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# ON SOME MICROLOCAL PROPERTIES OF THE RANGE OF A PSEUDODIFFERENTIAL OPERATOR OF PRINCIPAL TYPE 

Jens Wittsten


#### Abstract

We obtain microlocal analogues of results by L. Hörmander about inclusion relations between the ranges of first order differential operators with coefficients in $C^{\infty}$ that fail to be locally solvable. Using similar techniques, we study the properties of the range of classical pseudodifferential operators of principal type that fail to satisfy condition $(\Psi)$.


## 1. Introduction

We shall study the properties of the range of a classical pseudodifferential operator $P \in \Psi_{\mathrm{cl}}^{m}(X)$ that is not locally solvable, where $X$ is a $C^{\infty}$ manifold of dimension $n$. Here, classical means that the total symbol of $P$ is an asymptotic sum of homogeneous terms,

$$
\sigma_{P}(x, \xi)=p_{m}(x, \xi)+p_{m-1}(x, \xi)+\cdots,
$$

where $p_{k}$ is homogeneous of degree $k$ in $\xi$ and $p_{m}$ denotes the principal symbol of $P$. When no confusion can occur we will simply refer to $\sigma_{P}$ as the symbol of $P$. We shall restrict our study to operators of principal type, which means that the Hamilton vector field $H_{p_{m}}$ and the radial vector field are linearly independent when $p_{m}=0$. We shall also assume that all operators are properly supported, that is, both projections from the support of the kernel in $X \times X$ to $X$ are proper maps. For such operators, local solvability at a compact set $M \subset X$ means that for every $f$ in a subspace of $C^{\infty}(X)$ of finite codimension there is a distribution $u$ in $X$ such that

$$
\begin{equation*}
P u=f \tag{1-1}
\end{equation*}
$$

in a neighborhood of $M$. We can also define microlocal solvability at a set in the cosphere bundle, or equivalently, at a conic set in $T^{*}(X) \backslash 0$, the cotangent bundle of $X$ with the zero section removed. By a conic set $K \subset T^{*}(X) \backslash 0$ we mean a set that is conic in the fiber, that is,

$$
(x, \xi) \in K \quad \text { implies } \quad(x, \lambda \xi) \in K \quad \text { for all } \lambda>0 .
$$

If, in addition, $\pi_{x}(K)$ is compact in $X$, where $\pi_{x}: T^{*}(X) \rightarrow X$ is the projection, then $K$ is said to be compactly based. Thus, we say that $P$ is solvable at the compactly based cone $K \subset T^{*}(X) \backslash 0$ if there

[^0]is an integer $N$ such that for every $f \in H_{(N)}^{\text {loc }}(X)$ there exists a $u \in \mathscr{D}^{\prime}(X)$ with $K \cap W F(P u-f)=\varnothing$ (see Definition 2.1).

The famous example due to Hans Lewy [1957] of the existence of functions $f \in C^{\infty}$ ( $\mathbb{R}^{3}$ ) such that the equation

$$
\partial_{x_{1}} u+i \partial_{x_{2}} u-2 i\left(x_{1}+i x_{2}\right) \partial_{x_{3}} u=f
$$

does not have any solution $u \in \mathscr{D}^{\prime}(\Omega)$ in any open nonvoid subset $\Omega \subset \mathbb{R}^{3}$ contradicted the assumption that partial differential equations with smooth coefficients behave as analytic partial differential equations, for which existence of analytic solutions is guaranteed by the Cauchy-Kovalevsky theorem. This example led to an extension due to Hörmander [1960b; 1960a] in the sense of a necessary condition for a differential equation $P(x, D) u=f$ to have a solution locally for every $f \in C^{\infty}$. In fact (see [Hörmander 1963, Theorem 6.1.1]), if $\Omega$ is an open set in $\mathbb{R}^{n}$, and $P$ is a differential operator of order $m$ with coefficients in $C^{\infty}(\Omega)$ such that the differential equation $P(x, D) u=f$ has a solution $u \in \mathscr{D}^{\prime}(\Omega)$ for every $f \in C_{0}^{\infty}(\Omega)$, then $\left\{p_{m}, \bar{p}_{m}\right\}$ must vanish at every point $(x, \xi) \in \Omega \times \mathbb{R}^{n}$ for which $p_{m}(x, \xi)=0$, where

$$
\{a, b\}=\sum_{j=1}^{n} \partial_{\xi_{j}} a \partial_{x_{j}} b-\partial_{x_{j}} a \partial_{\xi_{j}} b
$$

denotes the Poisson bracket.
In addition to his example, Lewy conjectured that differential operators that fail to have local solutions are essentially uniquely determined by the range. Later Hörmander [1963, Chapter 6.2] proved that if $P$ and $Q$ are two first order differential operators with coefficients in $C^{\infty}(\Omega)$ and in $C^{1}(\Omega)$, respectively, such that the equation $P(x, D) u=Q(x, D) f$ has a solution $u \in \mathscr{D}^{\prime}(\Omega)$ for every $f \in C_{0}^{\infty}(\Omega)$, and $x$ is a point in $\Omega$ such that

$$
\begin{equation*}
p_{1}(x, \xi)=0 \quad \text { and } \quad\left\{p_{1}, \bar{p}_{1}\right\}(x, \xi) \neq 0 \tag{1-2}
\end{equation*}
$$

for some $\xi \in \mathbb{R}^{n}$, then there is a constant $\mu$ such that (at the fixed point $x$ )

$$
{ }^{t} Q(x, D)=\mu^{t} P(x, D)
$$

where ${ }^{t} Q$ and ${ }^{t} P$ are the formal adjoints of $Q$ and $P$. If (1-2) holds for a dense set of points $x$ in $\Omega$ and if the coefficients of $p_{1}(x, D)$ do not vanish simultaneously in $\Omega$, then there is a function $\mu \in C^{1}(\Omega)$ such that

$$
\begin{equation*}
Q(x, D) u=P(x, D)(\mu u) \tag{1-3}
\end{equation*}
$$

Furthermore, for such an operator $P$ and function $\mu$, the equation $P(x, D) u=\mu P(x, D) f$ has a solution $u \in \mathscr{D}^{\prime}(\Omega)$ for every $f \in C_{0}^{\infty}(\Omega)$ if and only if $p_{1}(x, D) \mu=0$.

Hörmander also showed that this result extends to operators of higher order in the following way (see [1963, Theorem 6.2.4]). If $P$ is a differential operator of order $m$ with coefficients in $C^{\infty}(\Omega)$ and $\mu$ is a function in $C^{m}(\Omega)$ such that the equation

$$
P(x, D) u=\mu P(x, D) f
$$

has a solution $u \in \mathscr{D}^{\prime}(\Omega)$ for every $f \in C_{0}^{\infty}(\Omega)$, then it follows that

$$
\sum_{j=1}^{n} \partial_{\xi_{j}} p_{m}(x, \xi) \partial_{x_{j}} \mu(x)=0
$$

for all $x \in \Omega$ and $\xi \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
\left\{p_{m}, \bar{p}_{m}\right\}(x, \xi) \neq 0 \quad \text { and } \quad p_{m}(x, \xi)=0 \tag{1-4}
\end{equation*}
$$

This means that the derivative of $\mu$ must vanish along every bicharacteristic element with initial data $(x, \xi)$ satisfying (1-4), that is, giving rise to nonexistence of solutions.

If $P$ is a pseudodifferential operator such that $P$ is microlocally elliptic near $\left(x_{0}, \xi_{0}\right)$, then there exists a microlocal inverse, called a parametrix $P^{-1}$ of $P$, such that in a conic neighborhood of $\left(x_{0}, \xi_{0}\right)$ we have $P P^{-1}=P^{-1} P=$ Identity modulo smoothing operators. $P$ is then trivially seen to be microlocally solvable near $\left(x_{0}, \xi_{0}\right)$, and for any pseudodifferential operator $Q$ we can write $Q=P P^{-1} Q+R=$ $P E+R$, where $R$ is a smoothing operator. When the range of $Q$ is microlocally contained in the range of $P$, we will show the existence of this type of representation for $Q$ in the case when $P$ is a nonsolvable pseudodifferential operator of principal type, although we will have to content ourselves with a weaker statement concerning the Taylor coefficients of the symbol of the operator $R$ (see Theorem 2.19 for the precise formulation of the result). Note that when $P$ is solvable but nonelliptic we cannot hope to obtain such a representation in general; see the remark on page 440.

For pseudodifferential operators of principal type, Hörmander [1985b] proved that local solvability in the sense of (1-1) implies that $M$ has an open neighborhood $Y$ in $X$ where $p_{m}$ satisfies condition $(\Psi)$, which means that

$$
\begin{equation*}
\text { Im } a p_{m} \text { does not change sign from }- \text { to }+ \text { along the oriented bicharacteristics of } \operatorname{Re} a p_{m} \tag{1-5}
\end{equation*}
$$

over $Y$ for any $0 \neq a \in C^{\infty}\left(T^{*}(Y) \backslash 0\right)$. The oriented bicharacteristics are the positive flow-outs of the Hamilton vector field $H_{\operatorname{Re} a p_{m}}$ on $\operatorname{Re} a p_{m}=0$. The proof relies on an idea due to Moyer [1978], and uses the fact that condition (1-5) is invariant under multiplication of $p_{m}$ with nonvanishing factors, and conjugation of $P$ with elliptic Fourier integral operators.

Rather recently Dencker [2006] proved that condition $(\Psi)$ is also sufficient for local and microlocal solvability for operators of principal type. To get local solvability at a point $x_{0}$, Dencker assumed the strong form of the nontrapping condition at $x_{0}$,

$$
\begin{equation*}
p_{m}=0 \quad \text { implies } \quad \partial_{\xi} p_{m} \neq 0 \tag{1-6}
\end{equation*}
$$

This was the original condition for principal type of Nirenberg and Treves [1970a; 1970b; 1971], which is always obtainable microlocally after a canonical transformation. Thus, we shall study pseudodifferential operators that fail to satisfy condition $(\Psi)$ in place of the condition given by (1-4), and show that such operators are, in analogue with the inclusion relations between the ranges of differential operators that fail to be locally solvable, essentially uniquely determined by the range. However, even though (1-4) is a microlocal condition, we get the mentioned local results for differential operators because of analyticity
in $\xi$ of the corresponding symbol. Since this is not generally true for pseudodifferential operators, our results will be inherently microlocal. We will combine the techniques used in [Hörmander 1963] to prove the inclusion relations for differential operators with the approach used in [Hörmander 1985b] to prove the necessity of condition $(\Psi)$ for local solvability of pseudodifferential operators of principal type.

It is possible to extend these results to certain systems of pseudodifferential operators, which will be addressed in a forthcoming joint paper with Nils Dencker.

## 2. Nonsolvable operators of principal type

Let $X$ be a $C^{\infty}$ manifold of dimension $n$. In what follows, $C$ will be taken to be a new constant every time unless stated otherwise. We let $\mathbb{N}=\{0,1,2, \ldots\}$, and if $\alpha \in \mathbb{N}^{n}$ is a multiindex $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right)$, we let

$$
D_{x}^{\alpha}=D_{x_{1}}^{\alpha_{1}} \cdots D_{x_{n}}^{\alpha_{n}}
$$

where $D_{x_{j}}=-i \partial_{x_{j}}$. We shall also employ the standard notation $f_{(\alpha)}^{(\beta)}(x, \xi)=\partial_{x}^{\alpha} \partial_{\xi}^{\beta} f(x, \xi)$ for multiindices $\alpha, \beta$.

In this section we will follow the outline of [Hörmander 1985b, Chapter 26, Section 4]. Recall that the Sobolev space $H_{(s)}(X)$ for $s \in \mathbb{R}$ is a local space, that is, if $\varphi \in C_{0}^{\infty}(X)$ and $u \in H_{(s)}(X)$, then $\varphi u \in H_{(s)}(X)$, and the corresponding operator of multiplication is continuous. Thus we can define

$$
H_{(s)}^{\mathrm{loc}}(X)=\left\{u \in \mathscr{D}^{\prime}(X): \varphi u \in H_{(s)}(X) \text { for all } \varphi \in C_{0}^{\infty}(X)\right\}
$$

This is a Fréchet space, and its dual with respect to the inner product on $L^{2}$ is

$$
H_{(-s)}^{\mathrm{comp}}(X)=H_{(-s)}^{\mathrm{loc}}(X) \cap \mathscr{E}^{\prime}(X)
$$

Definition 2.1. If $K \subset T^{*}(X) \backslash 0$ is a compactly based cone, we shall say that the range of $Q \in \Psi_{\mathrm{cl}}^{m}(X)$ is microlocally contained in the range of $P \in \Psi_{\mathrm{cl}}^{k}(X)$ at $K$ if there exists an integer $N$ such that for every $f \in H_{(N)}^{\text {loc }}(X)$, there exists a $u \in \mathscr{D}^{\prime}(X)$ with $W F(P u-Q f) \cap K=\varnothing$.

If $I \in \Psi_{\mathrm{cl}}^{0}(X)$ is the identity on $X$, we obtain from Definition 2.1 the definition of microlocal solvability for a pseudodifferential operator (see [Hörmander 1985b, Definition 26.4.3]) by setting $Q=I$. Thus, the range of the identity is microlocally contained in the range of $P$ at $K$ if and only if $P$ is microlocally solvable at $K$. Note also that if $P$ and $Q$ satisfy Definition 2.1 for some integer $N$, then due to the inclusion

$$
H_{(t)}^{\mathrm{loc}}(X) \subset H_{(s)}^{\mathrm{loc}}(X) \quad \text { if } s<t
$$

the statement also holds for any integer $N^{\prime} \geq N$. Hence $N$ can always be assumed to be positive. Furthermore, the property is preserved if $Q$ is composed with a properly supported pseudodifferential operator $Q_{1} \in \Psi_{\mathrm{cl}}^{m^{\prime}}(X)$ from the right. Indeed, let $g$ be an arbitrary function in $H_{\left(N+m^{\prime}\right)}^{\text {loc }}(X)$. Then $f=Q_{1} g \in H_{(N)}^{\text {loc }}(X)$ since the map

$$
Q_{1}: H_{(s)}^{\mathrm{loc}}(X) \rightarrow H_{\left(s-m^{\prime}\right)}^{\mathrm{loc}}(X)
$$

is continuous for every $s \in \mathbb{R}$, so by Definition 2.1 there exists a $u \in \mathscr{D}^{\prime}(X)$ with $W F(P u-Q f) \cap K=\varnothing$. Hence the range of $Q Q_{1}$ is microlocally contained in the range of $P$ at $K$ with the integer $N$ replaced by $N+m^{\prime}$.

The property given by Definition 2.1 is also preserved under composition of both $P$ and $Q$ with a properly supported pseudodifferential operator from the left. This follows immediately from the fact that properly supported pseudodifferential operators are microlocal, that is,

$$
W F(A u) \subset W F(u) \cap W F(A) \quad \text { for } u \in \mathscr{D}^{\prime}(X)
$$

Remark. In Definition 2.1 we may always assume that $f \in H_{(N)}^{\text {comp }}(X)$ and $u \in \mathscr{E}^{\prime}(X)$ when considering a fixed cone $K$. In fact, assume

$$
Q f=P u+g
$$

where $f \in H_{(N)}^{\text {loc }}(X)$ and $u, g \in \mathscr{D}^{\prime}(X)$ with $W F(g) \cap K=\varnothing$, and let $Y \Subset X$ satisfy $K \subset T^{*}(Y) \backslash 0$. (We write $Y \Subset X$ when $\bar{Y}$ is compact and contained in $X$.) Since $P$ and $Q$ are properly supported we can find $Z_{1}, Z_{2} \subset X$ such that $P v=0$ in $Y$ if $v=0$ in $Z_{1}$, and $Q v=0$ in $Y$ if $v=0$ in $Z_{2}$. We may of course assume that $Y \Subset Z_{j}$ for $j=1,2$. Fix $\phi_{j} \in C_{0}^{\infty}(X)$ with $\phi_{j}=1$ on $Z_{j}$. Then we have $P u=P\left(\phi_{1} u\right)$ and $Q f=Q\left(\phi_{2} f\right)$ in $Y$, so

$$
\varnothing=W F(Q f-P u) \cap K=W F\left(Q\left(\phi_{2} f\right)-P\left(\phi_{1} u\right)\right) \cap K
$$

where $\phi_{1} u$ and $\phi_{2} f$ have compact support. Hence we may assume that $u \in \mathscr{E}^{\prime}(X)$ and $f \in H_{(N)}^{\text {comp }}(X)=$ $H_{(N)}^{\text {loc }}(X) \cap \mathscr{E}^{\prime}(X)$ to begin with. Note that this also implies $g=Q f-P u \in \mathscr{E}^{\prime}(X)$ since $P$ and $Q$ are properly supported.

The following easy example will prove useful when discussing inclusion relations between the ranges of solvable but nonelliptic operators.

Example 2.2. If $X \subset \mathbb{R}^{n}$ is open, and $K \subset T^{*}(X) \backslash 0$ is a compactly based cone, then the range of $D_{1}=-i \partial / \partial x_{1}$ is microlocally contained in the range of $D_{2}$ at $K$. In fact, this is trivially true since both operators are surjective $\mathscr{D}^{\prime}(X) \rightarrow \mathscr{D}^{\prime}(X) / C^{\infty}(X)$. To see that for example $D_{1}$ is surjective, we note that by the remark on page 427 it suffices to show that there exists a number $N \in \mathbb{Z}$ such that the equation $D_{1} u=f$ has a solution $u \in \mathscr{D}^{\prime}(X)$ for every $f \in H_{(N)}^{\text {comp }}(X)=H_{(N)}^{\text {loc }}(X) \cap \mathscr{E}^{\prime}(X)$. By [Hörmander 1983b, Theorem 10.3.1] this is satisfied for every $N \in \mathbb{Z}$ if $u \in H_{(N+1)}^{\mathrm{loc}}(X)$ is given by $E * f$, where $E$ is the regular fundamental solution of $D_{1}$.

Just as the microlocal solvability of a pseudodifferential operator $P$ gives an a priori estimate for the adjoint $P^{*}$, we have the following result for operators satisfying Definition 2.1.

Lemma 2.3. Let $K \subset T^{*}(X) \backslash 0$ be a compactly based cone. Let $Q \in \Psi_{\mathrm{cl}}^{m}(X)$ and $P \in \Psi_{\mathrm{cl}}^{k}(X)$ be properly supported pseudodifferential operators such that the range of $Q$ is microlocally contained in the range of $P$ at $K$. If $Y \Subset X$ satisfies $K \subset T^{*}(Y)$ and if $N$ is the integer in Definition 2.1, then for every positive integer $\kappa$ we can find a constant $C$, a positive integer $v$ and a properly supported pseudodifferential
operator $A$ with $W F(A) \cap K=\varnothing$ such that

$$
\begin{equation*}
\left\|Q^{*} v\right\|_{(-N)} \leq C\left(\left\|P^{*} v\right\|_{(\nu)}+\|v\|_{(-N-\kappa-n)}+\|A v\|_{(0)}\right) \tag{2-1}
\end{equation*}
$$

for all $v \in C_{0}^{\infty}(Y)$.
Since (2-1) holds for any $\kappa$, it is actually superfluous to include the dimension $n$ in the norm $\|v\|_{(-N-\kappa-n)}$. However, for our purposes, it turns out that this is the most convenient formulation.

Proof. We shall essentially adapt the proof of [Hörmander 1985b, Lemma 26.4.5]. Let $\|\cdot\|_{(s)}$ denote a norm in $H_{(s)}^{\text {comp }}(X)$ that defines the topology in $H_{(s)}^{c}(M)=H_{(s)}^{\text {loc }}(X) \cap \mathscr{E}^{\prime}(M)$ for every compact set $M \subset X$. (The reason we change notation from $H_{(s)}^{\text {comp }}(M)$ to $H_{(s)}^{c}(M)$ when $M$ is compact is to signify that $H_{(s)}^{c}(M)$ is a Hilbert space for each fixed compact set $M$.) Let $Y \Subset Z \Subset X$, and take $\chi \in C_{0}^{\infty}(X)$ with supp $\chi=\bar{Z}$ to be a real-valued cutoff function identically equal to 1 in a neighborhood of $Y$. Then $\chi Q f \in H_{(N-m)}^{c}(\bar{Z})$ for all $f \in H_{(N)}^{\text {comp }}(X)$ since $Q$ is properly supported, and we claim that for fixed $f \in H_{(N)}^{\text {comp }}(X)$ we have for some $C, v$ and $A$ as in the statement of the lemma

$$
\begin{equation*}
|(\chi Q f, v)| \leq C\left(\left\|P^{*} v\right\|_{(v)}+\|v\|_{(-N-\kappa-n)}+\|A v\|_{(0)}\right) \tag{2-2}
\end{equation*}
$$

for all $v \in C_{0}^{\infty}(Y)$. Indeed, by hypothesis and the remark on page 427 we can find $u$ and $\tilde{g}$ in $\mathscr{E}^{\prime}(X)$ with $W F(\tilde{g}) \cap K=\varnothing$ such that

$$
\chi Q f=Q f-(1-\chi) Q f=P u+\tilde{g}-(1-\chi) Q f
$$

Since $K \subset T^{*}(Y)$ and $\chi \equiv 1$ near $Y$ we get $W F((1-\chi) Q f) \cap K=\varnothing$, so $\chi Q f=P u+g$ for some $g \in \mathscr{E}^{\prime}(X)$ with $W F(g) \cap K=\varnothing$. Thus

$$
(\chi Q f, v)=\left(u, P^{*} v\right)+(g, v) \quad \text { for } v \in C_{0}^{\infty}(Y)
$$

Now choose properly supported pseudodifferential operators $B_{1}$ and $B_{2}$ of order 0 with $I=B_{1}+B_{2}$ and $W F\left(B_{1}\right) \cap W F(g)=\varnothing$ and $W F\left(B_{2}\right) \cap K=\varnothing$, which is possible since $W F(g) \cap K=\varnothing$. Since $g \in \mathscr{E}^{\prime}(X)$ and $B_{1}: \mathscr{E}^{\prime}(X) \rightarrow \mathscr{E}^{\prime}(X)$ is continuous and microlocal we get $B_{1} g \in C_{0}^{\infty}(X)$, so ( $\left.B_{1} g, v\right)$ can be estimated by $C\|v\|_{(-N-\kappa-n)}$. Also, $g \in H_{(-\mu)}^{\text {loc }}(X)$ for some $\mu>0$ so if $B$ is properly supported and elliptic of order $\mu$, and $B^{\prime} \in \Psi_{\mathrm{cl}}^{-\mu}(X)$ is a properly supported parametrix of $B$, then

$$
\begin{equation*}
B_{2}^{*} v=B^{\prime} B B_{2}^{*} v+L B_{2}^{*} v \tag{2-3}
\end{equation*}
$$

where $L \in \Psi^{-\infty}(X)$ and both $B^{\prime}$ and $L$ are continuous $H_{(s)}^{\text {comp }}(X) \rightarrow H_{(s+\mu)}^{\text {comp }}(X)$. Hence

$$
\left|\left(B_{2} g, v\right)\right| \leq C\left\|B_{2}^{*} v\right\|_{(\mu)} \leq C\left(\left\|B B_{2}^{*} v\right\|_{(0)}+\left\|B_{2}^{*} v\right\|_{(0)}\right)
$$

and if we apply the identity (2-3) to $\left\|B_{2}^{*} v\right\|_{(0)},\left\|B_{2}^{*} v\right\|_{(-\mu)}, \ldots$ sufficiently many times, and then recall that $B_{2}^{*}$ is properly supported and of order 0 , we obtain

$$
\left|\left(B_{2} g, v\right)\right| \leq C\left(\left\|B B_{2}^{*} v\right\|_{(0)}+\|v\|_{(-N-\kappa-n)}\right)
$$

Since we chose $B$ to be properly supported this gives (2-2) with $A=B B_{2}^{*}$.

