\left(1+\left|x^{\prime}\right|\right)^{l-2 k}
$$

for every integer $k \geqslant 0$, it follows that for every integer $j \geqslant 0$

$$
\left|I_{\varepsilon}\left(x, \xi^{\prime}\right)\right| \leqslant C_{j}^{\prime}\left(1+\left|\xi^{\prime}\right|\right)^{m+2(\delta-1) j}
$$

with a constant $C_{j}^{\prime}$ independent of $\varepsilon$. Choosing $j$ so that $m+2(\delta-1) j<-n$ and using the Lebesgue dominated convergence theorem, we see that

$$
\begin{align*}
& \lim _{\mathrm{g} \rightarrow \infty} \int e^{i\left\langle x, \xi^{\prime}\right\rangle} I_{\mathrm{e}}\left(x, \xi^{\prime}\right) d x^{\prime}  \tag{6.16}\\
= & \int e^{i\left\langle x, \xi^{\prime}\right\rangle}\left(\int e^{-i\left\langle x^{\prime}, \xi^{\prime}\right\rangle} \bar{a}\left(x^{\prime}, \xi^{\prime}\right) v\left(x^{\prime}\right) d x^{\prime}\right) d \xi^{\prime} .
\end{align*}
$$

(Recall (6.14).) By Lemma 6.1, the right side of (6.16) equals $\left(A_{0}^{*} v\right)(x)$. Thus, combining (6.16) with (6.12), we obtain (6.10).

Now we give a theorem.
Theorem 3. Let hypothesis $(A)$ be fulfilled. The operator $A_{0}$ is symmetric if and only if

$$
\begin{equation*}
\mathscr{F}_{x}[a](\eta, \xi)=\mathscr{F}_{x}[a](\eta, \xi+\eta) \quad \text { in } \mathcal{S}^{\prime}\left(\boldsymbol{R}_{\eta}^{n} \times \boldsymbol{R}_{\xi}^{\eta}\right) \tag{6.17}
\end{equation*}
$$

Here $\mathscr{F}_{x}[a]$ is the Fourier transform of a with respect to $x$.
Proof. Let $b$ be as in (6.1). If

$$
\begin{equation*}
\mathscr{F}_{x}[b](\eta, \xi)=\mathscr{F}_{x}[a](\eta, \xi+\eta) \quad \text { in } \mathcal{S}^{\prime}\left(\boldsymbol{R}_{\eta}^{n} \times \boldsymbol{R}_{\xi}^{n}\right), \tag{6.18}
\end{equation*}
$$

then (6.17) is equivalent to

$$
\begin{equation*}
\mathscr{F}_{x}[a](\eta, \xi)=\mathscr{F}_{x}[b](\eta, \xi) \quad \text { in } \mathcal{S}^{\prime}\left(\boldsymbol{R}_{\eta}^{n} \times \boldsymbol{R}_{\xi}^{n}\right), \tag{6.19}
\end{equation*}
$$

or what is the same thing,

$$
\begin{equation*}
a(x, \xi)=b(x, \xi) \tag{6.20}
\end{equation*}
$$

Hence, by Lemma 6.3, (6.17) is equivalent to the fact that $A_{0}$ is symmetric. Thus, it suffices to show (6.18)
(6.18) means that for every $\varphi(\eta, \xi) \in \mathcal{S}\left(\boldsymbol{R}_{\eta}^{n} \times \boldsymbol{R}_{\xi}^{n}\right)$

$$
\begin{align*}
& \iint b(x, \xi) \mathscr{F}_{\eta}[\varphi](x, \xi) d x d \xi  \tag{6.21}\\
= & \iint a(x, \xi)\left(\int e^{-i\langle x, \eta\rangle} \varphi(\eta, \xi-\eta) d \eta\right) d x d \xi .
\end{align*}
$$

Choose $\chi_{0} \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$ so that $\chi_{0}(0)=1$ and replace $\chi\left(\varepsilon y, \varepsilon_{\eta}\right)$ in (6.3) by $\chi_{0}(\varepsilon y) \chi_{0}\left(\varepsilon_{\eta}\right)$. Then $b_{\varepsilon}$ has the same properties as before. By (6.6) and the fact that $b_{\mathrm{e}} \rightarrow b$ as $\varepsilon \rightarrow 0$, the left-hand side of (6.21) is equal to the limit of

$$
\begin{equation*}
\iint b_{\varepsilon}(x, \xi) \mathscr{F}_{\eta}[\varphi](x, \xi) d x d \xi \tag{6.22}
\end{equation*}
$$

which, by the definition of $b_{\mathrm{e}}$, equals

$$
\begin{equation*}
\iiint \int e^{-i\left\langle y, \eta^{\prime}\right\rangle} \chi_{0}(\varepsilon y) \chi_{0}\left(\varepsilon_{\eta^{\prime}}\right) \bar{a}\left(x+y, \xi+\eta^{\prime}\right) \mathscr{F}_{\eta}[\rho](x, \xi) d y d \eta^{\prime} d x d \xi \tag{6.23}
\end{equation*}
$$

By a change of variables, (6.23) equals

$$
\begin{align*}
& \iint a(x, \xi)\left(\iint e^{-i<y, \eta^{\prime}>} \chi_{0}(\varepsilon y) \chi_{0}\left(\varepsilon \eta^{\prime}\right) \mathscr{F}_{\eta}[\varphi]\left(x-y, \xi-\eta^{\prime}\right) d y d_{\eta^{\prime}}\right) d x d \xi  \tag{6.24}\\
\equiv & \iint a(x, \xi) J_{\mathrm{e}}(x, \xi) d x d \xi
\end{align*}
$$

Finally we should examine $J_{\varepsilon}$. Using the Fubini theorem, we see that

$$
\begin{equation*}
J_{\mathrm{e}}(x, \xi)=\iint e^{-i\langle x, \eta\rangle} \varphi\left(\eta, \xi-\eta^{\prime}\right) \chi_{0}\left(\varepsilon \eta^{\prime}\right)\left(\int e^{-i\left\langle y, \eta^{\prime}-\eta\right\rangle} \chi_{0}(\varepsilon y) d y\right) d \eta d \eta^{\prime} \tag{6.25}
\end{equation*}
$$

Writing $y^{\prime}=\varepsilon y$, we see that the right side of (6.25) is equal to

$$
\iint e^{-i\langle x, \eta\rangle} \varphi\left(\eta, \xi-\eta^{\prime}\right) \chi_{0}\left(\varepsilon \eta^{\prime}\right) \hat{X}_{0}\left(\varepsilon^{-1}\left(\eta^{\prime}-\eta\right)\right) \varepsilon^{-n} d \eta d \eta^{\prime}
$$

Thus a change of variables gives

$$
J_{\varepsilon}(x, \xi)=\iint e^{-i\left\langle x, \eta^{\prime}-\varepsilon \theta\right\rangle} \varphi\left(\eta^{\prime}-\varepsilon \theta, \xi-\eta^{\prime}\right) \chi_{0}\left(\varepsilon \eta^{\prime}\right) \hat{\chi}_{0}(\theta) d \theta d \eta^{\prime}
$$

Hence

$$
\begin{align*}
\lim _{\varepsilon \rightarrow 0} J_{\varepsilon}(x, \xi) & =\int e^{-i\left\langle x, \eta^{\prime}\right\rangle} \varphi\left(\eta^{\prime}, \xi-\eta^{\prime}\right)\left(\int \hat{X}_{0}(\theta) d \theta\right) d \eta^{\prime}  \tag{6.26}\\
& =\int e^{-i\left\langle x, \eta^{\prime}\right\rangle} \varphi\left(\eta^{\prime}, \xi-\eta^{\prime}\right) d \eta^{\prime}
\end{align*}
$$