For fixed $\kappa$, suppose $V$ is the space $C_{0}^{\infty}(Y)$ equipped with the topology defined by the seminorms $\|v\|_{(-N-\kappa-n)},\left\|P^{*} v\right\|_{(\nu)}$ for $v=1,2, \ldots$, and $\|A v\|_{(0)}$, where $A$ is a properly supported pseudodifferential operator with $K \cap W F(A)=\varnothing$. It suffices to use a countable sequence $A_{1}, A_{2}, \ldots$, where $A_{v}$ is noncharacteristic of order $v$ in a set that increases to $\left(T^{*}(X) \backslash 0\right) \backslash K$ as $v \rightarrow \infty$. Thus $V$ is a metrizable space. The sesquilinear form $(\chi Q f, v)$ in the product of the Hilbert space $H_{(N-m)}^{c}(\bar{Z})$ and the metrizable space $V$ is obviously continuous in $\chi Q f$ for fixed $v$, and by (2-2) it is also continuous in $v$ for fixed $f$. Hence it is continuous, which means that for some $v$ and $C$

$$
|(\chi Q f, v)| \leq C\|Q f\|_{(N-m)}\left(\left\|P^{*} v\right\|_{(v)}+\|v\|_{(-N-\kappa-n)}+\|A v\|_{(0)}\right)
$$

for all $f \in H_{(N)}^{\text {comp }}(X)$ and $v \in C_{0}^{\infty}(Y)$. Now $Q$ is continuous from $H_{(N)}^{\text {comp }}(X)$ to $H_{(N-m)}^{\text {comp }}(X)$ so we have $\|Q f\|_{(N-m)} \leq C\|f\|_{(N)}$. Since $\chi \equiv 1$ near $Y$ and $(\chi Q)^{*}=Q^{*} \chi$, this yields the estimate

$$
\begin{equation*}
\left|\left(f, Q^{*} v\right)\right| \leq C\|f\|_{(N)}\left(\left\|P^{*} v\right\|_{(\nu)}+\|v\|_{(-N-\kappa-n)}+\|A v\|_{(0)}\right) . \tag{2-4}
\end{equation*}
$$

For $v \in C_{0}^{\infty}(Y)$ and $Q^{*}$ properly supported we have $Q^{*} v \in C_{0}^{\infty}(X)$, and therefore also $Q^{*} v \in H_{(-N)}^{\text {loc }}(X)$. Viewing $Q^{*} v$ as a functional on $H_{(N)}^{\text {comp }}(X)$, the dual of $H_{(-N)}^{\text {loc }}(X)$ with respect to the standard inner product on $L^{2}$, we obtain (2-1) after taking the supremum over all $f \in H_{(N)}^{\text {comp }}(X)$ with $\|f\|_{(N)}=1$.

We will need the following analogue of [Hörmander 1985b, Proposition 26.4.4]. Recall that $\mathscr{H}$ : $T^{*}(Y) \backslash 0 \rightarrow T^{*}(X) \backslash 0$ is a canonical transformation if and only if its graph $C_{\mathscr{H}}$ in the product $\left(T^{*}(X) \backslash 0\right) \times\left(T^{*}(Y) \backslash 0\right)$ is Lagrangian with respect to the difference $\sigma_{X}-\sigma_{Y}$ of the symplectic forms of $T^{*}(X)$ and $T^{*}(Y)$ lifted to $T^{*}(X) \times T^{*}(Y)=T^{*}(X \times Y)$. This differs in sign from the symplectic form $\sigma_{X}+\sigma_{Y}$ of $T^{*}(X \times Y)$ so it is the twisted graph

$$
C_{\mathscr{H}}^{\prime}=\left\{(x, \xi, y,-\eta):(x, \xi, y, \eta) \in C_{\mathscr{H}}\right\},
$$

which is Lagrangian with respect to the standard symplectic structure in $T^{*}(X \times Y)$.
Proposition 2.4. Let $K \subset T^{*}(X) \backslash 0$ and $K^{\prime} \subset T^{*}(Y) \backslash 0$ be compactly based cones and let $\chi$ be a homogeneous symplectomorphism from a conic neighborhood of $K^{\prime}$ to one of $K$ such that $\chi\left(K^{\prime}\right)=K$. Let $A \in I^{m^{\prime}}\left(X \times Y, \Gamma^{\prime}\right)$ and $B \in I^{m^{\prime \prime}}\left(Y \times X,\left(\Gamma^{-1}\right)^{\prime}\right)$, where $\Gamma$ is the graph of $\chi$, and assume that $A$ and $B$ are properly supported and noncharacteristic at the restriction of the graphs of $\chi$ and $\chi^{-1}$ to $K^{\prime}$ and to $K$ respectively, while $W F^{\prime}(A)$ and $W F^{\prime}(B)$ are contained in small conic neighborhoods. Then the range of the pseudodifferential operator $Q$ in $X$ is microlocally contained in the range of the pseudodifferential operator $P$ in $X$ at $K$ if and only if the range of the pseudodifferential operator $B Q A$ in $Y$ is microlocally contained in the range of the pseudodifferential operator $B P A$ in $Y$ at $K^{\prime}$.

Proof. Choose $A_{1} \in I^{-m^{\prime \prime}}\left(X \times Y, \Gamma^{\prime}\right)$ and $B_{1} \in I^{-m^{\prime}}\left(Y \times X,\left(\Gamma^{-1}\right)^{\prime}\right)$ properly supported such that

$$
\begin{array}{ll}
K^{\prime} \cap W F\left(B A_{1}-I\right)=\varnothing, & K \cap W F\left(A_{1} B-I\right)=\varnothing \\
K^{\prime} \cap W F\left(B_{1} A-I\right)=\varnothing, & K \cap W F\left(A B_{1}-I\right)=\varnothing
\end{array}
$$

Assume that the range of $Q$ is microlocally contained in the range of $P$ at $K$ and choose $N$ as in Definition 2.1. Let $g \in H_{\left(N+m^{\prime}\right)}^{\text {loc }}(Y)$ and set $f=A g \in H_{(N)}^{\text {loc }}(X)$. Then we can find $u \in \mathscr{D}^{\prime}(X)$ such that
$K \cap W F(P u-Q f)=\varnothing$. Let $v=B_{1} u \in \mathscr{D}^{\prime}(Y)$. Then

$$
W F(A v-u)=W F\left(\left(A B_{1}-I\right) u\right)
$$

does not meet $K$, so $K \cap W F(P A v-Q f)=\varnothing$. Recalling that $f=A g$ this implies

$$
K^{\prime} \cap W F(B P A v-B Q A g)=\varnothing
$$

so the range of $B Q A$ is microlocally contained in the range of $B P A$ at $K^{\prime}$. Conversely, if the range of $B Q A$ is microlocally contained in the range of $B P A$ at $K^{\prime}$, it follows that the range of $A_{1} B Q A B_{1}$ is microlocally contained in the range of $A_{1} B P A B_{1}$ at $K$. Since

$$
K \cap W F\left(A_{1} B P A B_{1} u-A_{1} B Q A B_{1} f\right)=K \cap W F(P u-Q f),
$$

this means that the range of $Q$ is microlocally contained in the range of $P$ at $K$, which proves the proposition.

Before we can state our main theorem, we need to study the geometric situation that occurs when $p$ fails to satisfy condition $(\Psi)$. Recall that by [Hörmander 1985b, Theorem 26.4.12] we may always assume that the nonvanishing factor in condition (1-5) is a homogeneous function. We begin with a lemma concerning a reduction of the general case.

Lemma 2.5. Let $p$ and $q$ be homogeneous smooth functions on $T^{*}(X) \backslash 0$, and let $t \mapsto \gamma(t)$, for $a \leq t \leq b$, be a bicharacteristic interval of $\operatorname{Re} q p$ such that $q(\gamma(t)) \neq 0$ for $a \leq t \leq b$. If

$$
\begin{equation*}
\operatorname{Im} q p(\gamma(a))<0<\operatorname{Im} q p(\gamma(b)) \tag{2-5}
\end{equation*}
$$

then there exists a proper subinterval $\left[a^{\prime}, b^{\prime}\right] \subset[a, b]$, possibly reduced to a point, such that
(i) $\operatorname{Im} q p(\gamma(t))=0$ for $a^{\prime} \leq t \leq b^{\prime}$,
(ii) for every $\varepsilon>0$ there exist $a^{\prime}-\varepsilon<s_{-}<a^{\prime}$ and $b^{\prime}<s_{+}<b^{\prime}+\varepsilon$ such that $\operatorname{Im} q p\left(\gamma\left(s_{-}\right)\right)<0<$ $\operatorname{Im} q p\left(\gamma\left(s_{+}\right)\right)$.

If $\gamma(t)$ is defined for $a \leq t \leq b$ we shall in the sequel say that $\operatorname{Im} q p$ changes sign from - to + on $\gamma$ if (2-5) holds. If $\left.\gamma\right|_{\left[a^{\prime}, b^{\prime}\right]}$ is the restriction of $\gamma$ to $\left[a^{\prime}, b^{\prime}\right]$ and (i) and (ii) hold we shall say that $\operatorname{Im} q p$ strongly changes sign from - to + on $\left.\gamma\right|_{\left[a^{\prime}, b^{\prime}\right]}$.
Proof. It suffices to regard the case that $q=1, X=\mathbb{R}^{n}, p$ is homogeneous of degree 1 with $\operatorname{Re} p=\xi_{1}$, and the bicharacteristic of $\operatorname{Re} p$ is given by

$$
\begin{equation*}
a \leq x_{1} \leq b, \quad x^{\prime}=\left(x_{2}, \ldots, x_{n}\right)=0, \quad \xi=\varepsilon_{n} \tag{2-6}
\end{equation*}
$$

Here $\varepsilon_{n}=(0, \ldots, 0,1) \in \mathbb{R}^{n}$, and we shall in what follows write $\xi^{0}$ in place of $\varepsilon_{n}^{\prime}$. The proof of this fact is taken from [Hörmander 1985b, page 97] and is given here for the purpose of reference later, in particular in connection with Definition 2.11 below.

Choose a pseudodifferential operator $Q$ with principal symbol $q$. If we let $P_{1}=Q P$, then the principal symbol of $P_{1}$ is $p_{1}=q p$ so $\operatorname{Im} p_{1}$ changes sign from - to + on the bicharacteristic $\gamma$ of Re $p_{1}$. Now choose $Q_{1}$ to be of order $1-$ degree $P_{1}$ with positive, homogeneous principal symbol. If $p_{2}$ is the
principal symbol of $P_{2}=Q_{1} P_{1}$, it follows that $\operatorname{Re} p_{1}$ and $\operatorname{Re} p_{2}$ have the same bicharacteristics, including orientation, and since $p_{2}$ is homogeneous of degree 1 these can be considered to be curves on the cosphere bundle $S^{*}(X)$. Moreover, $\operatorname{Im} p_{1}$ and $\operatorname{Im} p_{2}$ have the same sign, so $\operatorname{Im} p_{2}$ changes sign from - to + along $\gamma \subset S^{*}(X)$. If $\gamma$ is a closed curve on $S^{*}(X)$ we can pick an arc that is not closed where the sign change still occurs. If we assume this to be done, then [Hörmander 1985b, Proposition 26.1.6] states that there exists a $C^{\infty}$ homogeneous canonical transformation $\chi$ from an open conic neighborhood of (2-6) to one of $\gamma$ such that $\chi\left(x_{1}, 0, \varepsilon_{n}\right)=\gamma\left(x_{1}\right)$ and $\chi^{*}\left(\operatorname{Re} p_{2}\right)=\xi_{1}$. Since the Hamilton field is symplectically invariant it follows that the equations of a bicharacteristic are invariant under the action of canonical transformations, that is, $\tilde{\gamma}$ is a bicharacteristic of $\chi^{*}\left(\operatorname{Re} p_{2}\right)$ if and only if $\chi(\tilde{\gamma})$ is a bicharacteristic of $\operatorname{Re} p_{2}$. This proves the claim.

In accordance with the notation in [Hörmander 1985b, page 97], let $\left(x^{\prime}, \xi^{\prime}\right)=\left(0, \xi^{0}\right)$ and consider

$$
L\left(0, \xi^{0}\right)=\inf \left\{t-s: a<s<t<b, \operatorname{Im} p\left(s, 0, \varepsilon_{n}\right)<0<\operatorname{Im} p\left(t, 0, \varepsilon_{n}\right)\right\}
$$

For every small $\delta>0$ there exist $s_{\delta}$ and $t_{\delta}$ such that $a<s_{\delta}<t_{\delta}<b, \operatorname{Im} p\left(s_{\delta}, 0, \varepsilon_{n}\right)<0<\operatorname{Im} p\left(t_{\delta}, 0, \varepsilon_{n}\right)$ and $t_{\delta}-s_{\delta}<L\left(0, \xi^{0}\right)+\delta$. Choose a sequence $\delta_{j} \rightarrow 0$ such that the limits $a^{\prime}=\lim s_{\delta_{j}}$ and $b^{\prime}=\lim t_{\delta_{j}}$ exist. Then $b^{\prime}-a^{\prime}=L\left(0, \xi^{0}\right)$ and in view of (2-5) we have $a<a^{\prime} \leq b^{\prime}<b$ by continuity. Moreover, $\operatorname{Im} p\left(t, 0, \varepsilon_{n}\right)=0$ for $a^{\prime} \leq t \leq b^{\prime}$. This is clear if $a^{\prime}=b^{\prime}$. If on the other hand $\operatorname{Im} p\left(t, 0, \varepsilon_{n}\right)$ is, say, strictly positive for some $a^{\prime}<t<b^{\prime}$, then $L\left(0, \xi^{0}\right) \leq t-s_{\delta_{j}} \rightarrow t-a^{\prime}<b^{\prime}-a^{\prime}$, a contradiction. Thus (i) holds.

To prove (ii), let $\varepsilon>0$. After possibly reducing to a subsequence we may assume that the sequences $\left\{s_{\delta_{j}}\right\}$ and $\left\{t_{\delta_{j}}\right\}$ given above are monotone increasing and decreasing, respectively. It then follows by (i) that $s_{\delta_{j}}<a^{\prime} \leq b^{\prime}<t_{\delta_{j}}$ for all $j$. Since $s_{\delta_{j}} \rightarrow a^{\prime}$ and $t_{\delta_{j}} \rightarrow b^{\prime}$ we can choose $j$ so that $a^{\prime}-\varepsilon<s_{\delta_{j}}<a^{\prime}$ and $b^{\prime}<t_{\delta_{j}}<b^{\prime}+\varepsilon$. By construction we have $\operatorname{Im} p\left(s_{\delta_{j}}, 0, \varepsilon_{n}\right)<0<\operatorname{Im} p\left(t_{\delta_{j}}, 0, \varepsilon_{n}\right)$. This completes the proof.

Although it will not be needed here, we note that if [ $a^{\prime}, b^{\prime}$ ] is the interval given by Lemma 2.5 and $a^{\prime}<b^{\prime}$, then in addition to (i) and (ii) we also have
(iii) there exists a $\delta>0$ such that $\operatorname{Im} q p(\gamma(s)) \leq 0 \leq \operatorname{Im} q p(\gamma(t))$ for all $a^{\prime}-\delta<s<a^{\prime}$ and $b^{\prime}<t<b^{\prime}+\delta$. Indeed, the infimum $L\left(0, \xi^{0}\right)=b^{\prime}-a^{\prime}$ would otherwise satisfy $L\left(0, \xi^{0}\right)<\delta$ for every $\delta$ in view of (ii), which is a contradiction when $a^{\prime}<b^{\prime}$.

We next recall the definition of a one-dimensional bicharacteristic.
Definition 2.6. A one-dimensional bicharacteristic of the pseudodifferential operator with homogeneous principal symbol $p$ is a $C^{1}$ map $\gamma: I \rightarrow T^{*}(X) \backslash 0$, where $I$ is an interval on $\mathbb{R}$, such that
(i) $p(\gamma(t))=0$ for $t \in I$,
(ii) $0 \neq \gamma^{\prime}(t)=c(t) H_{p}(\gamma(t))$ if $t \in I$
for some continuous function $c: I \rightarrow \mathbb{C}$.
Let $P$ be an operator of principal type on a $C^{\infty}$ manifold $X$ with principal symbol $p$, and suppose $p$ fails to satisfy condition $(\Psi)$ in $X$. By (1-5) there is a function $q$ in $C^{\infty}\left(T^{*}(X) \backslash 0\right)$ such that $\operatorname{Im} q p$
changes sign from - to + on a bicharacteristic $\gamma$ of $\operatorname{Re} q p$, where $q \neq 0$. As can be seen in [Hörmander 1985b, pages 96-97], we can then find a compact one-dimensional bicharacteristic interval $\Gamma \subset \gamma$ or a characteristic point $\Gamma \in \gamma$ such that the sign change occurs on bicharacteristics of $\operatorname{Re} q p$ arbitrarily close to $\Gamma$. What we mean by this will be clear from the following discussion, although we will not use this terminology in the sequel. By the proof of Lemma 2.5 it suffices to regard the case that $q=1, X=\mathbb{R}^{n}$, $p$ is homogeneous of degree 1 with $\operatorname{Re} p=\xi_{1}$, and the bicharacteristic of $\operatorname{Re} p$ is given by (2-6).

We shall now study a slightly more general situation is some detail. If $\gamma=I \times\left\{w_{0}\right\}$, where $I=[a, b]$, we shall by $|\gamma|$ denote the usual arc length in $\mathbb{R}^{2 n}$, so that $|\gamma|=b-a$. Furthermore, we will assume that all curves are bicharacteristics of $\operatorname{Re} p=\xi_{1}$, that is, $w_{0}=\left(x^{\prime}, 0, \xi^{\prime}\right) \in \mathbb{R}^{2 n-1}$. We owe parts of this exposition to Nils Dencker.

Lemma 2.7. Assume that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma=[a, b] \times\left\{w_{0}\right\}$. Then for any $\delta>0$ there exist $\varepsilon>0, a-\delta<s_{-}<a$ and $b<s_{+}<b+\delta$ such that $\pm \operatorname{Im} p\left(s_{ \pm}, w\right)>0$ for any $\left|w-w_{0}\right|<\varepsilon$.

Proof. Since $t \mapsto \operatorname{Im} p\left(t, w_{0}\right)$ strongly changes sign on $[a, b]$ we can find $s_{ \pm}$satisfying the conditions so that $\pm \operatorname{Im} p\left(s_{ \pm}, w_{0}\right)>0$. By continuity we can find $\varepsilon_{ \pm}>0$ so that $\pm \operatorname{Im} p\left(s_{ \pm}, w\right)>0$ for any $\left|w-w_{0}\right|<\varepsilon_{ \pm}$. The lemma now follows if we take $\varepsilon=\min \left(\varepsilon_{-}, \varepsilon_{+}\right)$.
Definition 2.8. Let $\gamma=[a, b] \times\left\{w_{0}\right\}$ and $\gamma_{j}=\left[a_{j}, b_{j}\right] \times\left\{w_{j}\right\}$. If $\liminf _{j \rightarrow \infty} a_{j} \geq a, \lim \sup _{j \rightarrow \infty} b_{j} \leq b$ and $\lim _{j \rightarrow \infty} w_{j}=w_{0}$, then we shall write $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$. If in addition $\lim _{j \rightarrow \infty} a_{j}=a$ and $\lim _{j \rightarrow \infty} b_{j}=b$ then we shall write $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$.
Definition 2.9. If $\gamma$ is a bicharacteristic of $\operatorname{Re} p=\xi_{1}$ and there exists a sequence $\left\{\gamma_{j}\right\}$ of bicharacteristics of $\operatorname{Re} p$ such that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{j}$ for all $j$ and $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$, we set

$$
\begin{equation*}
L_{p}(\gamma)=\inf _{\left\{\gamma_{j}\right\}}\left\{\liminf _{j \rightarrow \infty}\left|\gamma_{j}\right|: \gamma_{j} \rightarrow \gamma \text { as } j \rightarrow \infty\right\} \tag{2-7}
\end{equation*}
$$

where the infimum is taken over all such sequences. We shall write $L_{p}(\gamma) \geq 0$ to signify the existence of such a sequence $\left\{\gamma_{j}\right\}$.
Remark. The definition of $L_{p}(\gamma)$ corresponds to what is denoted by $L_{0}$ in [Hörmander 1985b, page 97], when $\gamma=[a, b] \times\left\{w_{0}\right\}$ is given by (2-6) and

$$
\begin{equation*}
\operatorname{Im} p\left(a, w_{0}\right)<0<\operatorname{Im} p\left(b, w_{0}\right) \tag{2-8}
\end{equation*}
$$

To prove this claim, we begin by showing that $L_{p}(\gamma) \leq L_{0}$, after having properly defined $L_{0}$. To this end, let $\tilde{\gamma}=[\tilde{a}, \tilde{b}] \times\{\tilde{w}\}$ be a bicharacteristic of $\operatorname{Re} p$ such that $\operatorname{Im} p$ changes sign on $\tilde{\gamma}$. For $w$ close to $w_{0}$ we set

$$
\mathscr{L}_{p}(\tilde{\gamma}, w)=\inf \{t-s: \tilde{a}<s<t<\tilde{b}, \operatorname{Im} p(\tilde{a}, w)<0<\operatorname{Im} p(\tilde{b}, w)\}
$$

(Using the notation in [Hörmander 1985b, page 97] we would have $\mathscr{L}_{p}(\gamma, w)=L\left(x^{\prime}, \xi^{\prime}\right)$ if $w=$ $\left(x^{\prime}, 0, \xi^{\prime}\right)$.) Then

$$
L_{0}=\liminf _{w \rightarrow w_{0}} \mathscr{L}_{p}(\gamma, w)
$$

By an adaptation of the arguments in [Hörmander 1985b, page 97] it follows from the definition of $L_{0}$ that we can find a sequence $\left\{\gamma_{j}\right\}$ of bicharacteristics of $\operatorname{Re} p$ with $\gamma_{j}=\left[a_{j}, b_{j}\right] \times\left\{w_{j}\right\}$ such that

$$
\operatorname{Im} p\left(a_{j}, w_{j}\right)<0<\operatorname{Im} p\left(b_{j}, w_{j}\right) \quad \text { for all } j
$$

where $\lim w_{j}=w_{0}$ and the limits $a_{0}=\lim a_{j}$ and $b_{0}=\lim b_{j}$ exist, belong to the interval $(a, b)$ and satisfy $b_{0}-a_{0}=L_{0}$. If we for each $j$ apply Lemma 2.5 to $\gamma_{j}$ we obtain a sequence of bicharacteristics $\Gamma_{j} \subset \gamma_{j}$ of Re $p$ such that $\operatorname{Im} p$ strongly changes sign from - to + on $\Gamma_{j}$, where $\left|\Gamma_{j}\right|=\mathscr{L}_{p}\left(\gamma_{j}, w_{j}\right)<\left|\gamma_{j}\right|$. Clearly $\Gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$. Since $a<a_{j} \leq b_{j}<b$ if $j$ is sufficiently large it follows that for such $j$ we have $\mathscr{L}_{p}\left(\gamma, w_{j}\right) \leq \mathscr{L}_{p}\left(\gamma_{j}, w_{j}\right)$ by definition. This implies

$$
\begin{align*}
L_{0} & =\liminf _{w \rightarrow w_{0}} \mathscr{L}_{p}(\gamma, w) \leq \liminf _{j \rightarrow \infty} \mathscr{L}_{p}\left(\gamma, w_{j}\right) \\
& \leq \liminf _{j \rightarrow \infty}\left|\Gamma_{j}\right| \leq \limsup _{j \rightarrow \infty}\left|\Gamma_{j}\right| \leq \lim _{j \rightarrow \infty}\left|\gamma_{j}\right|=L_{0} \tag{2-9}
\end{align*}
$$

so $\left|\Gamma_{j}\right| \rightarrow L_{0}$ as $j \rightarrow \infty$. Thus $L_{p}(\gamma) \leq L_{0}$.
For the reversed inequality, suppose $\left\{\tilde{\gamma}_{j}\right\}$ is any sequence satisfying the properties of Definition 2.9, with $\tilde{\gamma}_{j}=\left[\tilde{a}_{j}, \tilde{b}_{j}\right] \times\left\{\tilde{w}_{j}\right\}$. By assumption we have $\operatorname{Im} p\left(\tilde{a}_{j}, \tilde{w}_{j}\right)=\operatorname{Im} p\left(\tilde{b}_{j}, \tilde{w}_{j}\right)=0$ for all $j$, which together with (2-8) and a continuity argument implies the existence of a positive integer $j_{0}$ such that

$$
a<\tilde{a}_{j} \leq \tilde{b}_{j}<b \quad \text { for all } j \geq j_{0}
$$

If $\tilde{\gamma}_{j, \delta}=\left[\tilde{a}_{j}-\delta, \tilde{b}_{j}+\delta\right] \times\left\{\tilde{w}_{j}\right\}$, this means that for small $\delta>0$ and sufficiently large $j$ we have

$$
\mathscr{L}_{p}\left(\gamma, \tilde{w}_{j}\right) \leq \mathscr{L}_{p}\left(\tilde{\gamma}_{j, \delta}, \tilde{w}_{j}\right)
$$

Since $\operatorname{Im} p$ strongly changes sign from - to + on $\tilde{\gamma}_{j}$, the infimum in the right side exists for every $\delta>0$ and is bounded from above by $\tilde{b}_{j}-\tilde{a}_{j}+2 \delta$. Taking the limit as $\delta \rightarrow 0$ yields $\mathscr{L}_{p}\left(\gamma, \tilde{w}_{j}\right) \leq\left|\tilde{\gamma}_{j}\right|$. Since $\tilde{w}_{j} \rightarrow w_{0}$ as $j \rightarrow \infty$ the definition of $L_{0}$ now gives

$$
\begin{equation*}
L_{0} \leq \liminf _{j \rightarrow \infty} \mathscr{L}_{p}\left(\gamma, \tilde{w}_{j}\right) \leq \liminf _{j \rightarrow \infty}\left|\tilde{\gamma}_{j}\right| \tag{2-10}
\end{equation*}
$$

and since the sequence $\left\{\tilde{\gamma}_{j}\right\}$ was arbitrary, we obtain $L_{0} \leq L_{p}(\gamma)$ by Definition 2.9. This proves the claim.