where we use $\chi_{0}(0)=1$. On the other hand, since $\mathscr{F}_{n}[\varphi] \in \mathcal{S}\left(\boldsymbol{R}^{n} \times \boldsymbol{R}^{n}\right)$, it follows that for every integer $k \geqslant 0$

$$
\begin{equation*}
\left|J_{\mathrm{e}}(x, \xi)\right| \leqslant C_{k}(1+|x|)^{-k}(1+|\xi|)^{-k} \tag{6.27}
\end{equation*}
$$

with a constant $C_{k}$ independent of $\varepsilon$. Hence, by (6.22)-(6.24), (6.26), (6.27) and the Lebesgue dominated convergence theorem, (6.21) now follows. Q.E.D.

## 7. Self-adjoint extensions of $\boldsymbol{P}(\boldsymbol{D})+\boldsymbol{A}(\boldsymbol{X}, \boldsymbol{D})$

In this section we give sufficient conditions for $P(D)+A(X, D)$ with domain $\mathcal{S}\left(\boldsymbol{R}^{n}\right)$ to have self-adjoint extensions. Let $A_{0}$ be as in Section 6 and let $P_{0}=$ $P(D) \upharpoonright \mathcal{S}\left(\boldsymbol{R}^{n}\right)$. Throughout this section, we assume that $A_{0}$ is symmetric. The following is an extension of a proposition of [1].

Theorem 4. Let hypothesis $(A)$ be fulfilled and assume that $p$ is a realvalued continuous function so that

$$
|p(\xi)| \leqslant C(1+|\xi|)^{j}
$$

for some real numbers $C$ and $j$. If there is a symmetric and orthogonal matrix $M$ of type $n \times n$ such that

$$
\begin{equation*}
p(-M \xi)=p(\xi), \quad a(M x,-M \xi)=a(x, \xi) \tag{7.1}
\end{equation*}
$$

then $P_{0}+A_{0}$ has a self-adjoint extension.
Proof. For $u \in L^{2}\left(\boldsymbol{R}^{n}\right)$, define

$$
(U u)(x)=\overline{u(M x)} .
$$

Since $M$ is symmetric and orthogonal, $U$ is a conjugation in the sense of [4, Definition 17, p. 1231]. Using (7.1), we shall show that $P_{0}+A_{0}$ commutes with $U$. Let $u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$. Since $|M x|=|x|$, we see that $U u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$. By (7.1) and the fact that $\mathscr{F}[U u](\xi)=\overline{\mathscr{F}}[u](-M \xi)$,

$$
\left(P_{0}+A_{0}\right) U u(x)=\int e^{i\langle x, \xi\rangle}(p(-M \xi)+\bar{a}(M x,-M \xi)) \overline{\mathcal{F}[u](-M \bar{\xi}) d \xi . . . ~}
$$

A change of variables gives

$$
\begin{equation*}
\left(P_{0}+A_{0}\right) U u(x)=\int e^{i\langle-M x, \eta\rangle}(p(\eta)+\bar{a}(M x, \eta)) \overline{\hat{u}(\eta)} d \eta . \tag{7.2}
\end{equation*}
$$

Thus, by Theorem 18, p. 1231 of [4], $P_{0}+A_{0}$ has equal deficiency indices, which implies that $P_{0}+A_{0}$ admits a self-adjoint extension.
Q.E.D.

Next, we shall give a sufficient condition for $P_{0}+A_{0}$ to be essentially selfadjoint. We will need a stronger condition on $A(X, D)$.