When no confusion can occur we will omit the dependence on $p$ in Definition 2.9. We note that if $L_{p}(\gamma)$ exists, then $L_{p}(\gamma) \leq|\gamma|$ by definition. Also, if $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma$ then Lemma 2.7 implies that the conditions of Definition 2.9 are satisfied. This proves the first part of the following result.
Corollary 2.10. Let $\gamma=[a, b] \times\left\{w_{0}\right\}$ be a bicharacteristic of $\operatorname{Re} p=\xi_{1}$. If $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma$, then $0 \leq L_{p}(\gamma) \leq|\gamma|$. Moreover, for every $\delta, \varepsilon>0$ there exists a bicharacteristic $\tilde{\gamma}=\tilde{\gamma}_{\delta, \varepsilon}$ of $\operatorname{Re} p$ with

$$
\tilde{\gamma}=[\tilde{a}, \tilde{b}] \times\{\tilde{w}\}, \quad a-\varepsilon<\tilde{a} \leq \tilde{b}<b+\varepsilon, \quad\left|\tilde{w}-w_{0}\right|<\varepsilon
$$

such that $\operatorname{Im} p$ strongly changes sign from - to + on $\tilde{\gamma}$ and $|\tilde{\gamma}|<L_{p}(\gamma)+\delta$.

Proof. The existence of the sequence $\left\{\Gamma_{j}\right\}$ in the preceding remark can after some adjustments be used to prove the second part of Corollary 2.10, but we prefer the following direct proof.

Given $\delta>0$ we can by Definition 2.9 find a sequence $\gamma_{j}=\left[a_{j}, b_{j}\right] \times\left\{w_{j}\right\}$ of bicharacteristics of $\operatorname{Re} p$ such that $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty, \operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{j}$ and $\liminf \operatorname{inc}_{j \rightarrow \infty}\left|\gamma_{j}\right|<$ $L(\gamma)+\delta$. After reducing to a subsequence we may assume $\left|\gamma_{j}\right|<L(\gamma)+\delta$ for all $j$. We have $\liminf _{j \rightarrow \infty} a_{j} \geq a$, so for every $\varepsilon$ there exists a $j_{1}(\varepsilon)$ such that $a_{j}>a-\varepsilon$ for all $j \geq j_{1}$. Similarly there exists a $j_{2}(\varepsilon)$ such that $b_{j}<b+\varepsilon$ for all $j \geq j_{2}$. Also, $w_{j} \rightarrow w_{0}$ as $j \rightarrow \infty$, so there exists a $j_{3}(\varepsilon)$ such that $\left|w_{j}-w_{0}\right|<\varepsilon$ for all $j \geq j_{3}$. Hence we can take $\tilde{\gamma}=\gamma_{j_{0}}$, where $j_{0}=\max \left(j_{1}, j_{2}, j_{3}\right)$.

Consider now the general case when $\operatorname{Im} q p$ changes sign from - to + on a bicharacteristic $\gamma \subset$ $T^{*}(X) \backslash 0$ of $\operatorname{Re} q p$, where $q \neq 0$, that is, (2-5) holds. In view of the proof of Lemma 2.5 we can by means of (2-7) define a minimality property of a subset of the curve $\gamma$ in the following sense.

Definition 2.11. Let $I \subset \mathbb{R}$ be a compact interval possibly reduced to a point and let $\tilde{\gamma}: I \rightarrow T^{*}(X) \backslash 0$ be a characteristic point or a compact one-dimensional bicharacteristic interval of the homogeneous function $p \in C^{\infty}\left(T^{*}(X) \backslash 0\right)$. Suppose that there exists a function $q \in C^{\infty}\left(T^{*}(X) \backslash 0\right)$ and a $C^{\infty}$ homogeneous canonical transformation $\chi$ from an open conic neighborhood $V$ of

$$
\Gamma=\left\{\left(x_{1}, 0, \varepsilon_{n}\right): x_{1} \in I\right\} \subset T^{*}\left(\mathbb{R}^{n}\right)
$$

to an open conic neighborhood $\chi(V) \subset T^{*}(X) \backslash 0$ of $\tilde{\gamma}(I)$ such that
(i) $\chi\left(x_{1}, 0, \varepsilon_{n}\right)=\tilde{\gamma}\left(x_{1}\right)$ and $\operatorname{Re} \chi^{*}(q p)=\xi_{1}$ in $V$,
(ii) $L_{\chi^{*}(q p)}(\Gamma)=|\Gamma|$.

Then we say that $\tilde{\gamma}(I)$ is a minimal characteristic point or a minimal bicharacteristic interval if $|I|=0$ or $|I|>0$, respectively.

The definition of the arclength is of course dependent of the choice of Riemannian metric on $T^{*}\left(\mathbb{R}^{n}\right)$. However, since we are only using the arclength to compare curves where one is contained within the other and both are parametrizable through condition (i), the results here and Definition 2.11 in particular are independent of the chosen metric. By choosing a Riemannian metric on $T^{*}(X)$, one could therefore define the minimality property given by Definition 2.11 through the corresponding arclength in $T^{*}(X)$ directly, although there, the notion of convergence of curves is somewhat trickier. We shall not pursue this any further.

Note that condition (i) implies that $q \neq 0$ and $\operatorname{Re} H_{q p} \neq 0$ on $\tilde{\gamma}$, and that by definition, a minimal bicharacteristic interval is a compact one-dimensional bicharacteristic interval. Moreover, if $\operatorname{Im} q p$ changes sign from - to + on a bicharacteristic $\gamma \subset T^{*}(X) \backslash 0$ of $\operatorname{Re} q p$, where $q \neq 0$, then we can always find a minimal characteristic point $\tilde{\gamma} \in \gamma$ or a minimal bicharacteristic interval $\tilde{\gamma} \subset \gamma$. In view of the proof of Lemma 2.5, this follows from the conclusion of the extensive remark beginning on page 432 together with (2-9). The following proposition shows that this continues to hold even when the assumption (2-5) is relaxed in the sense of Definition 2.9. We will state this result only in the (very weak) generality needed here.

Proposition 2.12. Let $\gamma=[a, b] \times\left\{w_{0}\right\}$ be a bicharacteristic of $\operatorname{Re} p=\xi_{1}$, and assume that $L(\gamma) \geq 0$. Then there exists a minimal characteristic point $\Gamma \in \gamma$ of $p$ or a minimal bicharacteristic interval $\Gamma \subset \gamma$ of $p$ of length $L(\gamma)$ if $L(\gamma)=0$ or $L(\gamma)>0$, respectively. If $\Gamma=\left[a_{0}, b_{0}\right] \times\left\{w_{0}\right\}$ and $a_{0}<b_{0}$, that is, $L(\gamma)>0$, then

$$
\begin{equation*}
\operatorname{Im} p_{(\alpha)}^{(\beta)}\left(t, w_{0}\right)=0 \tag{2-11}
\end{equation*}
$$

for all $\alpha, \beta$ with $\beta_{1}=0$ if $a_{0} \leq t \leq b_{0}$. Conversely, if $\gamma$ is a minimal characteristic point or a minimal bicharacteristic interval then $L(\gamma)=|\gamma|$.

Lemma 2.13. Let $\gamma$ and $\gamma_{j}$ for $j \geq 1$ be bicharacteristics of $\operatorname{Re} p=\xi_{1}$, and assume that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{j}$ for each $j$. If $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$ then $L(\gamma) \leq \liminf _{j \rightarrow \infty} L\left(\gamma_{j}\right)$.

Proof. Let $\gamma_{j}=\left[a_{j}, b_{j}\right] \times\left\{w_{j}\right\}$ and $\gamma=[a, b] \times\left\{w_{0}\right\}$. Since $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{j}$ we can by Corollary 2.10 for each $j$ find a bicharacteristic $\tilde{\gamma}_{j}=\left[\tilde{a}_{j}, \tilde{b}_{j}\right] \times\left\{\tilde{w}_{j}\right\}$ of $\operatorname{Re} p$ with

$$
a_{j}-1 / j<\tilde{a}_{j} \leq \tilde{b}_{j}<b_{j}+1 / j \quad \text { and } \quad\left|\tilde{w}_{j}-w_{j}\right|<1 / j
$$

such that $\operatorname{Im} p$ strongly changes sign from - to + on $\tilde{\gamma}_{j}$ and $\left|\tilde{\gamma}_{j}\right|<L\left(\gamma_{j}\right)+1 / j$. Now $\left|\tilde{w}_{j}-w_{0}\right| \leq$ $\left|\tilde{w}_{j}-w_{j}\right|+\left|w_{j}-w_{0}\right|$, and since $\liminf _{j \rightarrow \infty} \tilde{a}_{j} \geq \liminf _{j \rightarrow \infty}\left(a_{j}-1 / j\right) \geq a$ and correspondingly for $\tilde{b}_{j}$, we find that $\tilde{\gamma}_{j} \rightarrow \gamma$ as $j \rightarrow \infty$. Thus

$$
L(\gamma) \leq \liminf _{j \rightarrow \infty}\left|\tilde{\gamma}_{j}\right| \leq \liminf _{j \rightarrow \infty}\left(L\left(\gamma_{j}\right)+1 / j\right)
$$

Proof of Proposition 2.12. We may without loss of generality assume that $w_{0}=\left(0, \varepsilon_{n}\right) \in \mathbb{R}^{2 n-1}$. The last statement is then an immediate consequence of Definition 2.11. To prove the theorem it then also suffices to show that we can find a characteristic point $\Gamma \in \gamma$ of $p$, or a compact one-dimensional bicharacteristic interval $\Gamma \subset \gamma$ of $p$ of length $L(\gamma)$, with the property that in any neighborhood of $\Gamma$ there is a bicharacteristic of $\operatorname{Re} p$ where $\operatorname{Im} p$ strongly changes sign from - to + . This is done by adapting the arguments in [Hörmander 1985b, page 97], which also yields (2-11).

For small $\delta>0$ we can find $\varepsilon(\delta)$ with $0<\varepsilon<\delta$ such that $L(\tilde{\gamma})>L(\gamma)-\delta / 2$ for any bicharacteristic $\tilde{\gamma}=[\tilde{a}, \tilde{b}] \times\{\tilde{w}\}$ with $a-\varepsilon<\tilde{a} \leq \tilde{b}<b+\varepsilon$ and $\left|\tilde{w}-w_{0}\right|<\varepsilon$ such that $\operatorname{Im} p$ strongly changes sign from - to + on $\tilde{\gamma}$. Indeed, otherwise there would exist a $\delta>0$ such that for each (sufficiently large) $k$ we can find a bicharacteristic $\gamma_{k}=\left[a_{k}, b_{k}\right] \times\left\{w_{k}\right\}$ with $a-1 / k<a_{k} \leq b_{k}<b+1 / k$ and $\left|w_{k}-w_{0}\right|<1 / k$ such that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{k}$ and $L\left(\gamma_{k}\right) \leq L(\gamma)-\delta / 2$. This implies that $\gamma_{k} \rightarrow \gamma$ as $k \rightarrow \infty$, so by Lemma 2.13 we obtain

$$
L(\gamma) \leq \liminf _{k \rightarrow \infty} L\left(\gamma_{k}\right) \leq L(\gamma)-\delta / 2
$$

a contradiction. Since $L(\gamma) \geq 0$ we have by Corollary 2.10 for some $\left|w_{\delta}-w_{0}\right|<\varepsilon$ and $a-\varepsilon<a_{\delta} \leq$ $b_{\delta}<b+\varepsilon$ with $w_{\delta}=\left(x_{\delta}^{\prime}, 0, \xi_{\delta}^{\prime}\right)$ that $\operatorname{Im} p$ strongly changes sign from - to + on the bicharacteristic $\gamma_{\delta}=\left[a_{\delta}, b_{\delta}\right] \times\left\{w_{\delta}\right\}$, and $\left|\gamma_{\delta}\right|<L(\gamma)+\delta / 4$. Thus,

$$
\begin{equation*}
L(\gamma)-\delta / 2<\left|\gamma_{\delta}\right|<L(\gamma)+\delta / 4 \tag{2-12}
\end{equation*}
$$

We claim that $\operatorname{Im} p$ and all derivatives with respect to $x^{\prime}$ and $\xi^{\prime}$ must vanish at $\left(t, w_{\delta}\right)$ if $a_{\delta}+\delta<t<b_{\delta}-\delta$. Indeed, by Lemma 2.7 we can find a $\rho>0, a_{\delta}-\delta / 4<s_{-}<a_{\delta}$ and $b_{\delta}<s_{+}<b_{\delta}+\delta / 4$ such that

$$
\operatorname{Im} p\left(s_{-}, w\right)<0<\operatorname{Im} p\left(s_{+}, w\right) \quad \text { for all }\left|w-w_{\delta}\right|<\rho
$$

If $\operatorname{Im} p$ and all derivatives with respect to $x^{\prime}$ and $\xi^{\prime}$ do not vanish at $\left(t, w_{\delta}\right)$ if $a_{\delta}+\delta<t<b_{\delta}-\delta$, then we can choose $w=\left(x^{\prime}, 0, \xi^{\prime}\right)$ so that $\left|w-w_{\delta}\right|<\rho,\left|w-w_{0}\right|<\varepsilon$ and $\operatorname{Im} p(t, w) \neq 0$ for some $a_{\delta}+\delta<t<b_{\delta}-\delta$. It follows that the required sign change of $\operatorname{Im} p\left(x_{1}, w\right)$ must occur on one of the intervals ( $\left.s_{-}, t\right)$ and $\left(t, s_{+}\right)$, which are shorter than $L(\gamma)-\delta / 2$. This contradiction proves the claim.

Now choose a sequence $\delta_{j} \rightarrow 0$ as $j \rightarrow \infty$ such that $\lim a_{\delta_{j}}$ and $\lim b_{\delta_{j}}$ exist. If we denote these limits by $a_{0}$ and $b_{0}$, respectively, then $L(\gamma)=b_{0}-a_{0}$ by (2-12), and (2-11) holds if $a_{0}<b_{0}$. In particular, if $a_{0}<b_{0}$ then

$$
H_{p}(\gamma(t))=\left(1+i \partial \operatorname{Im} p(\gamma(t)) / \partial \xi_{1}\right) \gamma^{\prime}(t) \quad \text { for } a_{0} \leq t \leq b_{0}
$$

so if $\Gamma=\left\{\left(t, w_{0}\right): t \in\left[a_{0}, b_{0}\right]\right\}$ then $\Gamma$ is a compact one-dimensional bicharacteristic interval of $p$ with the function $c$ in Definition 2.6 given by

$$
c(t)=\left(1+i \partial \operatorname{Im} p(\Gamma(t)) / \partial \xi_{1}\right)^{-1}
$$

Proposition 2.12 allows us to make some additional comments on the implications of Definition 2.11. With the notation in the definition, we note that condition (ii) implies that there exists a sequence $\left\{\Gamma_{j}\right\}$ of bicharacteristics of $\operatorname{Re} \chi^{*}(q p)$ on which $\operatorname{Im} \chi^{*}(q p)$ strongly changes sign from - to + , such that $\Gamma_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$. By our choice of terminology, the sequence $\left\{\Gamma_{j}\right\}$ may simply be a sequence of points when $L(\Gamma)=0$. Conversely, if $\left\{\Gamma_{j}\right\}$ is a point sequence then $L(\Gamma)=0$. Also note that if $\tilde{\gamma}(I)$ is minimal, and condition (i) in Definition 2.11 is satisfied for some other choice of maps $q^{\prime}$ and $\chi^{\prime}$, then condition (ii) also holds for $q^{\prime}$ and $\chi^{\prime}$; in other words,

$$
L_{\chi^{*}(q p)}(\Gamma)=|\Gamma|=L_{\left(\chi^{\prime}\right)^{*}\left(q^{\prime} p\right)}(\Gamma)
$$

This follows by an application of Proposition 2.12 together with [Hörmander 1985b, Lemma 26.4.10]. It is then also clear that $\tilde{\gamma}(I)$ is a minimal characteristic point or a minimal bicharacteristic interval of the homogeneous function $p \in C^{\infty}\left(T^{*}(X) \backslash 0\right)$ if and only if $\Gamma(I)$ is a minimal characteristic point or a minimal bicharacteristic interval of $\chi^{*}(q p) \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$ for any maps $q$ and $\chi$ satisfying condition (i) in Definition 2.11.

The proof of [Hörmander 1985b, Theorem 26.4.7] stating that condition $(\Psi)$ is necessary for local solvability relies on the imaginary part of the principal symbol satisfying (2-11). By Proposition 2.12, it is clear that $(2-11)$ holds on a minimal bicharacteristic interval $\Gamma$ in the case $q=1$ and $\operatorname{Re} p=\xi_{1}$. However, we shall require that we can find bicharacteristics arbitrarily close to $\Gamma$ for which the following stronger result is applicable, at least if $\operatorname{Im} p$ does not depend on $\xi_{1}$ as is the case for the standard normal form. This will be made precise below.

Proposition 2.14. Let $p=\xi_{1}+i \operatorname{Im} p$. Assume that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma=$ $[a, b] \times\{w\}$ and that $L(\gamma) \geq|\gamma|-\varrho$ for some $0<\varrho<|\gamma| / 2$. If $\operatorname{Im} p$ does not depend on $\xi_{1}$ then for any $\kappa>\varrho$ we find that $\operatorname{Im} p$ vanishes identically in a neighborhood of $I_{\kappa} \times\{w\}$, where $I_{\kappa}=[a+\kappa, b-\kappa]$.

The statement would of course be void if the hypotheses hold only for $\varrho \geq|\gamma| / 2$, for then $I_{\kappa}=\varnothing$.
Proof. If the statement is false, there exists a $\kappa>0$ such that $\operatorname{Im} p \not \equiv 0$ near $I_{\kappa} \times\{w\}$. Thus there exists a sequence $\left(s_{j}, w_{j}\right) \rightarrow I_{\kappa} \times\{w\}$ such that $\operatorname{Im} p\left(s_{j}, w_{j}\right) \neq 0$ for all $j$. Since $\operatorname{Im} p$ does not depend on $\xi_{1}$ we can choose $w_{j}$ to have $\xi_{1}$ coordinate equal to zero for all $j$, so that $\left(s_{j}, w_{j}\right)$ is contained in a bicharacteristic of Re $p$. We may choose a subsequence so that for some $s \in I_{\kappa}$ we have $\left|s_{j}-s\right| \rightarrow 0$ and $\left|w_{j}-w\right| \rightarrow 0$ monotonically, and either $\operatorname{Im} p\left(s_{j}, w_{j}\right)>0$ or $-\operatorname{Im} p\left(s_{j}, w_{j}\right)>0$ for all $j$. We shall consider the case with positive sign, the negative case works similarly.

Choose $\delta<(\kappa-\varrho) / 3$ and use Lemma 2.7. We find that there exists $a-\delta<s_{-}<a$ and $\varepsilon>0$ such that $\operatorname{Im} p\left(s_{-}, v\right)<0$ for any $|v-w|<\varepsilon$. Choose $k>0$ so that $\left|s_{j}-s\right|<\delta$ and $\left|w_{j}-w\right|<\varepsilon$ when $j>k$. Then $t \mapsto \operatorname{Im} p\left(t, w_{j}\right)$ changes sign from - to + on $I_{j}=\left[s_{-}, s_{j}\right]$, which has length

$$
\left|I_{j}\right|=s_{j}-s_{-} \leq\left|s_{j}-s\right|+s-a+a-s_{-}<|\gamma|-\kappa+2 \delta<|\gamma|-\varrho-\delta .
$$

If we for each $j$ apply Lemma 2.5 to $I_{j} \times\left\{w_{j}\right\}$ and let $j \rightarrow \infty$ we obtain a contradiction to the hypothesis $L(\gamma) \geq|\gamma|-\varrho$.

One could state Proposition 2.14 without the condition that the imaginary part is independent of $\xi_{1}$. The invariant statement would then be that the restriction of the imaginary part to the characteristic set of the real part vanishes in a neighborhood of $\gamma$.

The fact that Proposition 2.14 assumes that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma$ means that the conditions are not in general satisfied when $\gamma$ is a minimal bicharacteristic interval. As mentioned above, we will instead show that arbitrarily close to a minimal bicharacteristic interval one can always find bicharacteristics for which Proposition 2.14 is applicable. Before we state the results we introduce a helpful definition together with some (perhaps contrived but illustrative) examples.

Definition 2.15. A minimal bicharacteristic interval $\Gamma=\left[a_{0}, b_{0}\right] \times\left\{w_{0}\right\} \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ of the homogeneous function $p=\xi_{1}+i \operatorname{Im} p$ of degree 1 is said to be $\varrho$-minimal if there exists a $\varrho \geq 0$ such that $\operatorname{Im} p$ vanishes in a neighborhood of $\left[a_{0}+\kappa, b_{0}-\kappa\right] \times\left\{w_{0}\right\}$ for any $\kappa>\varrho$.

By a 0-minimal bicharacteristic interval $\Gamma$ we thus mean a minimal bicharacteristic interval such that the imaginary part vanishes in a neighborhood of any proper closed subset of $\Gamma$. Note that this does not hold for minimal bicharacteristic intervals in general.

Example 2.16. Let $f \in C^{\infty}(\mathbb{R})$ be given by

$$
f(t)= \begin{cases}-e^{-1 / t^{2}} & \text { if } t<0  \tag{2-13}\\ 0 & \text { if } 0 \leq t \leq 2 \\ e^{-1 /(t-2)^{2}} & \text { if } t>2\end{cases}
$$

and let $\phi \in C^{\infty}(\mathbb{R})$ be a smooth cutoff function with $\operatorname{supp} \phi=[0,2]$ such that $\phi>0$ on $(0,2)$. If $\xi=\left(\xi_{1}, \xi^{\prime}\right)$ then

$$
p_{1}(x, \xi)=\xi_{1}+i\left|\xi^{\prime}\right|\left(f\left(x_{1}\right)+x_{2} \phi\left(x_{1}\right)\right)
$$

is homogeneous of degree 1 . If we write $x=\left(x_{1}, x_{2}, x^{\prime \prime}\right)$ then for any fixed $\left(x^{\prime \prime}, \xi^{\prime}\right) \in \mathbb{R}^{n-2} \times \mathbb{R}^{n-1}$ with $\xi^{\prime} \neq 0$ we find that $\left\{\left(x_{1}, x_{2}, x^{\prime \prime}, 0, \xi^{\prime}\right): x_{1}=a, x_{2}=c\right\}$ is a minimal characteristic point of $p_{1}$ if $c \geq 0$ and $a=0$ or if $c \leq 0$ and $a=2$. Note that if $\xi^{\prime} \neq 0$ then $\operatorname{Im} p_{1}$ changes sign from - to + on the bicharacteristic $\gamma\left(x_{1}\right)=\left\{\left(x_{1}, 0, x^{\prime \prime}, 0, \xi^{\prime}\right)\right\}$ of Re $p_{1}$, but that none of the points $\left\{\gamma\left(x_{1}\right): 0<x_{1}<2\right\}$ are minimal characteristic points. ${ }^{1}$ On the other hand, if $f$ is given by (2-13) let

$$
h\left(x, \xi^{\prime}\right)= \begin{cases}\left|\xi^{\prime}\right| f\left(x_{1}-1\right) e^{1 / x_{2}} & \text { if } x_{2}<0 \\ 0 & \text { if } x_{2}=0 \\ \left|\xi^{\prime}\right| f\left(x_{1}\right) e^{-1 / x_{2}} & \text { if } x_{2}>0\end{cases}
$$

be the imaginary part of $p_{2}(x, \xi)$. If $\operatorname{Re} p_{2}=\xi_{1}$ then $p_{2}$ is homogeneous of degree 1 and

$$
\Gamma_{c}=\left\{\left(x_{1}, x_{2}, x^{\prime \prime}, 0, \xi^{\prime}\right): x_{2}=c, x_{1} \in I_{c}\right\}
$$

is a minimal bicharacteristic interval of $p_{2}$ for any $\left(x^{\prime \prime}, \xi^{\prime}\right) \in \mathbb{R}^{n-2} \times \mathbb{R}^{n-1}$ with $\xi^{\prime} \neq 0$ if $c \geq 0$ and $I_{c}=[0,2]$ or if $c \leq 0$ and $I_{c}=[1,3]$. Moreover, if $c \lessgtr 0$ then $\Gamma_{c}$ is a 0 -minimal bicharacteristic interval. However, there is no $\varrho>0$ such that the minimal bicharacteristic interval $\Gamma=\left\{\left(x_{1}, 0, x^{\prime \prime}, 0, \xi^{\prime}\right): 0 \leq x_{1} \leq 2\right\}$ is $\varrho$-minimal. The same holds for the minimal bicharacteristic interval $\tilde{\Gamma}=\left\{\left(x_{1}, 0, x^{\prime \prime}, 0, \xi^{\prime}\right): 1 \leq x_{1} \leq 3\right\}$. Figure 1 shows a cross-section of the characteristic sets of $\operatorname{Im} p_{1}$ and $\operatorname{Im} p_{2}$.

Lemma 2.17. Let $p=\xi_{1}+i \operatorname{Im} p$, and assume that $L(\gamma)>0$ and that $\operatorname{Im} p$ does not depend on $\xi_{1}$. Then one can find $\tilde{\gamma}_{j} \subset \gamma_{j} \rightarrow \gamma$ such that $\left|\tilde{\gamma}_{j}\right| \rightarrow L(\gamma)$, Im $p$ strongly changes sign from - to + on $\gamma_{j}$ and $\operatorname{Im} p$ vanishes in a neighborhood of $\tilde{\gamma}_{j}$.

[^1]

Figure 1. Cross-sections of the characteristic sets of $\operatorname{Im} p_{1}$ and $\operatorname{Im} p_{2}$, respectively.

Note that the conditions imply that $\tilde{\gamma}_{j} \rightarrow \gamma$ as $j \rightarrow \infty$.
Proof. Choose $\gamma_{j} \rightarrow \gamma$ when $j \rightarrow \infty$ as in the proof of Proposition 2.12, so that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{j}$ and $L(\gamma)=\lim _{j \rightarrow \infty}\left|\gamma_{j}\right|$. By Lemma 2.13 and Corollary 2.10 we have

$$
L(\gamma) \leq \liminf _{j \rightarrow \infty} L\left(\gamma_{j}\right) \leq \liminf _{j \rightarrow \infty}\left|\gamma_{j}\right|=L(\gamma)
$$

Thus we can for every $\varepsilon>0$ choose $j$ so that $\left|L(\gamma)-\left|\gamma_{j}\right|\right|<\varepsilon$ and $\left|L\left(\gamma_{j}\right)-\left|\gamma_{j}\right|\right|<\varepsilon$. If we choose $\varepsilon<L(\gamma) / 5$ then

$$
2 \varepsilon<(L(\gamma)-\varepsilon) / 2<\left|\gamma_{j}\right| / 2
$$

Hence, if $\gamma_{j}=\left[a_{j}, b_{j}\right] \times w_{j}$ then by using Proposition 2.14 on $\gamma_{j}$ we find that $\operatorname{Im} p$ vanishes identically in a neighborhood of $\tilde{\gamma}_{j}=\left[a_{j}+2 \varepsilon, b_{j}-2 \varepsilon\right] \times\left\{w_{j}\right\}$. Now choose a sequence $\varepsilon_{k} \rightarrow 0$ as $k \rightarrow \infty$. Then $\tilde{\gamma}_{j(k)} \subset \gamma_{j(k)}$ and assuming as we may that $j(k)>j\left(k^{\prime}\right)$ if $k>k^{\prime}$ we obtain $\left|\tilde{\gamma}_{j(k)}\right| \rightarrow L(\gamma)$ as $k \rightarrow \infty$, which completes the proof.

If $\Gamma \subset \gamma$ is a minimal bicharacteristic interval in $T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ of the homogeneous function $p=$ $\xi_{1}+i \operatorname{Im} p$ of degree 1 , where the imaginary part is independent of $\xi_{1}$, then by Definition 2.11 and Proposition 2.12 we have $0<|\Gamma|=L(\Gamma)$. By the proof of Lemma 2.17 there exists a sequence $\gamma_{j} \rightarrow \Gamma$ of bicharacteristics of $\operatorname{Re} p$ such that $\operatorname{Im} p$ strongly changes sign from - to + on $\gamma_{j}$ and vanishes identically in a neighborhood of a subinterval $\tilde{\gamma}_{j} \subset \gamma_{j}$. Moreover, $\tilde{\gamma}_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$. By Lemma 2.13 we have $L\left(\gamma_{j}\right)>0$ for sufficiently large $j$, so according to Proposition 2.12 we can for each such $j$ find a minimal bicharacteristic interval $\Gamma_{j} \subset \gamma_{j}$. We have $\gamma_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$ and since

$$
|\Gamma|=L(\gamma) \leq \liminf _{j \rightarrow \infty} L\left(\gamma_{j}\right)=\liminf _{j \rightarrow \infty}\left|\Gamma_{j}\right| \leq \limsup _{j \rightarrow \infty}\left|\Gamma_{j}\right| \leq \lim _{j \rightarrow \infty}\left|\gamma_{j}\right|=|\Gamma|
$$

it follows that $\Gamma_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$. Since also $\tilde{\gamma}_{j} \subset \gamma_{j}$ and $\tilde{\gamma}_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$, the intersection $\tilde{\gamma}_{j} \cap \Gamma_{j}$ must be nonempty for large $j$. For such $j$ it follows that $\tilde{\gamma}_{j}$ must be a proper subinterval of $\Gamma_{j}$, for if not, this would contradict the fact that $\Gamma_{j}$ is a minimal bicharacteristic interval. Hence we can find a sequence $\left\{\varrho_{j}\right\}$ of positive numbers with $\varrho_{j} \rightarrow 0$ as $j \rightarrow \infty$, such that $\Gamma_{j}$ is a $\varrho_{j}$-minimal bicharacteristic interval. We have thus proved the following theorem, which concludes our study of the bicharacteristics.

Theorem 2.18. If $\Gamma$ is a minimal bicharacteristic interval in $T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ of the homogeneous function $p=\xi_{1}+i \operatorname{Im} p$ of degree 1 , where the imaginary part is independent of $\xi_{1}$, then there exists a sequence $\left\{\Gamma_{j}\right\}$ of $\varrho_{j}$-minimal bicharacteristic intervals of $p$ such that $\Gamma_{j} \rightarrow \Gamma$ and $\varrho_{j} \rightarrow 0$ as $j \rightarrow \infty$.

We can now state our main theorem, which yields necessary conditions for inclusion relations between the ranges of operators that fail to be microlocally solvable.
Theorem 2.19. Let $K \subset T^{*}(X) \backslash 0$ be a compactly based cone. Let $P \in \Psi_{\mathrm{cl}}^{k}(X)$ and $Q \in \Psi_{\mathrm{cl}}^{k^{\prime}}(X)$ be properly supported pseudodifferential operators such that the range of $Q$ is microlocally contained in the range of $P$ at $K$, where $P$ is an operator of principal type in a conic neighborhood of $K$. Let $p_{k}$ be the homogeneous principal symbol of $P$, and let $I=\left[a_{0}, b_{0}\right] \subset \mathbb{R}$ be a compact interval possibly reduced to a point. Suppose that $K$ contains a conic neighborhood of $\gamma(I)$, where $\gamma: I \rightarrow T^{*}(X) \backslash 0$ is either
(a) a minimal characteristic point of $p_{k}$, or
(b) a minimal bicharacteristic interval of $p_{k}$ with injective regular projection in $S^{*}(X)$.

Then there exists a pseudodifferential operator $E \in \Psi_{\mathrm{cl}}^{k^{\prime}-k}(X)$ such that the terms in the asymptotic sum of the symbol of $Q-P E$ have vanishing Taylor coefficients at $\gamma(I)$.

The hypotheses of Theorem 2.19 imply that $P$ is not solvable at the cone $K$. Indeed, solvability at $K \subset T^{*}(X) \backslash 0$ implies solvability at any smaller closed cone, and in view of Definition 2.11 it follows by [Hörmander 1985b, Theorem 26.4.7'] together with [Hörmander 1985b, Proposition 26.4.4] that $P$ is not solvable at the cone generated by $\gamma(I)$. Conversely, suppose that $P$ is an operator of principal type that is not microlocally solvable in any neighborhood of a point $\left(x_{0}, \xi_{0}\right) \in T^{*}(X) \backslash 0$. Then the principal symbol $p_{k}$ fails to satisfy condition (1-5) in every neighborhood of ( $x_{0}, \xi_{0}$ ) by [Dencker 2006, Theorem 1.1]. In view of the alternative version of condition (1-5) given by [Hörmander 1985b, Theorem 26.4.12], it is then easy to see using [Hörmander 1985a, Theorem 21.3.6] and [Hörmander 1985b, Lemma 26.4.10] that $\left(x_{0}, \xi_{0}\right)$ is a minimal characteristic point of $p_{k}$, so Theorem 2.19 applies there.

We also mention that if $P$ is of principal type and $\gamma$ is a minimal bicharacteristic interval of the principal symbol $p_{k}$ contained in a curve along which $p_{k}$ fails to satisfy condition (1-5), then $\gamma$ has injective regular projection in $S^{*}(X)$ by the proof of [Hörmander 1985b, Theorem 26.4.12].

Remark. As pointed out in the introduction, we cannot hope to obtain a result such as Theorem 2.19 for solvable nonelliptic operators in general. Indeed, Example 2.2 shows that if $X \subset \mathbb{R}^{n}$ is open, and $K \subset T^{*}(X) \backslash 0$ is a compactly based cone, then the range of $D_{2}$ is microlocally contained in the range of $D_{1}$ at $K$. If there were to exist a pseudodifferential operator $e(x, D) \in \Psi_{\mathrm{cl}}^{0}(X)$ such that all the terms in the symbol of $R(x, D)=D_{2}-D_{1} \circ e(x, D)$ have vanishing Taylor coefficients at a point $\left(x_{0}, \xi_{0}\right) \in K$ contained in a bicharacteristic of the principal symbol $\sigma\left(D_{1}\right)=\xi_{1}$ of $D_{1}$, then in particular this would hold for the principal symbol

$$
\sigma(R)(x, \xi)=\xi_{2}-\xi_{1} e_{0}(x, \xi)
$$

if $e_{0}$ denotes the principal symbol of $e(x, D)$. However, taking the $\xi_{2}$ derivative of the equation above and evaluating at $\left(x_{0}, \xi_{0}\right)$ then immediately yields the contradiction $0=1$ since $\left(x_{0}, \xi_{0}\right)$ belongs to the hypersurface $\xi_{1}=0$.

In the proof of the theorem we may assume that $P$ and $Q$ are operators of order 1. In fact, the discussion following Definition 2.1 shows that if the conditions of Theorem 2.19 hold and $Q_{1} \in \Psi_{\mathrm{cl}}^{k-k^{\prime}}(X)$ and $Q_{2} \in \Psi_{\mathrm{cl}}^{1-k}(X)$ are properly supported, then the range of $Q_{2} Q Q_{1} \in \Psi_{\mathrm{cl}}^{1}(X)$ is microlocally contained in the range of $Q_{2} P \in \Psi_{\mathrm{cl}}^{1}(X)$ at $K$. If the theorem holds for operators of the same order $k$ then there exists an operator $E \in \Psi_{\mathrm{cl}}^{0}(X)$ such that all the terms in the asymptotic expansion of the symbol of $Q Q_{1}-P E$ have vanishing Taylor coefficients at $\gamma(I)$. If we choose $Q_{1}$ to be elliptic, then we can find a parametrix $Q_{1}^{-1}$ of $Q_{1}$ such that

$$
Q-P E Q_{1}^{-1} \equiv\left(Q Q_{1}-P E\right) \circ Q_{1}^{-1} \quad \bmod \Psi^{-\infty}(X)
$$

has symbol

$$
\begin{equation*}
\sigma_{A \circ Q_{1}^{-1}}(x, \xi) \sim \sum \partial_{\xi}^{\alpha} \sigma_{A}(x, \xi) D_{x}^{\alpha} \sigma_{Q_{1}^{-1}}(x, \xi) / \alpha! \tag{2-14}
\end{equation*}
$$

with $A=Q Q_{1}-P E$. Clearly, all the terms in the asymptotic expansion of the symbol of $Q-P E Q_{1}^{-1}$ then have vanishing Taylor coefficients at $\gamma(I)$, and $E_{1}=E Q_{1}^{-1} \in \Psi_{\mathrm{cl}}^{k^{\prime}-k}(X)$, so the theorem holds with $E$ replaced by $E_{1}$. If the theorem holds for operators of order 1 we can choose $Q_{2}$ elliptic and use the same argument to show that if all the terms in the asymptotic expansion of the symbol of $Q_{2} Q Q_{1}-Q_{2} P E$ have vanishing Taylor coefficients at $\gamma(I)$, then the same holds for

$$
Q-P E Q_{1}^{-1} \equiv Q_{2}^{-1} \circ\left(Q_{2} Q Q_{1}-Q_{2} P E\right) \circ Q_{1}^{-1} \quad \bmod \Psi^{-\infty}(X)
$$

where $Q_{2}^{-1}$ is a parametrix of $Q_{2}$. Here we use the fact that if $\gamma(I)$ is a minimal characteristic point or a minimal bicharacteristic interval of the principal symbol of $P$, then this also holds for the principal symbol of $Q_{2} P$ by Definition 2.11.

For pseudodifferential operators, the property that all terms in the asymptotic expansion of the total symbol have vanishing Taylor coefficients is preserved under conjugation with Fourier integral operators associated with a canonical transformation (see Lemma A.1). Thus we will be able to prove Theorem 2.19 by local arguments and an application of Proposition 2.4.

Let $\gamma: I \rightarrow T^{*}(X) \backslash 0$, with $I=\left[a_{0}, b_{0}\right] \subset \mathbb{R}$, be the map given by Theorem 2.19. By using [Hörmander 1985a, Theorem 21.3.6] or [Hörmander 1985b, Theorem 26.4.13], when $\gamma$ is a characteristic point or a one-dimensional bicharacteristic, respectively, we can find a $C^{\infty}$ canonical transformation $\chi$ from a conic neighborhood of $\Gamma=\left\{\left(x, \varepsilon_{n}\right): x_{1} \in I, x^{\prime}=0\right\}$ in $T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ to a conic neighborhood of $\gamma(I)$ in $T^{*}(X) \backslash 0$ and a $C^{\infty}$ homogeneous function $b$ of degree 0 with no zero on $\gamma(I)$ such that $\chi\left(x_{1}, 0, \varepsilon_{n}\right)=\gamma\left(x_{1}\right)$ for $x_{1} \in I$ and

$$
\begin{equation*}
\chi^{*}\left(b p_{1}\right)=\xi_{1}+i f\left(x, \xi^{\prime}\right) \tag{2-15}
\end{equation*}
$$

where $f$ is real-valued, homogeneous of degree 1 and independent of $\xi_{1}$. Thus, by the hypotheses of Theorem 2.19 one can in any neighborhood of $\Gamma$ find an interval in the $x_{1}$ direction where $f$ changes sign from - to + for increasing $x_{1}$. Also, if $I$ is an interval then $f$ vanishes of infinite order on $\Gamma$ by (2-11), and by Theorem 2.18 there exists a sequence $\left\{\Gamma_{j}\right\}$ of $\varrho_{j}$-minimal bicharacteristics of $\chi^{*}\left(b p_{1}\right)$ such that $\varrho_{j} \rightarrow 0$ and $\Gamma_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$.

The existence of the canonical transformation $\chi$ together with Proposition 2.4 implies that we can find Fourier integral operators $A$ and $B$ such that the range of $B Q A$ is microlocally contained in the range of $B P A$ at a cone $K^{\prime}$ containing $\Gamma$, where the principal symbol of $B P A$ is given by (2-15). In view of Lemma A. 1 we may therefore reduce the proof to the case that $P, Q \in \Psi_{\mathrm{cl}}^{1}\left(\mathbb{R}^{n}\right)$ and the principal symbol $p$ of $P$ is given by (2-15). In accordance with the notation in Proposition 2.4 we will assume that the range of $Q$ is microlocally contained in the range of $P$ at a cone $K$ containing $\Gamma$, thus renaming $K^{\prime}$ to $K$. If

$$
\sigma_{Q}=q_{1}+q_{0}+\cdots
$$

is the asymptotic sum of homogeneous terms of the symbol of $Q$, we can then use the Malgrange preparation theorem (see [Hörmander 1983a, Theorem 7.5.6]) to find $e_{0}, r_{1} \in C^{\infty}$ near $\Gamma$ such that

$$
q_{1}(x, \xi)=\left(\xi_{1}+i f\left(x, \xi^{\prime}\right)\right) e_{0}(x, \xi)+r_{1}\left(x, \xi^{\prime}\right)
$$

where $r_{1}$ is independent of $\xi_{1}$. Restricting to $|\xi|=1$ and extending by homogeneity we can make $e_{0}$ and $r_{1}$ homogeneous of degree 0 and 1 , respectively. The term of degree 1 in the symbol of $Q-P \circ e_{0}(x, D)$ is $r_{1}\left(x, \xi^{\prime}\right)$. Again, by Malgrange's preparation theorem we can find $e_{-1}, r_{0} \in C^{\infty}$ near $\Gamma$ such that

$$
q_{0}(x, \xi)-\sigma_{0}\left(P \circ e_{0}(x, D)\right)(x, \xi)=\left(\xi_{1}+i f\left(x, \xi^{\prime}\right)\right) e_{-1}(x, \xi)+r_{0}\left(x, \xi^{\prime}\right)
$$

where $e_{-1}$ and $r_{0}$ are homogeneous of degree -1 and 0 , respectively, and $r_{0}$ is independent of $\xi_{1}$. The term of degree 0 in the symbol of

$$
Q-P \circ e_{0}(x, D)-P \circ e_{-1}(x, D)
$$

is $r_{0}\left(x, \xi^{\prime}\right)$. Repetition of the argument allows us to write

$$
\begin{equation*}
Q=P \circ E+R\left(x, D_{x^{\prime}}\right) \tag{2-16}
\end{equation*}
$$

where $\sigma_{R}\left(x, \xi^{\prime}\right)=r_{1}\left(x, \xi^{\prime}\right)+r_{0}\left(x, \xi^{\prime}\right)+\cdots$ is an asymptotic sum of homogeneous terms, all independent of $\xi_{1}$. Thus $R\left(x, D_{x^{\prime}}\right)$ is a pseudodifferential operator in the $n-1$ variables $x^{\prime}$ depending on $x_{1}$ as a parameter. Furthermore, the range of $R\left(x, D_{x^{\prime}}\right)$ is microlocally contained in the range of $P$ at $K$. Indeed, suppose $N$ is the integer given by Definition 2.1. If $g \in H_{(N)}^{\mathrm{loc}}\left(\mathbb{R}^{n}\right)$, then $R g=P E g-Q g=P v-Q g$ for some $v \in \mathscr{D}^{\prime}\left(\mathbb{R}^{n}\right)$, and there exists a $u \in \mathscr{D}^{\prime}\left(\mathbb{R}^{n}\right)$ such that

$$
K \cap W F(Q g-P u)=\varnothing
$$

Hence, $W F(P(v-u)-R g)$ does not meet $K$, so the range of $R$ is microlocally contained in the range of $P$ at $K$. We claim that under the assumptions of Theorem 2.19, this implies that all terms in the asymptotic sum of the symbol of the operator $R\left(x, D_{x^{\prime}}\right)$ in (2-16) have vanishing Taylor coefficients at $\Gamma$, thus proving Theorem 2.19. The proof of this claim will be based on the two theorems stated below. As we have seen, the principal symbol $p$ of $P$ may be assumed to have the normal form given by (2-15). By means of Theorem 2.20 below, we shall also use the fact that an even simpler normal form exists near a point where $p=0$ and $\{\operatorname{Re} p, \operatorname{Im} p\} \neq 0$. To prove these two theorems, we will use techniques that actually require the lower order terms of $P$ to be independent of $\xi_{1}$ near $\Gamma$. However, we claim that this may always be assumed. In fact, Malgrange's preparation theorem implies that

$$
p_{0}(x, \xi)=a(x, \xi)\left(\xi_{1}+i f\left(x, \xi^{\prime}\right)\right)+b\left(x, \xi^{\prime}\right)
$$

where $a$ is homogeneous of degree -1 and $b$ homogeneous of degree 0 , as demonstrated in the construction of the operators $E$ and $R$ above. The term of degree 0 in the symbol of $(I-a(x, D)) P$ is equal to $b\left(x, \xi^{\prime}\right)$. Repetition of the argument implies that there exists a classical operator $\tilde{a}(x, D)$ of order -1 such that $(I-\tilde{a}(x, D)) P$ has principal symbol $\xi_{1}+i f\left(x, \xi^{\prime}\right)$ and all lower order terms are independent of $\xi_{1}$. The microlocal property of pseudodifferential operators immediately implies that the
range of $(I-\tilde{a}(x, D)) Q$ is microlocally contained in the range of $(I-\tilde{a}(x, D)) P$ at $K$. Hence, if there are operators $E$ and $R$ with

$$
R=(I-\tilde{a}(x, D)) Q-(I-\tilde{a}(x, D)) P E
$$

such that all terms in the asymptotic expansion of the symbol of $R$ have vanishing Taylor coefficients at $\Gamma$, then this also holds for the symbol of $Q-P E \equiv(I-\tilde{a}(x, D))^{-1} R \bmod \Psi^{-\infty}$, since this property is preserved under composition with elliptic pseudodifferential operators by (2-14).

Theorem 2.20. Suppose that in a conic neighborhood $\Omega$ of

$$
\Gamma^{\prime}=\left\{\left(0, \varepsilon_{n}\right)\right\} \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0
$$

$P$ has the form $P=D_{1}+i x_{1} D_{n}$ and the symbol of $R\left(x, D_{x^{\prime}}\right)$ is given by the asymptotic sum

$$
\sigma_{R}=\sum_{j=0}^{\infty} r_{1-j}\left(x, \xi^{\prime}\right)
$$

with $r_{1-j}$ homogeneous of degree $1-j$ and independent of $\xi_{1}$. If there exists a compactly based cone $K \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ containing $\Omega$ such that the range of $R$ is microlocally contained in the range of $P$ at $K$, then all the terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients at $\Gamma^{\prime}$.

Theorem 2.21. Suppose that in a conic neighborhood $\Omega$ of

$$
\Gamma^{\prime}=\left\{\left(x_{1}, x^{\prime}, 0, \xi^{\prime}\right): a \leq x_{1} \leq b\right\} \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0
$$

the principal symbol of $P$ has the form

$$
p(x, \xi)=\xi_{1}+i f\left(x, \xi^{\prime}\right)
$$

where $f$ is real-valued and homogeneous of degree 1 , and suppose that if $b>a$ then $f$ vanishes of infinite order on $\Gamma^{\prime}$ and there exists $a \varrho \geq 0$ such that for any $\varepsilon>\varrho$ one can find a neighborhood of

$$
\begin{equation*}
\Gamma_{\varepsilon}^{\prime}=\left\{\left(x_{1}, x^{\prime}, 0, \xi^{\prime}\right): a+\varepsilon \leq x_{1} \leq b-\varepsilon\right\} \tag{2-17}
\end{equation*}
$$

where $f$ vanishes identically. Suppose also that

$$
\begin{equation*}
f\left(x, \xi^{\prime}\right)=0 \quad \text { implies } \quad \partial f\left(x, \xi^{\prime}\right) / \partial x_{1} \leq 0 \tag{2-18}
\end{equation*}
$$

in $\Omega$ and that in any neighborhood of $\Gamma^{\prime}$ one can find an interval in the $x_{1}$ direction where $f$ changes sign from - to + for increasing $x_{1}$. Furthermore, suppose that in $\Omega$ the symbol of $R\left(x, D_{x^{\prime}}\right)$ is given by the asymptotic sum

$$
\sigma_{R}=\sum_{j=0}^{\infty} r_{1-j}\left(x, \xi^{\prime}\right)
$$

with $r_{1-j}$ homogeneous of degree $1-j$ and independent of $\xi_{1}$. If the lower order terms $p_{0}, p_{-1}, \ldots$ in the symbol of $P$ are independent of $\xi_{1}$ near $\Gamma^{\prime}$, and there exists a compactly based cone $K \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ containing $\Omega$ such that the range of $R$ is microlocally contained in the range of $P$ at $K$, then all the
terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients on $\Gamma_{\varrho}^{\prime}$ if $a<b$, and at $\Gamma^{\prime}$ if $a=b$.