Theorem 5. Let $m, \rho, \delta$ be real numbers with $0 \leqslant \delta \leqslant \rho \leqslant 1, \delta<1$ and assume that a is a $C^{\infty}$ function so that

$$
\begin{equation*}
\left|(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta} a(x, \xi)\right| \leqslant C_{\alpha \beta}(|\xi|)(1+|\xi|)^{m+\delta|\beta|-\rho|\alpha|} \tag{7.3}
\end{equation*}
$$

for all multi-indices $\alpha, \beta$. Here $C_{\alpha \beta}(|\xi|)$ is a continuous function on $\boldsymbol{R}^{n}$ such that $C_{a \beta}(|\xi|) \rightarrow 0$ as $|\xi| \rightarrow \infty$. Moreover, assume that $p$ is a real-valued continuous function so that

$$
\begin{equation*}
C_{1}(1+|\xi|)^{m} \leqslant|p(\xi)| \leqslant C_{2}(1+|\xi|)^{\bar{m}} \tag{7.4}
\end{equation*}
$$

for some $C_{1}, C_{2}$ and $\bar{m}$. Then $P_{0}+A_{0}$ is essentially self-adjoint.
Proof. For every $R>0$, choose a function $\varphi_{R} \in C^{\infty}\left(\boldsymbol{R}^{n}\right)$ with $\varphi_{R}(\xi)=1$ when $|\xi|<R$ and $\varphi_{R}(\xi)=0$ when $|\xi|>R+1$. Let $\Phi_{R}(D)$ be a pseudo-differential operator with the symbol $\varphi_{R}(\xi)$. We make a decomposition

$$
\begin{align*}
& A(X, D)(1-\Delta)^{-m / 2}  \tag{7.5}\\
= & A(X, D)\left(1-\Phi_{R}(D)\right)(1-\Delta)^{-m / 2}+A(X, D) \Phi_{R}(D)(1-\Delta)^{-m / 2} .
\end{align*}
$$

Note that the symbol of the first term on the right side is

$$
a(x, \xi)\left(1-\varphi_{R}(\xi)\right)\left(1+|\xi|^{2}\right)^{-m / 2} .
$$

Since $1-\varphi_{R}(\xi)=0$ for $|\xi|<R$, we obtain for all multi-indices $\alpha, \beta$

$$
\begin{aligned}
& \left|(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta}\left\{a(x, \xi)\left(1-\varphi_{R}(\xi)\right)\left(1+|\xi|^{2}\right)^{-m / 2}\right\}\right| \\
\leqslant & M_{\alpha \beta}\left(\max _{\mu \leqslant \alpha} \sup _{\nu \leqslant \beta} \sup _{\mid \xi \geqslant 2} C_{\mu_{\nu}}(|\xi|)\right)(1+|\xi|)^{\delta|\beta|-\rho|\alpha|}
\end{aligned}
$$

with a constant $M_{\alpha \beta}$. Here we have used (7.3). From a theorem due to Calderon and Vaillancourt ([8], p. 215) it follows that

$$
\begin{equation*}
\left\|A(X, D)\left(1-\Phi_{R}(D)\right)(1-\Delta)^{-m / 2} v\right\|_{L^{2}} \leqslant C_{R}\|v\|_{L^{2}}, \quad v \in L^{2}\left(\boldsymbol{R}^{n}\right) \tag{7.6}
\end{equation*}
$$

where

$$
\begin{align*}
C_{R}=C & \max _{\substack{|\alpha| \leqslant M \\
|\beta|<N}} \sup _{x, \xi}\left[\mid(\partial / \partial \xi)^{\alpha}(\partial / \partial x)^{\beta}\{a(x, \xi)\right.  \tag{7.7}\\
& \left.\left.\quad \times\left(1-\varphi_{R}(\xi)\right)\left(1+|\xi|^{2}\right)^{(-m-\delta|\beta|+\rho|\alpha|) / 2}\right\} \mid\right]
\end{align*}
$$

and $M=2[n / 2+1], N=2[n /(2(1-\delta))+1]$. The constant $C_{R}$ can be estimated by

$$
\begin{equation*}
\max _{\substack{|\alpha| \leqslant M \\|\beta|<N}}\left(\sup _{|\xi| \geqslant R} C_{a \beta}(|\xi|)\right) . \tag{7.8}
\end{equation*}
$$

Similarly, we note that the symbol of the second term on the right side of (7.5) is

$$
a(x, \xi) \varphi_{R}(\xi)\left(1+|\xi|^{2}\right)^{-m / 2}
$$

which vanishes when $|\xi|>R$. Again we can apply a theorem due to Calderon and Vaillancourt which shows that

$$
\begin{equation*}
\left\|A(X, D) \Phi_{R}(D)(1-\Delta)^{-m / 2} v\right\|_{L^{2}} \leqslant C_{R}^{\prime}\|v\|_{-m}, \quad v \in H_{-m} \tag{7.9}
\end{equation*}
$$

with $C_{R}^{\prime}$ depending on $R$.
Now, let $u \in \mathcal{S}\left(\boldsymbol{R}^{n}\right)$ and put $v=(1-\Delta)^{m / 2} u$. Then,

$$
\begin{equation*}
\|v\|_{-m}=\|u\|_{L^{2}} \tag{7.10}
\end{equation*}
$$

and

$$
\begin{align*}
\|v\|_{L^{2}}^{2} & =\left(\int_{|\xi|<R}+\int_{|\xi|>R}\right)\left(1+|\xi|^{2}\right)^{m}|\hat{u}(\xi)|^{2} d \xi  \tag{7.11}\\
& \leqslant\left(1+R^{2}\right)^{|m|}\|u\|_{L^{2}}^{2}+C\left\|P_{0} u\right\|_{L^{2}}^{2}
\end{align*}
$$

where we use (7.4). Hence, by (7.5)-(7.11),

$$
\begin{align*}
\left\|A_{0} u\right\|_{L^{2}} & =\left\|A(X, D)(1-\Delta)^{-m / 2} v\right\|_{L^{2}}  \tag{7.12}\\
& \leqslant C\left(\max _{\substack{|\alpha| \leqslant M \\
|\beta|<N}} \sup _{|\xi| \geqslant R} C_{\alpha \beta}(|\xi|)\right)\left\|P_{0} u\right\|_{L^{2}}+C_{R}^{\prime \prime}\|u\|_{L^{2}}
\end{align*}
$$

with $C_{R}^{\prime \prime}$ depending on $R$ and $C$ independent of $R$. Since $C_{\alpha \beta}(|\xi|) \rightarrow 0$ as $|\xi| \rightarrow$ $\infty$, (7.12) means that $A_{0}$ is $P_{0}$-bounded with relative bound 0 . Note that $P_{0}$ is essentially self-adjoint. Thus, by Theorem 4.4, p. 288 of [7], we conclude that $P_{0}+A_{0}$ is also essentially self-adjoint.
Q.E.D.

## Appendix

In [6], Hörmander systematized the method of stationary phase. In this appendix, we shall reproduce his results in a somewhat different form.

Lemma I. Let $K \subset \boldsymbol{R}^{n}$ be a compact set, $\Omega$ a neighborhood of $K$ and let $\Phi$ be a subset of $C^{k+1}(\Omega, \boldsymbol{R})$ with the following properties:
(a) There is a constant $C_{1}$ such that for every $f \in \Phi$ and every multi-index $\alpha$ with $1 \leqslant|\alpha| \leqslant k+1$ the estimate

$$
\left|(\partial / \partial x)^{\alpha} f(x)\right| \leqslant C_{1}, \quad x \in \Omega
$$

is valid.
(b) There is a positive constant $C_{2}$ such that for every $f \in \Phi$

$$
\left|f^{\prime}(x)\right| \geqslant C_{2}, \quad x \in K
$$

Then we have for $\omega \in \boldsymbol{R}$ and all $f \in \Phi$

$$
\left|\int e^{i \omega f(x)} u(x) d x\right| \leqslant C(1+|\omega|)^{-k}|u|_{k}, \quad u \in C_{0}^{k}(K)
$$

If one looks carefully into the proof of Lemma A. 1 of [6] one finds easily that the conclusion of Lemma I remains valid. So we omit the proof.