Assuming these results for the moment, we can now show how Theorem 2.19 follows.
End of proof of Theorem 2.19. Recall that

$$
\Gamma=\left\{\left(x_{1}, 0, \varepsilon_{n}\right): a_{0} \leq x_{1} \leq b_{0}\right\} \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0
$$

By what we have shown, it suffices to regard the case $Q=P E+R$, where we may assume that the conditions of Theorem 2.21 are all satisfied in a conic neighborhood $\Omega$ of $\Gamma$, with the exception of (2-18) and the condition concerning the existence of a neighborhood of (2-17) in which $f$ vanishes identically when $a_{0}<b_{0}$. We consider three cases.
(i) $\Gamma$ is an interval. We then claim that condition (2-18) imposes no restriction. Indeed, if there is no neighborhood of $\Gamma$ in which (2-18) holds, then there exists a sequence $\left\{\gamma_{j}\right\}=\left\{\left(t_{j}, x_{j}^{\prime}, 0, \xi_{j}^{\prime}\right)\right\}$ such that $a_{0} \leq \liminf t_{j} \leq \limsup t_{j} \leq b_{0},\left(x_{j}^{\prime}, \xi_{j}^{\prime}\right) \rightarrow\left(0, \xi^{0}\right) \in \mathbb{R}^{2 n-2}$,

$$
\begin{equation*}
f\left(t_{j}, x_{j}^{\prime}, \xi_{j}^{\prime}\right)=0 \quad \text { and } \quad \partial f\left(t_{j}, x_{j}^{\prime}, \xi_{j}^{\prime}\right) / \partial x_{1}>0 \tag{2-19}
\end{equation*}
$$

for each $j$. By (2-19) we can choose a sequence $0<\delta_{j} \rightarrow 0$ such that

$$
f\left(t_{j}-\delta_{j}, x_{j}^{\prime}, \xi_{j}^{\prime}\right)<0<f\left(t_{j}+\delta_{j}, x_{j}^{\prime}, \xi_{j}^{\prime}\right)
$$

In view of Definition 2.9 we must therefore have $L(\Gamma)=0$. Since $\Gamma$ is minimal, this implies that $|\Gamma|=0$, so $\gamma_{j} \rightarrow \Gamma$. Thus, if there is no neighborhood of $\Gamma$ in which (2-18) holds, then $\Gamma$ is a point, and we will in this case use the existence of the sequence $\left\{\gamma_{j}\right\}$ satisfying (2-19) to reduce the proof of Theorem 2.19 to Theorem 2.20, as demonstrated in case (iii) below. In the present case however, $\Gamma$ is assumed to be an interval, so there exists a neighborhood $U$ of $\Gamma$ in which (2-18) holds. We may assume that $\cup \subset \Omega$ and since $f$ is homogeneous of degree 1 we may also assume that $U$ is conic.

By Theorem 2.18, there exists a sequence $\left\{\Gamma_{j}\right\}$ of $\varrho_{j}$-minimal bicharacteristic intervals such that $\varrho_{j} \rightarrow 0$ and $\Gamma_{j} \rightarrow \Gamma$ as $j \rightarrow \infty$. For sufficiently large $j$ we have $\Gamma_{j} \subset \mathcal{U}$. Hence, if

$$
\Gamma_{j}=\left\{\left(x_{1}, x_{j}^{\prime}, 0, \xi_{j}^{\prime}\right): a_{j} \leq x_{1} \leq b_{j}\right\}
$$

then all the terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients on

$$
\Gamma_{\varrho_{j}}=\left\{\left(x_{1}, x_{j}^{\prime}, 0, \xi_{j}^{\prime}\right): a_{j}+\varrho_{j} \leq x_{1} \leq b_{j}-\varrho_{j}\right\}
$$

by Theorem 2.21. Since $\Gamma_{\varrho_{j}} \rightarrow \Gamma$ as $j \rightarrow \infty$, and all the terms in the asymptotic sum of the symbol of $R$ are smooth functions, it follows that all the terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients on $\Gamma$. This proves Theorem 2.19 in this case.
(ii) $\Gamma$ is a point and condition (2-18) holds. Then all the terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients on $\Gamma$ by Theorem 2.21, so Theorem 2.19 follows.
(iii) $\Gamma$ is a point and (2-18) is false. Let $\left\{\gamma_{j}\right\}$ be the sequence satisfying (2-19). Then $\{\operatorname{Re} p, \operatorname{Im} p\}\left(\gamma_{j}\right)>0$ and $p\left(\gamma_{j}\right)=0$ for each $j$ since $\gamma_{j}=\left(t_{j}, x_{j}^{\prime}, 0, \xi_{j}^{\prime}\right)$. For fixed $j$ we may assume that $\gamma_{j}=(0, \eta)$ and use
[Hörmander 1985a, Theorem 21.3.3] to find a canonical transformation $\chi$ together with Fourier integral operators $A, B, A_{1}$ and $B_{1}$ as in Proposition 2.4 such that $\chi\left(0, \varepsilon_{n}\right)=\gamma_{j}$, and $B P A=D_{1}+i x_{1} D_{n}$ in a conic neighborhood $\Omega$ of $\left\{\left(0, \varepsilon_{n}\right)\right\}$. Repetition of the arguments above allows us to write

$$
\begin{equation*}
B Q A=B P A E+R\left(x, D_{x^{\prime}}\right), \tag{2-20}
\end{equation*}
$$

where the range of $R$ is microlocally contained in the range of $B P A$ at some compactly based cone $K^{\prime}$ containing $\Omega$ with $\chi\left(K^{\prime}\right)=K$. As before, $E$ and $R$ have classical symbols. Then all the terms in the asymptotic expansion of the symbol of $R$ have vanishing Taylor coefficients at $\left\{\left(0, \varepsilon_{n}\right)\right\}$ by Theorem 2.20, and therefore all the terms in the asymptotic expansion of the symbol of $A_{1} R B_{1}$ have vanishing Taylor coefficients at $\gamma_{j}$ by Lemma A. 1 in the appendix. Since the Fourier integral operators are chosen so that

$$
K \cap W F\left(A_{1} B-I\right)=\varnothing \quad \text { and } \quad K \cap W F\left(A B_{1}-I\right)=\varnothing
$$

we have

$$
\varnothing=K \cap W F\left(A_{1} B Q A B_{1}-A_{1} B P A E B_{1}-A_{1} R B_{1}\right)=K \cap W F\left(Q-P A E B_{1}-A_{1} R B_{1}\right)
$$

in view of (2-20). Hence, all the terms in the asymptotic expansion of the symbol of

$$
\begin{equation*}
Q-P E_{1}=A_{1} R B_{1}+S, \quad \text { where } \quad W F(S) \cap K=\varnothing, \tag{2-21}
\end{equation*}
$$

have vanishing Taylor coefficients at $\gamma_{j}$ if $E_{1}=A E B_{1}$. (Strictly speaking, the change of base variables $\gamma_{j} \mapsto(0, \eta)$ should be represented in (2-21) by conjugation of a linear transformation $\kappa: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, but this could be integrated in the Fourier integral operators $A_{1}$ and $B_{1}$ so it has been left out since it will not affect the arguments below.) It is clear that $E_{1} \in \Psi_{\mathrm{cl}}^{0}\left(\mathbb{R}^{n}\right)$.

We have now shown that for each $j$ there exists an operator $E_{j} \in \Psi_{\mathrm{cl}}^{0}\left(\mathbb{R}^{n}\right)$ such that all the terms in the asymptotic expansion of the symbol of $Q-P E_{j}$ have vanishing Taylor coefficients at $\gamma_{j}$. To construct the operator $E$ in Theorem 2.19, we do the following. For each $j$, denote the symbol of $E_{j}$ by

$$
e^{j}(x, \xi) \sim \sum_{l=0}^{\infty} e_{-l}^{j}(x, \xi)
$$

where $e_{0}^{j}(x, \xi)$ is the principal part, and $e_{-l}^{j}(x, \xi)$ is homogeneous of degree $-l$. If $q$ is the principal symbol of $Q$, then by Proposition A. 3 there exists a function $e_{0} \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$, homogeneous of degree 0 , such that $q-p e_{0}$ has vanishing Taylor coefficients at $\Gamma$.

This argument can be repeated for lower order terms. Indeed, if $\sigma_{Q}=q+q_{0}+\cdots$, then the term of degree 0 in $\sigma_{Q-P E_{j}}$ is

$$
\sigma_{0}\left(Q-P E_{j}\right)=\tilde{q}_{j}-p e_{-1}^{j}
$$

where (see (2-25) below)

$$
\tilde{q}_{j}(x, \xi)=q_{0}(x, \xi)-p_{0}(x, \xi) e_{0}^{j}(x, \xi)-\sum_{k} \partial_{\xi_{k}} p(x, \xi) D_{x_{k}} e_{0}^{j}(x, \xi)
$$

We can write

$$
p(x, \xi) e_{-1}^{j}(x, \xi)=p(x, \xi /|\xi|) e_{-1}^{j}(x, \xi /|\xi|)
$$

so that $\tilde{q}_{j}(x, \xi), p(x, \xi /|\xi|)$ and $e_{-1}^{j}(x, \xi /|\xi|)$ are all homogeneous of degree 0 . Since

$$
\partial_{x}^{\alpha} \partial_{\xi}^{\beta} e_{0}(\Gamma)=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} e_{0}^{j}\left(\gamma_{j}\right)
$$

it follows by Proposition A. 3 that there is a function $g \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$, homogeneous of degree 0 , such that

$$
q_{0}(x, \xi)-p_{0}(x, \xi) e_{0}(x, \xi)-\sum_{k} \partial_{\xi_{k}} p(x, \xi) D_{x_{k}} e_{0}(x, \xi)-p(x, \xi /|\xi|) g(x, \xi)
$$

has vanishing Taylor coefficients at $\Gamma$. Putting $e_{-1}(x, \xi)=|\xi|^{-1} g(x, \xi)$ we find that

$$
\partial_{x}^{\alpha} \partial_{\xi}^{\beta} e_{-1}(\Gamma)=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} e_{-1}^{j}\left(\gamma_{j}\right)
$$

and that

$$
\sigma_{0}\left(Q-P \circ e_{0}(x, D)-P \circ e_{-1}(x, D)\right)
$$

has vanishing Taylor coefficients at $\Gamma$. Continuing this way we successively obtain functions $e_{m}(x, \xi) \in$ $C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$, homogeneous of degree $m$ for $m \leq 0$, such that

$$
\sigma_{Q}-\left(\sum_{m=0}^{M} e_{-m}\right) \sigma_{P} \quad \bmod S_{\mathrm{cl}}^{-M}
$$

has vanishing Taylor coefficients at $\Gamma$. If we let $E$ have symbol

$$
\sigma_{E}(x, \xi) \sim \sum_{m=0}^{\infty}(1-\phi(\xi)) e_{-m}(x, \xi)
$$

with $\phi \in C_{0}^{\infty}$ equal to 1 for $\xi$ close to 0 , then $E \in \Psi_{\mathrm{cl}}^{0}\left(\mathbb{R}^{n}\right)$ and all terms in the asymptotic expansion of the symbol of $Q-P E$ have vanishing Taylor coefficients at $\Gamma$. We have proved Theorem 2.19.
Remark. Instead of reducing to the study of the normal form $P=D_{x_{1}}+i x_{1} D_{x_{n}}$ when condition (2-18) does not hold, as in case (iii) above, one could show that the terms in the asymptotic expansion of the operator $R$ given by (2-16) has vanishing Taylor coefficients at every point in the sequence $\left\{\gamma_{j}\right\}$ satisfying (2-19) using techniques very similar to those used to prove Theorem 2.21. Theorem 2.19 would then follow by continuity, but the proof of the analogue of Theorem 2.20 would be more involved. In particular, we would have to construct a phase function $w$ solving the eikonal equation

$$
\partial w / \partial x_{1}-i f\left(x, \partial w / \partial x^{\prime}\right)=0
$$

approximately instead of explicitly (confer the proofs of Theorems 2.21 and 2.20, respectively). For fixed $j$ this could be accomplished by adapting the approach in [Hörmander 1963; Hörmander 1966] (for a brief discussion, see [Hörmander 1981, p. 83]), where one has $f=0$ and $\partial f / \partial x_{1}>0$ at $\left(0, \xi^{0}\right)$ instead of at $\gamma_{j}$.

We shall now show how our results relate to the ones referred to in the introduction, beginning with (1-3). There, it sufficed to have the coefficients of $P$ and $Q$ in $C^{\infty}$ and $C^{1}$, respectively. However, in order for Theorem 2.19 to qualify, we must require both $P$ and $Q$ to have smooth coefficients. On the other hand, we shall only require the equation $P u=Q f$ to be microlocally solvable (at an appropriate cone $K$ ) as given by Definition 2.1. Note that if $P$ is a first order differential operator on an open set $\Omega \subset \mathbb{R}^{n}$, such that the principal symbol $p$ of $P$ satisfies condition (1-4) at a point $(x, \xi) \in T^{*}(\Omega) \backslash 0$, then either $\{\operatorname{Re} p, \operatorname{Im} p\}>0$ at $(x, \xi)$, or $\{\operatorname{Re} p, \operatorname{Im} p\}>0$ at $(x,-\xi)$. (The order of the operator is not important; the statement is still true for a differential operator of order $m$, since the Poisson bracket is then homogeneous of order $2 m-1$.) Assuming the former, this implies that ( $x, \xi$ ) satisfies condition (a) in Theorem 2.19 by an application of [Hörmander 1985a, Theorem 21.3.3] and Lemma 2.7. In order to keep the formulation of the following result as simple as possible, we will assume that there exists a compactly based cone $K \subset T^{*}(\Omega) \backslash 0$ with nonempty interior such that $K$ contains the appropriate point $(x, \pm \xi)$, and such that the equation $P u=Q f$ is microlocally solvable at $K$. This is clearly the case if the equation $P u=Q f$ is locally solvable in $\Omega$ in the weak sense suggested by (1-1).

Corollary 2.22. Let $\Omega \subset \mathbb{R}^{n}$ be open, and let $P(x, D)$ and $Q(x, D)$ be two first order differential operators with coefficients in $C^{\infty}(\Omega)$. Let $p$ be the principal symbol of $P$, and let $x_{0}$ be a point in $\Omega$ such that

$$
\begin{equation*}
p\left(x_{0}, \xi_{0}\right)=0 \quad \text { and } \quad\{\operatorname{Re} p, \operatorname{Im} p\}\left(x_{0}, \xi_{0}\right)>0 \tag{2-22}
\end{equation*}
$$

for some $\xi_{0} \in \mathbb{R}^{n}$. If $K \subset T^{*}(\Omega) \backslash 0$ is a compactly based cone containing $\left(x_{0}, \xi_{0}\right)$ such that the range of $Q$ is microlocally contained in the range of $P$ at $K$, then there exists a constant $\mu$ such that (at the fixed point $x_{0}$ )

$$
\begin{equation*}
Q^{*}\left(x_{0}, D\right)=\mu P^{*}\left(x_{0}, D\right) \tag{2-23}
\end{equation*}
$$

where $Q^{*}$ and $P^{*}$ are the adjoints of $Q$ and $P$.
Proof. By (2-22), $P \in \Psi_{\mathrm{cl}}^{1}(\Omega)$ is an operator of principal type microlocally near $\left(x_{0}, \xi_{0}\right) . P$ and $Q$ therefore satisfy the hypotheses of Theorem 2.19 , and in view of the discussion above regarding the point $(x, \xi)$ we find that there exists an operator $E \in \Psi_{\mathrm{cl}}^{0}(\Omega)$ such that all the terms in the asymptotic expansion of the symbol of $Q-P E$ have vanishing Taylor coefficients at $\left(x_{0}, \xi_{0}\right)$. By the discussion following (3-7) on page 452 below, it follows that the same must hold for the adjoint $Q^{*}-E^{*} P^{*}$. If we let $Q^{*}$ and $P^{*}$ have symbols $\sigma_{Q^{*}}(x, \xi)=q_{1}(x, \xi)+q_{0}(x)$ and $\sigma_{P^{*}}(x, \xi)=p_{1}(x, \xi)+p_{0}(x)$, then $E^{*} P^{*}$ has principal symbol $e_{0} p_{1}$ if $\sigma_{E^{*}}=e_{0}+e_{-1}+\ldots$ denotes the symbol of $E^{*}$. Hence

$$
\partial q_{1}\left(x_{0}, \xi_{0}\right) / \partial \xi_{k}=e_{0}\left(x_{0}, \xi_{0}\right) \partial p_{1}\left(x_{0}, \xi_{0}\right) / \partial \xi_{k} \quad \text { for } 1 \leq k \leq n
$$

and $p_{1}\left(x_{0}, \xi_{0}\right)=\overline{p\left(x_{0}, \xi_{0}\right)}=0$. Since $q_{1}$ and $p_{1}$ are polynomials in $\xi$ of degree 1 , this means that at the fixed point $x_{0}$ we have $q_{1}\left(x_{0}, \xi\right)=\mu p_{1}\left(x_{0}, \xi\right)$ for $\xi \in \mathbb{R}^{n}$, where the constant $\mu$ is given by the value of $e_{0}$ at $\left(x_{0}, \xi_{0}\right)$. Moreover,

$$
\begin{equation*}
0=\partial_{\xi_{j}} \partial_{\xi_{k}} q_{1}\left(x_{0}, \xi_{0}\right)=\partial_{\xi_{j}} e_{0}\left(x_{0}, \xi_{0}\right) \partial_{\xi_{k}} p_{1}\left(x_{0}, \xi_{0}\right)+\partial_{\xi_{k}} e_{0}\left(x_{0}, \xi_{0}\right) \partial_{\xi_{j}} p_{1}\left(x_{0}, \xi_{0}\right) \tag{2-24}
\end{equation*}
$$

By assumption, the coefficients of $p(x, D)$ do not vanish simultaneously, so the same is true for $p_{1}(x, D)$. Hence $\partial_{\xi_{j}} p_{1}\left(x_{0}, \xi_{0}\right) \neq 0$ for some $j$. Assuming this holds for $j=1$, we find by choosing $j=k=1$ in (2-24) that $\partial_{\xi_{1}} e_{0}\left(x_{0}, \xi_{0}\right)=0$. But this immediately yields

$$
\partial_{\xi_{k}} e_{0}\left(x_{0}, \xi_{0}\right)=-\partial_{\xi_{1}} e_{0}\left(x_{0}, \xi_{0}\right) \partial_{\xi_{k}} p_{1}\left(x_{0}, \xi_{0}\right) / \partial_{\xi_{1}} p_{1}\left(x_{0}, \xi_{0}\right)=0
$$

for $2 \leq k \leq n$. Now

$$
\sigma_{E^{*} P^{*}}(x, \xi) \sim \sum \frac{1}{\alpha!} \partial_{\xi}^{\alpha} \sigma_{E^{*}} D_{x}^{\alpha}\left(p_{1}(x, \xi)+p_{0}(x)\right)
$$

and since we have a bilinear map

$$
S_{\mathrm{cl}}^{m^{\prime}} / S^{-\infty} \times S_{\mathrm{cl}}^{m^{\prime \prime}} / S^{-\infty} \ni(a, b) \mapsto a \# b \in S_{\mathrm{cl}}^{m^{\prime}+m^{\prime \prime}} / S^{-\infty}
$$

with

$$
(a \# b)(x, \xi) \sim \sum \frac{1}{\alpha!} \partial_{\xi}^{\alpha} a(x, \xi) D_{x}^{\alpha} b(x, \xi)
$$

we find that the term of order 0 in the symbol of $E^{*} P^{*}$ is

$$
\begin{equation*}
\sigma_{0}\left(E^{*} P^{*}\right)(x, \xi)=e_{-1}(x, \xi) p_{1}(x, \xi)+e_{0}(x, \xi) p_{0}(x)+\sum_{k=1}^{n} \partial_{\xi_{k}} e_{0}(x, \xi) D_{k} p_{1}(x, \xi) \tag{2-25}
\end{equation*}
$$

Since $\partial_{\xi_{k}} e_{0}$ and $p_{1}$ vanish at $\left(x_{0}, \xi_{0}\right)$ we find that $q_{0}\left(x_{0}\right)=\mu p_{0}\left(x_{0}\right)$ at the fixed point $x_{0}$, which completes the proof.

Having proved this result, we immediately obtain the following after making the obvious adjustments to [Hörmander 1963, Theorem 6.2.2]. The fact that we require higher regularity on the coefficients of $Q$ then yields higher regularity on the propertionality factor. Since the proof remains the same, it is omitted.

Corollary 2.23. Let $\Omega \subset \mathbb{R}^{n}$ be open, and let $P(x, D)$ and $Q(x, D)$ be two first order differential operators with coefficients in $C^{\infty}(\Omega)$. Let $p$ be the principal symbol of $P$, and assume that the coefficients of $p(x, D)$ do not vanish simultaneously in $\Omega$. If for a dense set of points $x$ in $\Omega$ one can find $\xi \in \mathbb{R}^{n}$ such that (2-22) is fulfilled, and iffor each $(x, \xi)$ there is a compactly based cone $K \subset T^{*}(\Omega) \backslash 0$ containing $(x, \xi)$ such that the range of $Q$ is microlocally contained in the range of $P$ at $K$, then there exists $a$ function $e \in C^{\infty}(\Omega)$ such that

$$
\begin{equation*}
Q(x, D) u \equiv P(x, D)(e u) \tag{2-26}
\end{equation*}
$$

In stating Corollary 2.23 we could replace the assumption that the coefficients of $p(x, D)$ do not vanish simultaneously in $\Omega$ with the condition that $P$ is of principal type. Indeed, if $d p \neq 0$ then by a canonical transformation we find that condition (1-6) holds. Since $p \neq 0$ implies $\partial_{\xi} p \neq 0$ by the Euler homogeneity equation we then have $\partial_{\xi} p \neq 0$ everywhere, that is, the coefficients of $p(x, D)$ do not vanish simultaneously in $\Omega$. The converse is obvious.

As shown in Example 2.25 below, we also recover the result for higher order differential operators mentioned in the introduction as a special case of the following corollary to Theorem 2.19, although we again need to assume higher regularity in order to apply our results.

Proposition 2.24. Let $X$ be a smooth manifold, and let $P \in \Psi_{\mathrm{cl}}^{k}(X)$ and $Q \in \Psi_{\mathrm{cl}}^{k^{\prime}}(X)$ be properly supported such that the range of $Q \circ P$ is microlocally contained in the range of $P$ at a compactly based cone $K \subset T^{*}(X) \backslash 0$. Let $p$ and $q$ be the principal symbols of $P$ and $Q$, respectively, and assume that $P$ is of principal type microlocally near $K$. If $\gamma: I \rightarrow T^{*}(X) \backslash 0$ is a minimal characteristic point or a minimal bicharacteristic interval of $p$ contained in $K$ then it follows that

$$
H_{p}^{m}(q)=0 \quad \text { for all }(x, \xi) \in \gamma(I) \text { and } m \geq 1
$$

Here $H_{p}^{m}(q)$ is defined recursively by $H_{p}(q)=\{p, q\}$ and $H_{p}^{m}(q)=\left\{p, H_{p}^{m-1}(q)\right\}$ for $m \geq 2$.
Proof. First note that if the range of $Q \in \Psi_{\mathrm{cl}}^{k^{\prime}}(X)$ is microlocally contained in the range of $P \in \Psi_{\mathrm{cl}}^{k}(X)$ at $K$ and both operators are properly supported, then it follows that the range of $Q \circ P$ is microlocally contained in the range of $P$ at $K$. (The converse is not true in general.) Indeed, let $N$ be the integer given by Definition 2.1, and let $f \in H_{(N+k)}^{\text {loc }}(X)$. Since $P: H_{(N+k)}^{\text {loc }}(X) \rightarrow H_{(N)}^{\text {loc }}(X)$ is continuous, we have $g=P f \in H_{(N)}^{\text {loc }}(X)$. Thus, there exists a $u \in \mathscr{D}^{\prime}(X)$ such that

$$
\varnothing=K \cap W F(Q g-P u)=K \cap W F(Q P f-P u)
$$

so the conditions of Definition 2.1 are satisfied with $N$ replaced with $N+k$.
Let $(x, \xi) \in \gamma(I)$. The range of $P Q$ is easily seen to be microlocally contained in the range of $P$ for any properly supported pseudodifferential operator $Q$. The assumptions of the proposition therefore imply that the range of the commutator

$$
\begin{equation*}
R_{1}=P \circ Q-Q \circ P \in \Psi_{\mathrm{cl}}^{k+k^{\prime}-1}(X) \tag{2-27}
\end{equation*}
$$

is microlocally contained in the range of $P$ at $K$. Hence, by Theorem 2.19 there exists an operator $E \in \Psi_{\mathrm{cl}}^{k^{\prime}-1}(X)$ such that, in particular, the principal symbol of $R_{1}-P E$ vanishes at $(x, \xi)$. If $e$ is the principal symbol of $E$, homogeneous of degree $k^{\prime}-1$, then the principal symbol of $P E$ satisfies $p(x, \xi) e(x, \xi)=0$ since $p \circ \gamma=0$. Since the principal symbol of $R_{1}$ is

$$
\sigma_{k+k^{\prime}-1}\left(R_{1}\right)=\frac{1}{i}\{p, q\}
$$

the result follows for $m=1$.
Let $R_{m}$ be defined recursively by $R_{m}=\left[P, R_{m-1}\right]$ for $m \geq 2$ with $R_{1}$ given by (2-27). Arguing by induction, we conclude in view of the first paragraph of the proof that the range of $R_{m}$ is microlocally contained in the range of $P$ at $K$ for $m=1,2 \ldots$ since this holds for $R_{1}$. Assuming the proposition holds for some $m \geq 1$, we can repeat the arguments above to show that the principal symbol of $R_{m+1}$ must vanish at $(x, \xi)$. Since the principal symbol of $R_{m+1}$ equals $-i\left\{p, H_{p}^{m}(q)\right\}$, this completes the proof.

Example 2.25. Let $\Omega \subset \mathbb{R}^{n}$ be open, $P(x, D)$ be a differential operator of order $m$ with coefficients in $C^{\infty}(\Omega)$, and let $\mu$ be a function in $C^{\infty}(\Omega)$ such that the equation

$$
P(x, D) u=\mu P(x, D) f
$$

has a solution $u \in \mathscr{D}^{\prime}(\Omega)$ for every $f \in C_{0}^{\infty}(\Omega)$. If $p$ is the principal symbol of $P$ then it follows that

$$
\begin{equation*}
\sum_{j=1}^{n} \partial_{\xi_{j}} p(x, \xi) D_{x_{j}} \mu(x)=0 \tag{2-28}
\end{equation*}
$$

for all $x \in \Omega$ and $\xi \in \mathbb{R}^{n}$ such that

$$
\begin{equation*}
\{p, \bar{p}\}(x, \xi) \neq 0 \quad \text { and } \quad p(x, \xi)=0 \tag{2-29}
\end{equation*}
$$

Indeed, if $(x, \xi)$ satisfies (2-29) then we may assume that

$$
\{\operatorname{Re} p, \operatorname{Im} p\}(x, \xi)=-\frac{1}{2 i}\{p, \bar{p}\}(x, \xi)>0
$$

since otherwise we just regard $(x,-\xi)$ instead, as per the remarks preceding Corollary 2.22. By the same discussion it is also clear that $(x, \xi)$ is a minimal characteristic point of $p$. Now the conditions above imply that there exists a compactly based cone $K \subset T^{*}(\Omega) \backslash 0$ containing ( $x, \xi$ ) such that the range of $\mu P$ is microlocally contained in the range of $P$ at $K$. By condition (2-29), $P$ is of principal type near $(x, \xi)$, so Proposition 2.24 implies that $\{p, \mu\}=0$ at $(x, \xi)$, that is,

$$
\sum_{j=1}^{n} \partial_{\xi_{j}} p(x, \xi) \partial_{x_{j}} \mu(x)-\partial_{x_{j}} p(x, \xi) \partial_{\xi_{j}} \mu(x)=0
$$

Since $\mu$ is independent of $\xi$ we find that (2-28) holds at $(x, \xi)$. By homogeneity it then also holds at $(x,-\xi)$.

## 3. Proof of Theorem $\mathbf{2 . 2 0}$

Throughout this section we assume that the hypotheses of Theorem 2.20 hold. We shall prove the theorem by using Lemma 2.3 on approximate solutions of the equation $P^{*} v=0$ concentrated near $\Gamma^{\prime}=\left\{\left(0, \varepsilon_{n}\right)\right\}$. We take as starting point the construction on [Hörmander 1985b, page 103], but some modifications need to be made in particular to the amplitude function $\phi$, so the results there concerning the estimates for the right side of (2-1) cannot be used immediately. To obtain the desired estimates we will instead have to use [Hörmander 1985b, Lemma 26.4.15]. Set

$$
\begin{equation*}
v_{\tau}(x)=\phi(x) e^{i \tau w(x)} \tag{3-1}
\end{equation*}
$$

where

$$
w(x)=x_{n}+i\left(x_{1}^{2}+x_{2}^{2}+\cdots+x_{n-1}^{2}+\left(x_{n}+i x_{1}^{2} / 2\right)^{2}\right) / 2
$$

satisfies $P^{*} w=0$ and $\phi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$. By the Cauchy-Kovalevsky theorem we can solve $D_{1} \phi-i x_{1} D_{n} \phi=0$ in a neighborhood of 0 for any analytic initial data $\phi\left(0, x^{\prime}\right)=f\left(x^{\prime}\right) \in C^{\omega}\left(\mathbb{R}^{n-1}\right)$; in particular we are free to specify the Taylor coefficients of $f\left(x^{\prime}\right)$ at $x^{\prime}=0$. We take $\phi$ to be such a solution. If need be we can reduce the support of $\phi$ by multiplying by a smooth cutoff function $\chi$, where $\chi$ is equal to 1 in
some smaller neighborhood of 0 so that $\chi \phi$ solves the equation there. We assume this to be done and note that if $\operatorname{supp} \phi$ is small enough then

$$
\begin{equation*}
\operatorname{Im} w(x) \geq|x|^{2} / 4 \quad \text { for } x \in \operatorname{supp} \phi \tag{3-2}
\end{equation*}
$$

Since

$$
d \operatorname{Re} w(x)=-x_{1} x_{n} d x_{1}+\left(1-x_{1}^{2} / 2\right) d x_{n}
$$

we may similarly assume that $d \operatorname{Re} w(x) \neq 0$ in the support of $\phi$. We then have the following result.
Lemma 3.1. Suppose $P=D_{1}+i x_{1} D_{n}$ and let $v_{\tau}$ be defined by (3-1). Then $\phi$ and $w$ can be chosen so that for any $f \in C^{\omega}\left(\mathbb{R}^{n-1}\right)$ and any positive integers $k$ and $m$ we have $\phi\left(0, x^{\prime}\right)=f\left(x^{\prime}\right)$ in a neighborhood of $(0,0), \tau^{k}\left\|P^{*} v_{\tau}\right\|_{(m)} \rightarrow 0$ as $\tau \rightarrow \infty$, and

$$
\begin{equation*}
\left\|v_{\tau}\right\|_{(-m)} \leq C_{m} \tau^{-m} \tag{3-3}
\end{equation*}
$$

If $\tilde{\Gamma}$ is the cone generated by

$$
\left\{\left(x, w^{\prime}(x)\right): x \in \operatorname{supp} \phi, \operatorname{Im} w(x)=0\right\}
$$

then $\tau^{k} v_{\tau} \rightarrow 0$ in $\mathscr{D}_{\tilde{\Gamma}}^{\prime}$ as $\tau \rightarrow \infty$; hence $\tau^{k} A v_{\tau} \rightarrow 0$ in $C^{\infty}\left(\mathbb{R}^{n}\right)$ if $A$ is a pseudodifferential operator with $W F(A) \cap \tilde{\Gamma}=\varnothing$.

Here $\mathscr{D}_{\tilde{\Gamma}}^{\prime}(X)=\left\{u \in \mathscr{D}^{\prime}(X): W F(u) \subset \tilde{\Gamma}\right\}$, equipped with the topology given by all the seminorms on $\mathscr{D}^{\prime}(X)$ for the weak topology, together with all seminorms of the form

$$
P_{\phi, V, N}(u)=\sup _{\xi \in V}|\widehat{\phi u}(\xi)|(1+|\xi|)^{N}
$$

where $N \geq 0, \phi \in C_{0}^{\infty}(X)$, and $V \subset \mathbb{R}^{n}$ is a closed cone with $(\operatorname{supp} \phi \times V) \cap \tilde{\Gamma}=\varnothing$. Note that $u_{j} \rightarrow u$ in $\mathscr{D}_{\tilde{\Gamma}}^{\prime}(X)$ is equivalent to $u_{j} \rightarrow u$ in $\mathscr{D}^{\prime}(X)$ and $A u_{j} \rightarrow A u$ in $C^{\infty}$ for every properly supported pseudodifferential operator $A$ with $\tilde{\Gamma} \cap W F(A)=\varnothing$; see the remark following [Hörmander 1985a, Theorem 18.1.28].
Proof. We observe that $\tau^{k} P^{*} v_{\tau}=\tau^{k}\left(P^{*} \phi\right) e^{i \tau w} \rightarrow 0$ in $C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ for any $k$ as $\tau \rightarrow \infty$, if $w$ and $\phi$ are chosen in the way given above. Hence $\tau^{k}\left\|P^{*} v_{\tau}\right\|_{(m)} \rightarrow 0$ for any positive integers $k$ and $m$. In view of (3-2) and the fact that $d \operatorname{Re} w \neq 0$ in the support of $\phi$ we can apply [Hörmander 1985b, Lemma 26.4.15] to $v_{\tau}$. This immediately yields (3-3) and also that $\tau^{k} v_{\tau} \rightarrow 0$ in $\mathscr{D}_{\tilde{\Gamma}}^{\prime}$ as $\tau \rightarrow \infty$, proving the lemma.

We are now ready to proceed with a tool that will be instrumental in proving Theorem 2.21. The idea is based on techniques found in [Hörmander 1963].

Let $R$ be the operator given by Theorem 2.20. By assumption there exists a compactly based cone $K \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ such that the range of $R$ is microlocally contained in the range of $P$ at $K$. If $N$ is the integer given by Definition 2.1, let $H(x) \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ and set

$$
\begin{equation*}
h_{\tau}(x)=\tau^{-N} H(\tau x) \tag{3-4}
\end{equation*}
$$

Since $\hat{h}_{\tau}(\xi)=\tau^{-N-n} \hat{H}(\xi / \tau)$ it is clear that for $\tau \geq 1$ we have $h_{\tau} \in H_{(N)}\left(\mathbb{R}^{n}\right)$ and $\left\|h_{\tau}\right\|_{(N)} \leq C \tau^{-n / 2}$. In particular, $\left\|h_{\tau}\right\|_{(N)} \leq C$ for $\tau \geq 1$, where the constant depends on $H$ but not on $\tau$. Now denote by $I_{\tau}$ the integral

$$
\begin{equation*}
I_{\tau}=\tau^{n} \int H(\tau x) R^{*} v_{\tau}(x) d x=\tau^{N+n}\left(R^{*} v_{\tau}, \overline{h_{\tau}}\right) \tag{3-5}
\end{equation*}
$$

where $R^{*}$ is the adjoint of $R$. For any $\kappa$ we then have by the second equality and Lemma 2.3 that

$$
\left|I_{\tau}\right| \leq \tau^{N+n}\left\|h_{\tau}\right\|_{(N)}\left\|R^{*} v_{\tau}\right\|_{(-N)} \leq C_{\kappa} \tau^{N+n}\left(\left\|P^{*} v_{\tau}\right\|_{(\nu)}+\left\|v_{\tau}\right\|_{(-N-\kappa-n)}+\left\|A v_{\tau}\right\|_{(0)}\right)
$$

for some positive integer $v$ and properly supported pseudodifferential operator $A$ with $W F(A) \cap K=\varnothing$. By Lemma 3.1 this implies

$$
\begin{equation*}
\left|I_{\tau}\right| \leq C_{\kappa} \tau^{-\kappa} \tag{3-6}
\end{equation*}
$$

for any positive integer $\kappa$ if $\tau$ is sufficiently large.
Recall that $R\left(x, D_{x^{\prime}}\right)$ is a pseudodifferential operator in $x^{\prime}$ depending on $x_{1}$ as a parameter. Its symbol is given by the asymptotic sum

$$
\sigma_{R}\left(x, \xi^{\prime}\right)=r_{1}\left(x, \xi^{\prime}\right)+r_{0}\left(x, \xi^{\prime}\right)+\cdots,
$$

where $r_{-j}\left(x, \xi^{\prime}\right)$ is homogeneous of degree $-j$ in $\xi^{\prime}$. The symbol of $R^{*}$ has the asymptotic expansion

$$
\sigma_{R^{*}}=\sum \partial_{\xi}^{\alpha} D_{x}^{\alpha} \overline{\sigma_{R}\left(x, \xi^{\prime}\right)} / \alpha!
$$

which shows that $R^{*}$ is also a pseudodifferential operator in $x^{\prime}$ depending on $x_{1}$ as a parameter. If we sort the terms above with respect to homogeneity we can write

$$
\begin{equation*}
\sigma_{R^{*}}=q_{1}\left(x, \xi^{\prime}\right)+q_{0}\left(x, \xi^{\prime}\right)+\cdots \tag{3-7}
\end{equation*}
$$

where $q_{-j}$ is homogeneous of order $-j, q_{1}\left(x, \xi^{\prime}\right)=\overline{r_{1}\left(x, \xi^{\prime}\right)}$ and

$$
q_{0}\left(x, \xi^{\prime}\right)=\overline{r_{0}\left(x, \xi^{\prime}\right)}+\sum_{k=2}^{n} \partial_{\xi_{k}} D_{x_{k}} \overline{r_{1}\left(x, \xi^{\prime}\right)}
$$

A moment's reflection shows that if all the terms in (3-7) have vanishing Taylor coefficients at some point $\left(x, \xi^{\prime}\right)$, then the same must hold for $\sigma_{R}$.

Our goal is to show that if $q_{-j(\alpha)}^{(\beta)}\left(0, \xi^{0}\right)$ does not vanish for all $j \geq-1$ and all $\alpha, \beta \in \mathbb{N}^{n}$, then (3-6) cannot hold. For this purpose, we introduce a total well-ordering $>_{t}$ on the Taylor coefficients by means of an ordering of the indices $(j, \alpha, \beta)$ as follows.

Definition 3.2. Let $\alpha_{i}, \beta_{i} \in \mathbb{N}^{n}$ and $j_{i} \geq-1$ for $i=1,2$. We say that

$$
q_{-j_{1}\left(\alpha_{1}\right)}^{\left(\beta_{1}\right)}\left(0, \xi^{0}\right)>_{t} q_{-j_{2}\left(\alpha_{2}\right)}^{\left(\beta_{2}\right)}\left(0, \xi^{0}\right) \quad \text { if } j_{1}+\left|\alpha_{1}\right|+\left|\beta_{1}\right|>j_{2}+\left|\alpha_{2}\right|+\left|\beta_{2}\right|
$$

To "break ties", we say that if $j_{1}+\left|\alpha_{1}\right|+\left|\beta_{1}\right|=j_{2}+\left|\alpha_{2}\right|+\left|\beta_{2}\right|$, then

$$
q_{-j_{1}\left(\alpha_{1}\right)}^{\left(\beta_{1}\right)}\left(0, \xi^{0}\right)>_{t} q_{-j_{2}\left(\alpha_{2}\right)}^{\left(\beta_{2}\right)}\left(0, \xi^{0}\right) \quad \text { if }\left|\beta_{2}\right|>\left|\beta_{1}\right|
$$

Note the reversed order. If also $\left|\beta_{1}\right|=\left|\beta_{2}\right|$, then we use a monomial ordering on the $\beta$ index to break ties. Recall that this is any relation $>$ on $\mathbb{N}^{n}$ such that $>$ is a total well-ordering on $\mathbb{N}^{n}$ and $\beta_{1}>\beta_{2}$ and $\gamma \in \mathbb{N}^{n}$ implies $\beta_{1}+\gamma>\beta_{2}+\gamma$. Having come this far, the actual order turns out not to matter for the proof of Theorem 2.20, but it will have bearing on the proof of Theorem 2.21. Which monomial ordering we use on the $\beta$ index will not be important, but for completeness let us choose lexicographic order since this will be used at a later stage in the definition. Here we by lexicographic order refer to the usual one, corresponding to the variables being ordered $x_{1}>\cdots>x_{n}$. That is to say, if $\alpha_{i} \in \mathbb{N}^{n}$ for $i=1,2$, then $\alpha_{1}>_{\text {lex }} \alpha_{2}$ if, in the vector difference $\alpha_{1}-\alpha_{2} \in \mathbb{Z}^{n}$, the leftmost nonzero entry is positive. Thus, if $j_{1}+\left|\alpha_{1}\right|+\left|\beta_{1}\right|=j_{2}+\left|\alpha_{2}\right|+\left|\beta_{2}\right|$ and $\beta_{1}=\beta_{2}$, then we first say that

$$
\begin{equation*}
q_{-j_{1}\left(\alpha_{1}\right)}^{\left(\beta_{1}\right)}\left(0, \xi^{0}\right)>_{t} q_{-j_{2}\left(\alpha_{2}\right)}^{\left(\beta_{2}\right)}\left(0, \xi^{0}\right) \quad \text { if }\left|\alpha_{2}\right|>\left|\alpha_{1}\right| \tag{3-8}
\end{equation*}
$$

and then use lexicographic order on the $n$-tuples $\alpha$ to break ties at this stage. Using the lexicographic order on both multiindices (separately) we get

$$
q_{1}<_{t} q_{1}^{\left(\varepsilon_{n}\right)}<_{t} \cdots<_{t} q_{1}^{\left(\varepsilon_{1}\right)}<_{t} q_{1\left(\varepsilon_{n}\right)}<_{t} \cdots<_{t} q_{1\left(\varepsilon_{1}\right)}<_{t} q_{0}<_{t} \cdots
$$

As indicated above we will prove Theorem 2.20 by a contradiction argument, so in the sequel we let $\kappa$ denote an integer such that

$$
\begin{equation*}
j+|\alpha|+|\beta|<\kappa \tag{3-9}
\end{equation*}
$$

if $q_{-j(\alpha)}^{(\beta)}\left(0, \xi^{0}\right)$ is the first nonvanishing Taylor coefficient with respect to the ordering $>_{t}$. Since $j \geq-1$ we will thus have $\kappa \geq 0$.

To simplify notation, we shall in what follows write $t$ instead of $x_{1}$ and $x$ instead of $x^{\prime}$. Then $v_{\tau}$ takes the form

$$
v_{\tau}(t, x)=\phi(t, x) e^{i \tau w(t, x)}
$$

where

$$
\begin{equation*}
w(t, x)=x_{n-1}+i\left(t^{2}+x_{1}^{2}+\cdots+x_{n-2}^{2}+\left(x_{n-1}+i t^{2} / 2\right)^{2}\right) / 2 \tag{3-10}
\end{equation*}
$$

We shall as before use the notation $\xi^{0}=(0, \ldots, 0,1) \in \mathbb{R}^{n-1}$ when in this context. To interpret the integral $I_{\tau}$ we will need a formula for how $R^{*}(t, x, D)$ acts on the functions $v_{\tau}$. This is given by the following lemma, where the parameter $t$ has been suppressed to simplify notation.
Lemma 3.3 [Hörmander 1985b, Lemma 26.4.16]. Let $q(x, \xi) \in S^{\mu}\left(\mathbb{R}^{n-1} \times \mathbb{R}^{n-1}\right)$, let $\phi \in C_{0}^{\infty}\left(\mathbb{R}^{n-1}\right)$ and $w \in C^{\infty}\left(\mathbb{R}^{n-1}\right)$, and assume that $\operatorname{Im} w>0$ except at a point $y$ where $w^{\prime}(y)=\eta \in \mathbb{R}^{n-1} \backslash 0$ and $\operatorname{Im} w^{\prime \prime}$ is positive definite. Then

$$
\begin{equation*}
\left|q(x, D)\left(\phi e^{i \tau w}\right)-\sum_{|\alpha|<k} q^{(\alpha)}(x, \tau \eta)(D-\tau \eta)^{\alpha}\left(\phi e^{i \tau w}\right) / \alpha!\right| \leq C_{k} \tau^{\mu-k / 2} \tag{3-11}
\end{equation*}
$$

for $\tau>1$ and $k=1,2, \ldots$..
An inspection of the proof of [Hörmander 1985b, Lemma 26.4.16] shows that the result is still applicable if $\operatorname{Im} w>0$ everwhere. This is also used without mention in [Hörmander 1985b] when proving the necessity of condition $(\Psi)$. Thus the statement holds if $\operatorname{Im} w>0$ except possibly at a point $y$ where
$w^{\prime}(y)=\eta \in \mathbb{R}^{n-1} \backslash 0$ and $\operatorname{Im} w^{\prime \prime}$ is positive definite. We will also use this fact, but we have refrained from altering the statement of the lemma.