Lemma II. Let $f \in C^{\infty}(\Omega)$ be a real-valued function in a neighborhood $\Omega$ of 0 in $\boldsymbol{R}^{n}$. Assume that $f^{\prime}(0)=0$ and that $f^{\prime \prime}(0)$ is non-singular. Moreover, assume that $f^{\prime}(x) \neq 0$ for every $x \in \Omega \backslash\{0\}$. Then there exist differential operators $L_{f, j}$ of order $2 j$ such that for every integer $k>0$ and every $s>2 k+n / 2$, we have when $u \in C_{0}^{\infty}(K), K$ compact $\subset \Omega$,

$$
\begin{align*}
& \left|\int u(x) e^{i \omega f(x)} d x-\sum_{j=0}^{k-1}\right| \operatorname{det} f^{\prime \prime}(0) /\left.2 \pi\right|^{-1 / 2} e^{\pi i \sigma / 4} e^{i \omega f(0)} L_{f, j} u(0) \omega^{-j-n / 2} \mid  \tag{A.1}\\
& \leqslant C \omega^{-n / 2-k}|u|_{s}, \quad \omega>1 .
\end{align*}
$$

Here $\sigma=\operatorname{sgn} f^{\prime \prime}(0)$, i.e.,

$$
\sigma=\sum \operatorname{sgn} \lambda_{j}
$$

in which $\lambda_{j}, j=1, \cdots, n$, are the repeated eigenvalues of the symmetric matrix $f^{\prime \prime}(0)$.
Since we assume that $f^{\prime}(x) \neq 0$ for every $x \in \Omega \backslash\{0\}$, this lemma follows immediately from Lemma A. 1 and Lemma A. 4 of [6]. In fact, choosing $\alpha \in$ $C_{0}^{\infty}(\Omega)$ so that supp $\alpha$ is contained in a small neighborhood of the origin and using Lemma A. 4 of [6], we find that there exist differential operators $L_{f, j}$ of order $2 j$ such that

$$
\begin{align*}
& \left|\int u(x) \alpha(x) e^{i \omega f(x)} d x-\sum_{j=0}^{k-1}\right| \operatorname{det} f^{\prime \prime}(0) /\left.2 \pi\right|^{-1 / 2} e^{\pi i \sigma / 4} e^{i \omega f(0)} L_{f, j}(u \alpha)(0) \omega^{-j-n / 2} \mid  \tag{A.2}\\
& \leqslant C \omega^{-n / 2-k}|u \alpha|_{s} .
\end{align*}
$$

If $r_{0}>0$ is sufficiently small and $\alpha(x)=1$ when $|x| \leqslant r_{0}$, then $f$ has no critical point in $\operatorname{supp} u(1-\alpha)$. Applying Lemma A. 1 of [6] to the integral

$$
\int u(x)(1-\alpha(x)) e^{i \omega f(x)} d x
$$

we have

$$
\begin{align*}
& \left|\int u(x)(1-\alpha(x)) e^{i \omega f(x)} d x\right|  \tag{A.3}\\
\leqslant & C(1+|\omega|)^{-k-[n / 2]-1}|(1-\alpha) u|_{k+[n / 2]+1} .
\end{align*}
$$

Combining (A.2) and (A.3), we obtain the conclusion of Lemma II.

We remark that Lemma II is applicable uniformly to all functions in a small neighborhood of $f$ in $C^{\infty}(\Omega)$. Indeed, Hörmander [6] showed that (A.2) is applicable uniformly to all functions in a small neighborhood of $f$ in $C^{\infty}(\Omega)$. Remembering that $f^{\prime}(x) \neq 0$ for every $x \in \Omega \backslash\{0\}$, we find that for any compact set $K \subset \Omega \backslash\{0\}$ and any integer $k \geqslant 0$ the hypotheses of Lemma I are satisfied, provided that we take $\Phi$ to be a small neighborhood of $f$ in $C^{\infty}(\Omega)$. Thus (A.3) is also applicable uniformly to all functions in a small neighborhood of $f$ in $C^{\infty}(\Omega)$.

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