If $q$ is homogeneous of degree $\mu$, then the sum in (3-11) consists (apart from the factor $e^{i \tau w}$ ) of terms that are homogeneous in $\tau$ of degree $\mu, \mu-1, \ldots$. The terms of degree $\mu$ are those in

$$
\begin{equation*}
\phi \sum q^{(\alpha)}(x, \tau \eta)\left(\tau w^{\prime}(x)-\tau \eta\right)^{\alpha} / \alpha! \tag{3-12}
\end{equation*}
$$

which is the Taylor expansion at $\tau \eta$ of $q\left(x, \tau w^{\prime}\right)$. In this way one can give meaning to the expression $q\left(x, \tau w^{\prime}\right)$ even though $q(x, \xi)$ may not be defined for complex $\xi$. The terms of degree $\mu-1$ where $\phi$ is differentiated are similarly

$$
\sum_{k=1}^{n-1} q^{(k)}\left(x, \tau w^{\prime}(x)\right) D_{k} \phi
$$

where $q^{(k)}$ should be replaced by the Taylor expansion at $\tau \eta$ representing the value at $\tau w^{\prime}(x)$, as in (3-12). In the present case we have

$$
w_{x}^{\prime}(t, x)-\xi^{0}=i x-\left(t^{2} / 2\right) \xi^{0}
$$

so the expression $q_{-j}\left(t, x, w_{x}^{\prime}(t, x)\right)$ is given meaning if it is replaced by a finite Taylor expansion

$$
\sum_{\beta} q_{-j}^{(\beta)}\left(t, x, \xi^{0}\right)\left(w_{x}^{\prime}(t, x)-\xi^{0}\right)^{\beta} /|\beta|!
$$

of sufficiently high order.
Using the classicality of $R^{*}$ we have

$$
\sigma_{R^{*}}(t, x, \xi)-\sum_{j=-1}^{M} q_{-j}(t, x, \xi) \in \Psi_{\mathrm{cl}}^{-M-1}\left(\mathbb{R}^{n}\right)
$$

so there is a symbol $a \in S_{\mathrm{cl}}^{-M-1}\left(\mathbb{R}^{n} \times \mathbb{R}^{n-1}\right)$ such that

$$
a(t, x, D)=R^{*}(t, x, D)-\sum_{j=-1}^{M} q_{-j}(t, x, D) \quad \bmod \Psi^{-\infty}\left(\mathbb{R}^{n}\right)
$$

By (3-2) and (3-10) it is clear that $w$ satisfies the conditions of Lemma 3.3, so

$$
a(t, x, D) v_{\tau}=a\left(t, x, \tau \xi^{0}\right) v_{\tau}+\mathbb{O}\left(\tau^{-M-3 / 2}\right)=\tau^{-M-1} a\left(t, x, \xi^{0}\right) v_{\tau}+\mathbb{O}\left(\tau^{-M-3 / 2}\right)
$$

which implies that $\left|a(t, x, D) v_{\tau}\right| \leq C \tau^{-M-1}$. If we for each $-1 \leq j \leq M$ write

$$
\left|q_{-j}(t, x, D) v_{\tau}-\sum_{|\alpha|<k_{j}} q_{-j}^{(\alpha)}\left(t, x, \tau \xi^{0}\right)\left(D_{x}-\tau \xi^{0}\right)^{\alpha} v_{\tau} / \alpha!\right| \leq C_{k_{j}} \tau^{-j-k_{j} / 2}
$$

with $k_{j}=2 M-2 j+1$, then

$$
R^{*}(t, x, D) v_{\tau}=\sum_{j=-1}^{M} \sum_{|\alpha|<k_{j}} q_{-j}^{(\alpha)}\left(t, x, \tau \xi^{0}\right)\left(D_{x}-\tau \xi^{0}\right)^{\alpha} v_{\tau} / \alpha!+\mathbb{O}\left(\tau^{-M-1 / 2}\right)
$$

Now recall the discussion above regarding the homogeneity of the terms in (3-11), and choose $M \geq \kappa$, where $\kappa$ is an integer satisfying (3-9). Then

$$
\begin{aligned}
R^{*}(t, x, D) v_{\tau} & =e^{i \tau w} \sum_{j=-1}^{M} \sum_{|\alpha| \leq 2 M-2 j} q_{-j}^{(\alpha)}\left(t, x, \tau w_{x}^{\prime}(t, x)\right) D^{\alpha} \phi \\
& =e^{i \tau w} \sum_{j=-1}^{M} \sum_{|\alpha| \leq 2 M-2 j} \tau^{-j-|\alpha|} q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}(t, x)\right) D^{\alpha} \phi \\
& =e^{i \tau w} \sum_{J=-1}^{M} \tau^{-J} \lambda_{J}(t, x)
\end{aligned}
$$

with an error of order $\mathcal{O}\left(\tau^{-\kappa-1 / 2}\right)$, where

$$
\begin{equation*}
\lambda_{J}(t, x)=\sum_{j+|\alpha|=J} q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}(t, x)\right) D^{\alpha} \phi \quad \text { for } j \geq-1 . \tag{3-13}
\end{equation*}
$$

As before, $q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}(t, x)\right)$ should be replaced by a finite Taylor expansion at $\xi^{0}$ of sufficiently high order representing the value at $w_{x}^{\prime}(t, x)$. In view of (3-5), this yields

$$
I_{\tau}=\tau^{n} \int H(\tau t, \tau x) e^{i \tau w(t, x)}\left(\sum_{J=-1}^{\kappa} \tau^{-J} \lambda_{J}(t, x)+\mathcal{O}\left(\tau^{-\kappa-1 / 2}\right)\right) d t d x
$$

After the change of variables $(\tau t, \tau x) \mapsto(t, x)$ we find that

$$
\begin{equation*}
I_{\tau}=\int H(t, x) e^{i \tau w(t / \tau, x / \tau)}\left(\sum_{J=-1}^{\kappa} \tau^{-J} \lambda_{J}(t / \tau, x / \tau)+\mathbb{O}\left(\tau^{-\kappa-1 / 2}\right)\right) d t d x \tag{3-14}
\end{equation*}
$$

To illustrate how we will proceed to prove Theorem 2.20 by contradiction, let us for the moment assume that $q_{1}\left(0,0, \xi^{0}\right) \neq 0$, where $\xi^{0}=(0, \ldots, 0,1) \in \mathbb{R}^{n-1}$. Since

$$
\begin{equation*}
\lambda_{-1}(t / \tau, x / \tau)=\phi(t / \tau, x / \tau) \sum_{\beta} q_{1}^{(\beta)}\left(t / \tau, x / \tau, \xi^{0}\right)\left(w_{x}^{\prime}(t / \tau, x / \tau)-\xi^{0}\right)^{\beta} /|\beta|! \tag{3-15}
\end{equation*}
$$

where

$$
\begin{equation*}
w_{x}^{\prime}(t / \tau, x / \tau)-\xi^{0}=i x / \tau-\left(t^{2} /\left(2 \tau^{2}\right)\right) \xi^{0}=\mathbb{O}\left(\tau^{-1}\right) \tag{3-16}
\end{equation*}
$$

and (3-10) implies that $\tau w(t / \tau, x / \tau) \rightarrow x_{n-1}$ as $\tau \rightarrow \infty$, we obtain

$$
\lim _{\tau \rightarrow \infty} I_{\tau} / \tau=\int H(t, x) e^{i x_{n-1}} \phi(0,0) q_{1}\left(0,0, \xi^{0}\right) d t d x
$$

Since we may choose $\phi \neq 0$ at the origin, the limit above will then not be equal to 0 for a suitable choice of $H$. However, this contradicts (3-6).

Now assume that $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}\left(0,0, \xi^{0}\right)$ is the first nonvanishing Taylor coefficient with respect to the ordering $>_{t}$, and let

$$
\begin{equation*}
m=j_{0}+k_{0}+\left|\alpha_{0}\right|+\left|\beta_{0}\right| \tag{3-17}
\end{equation*}
$$

so that $m<\kappa$ by (3-9). Note that $\alpha_{0}, \beta_{0} \in \mathbb{N}^{n-1}$ and that the integer $k_{0}$ accounts for derivatives in $t$ while there is no corresponding term for derivatives in the Fourier transform of $t$ since the $q_{-j}$ are independent of this variable. Note also that since $j_{0}$ is permitted to be -1 , we have $0 \leq k_{0},\left|\alpha_{0}\right|,\left|\beta_{0}\right| \leq m+1$.

To use our assumption we will need for each term $q_{-j}^{(\beta+\gamma)}\left(t / \tau, x / \tau, \xi^{0}\right)$ in the Taylor expansion of $q_{-j}^{(\gamma)}\left(t / \tau, x / \tau, w_{x}^{\prime}(t / \tau, x / \tau)\right)$ (as it appears in (3-13)) at $\xi^{0}$ to consider Taylor expansions in $t$ and $x$ at the origin. Note that for given $j$ and $\gamma$, it suffices to consider finite Taylor expansions of $q_{-j}^{(\gamma)}$ of order $\kappa-j-|\gamma|$ by (3-14) and (3-16). For each $j$ and $\gamma$ we thus write

$$
\begin{aligned}
& q_{-j}^{(\gamma)}\left(t / \tau, x / \tau, w_{x}^{\prime}(t / \tau, x / \tau)\right) \\
& =\sum_{k+|\alpha|+|\beta| \leq \kappa-j-|\gamma|}\left(\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}\right)\left(0,0, \xi^{0}\right) \tau^{-k-|\alpha|} t^{k} x^{\alpha}\left(w_{x}^{\prime}(t / \tau, x / \tau)-\xi^{0}\right)^{\beta} /(k!|\alpha|!|\beta|!)+\mathbb{O}\left(\tau^{-\kappa-1+j+|\gamma|}\right),
\end{aligned}
$$

where $\left(w_{x}^{\prime}(t / \tau, x / \tau)-\xi^{0}\right)^{\beta}$ should be interpreted by means of (3-16). As we shall see, the term $\left(t^{2} /\left(2 \tau^{2}\right)\right) \xi^{0}$ will not pose any problem, since it is $\mathcal{O}\left(\tau^{-2}\right)$. We have
$\lambda_{J}(t / \tau, x / \tau)=\sum_{j+|\gamma|=J} \sum_{k+|\alpha|+|\beta| \leq \kappa-J}\left(\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}\right)\left(0,0, \xi^{0}\right) D^{\gamma} \phi(t / \tau, x / \tau), ~ \tau^{-k-|\alpha|} t^{k} x^{\alpha}\left(w_{x}^{\prime}(t / \tau, x / \tau)-\xi^{0}\right)^{\beta} /(k!|\alpha|!|\beta|!)+\mathcal{O}\left(\tau^{-\kappa-1+J}\right)$,
where $-1 \leq j \leq J$. If we are only interested in terms of order $\tau^{-m}$ in (3-14), we can use the assumption that $\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}\left(0,0, \xi^{0}\right)=0$ for all $-1 \leq j+k+|\alpha|+|\beta|+|\gamma|<m$ to let the term $\left(t^{2} /\left(2 \tau^{2}\right)\right) \xi^{0}$ from (3-16) be absorbed by the error term in the expression above. This yields

$$
\sum_{J=-1}^{m} \tau^{-J} \lambda_{J}(t / \tau, x / \tau)=\sum_{j+k+|\alpha|+|\beta|+|\gamma|=m} \begin{aligned}
&\left(\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}\right)\left(0,0, \xi^{0}\right) \\
& \times D^{\gamma} \phi(t / \tau, x / \tau) \tau^{-m} t^{k} x^{\alpha}(i x)^{\beta} /(k!|\alpha|!|\beta|!)+\mathbb{O}\left(\tau^{-m-1}\right)
\end{aligned}
$$

where we use $J=j+|\gamma|$ together with the fact that we get a factor $\tau^{-|\beta|}$ from $\left(w_{x}^{\prime}(t / \tau, x / \tau)-\xi^{0}\right)^{\beta}$ by (3-16). Thus,

$$
\begin{aligned}
\lim _{\tau \rightarrow \infty} \tau^{m} I_{\tau}=\int H(t, x) & e^{i x_{n-1}} \\
& \times\left(\sum_{j+k+|\alpha|+|\beta|+|\gamma|=m} t^{k} x^{\alpha}(i x)^{\beta}\left(\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}\right)\left(0,0, \xi^{0}\right) D^{\gamma} \phi(0,0) /(k!|\alpha|!|\beta|!)\right) d t d x
\end{aligned}
$$

Now choose $\phi$ such that $D^{\beta_{0}} \phi(0,0)=1$, but $D^{\gamma} \phi(0,0)=0$ for all other $\gamma$ such that $|\gamma| \leq\left|\beta_{0}\right|$. This is possible by the discussion following (3-1). By (3-17) and our choice of the ordering $>_{t}$, we have $\partial_{t}^{k} q_{-j(\alpha)}^{\left(\beta+\beta_{0}\right)}\left(0,0, \xi^{0}\right)=0$ for all $\beta$ such that $|\beta|>0$ as long as $j+k+|\alpha|+|\beta|+\left|\beta_{0}\right|=m$. Hence, with this choice of $\phi$, the last expression takes the form

$$
\begin{equation*}
\lim _{\tau \rightarrow \infty} \tau^{m} I_{\tau}=\int H(t, x) e^{i x_{n-1}}\left(\sum_{j+k+|\alpha|+\left|\beta_{0}\right|=m} t^{k} x^{\alpha}\left(\partial_{t}^{k} q_{-j(\alpha)}^{\left(\beta_{0}\right)}\right)\left(0,0, \xi^{0}\right) /(k!|\alpha|!)\right) d t d x \tag{3-18}
\end{equation*}
$$

where as usual $j$ is allowed to be -1 so that $j \in\left[-1, m-\left|\beta_{0}\right|\right]$ in (3-18). Now some of the Taylor coefficients in (3-18) may be zero, in particular, the expression may well contain Taylor coefficients that preceed $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}\left(0,0, \xi^{0}\right)$, and those are by assumption zero. However, we claim that if at least one of
the Taylor coefficients above are nonzero, then we may choose $H$ so that the limit is nonzero. Indeed, if that were not the case then the expression within brackets in (3-18) would be a polynomial with infinitely many zeros, and thus it would have to have vanishing coefficients. Since this violates our assumption, we conclude that the limit is nonzero. However, this contradicts (3-6), which proves Theorem 2.20.

## 4. Proof of Theorem 2.21

In this section we shall give the proof of Theorem 2.21, using ideas taken from [Hörmander 1963] together with the approach used to prove [Hörmander 1985b, Theorem 26.4.7']. As in the previous section, we aim to use Lemma 2.3 to estimate the operator $R\left(x, D_{x^{\prime}}\right)$ on approximate solutions of the equation $P^{*} v=0$, concentrated near

$$
\begin{equation*}
\Gamma^{\prime}=\left\{\left(x_{1}, x^{\prime}, 0, \xi^{\prime}\right): x_{1} \in I^{\prime}\right\} \subset T^{*}\left(\mathbb{R}^{n}\right) \backslash 0 \tag{4-1}
\end{equation*}
$$

The proofs will be similar, but the situation is more complicated now, which will affect the construction of the approximate solutions. We will also have to make some adjustments to the proof of [Hörmander 1985b, Theorem 26.4.7'] to make it work, so a lot of the details will have to be revisited. Our approximate solutions will also differ slightly from the ones used to prove [Hörmander 1985b, Theorem 26.4.7'], so although we will refer directly to results in [ibid.] whenever possible, the formulation of some of these results will be affected. For a more complete description of the approximate solutions, we refer the reader to [Hörmander 1981] or [Hörmander 1985b], where their construction is carried out in greater detail. When proving Theorem 2.21 we may without loss of generality assume that $x^{\prime}=0$ and $\xi^{\prime}=\xi^{0}$ in (4-1). In accordance with the notation in the proof of Theorem 2.19, we shall therefore throughout this section refer to $\Gamma^{\prime}$ simply by $\Gamma$, and we will let $I^{\prime}=\left[a_{0}, b_{0}\right]$.

To simplify notation we shall in what follows write $t$ instead of $x_{1}$ and $x$ instead of $x^{\prime}$. If $N$ is the integer given by Definition 2.1, and $n$ is the dimension, the approximate solutions $v_{\tau}$ will be taken of the form

$$
\begin{equation*}
v_{\tau}(t, x)=\tau^{N+n} e^{i \tau w(t, x)} \sum_{0}^{M} \phi_{j}(t, x) \tau^{-j} . \tag{4-2}
\end{equation*}
$$

Here $\phi_{0}, \phi_{1}, \ldots$ are amplitude functions, and $w$ is a phase function that should satisfy the eikonal equation

$$
\begin{equation*}
\partial w / \partial t-i f(t, x, \partial w / \partial x)=0 \tag{4-3}
\end{equation*}
$$

approximately, where $f$ is the imaginary part of the principal symbol of $P$. We take $w$ of the form

$$
\begin{equation*}
w(t, x)=w_{0}(t)+\langle x-y(t), \eta(t)\rangle+\sum_{2 \leq|\alpha| \leq M} w_{\alpha}(t)(x-y(t))^{\alpha} /|\alpha|! \tag{4-4}
\end{equation*}
$$

where $M$ is a large integer to be determined later and $x=y(t)$ is a smooth real curve. When discussing the functions $w_{\alpha}$ we shall permit ourselves to use the notation $\alpha=\left(\alpha_{1}, \ldots, \alpha_{s}\right)$ for a sequence of $s=|\alpha|$ indices between 1 and the dimension $n-1$ of the $x$ variable, and $w_{\alpha}$ will be symmetric in these indices. If we take $\eta(t)$ to be real-valued and make sure the matrix $\left(\operatorname{Im} w_{j k}\right)$ is positive definite, then $\operatorname{Im} w$ will have a strict minimum when $x=y(t)$ as a function of the $x$ variables.

On the curve $x=y(t)$ the eikonal equation (4-3) is reduced to

$$
\begin{equation*}
w_{0}^{\prime}(t)=\left\langle y^{\prime}(t), \eta(t)\right\rangle+i f(t, y(t), \eta(t)) \tag{4-5}
\end{equation*}
$$

which is the only equation where $w_{0}$ occurs. Hence it can be used to determine $w_{0}$ after $y$ and $\eta$ have been chosen. In particular

$$
\begin{equation*}
d \operatorname{Im} w_{0}(t) / d t=f(t, y(t), \eta(t)) \tag{4-6}
\end{equation*}
$$

In the proof of Theorem 2.20 we could solve the corresponding eikonal equation explicitly. Here this is not possible, so our goal will instead be to make (4-3) valid apart from an error of order $M+1$ in $x-y(t)$. Note that $f(t, x, \xi)$ is not defined for complex $\xi$, but since

$$
\partial w(t, x) / \partial x_{j}-\eta_{j}(t)=\sum w_{\alpha, j}(t)(x-y(t))^{\alpha} /|\alpha|!
$$

(4-3) is given meaning if $f(t, x, \partial w / \partial x)$ is replaced by the finite Taylor expansion

$$
\begin{equation*}
\sum_{|\beta| \leq M} f^{(\beta)}(t, x, \eta(t))(\partial w(t, x) / \partial x-\eta(t))^{\beta} /|\beta|!. \tag{4-7}
\end{equation*}
$$

To compute the coefficient of $(x-y(t))^{\alpha}$ in (4-7) we just have to consider the terms with $|\beta| \leq|\alpha|$. Since $\partial w / \partial t=w_{0}^{\prime}-\left\langle y^{\prime}, \eta\right\rangle+\left\langle x-y, \eta^{\prime}\right\rangle$

$$
+\sum_{2 \leq|\alpha| \leq M} w_{\alpha}^{\prime}(t)(x-y)^{\alpha} /|\alpha|!-\sum_{k} \sum_{1 \leq|\alpha| \leq M-1} w_{\alpha, k}(t)(x-y)^{\alpha} d y_{k} / d t /|\alpha|!
$$

the first order terms in the equation (4-3) give

$$
\begin{equation*}
d \eta_{j} / d t-\sum_{k} w_{j k}(t) d y_{k} / d t=i\left(f_{(j)}(t, y, \eta)+\sum_{k} f^{(k)}(t, y, \eta) w_{j k}(t)\right) \tag{4-8}
\end{equation*}
$$

Note that this is a system of $2 n$ equations

$$
\begin{gather*}
d \eta_{j} / d t-\sum_{k} \operatorname{Re} w_{j k}(t) d y_{k} / d t=-\sum_{k} \operatorname{Im} w_{j k}(t) f^{(k)}(t, y, \eta)  \tag{4-8}\\
\sum_{k} \operatorname{Im} w_{j k}(t) d y_{k} / d t=-f_{(j)}(t, y, \eta)-\sum_{k} \operatorname{Re} w_{j k}(t) f^{(k)}(t, y, \eta) \tag{4-8}
\end{gather*}
$$

since $y$ and $\eta$ are real, and under the assumption that $\operatorname{Im} w_{j k}$ is positive definite these equations can be solved for $d y / d t$ and $d \eta / d t$. We observe that at a point where $f=d f=0$ they just mean that $d y / d t=d \eta / d t=0$.

When $2 \leq|\alpha| \leq M$ we obtain a differential equation

$$
\begin{equation*}
d w_{\alpha} / d t-\sum_{k} w_{\alpha, k} d y_{k} / d t=F_{\alpha}\left(t, y, \eta,\left\{w_{\beta}\right\}\right) \tag{4-9}
\end{equation*}
$$

from (4-3). Here $F_{\alpha}$ is a linear combination of the derivatives of $f$ of order $|\alpha|$ or less, multiplied with polynomials in $w_{\beta}$ with $2 \leq|\beta| \leq|\alpha|+1$. Of course, when $|\alpha|=M$ the sum on the left side of (4-9) should be dropped, and $\beta$ should satisfy $|\beta| \leq|\alpha|$ instead. Altogether (4-8)', (4-8) " and (4-9) form a quasilinear system of differential equations with as many equations as unknowns. Hence we have local
solutions with prescribed initial data. According to [Hörmander 1985b, pages 105-106] we can find a $c>0$ such that the equations (4-8) and (4-9) with initial data

$$
\begin{align*}
& w_{j k}=i \delta_{j k}, \quad w_{\alpha}=0, \quad \text { when } 2<|\alpha| \leq M \text { and } t=\left(a_{0}+b_{0}\right) / 2,  \tag{4-10}\\
& y=x, \quad \eta=\xi, \quad \text { when } t=\left(a_{0}+b_{0}\right) / 2 \tag{4-11}
\end{align*}
$$

have a unique solution in $\left(a_{0}-c, b_{0}+c\right)$ for all $x$ and $\xi$ with $|x|+\left|\xi-\xi^{0}\right|<c$. (Here $\delta_{j k}$ is the Kronecker delta.) Moreover,
(i) $\left(\operatorname{Im} w_{j k}-\delta_{j k} / 2\right)$ is positive definite,
(ii) the map

$$
(x, \xi, t) \mapsto(y, \eta, t), \quad \text { where }|x|+\left|\xi-\xi^{0}\right|<c, a_{0}-c<t<b_{0}+c
$$

is a diffeomorphism.
In the range $X_{c}$ of the map (ii) we let $v$ denote the image of the vector field $\partial / \partial t$ under the map. Thus $v$ is the tangent vector field of the integral curves, and when $f=d f=0$ we have $v=\partial / \partial t$. By assumption $f=0$ implies $\partial f / \partial t \leq 0$ in a neighborhood of $\Gamma$ (see (2-18)), so if $c$ is small enough this also holds in $X_{c}$. An application of [Hörmander 1985b, Lemma 26.4.11] now yields that $f$ must have a change of sign from - to + along an integral curve of $v$ in $X_{c}$, for otherwise there would be no such sign change for increasing $t$ and fixed $(x, \xi)$, and that contradicts the hypothesis in Theorem 2.21. By (4-6) this means that $\operatorname{Im} w_{0}(t)$ will start decreasing and end increasing, so the minimum is attained at an interior point. We can normalize the minimum value to zero and have then for a suitable interval of $t$ that $\operatorname{Im} w_{0}>0$ at the end points and $\operatorname{Im} w_{0}=0$ at some interior point. Since $\operatorname{Re} w_{0}$ is given by (4-5) we can at this interior point also normalize the value of $\operatorname{Re} w_{0}$ to zero. This completes the proof of [Hörmander 1985b, Lemma 26.4.14]. However, in order to prove Theorem 2.21 when $a_{0}<b_{0}$ we shall need the following stronger result.

Lemma 4.1. Assume that the hypotheses of Theorem 2.21 are fulfilled, the variables being denoted $(t, x)$ now. Then given $M \in \mathbb{N}$ we can find
(i) a curve $t \mapsto(t, y(t), 0, \eta(t)) \in \mathbb{R}^{2 n}$, with $a^{\prime} \leq t \leq b^{\prime}$ as close to $\Gamma$ as desired,
(ii) $C^{\infty}$ functions $w_{\alpha}(t)$ for $2 \leq|\alpha| \leq M$, with $\left(\operatorname{Im} w_{j k}-\delta_{j k} / 2\right)$ positive definite when $a^{\prime} \leq t \leq b^{\prime}$,
(iii) a function $w_{0}(t)$ with $\operatorname{Im} w_{0}(t) \geq 0$ for $a^{\prime} \leq t \leq b^{\prime}, \operatorname{Im} w_{0}\left(a^{\prime}\right)>0$, $\operatorname{Im} w_{0}\left(b^{\prime}\right)>0$ and $\operatorname{Re} w_{0}\left(c^{\prime}\right)=$ $\operatorname{Im} w_{0}\left(c^{\prime}\right)=0$ for some $c^{\prime} \in\left(a^{\prime}, b^{\prime}\right)$
such that (4-4) is a formal solution to (4-3) with an error of order $\mathcal{O}\left(|x-y(t)|^{M+1}\right)$. If $a_{0}<b_{0}$ then (iii) can be improved in the sense that if $\varrho \geq 0$ is the number given by Theorem 2.21 , then we can for any $\varepsilon>\varrho$ find
(iii)' a function $w_{0}(t)$ with $\operatorname{Im} w_{0}(t) \geq 0, a^{\prime} \leq t \leq b^{\prime}, \operatorname{Im} w_{0}\left(a^{\prime}\right)>0, \operatorname{Im} w_{0}\left(b^{\prime}\right)>0$ and $\operatorname{Re} w_{0}(t)=$ $\operatorname{Im} w_{0}(t)=0$ for all $t \in\left[a_{0}+\varepsilon, b_{0}-\varepsilon\right]$.

Proof. In view of [Hörmander 1985b, Lemma 26.4.14] we only need to prove (iii)'.
Let $\varepsilon>\varrho$, and let $I_{\varepsilon}=\left[a_{0}+\varepsilon, b_{0}-\varepsilon\right]$. By the hypotheses of Theorem 2.21, there is a neighborhood U of

$$
\Gamma_{\varepsilon}=\left\{\left(t, 0,0, \xi^{0}\right): t \in I_{\varepsilon}\right\}
$$

where $f$ vanishes identically. Take $\delta>0$ sufficiently small so that

$$
t \in I_{\varepsilon},|x|+\left|\xi-\xi^{0}\right|<\delta \quad \text { implies } \quad(t, x, 0, \xi) \in U
$$

As above we can find $c>0$ such that the equations (4-8) and (4-9) with initial data (4-10) and (4-11) have a unique solution in $\left(a_{0}-c, b_{0}+c\right)$ for all $x$ and $\xi$ with $|x|+\left|\xi-\xi^{0}\right|<c$. Since the map

$$
(x, \xi, t) \mapsto(y, \eta, t), \quad \text { where }|x|+\left|\xi-\xi^{0}\right|<c, a_{0}-c<t<b_{0}+c
$$

is a diffeomorphism, we can choose $c$ small enough so that if $(y, \eta, t)$ is in the range $X_{c}$ of this map, then $|y|+\left|\eta-\xi^{0}\right|<\delta$. As we have seen, $f$ must change sign from - to + along an integral curve of $v$ in $X_{c}$ if $c$ is small enough, where in $X_{c}$ we denote by $v$ the image of the vector field $\partial / \partial t$ under the map. Let this integral curve be given by

$$
\gamma(t)=(t, y(t), 0, \eta(t)) \in \mathbb{R}^{2 n} \quad \text { for } a^{\prime} \leq t \leq b^{\prime}
$$

for some choice of $a^{\prime}$ and $b^{\prime}$ such that $a_{0}-c<a^{\prime}, b^{\prime}<b_{0}+c$ and

$$
f\left(a^{\prime}, y\left(a^{\prime}\right), \eta\left(a^{\prime}\right)\right)<0<f\left(b^{\prime}, y\left(b^{\prime}\right), \eta\left(b^{\prime}\right)\right)
$$

Recall that at a point where $f=d f=0$ the equations (4-8)' and (4-8)" imply that $d y / d t=d \eta / d t=0$. Since $f$ vanishes identically on $\gamma$ for $t \in I_{\varepsilon}$ and the function $w_{0}$ is determined by (4-5), this proves the lemma after a suitable normalization.

Note that if $\Gamma$ is a point then by Lemma 4.1 we can obtain a sequence $\left\{\gamma_{j}\right\}$ of curves

$$
\gamma_{j}(t)=\left(t, y_{j}(t), 0, \eta_{j}(t)\right) \quad \text { for } a_{j}^{\prime} \leq t \leq b_{j}^{\prime}
$$

approaching $\Gamma$, which implies that at $t=c_{j}^{\prime}$ we have

$$
\left(c_{j}^{\prime}, y_{j}\left(c_{j}^{\prime}\right), 0, \eta_{j}\left(c_{j}^{\prime}\right)\right) \rightarrow \Gamma \quad \text { as } j \rightarrow \infty
$$

in $T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$, where $c_{j}^{\prime}$ is the point where $\operatorname{Re} w_{0 j}=\operatorname{Im} w_{0 j}=0$. Similarly, if $\Gamma$ is an interval and $\varrho \geq 0$ is the number given by Theorem 2.21, then for any point $\omega$ in the interior of $\Gamma_{\varrho}$ we can use Lemma 4.1 to obtain a sequence $\left\{\gamma_{j}\right\}$ of curves approaching $\Gamma$ and a sequence $\left\{w_{0 j}\right\}$ of functions such that for each $j$ there exists a point $\omega_{j} \in \gamma_{j}$ with $\omega_{j}=\gamma_{j}\left(t_{j}\right)$ that can be chosen so that $\operatorname{Re} w_{0 j}\left(t_{j}\right)=\operatorname{Im} w_{0 j}\left(t_{j}\right)=0$ and $\omega_{j} \rightarrow \omega$ as $j \rightarrow \infty$. This will be crucial in proving Theorem 2.21. Our strategy is to show that all the terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients at $\omega_{j}$, or at $\left(c_{j}^{\prime}, y_{j}\left(c_{j}^{\prime}\right), 0, \eta_{j}\left(c_{j}^{\prime}\right)\right)$ when $\Gamma$ is a point. Theorem 2.21 will then follow by continuity. In what follows we will suppress the index $j$ to simplify notation.

Let $K$ and $\Omega$ be the cones given by Theorem 2.21, and suppose that the function $w$ given by (4-4) is a formal solution to (4-3) with an error of order $\mathcal{O}\left(|x-y(t)|^{M+1}\right)$ in a neighborhood $Y$ of

$$
\left\{(t, 0): a_{0} \leq t \leq b_{0}\right\} \subset \mathbb{R}^{n}
$$

with $K \subset T^{*}(Y)$, such that $\operatorname{Im} w>0$ in $Y$ except on a compact nonempty subset $T$ of the curve $x=y(t)$, with $\left(t_{0}, y\left(t_{0}\right)\right) \in T$ and $w=0$ on $T$. We want to show that all the terms in the asymptotic sum of the symbol of $R$ have vanishing Taylor coefficients at ( $\left.t_{0}, y\left(t_{0}\right), 0, \eta\left(t_{0}\right)\right)$. By part (i) of Lemma 4.1 we can choose $w$ so that

$$
\begin{equation*}
\Gamma_{0}=\{(t, x, \partial w(t, x) / \partial t, \partial w(t, x) / \partial x):(t, x) \in T\} \tag{4-12}
\end{equation*}
$$

is contained in $\Omega$. This is done to ensure that if $A$ is a given pseudodifferential operator with wavefront set contained in the complement of $K$, then $W F(A)$ does not meet the cone generated by $\Gamma_{0}$.

We now turn our attention to the amplitude functions $\phi_{j}$. With the exception of $\phi_{0}$, which will be of great interest to us, we will not be very thorough in describing them. Suffice it to say that these functions can be chosen so that if $P^{*}$ is the adjoint of $P$ then

$$
\begin{equation*}
\left\|P^{*} v_{\tau}\right\|_{(\nu)} \leq C \tau^{N+n+\nu+(1-M) / 2} \tag{4-13}
\end{equation*}
$$

where $M$ is the number given by (4-2). The procedure begins by setting

$$
\phi_{0}(t, x)=\sum_{|\alpha|<M} \phi_{0 \alpha}(t)(x-y(t))^{\alpha}
$$

with $y(t)$ as above, and having $\phi_{0 \alpha}$ satisfy the linear system of ordinary differential equations

$$
\begin{equation*}
D_{t} \phi_{0 \alpha}+\sum_{|\beta|<M} a_{\alpha \beta} \phi_{0 \beta}=0 \tag{4-14}
\end{equation*}
$$

In the same way we then successively choose $\phi_{j}$ and obtain (4-13). The precise details can be found in [Hörmander 1981, pages 87-89], or in [Hörmander 1985b, pages 107-110]. Note that we for any positive integer $J<M$ can solve the equations that determine $\phi_{0}$ so that at the point $\left(t_{0}, y\left(t_{0}\right)\right) \in T$ we have $D_{x}^{\alpha} \phi_{0}\left(t_{0}, y\left(t_{0}\right)\right)=0$ for all $|\alpha| \leq J$ except for one index $\alpha,|\alpha|=J$. This will be important later on. Note also that the estimate (4-13) is not affected if the functions $\phi_{j}$ are multiplied by a cutoff function in $C_{0}^{\infty}(Y)$ that is 1 in a neighborhood of $T$. Since the $\phi_{j}$ will be irrelevant outside of $Y$ for large $\tau$ by construction, we can in this way choose them to be supported in $Y$ so that $v_{\tau} \in C_{0}^{\infty}(Y)$.

Having completed the construction of the approximate solutions, we are now ready to start to follow the proof of Theorem 2.20. To get the estimates for the right side of (2-1) when $v$ is an approximate solution, we shall need the following two results. The first, corresponding to Lemma 3.1, is taken from [Hörmander 1985b]. Observe that here it is stated for our approximate solutions which differ from those in [ibid.] by a factor of $\tau^{N+n}$, which explains the difference in appearance. Note also that although we will not use the lower bound for the approximate solutions, that estimate is included so as not to alter the statement.

Lemma 4.2 [Hörmander 1985b, Lemma 26.4.15]. Let $X \subset \mathbb{R}^{n}$ be open, and let $v_{\tau}$ be defined by (4-2), where $w \in C^{\infty}(X), \phi_{j} \in C_{0}^{\infty}(X), \operatorname{Im} w \geq 0$ in $X$ and $d \operatorname{Re} w \neq 0$. For any positive integer $m$ we then have

$$
\begin{equation*}
\left\|v_{\tau}\right\|_{(-m)} \leq C \tau^{N+n-m} \quad \text { for } \tau>1 \tag{4-15}
\end{equation*}
$$

If $\operatorname{Im} w\left(t_{0}, x_{0}\right)=0$ and $\phi_{0}\left(t_{0}, x_{0}\right) \neq 0$ for some $\left(t_{0}, x_{0}\right) \in X$ then

$$
\left\|v_{\tau}\right\|_{(-m)} \geq c \tau^{N+n / 2-m} \quad \text { for } \tau>1
$$

and for some $c>0$. If $\tilde{\Gamma}$ is the cone generated by

$$
\left\{\left(t, x, \partial_{t} w(t, x), \partial_{x} w(t, x)\right):(t, x) \in \bigcup_{j} \operatorname{supp} \phi_{j}, \operatorname{Im} w(t, x)=0\right\}
$$

then $\tau^{k} v_{\tau} \rightarrow 0$ in $\mathscr{D}_{\tilde{\Gamma}}^{\prime}$ as $\tau \rightarrow \infty$; hence $\tau^{k} A v_{\tau} \rightarrow 0$ in $C^{\infty}\left(\mathbb{R}^{n}\right)$ if $A$ is a pseudodifferential operator with $W F(A) \cap \tilde{\Gamma}=\varnothing$ and $k$ is any real number.

Proposition 4.3. Assume that the hypotheses of Theorem 2.21 are fulfilled, the variables being denoted $(t, x)$ now, and let $v_{\tau}$ be given by (4-2), where $w \in C^{\infty}(Y), \phi_{j} \in C_{0}^{\infty}(Y), \operatorname{Im} w \geq 0$ in $Y$ and $d \operatorname{Re} w \neq 0$. Here $Y$ is a neighborhood of $\left\{(t, 0): a_{0} \leq t \leq b_{0}\right\}$ such that $K \subset T^{*}(Y)$. Let $H(t, x) \in C_{0}^{\infty}\left(\mathbb{R} \times \mathbb{R}^{n-1}\right)$ and set

$$
\begin{equation*}
h_{\tau}(t, x)=\tau^{-N} H\left(\tau\left(t-t_{0}\right), \tau(x-y(t))\right) \tag{4-16}
\end{equation*}
$$

where $N$ is the positive integer given by Definition 2.1 for the operators $R$ and $P$ in Theorem 2.21. Then $h_{\tau} \in H_{(N)}\left(\mathbb{R}^{n}\right)$ for all $\tau \geq 1$ and $\left\|h_{\tau}\right\|_{(N)} \leq C$, where the constant depends on $H$ but not on $\tau$. Furthermore, if $M$ is the integer given by the definition of $v_{\tau}$ in (4-2) so that (4-13) holds, and $I_{\tau}$ is the integral

$$
\begin{equation*}
I_{\tau}=\left(R^{*} v_{\tau}, \overline{h_{\tau}}\right) \tag{4-17}
\end{equation*}
$$

where $R^{*}$ is the adjoint of $R(t, x, D)$, then for any positive integer $\kappa$ there exists a constant $C$ such that $\left|I_{\tau}\right| \leq C \tau^{-\kappa}$ if $M=M(\kappa)$ is sufficiently large.

Proof. In Section 3, one easily obtains a formula for the Fourier transform of the corresponding function $h_{\tau}$ (see (3-4)), which yields the estimates needed to show that $h_{\tau} \in H_{(N)}$. Here we shall instead use the equality

$$
\iint\left|h_{\tau}(t, x)\right|^{2} d t d x=\tau^{-2 N} \iint\left|H\left(\tau\left(t-t_{0}\right), \tau(x-y(t))\right)\right|^{2} d t d x
$$

which shows that if $\tau \geq 1$ then $D_{t}^{j} D_{x}^{\alpha} h_{\tau} \in L^{2}\left(\mathbb{R}^{n}\right)$ for all $(j, \alpha) \in \mathbb{N} \times \mathbb{N}^{n-1}$ such that $j+|\alpha| \leq N+[n / 2]$. Hence, by using the equivalent norm on $H_{(N)}\left(\mathbb{R}^{n}\right)$ given by

$$
\left\|h_{\tau}\right\|_{(N)}=\sum_{j+|\alpha| \leq N}\left\|D_{t}^{j} D_{x}^{\alpha} h_{\tau}\right\|_{(0)}
$$

we find that $\left\{h_{\tau}\right\}_{\tau \geq 1}$ is a bounded one parameter family in $H_{(N)}\left(\mathbb{R}^{n}\right)$, which proves the first assertion of the proposition.

To prove the second part, let $\kappa$ be an arbitrary positive integer, and let $\nu$ be the positive integer given by Lemma 2.3 (applied to the operator $R$ instead of $Q$ ) so that (2-1) holds for the choice of seminorm $\left\|P^{*} v\right\|_{(v)}$ in the right side. If we choose

$$
\begin{equation*}
(1-M) / 2 \leq-N-n-v-\kappa, \tag{4-18}
\end{equation*}
$$

and recall (4-13), then

$$
\begin{equation*}
\left\|P^{*} v_{\tau}\right\|_{(\nu)} \leq C \tau^{-\kappa} \tag{4-19}
\end{equation*}
$$

Since supp $H$ is compact, we can find a bounded open ball containing $\operatorname{supp} h_{\tau}$ for all $\tau \geq 1$. Hence $h_{\tau} \in H_{(N)}\left(\mathbb{R}^{n}\right)$ has compact support and $v_{\tau} \in C_{0}^{\infty}(Y)$, so the result now follows by the estimate (2-4) together with Lemma 4.2.

To shorten the notation we will from now on assume that $t_{0}=0$, so that $w(0, y(0))=0$. As in the proof of Theorem 2.20 it suffices to show that all terms in the asymptotic expansion of the symbol of $R^{*}$, given by

$$
\sigma_{R^{*}}=q_{1}(t, x, \xi)+q_{0}(t, x, \xi)+\cdots
$$

with $q_{j}$ homogeneous of degree $j$ in $\xi$, have vanishing Taylor coefficients at $(0, y(0), \eta(0))$. The method will be to argue by contradiction that if not, then Proposition 4.3 does not hold. Therefore, let us assume that $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0))$ is the first nonvanishing Taylor coefficient with respect to the ordering $>_{t}$ given by Definition 3.2, and let

$$
\begin{equation*}
m=j_{0}+k_{0}+\left|\alpha_{0}\right|+\left|\beta_{0}\right| . \tag{4-20}
\end{equation*}
$$

Now let $\kappa$ be a positive integer such that $m<\kappa$, and sort the terms in $I_{\tau}$, given by (4-17), with respect to homogeneity degree in $\tau$. We can use Lemma 3.3 and the classicality of the symbol $\sigma_{R^{*}}$ to write

$$
\begin{aligned}
R^{*}(t, x, D) v_{\tau} & =\sum_{j=-1}^{M^{\prime}} q_{-j}(t, x, D) v_{\tau}+\mathbb{O}\left(\tau^{N+n-M^{\prime}-1}\right) \\
& =\sum_{j=-1}^{M^{\prime}} \sum_{l=0}^{M} \tau^{N+n-l} q_{-j}(t, x, D)\left(e^{i \tau w} \phi_{l}\right)+\mathcal{O}\left(\tau^{N+n-M^{\prime}-1}\right)
\end{aligned}
$$

for some large number $M^{\prime}$. Note that (4-18) implies a lower bound on $M$, but as we shall see below, we must also make sure to pick $M>2 M^{\prime}+1$. For each $j$ we then estimate $q_{-j}(t, x, D)\left(e^{i \tau w} \phi_{l}\right)$ using (3-11) with $k=M-1-2 j$, so that

$$
q_{-j}(t, x, D)\left(e^{i \tau w} \phi_{l}\right)=\sum_{|\alpha|<M-1-2 j} q_{-j}^{(\alpha)}(t, x, \tau \eta)(D-\tau \eta)^{\alpha}\left(\phi_{l} e^{i \tau w}\right) / \alpha!
$$

with an error of order $\mathbb{O}\left(\tau^{(1-M) / 2}\right)$. Recalling (4-18) and the discussion following Lemma 3.3 regarding the homogeneity of the terms in (3-11), this yields, for sufficiently large $M^{\prime}$,

$$
\begin{align*}
R^{*}(t, x, D) v_{\tau} & =\sum_{j=-1}^{M^{\prime}} \sum_{l=0}^{M} \tau^{N+n-l} e^{i \tau w} \sum_{|\alpha|<M-1-2 j} q_{-j}^{(\alpha)}\left(t, x, \tau w_{x}^{\prime}\right) D^{\alpha} \phi_{l}+\mathscr{O}\left(\tau^{-\kappa-1}\right) \\
& =\tau^{N+n} e^{i \tau w} \sum_{j=-1}^{M^{\prime}} \sum_{l=0}^{M} \sum_{|\alpha|<M-1-2 j} \tau^{-j-|\alpha|-l} q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}\right) D^{\alpha} \phi_{l}+\mathbb{O}\left(\tau^{-\kappa-1}\right) \tag{4-21}
\end{align*}
$$

Note that $\tau^{-j-|\alpha|-l} q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}\right) D^{\alpha} \phi_{l}$ is now homogeneous of order $-j-|\alpha|-l$ in $\tau$, and that as before, $q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}\right)$ should be replaced by a finite Taylor expansion at $\eta$ of sufficiently high order. For each $-1 \leq J \leq \kappa$, collect all terms of the form $\tau^{-j-|\alpha|-l} q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}\right) D^{\alpha} \phi_{l}$ in (4-21) that are homogeneous of order $-J$ in $\tau$, that is, all terms that satisfy $j+|\alpha|+l=J$ for $j \geq-1$, and $|\alpha|, l \geq 0$. If

$$
\lambda_{J}(t, x)=\sum_{j+|\alpha|+l=J} q_{-j}^{(\alpha)}\left(t, x, w_{x}^{\prime}(t, x)\right) D^{\alpha} \phi_{l}(t, x)
$$

for the permitted values of $j$ and $l$, then

$$
I_{\tau}=\tau^{n} \iint H(\tau t, \tau(x-y(t)))\left(e^{i \tau w(t, x)} \sum_{J=-1}^{\kappa} \tau^{-J} \lambda_{J}(t, x)+\mathbb{O}\left(\tau^{-\kappa-1}\right)\right) d t d x
$$

After the change of variables $(\tau t, \tau(x-y(t))) \mapsto(t, x)$ we obtain

$$
\begin{equation*}
I_{\tau}=\iint H(t, x)\left(e^{i \tau w(t / \tau, x / \tau+y(t / \tau))} \sum_{J=-1}^{\kappa} \tau^{-J} \lambda_{J}(t / \tau, x / \tau+y(t / \tau))+\mathcal{O}\left(\tau^{-\kappa-1}\right)\right) d t d x \tag{4-22}
\end{equation*}
$$

where

$$
\lambda_{J}(t / \tau, x / \tau+y(t / \tau))
$$

$$
\begin{equation*}
=\sum_{j+|\alpha|+l=J} D^{\alpha} \phi_{l}(t / \tau, x / \tau+y(t / \tau)) q_{-j}^{(\alpha)}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right) \tag{4-23}
\end{equation*}
$$

Recall that $w_{0}(0)=0$, which together with (4-4) implies

$$
i \tau w(t / \tau, x / \tau+y(t / \tau))=i t w_{0}^{\prime}(0)+i\langle x, \eta(t / \tau)\rangle+\mathbb{O}\left(\tau^{-1}\right)
$$

Hence

$$
\begin{equation*}
\lim _{\tau \rightarrow \infty} e^{i \tau w(t / \tau, x / \tau+y(t / \tau))}=e^{i t w_{0}^{\prime}(0)+i\langle x, \eta(0)\rangle} \tag{4-24}
\end{equation*}
$$

In the sequel we shall also need

$$
\begin{equation*}
\partial w / \partial x_{j}(t / \tau, x / \tau+y(t / \tau))-\eta_{j}(t / \tau)=\sum_{k=1}^{n-1} w_{j, k}(t / \tau)\left(x_{k} / \tau\right)+\mathbb{O}\left(\tau^{-2}\right) \tag{4-25}
\end{equation*}
$$

which follows from the definition of $w$ and the fact that $w_{\alpha}$ is symmetric in these special indices $\alpha$. In particular, $w_{j, k}(t)=w_{k, j}(t)$ for all $j, k \in[1, n-1]$.

Recall that we chose the integer $\kappa$ such that $m<\kappa$. By Proposition 4.3 there is a constant $C$ such that

$$
\begin{equation*}
\left|I_{\tau}\right| \leq C \tau^{-\kappa} \tag{4-26}
\end{equation*}
$$

and we shall now show that if $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0))$ is the first nonvanishing Taylor coefficient with respect to the ordering $>_{t}$, where $m=j_{0}+k_{0}+\left|\alpha_{0}\right|+\left|\beta_{0}\right|$, then (4-26) cannot hold. (Since we are denoting the variables by $(t, x)$ now, the index $\alpha$ in Definition 3.2 will be replaced by the pair $(k, \alpha) \in \mathbb{N} \times \mathbb{N}^{n-1}$.) We will do this by determining the limit of $\tau^{m} I_{\tau}$ as $\tau \rightarrow \infty$. To see what is needed, consider $\lambda_{-1}(t / \tau, x / \tau+y(t / \tau))$ and recall that this is

$$
q_{1}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right) \phi_{0}(t / \tau, x / \tau+y(t / \tau))
$$

which should be regarded as a Taylor expansion in $\xi$ of $q_{1}$ at $\eta(t / \tau)$ of finite order. The same applies to all the other terms of the form $q_{-j}^{(\alpha)}$. For given $j$ and $\alpha$, we only ever need to consider Taylor expansions of $q_{-j}^{(\alpha)}$ of order $\kappa-j-|\alpha|$ in view of (4-22) and (4-25). To keep things simple, we shall first only consider $q_{1}$; it will be clear by symmetry what the corresponding expressions for the other terms should be. Thus,

$$
\begin{align*}
& q_{1}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right) \\
& =\sum_{|\beta| \leq \kappa+1} q_{1}^{(\beta)}(t / \tau, x / \tau+y(t / \tau), \eta(t / \tau))\left(w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))-\eta(t / \tau)\right)^{\beta} /|\beta|!+\mathbb{O}\left(\tau^{-\kappa-2}\right), \tag{4-27}
\end{align*}
$$

which shows that to use our assumption regarding the Taylor coefficient $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0))$, we have to for each $\beta$ write $q_{1}^{(\beta)}(t / \tau, x / \tau+y(t / \tau), \eta(t / \tau))$ as a Taylor series at $\eta(0)$, in addition to having to expand each term as a Taylor series in $t$ and $x$. However, it is immediate from (4-25) that if $\beta$ is an ( $n-1$ )-tuple corresponding to a given differential operator $D_{\xi}^{\beta}$, then there is a sequence $\tilde{\beta}=\left(\tilde{\beta}_{1}, \ldots, \tilde{\beta}_{s}\right)$ of $s=|\beta|$ indices between 1 and the dimension $n-1$ of the $x$ variable, such that

$$
\begin{equation*}
g_{\tau}^{\beta}(t, x)=\left(w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))-\eta(t / \tau)\right)^{\beta} \tag{4-28}
\end{equation*}
$$

as it appears in (4-27), satisfies

$$
g_{\tau}^{\beta}(t, x)=c_{\beta}(t / \tau, x / \tau)+\mathbb{O}\left(\tau^{-|\beta|-1}\right)
$$

where

$$
c_{\beta}(t / \tau, x / \tau)=\prod_{j=1}^{s}\left(\sum_{k=1}^{n-1} w_{k, \tilde{\beta}_{j}}(t / \tau) x_{k} / \tau\right) \quad \text { and } \quad c_{\beta}(0, x / \tau)=\tau^{-|\beta|} c_{\beta}(0, x)
$$

These expressions make sense if we choose the sequence $\tilde{\beta}$ to be increasing, for then it is uniquely determined by $\beta$. If for instance $D_{\xi}^{\beta}=-\partial^{2} / \partial \xi_{i} \partial \xi_{j}$, then $\tilde{\beta}=(i, j)$ if $i \leq j$ (see the indices $\alpha$ used in connection with $w_{\alpha}$ in (4-4)). Thus (4-27) takes the form

$$
\begin{aligned}
q_{1}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau\right. & +y(t / \tau))) \\
& =\sum_{|\beta| \leq \kappa+1} q_{1}^{(\beta)}(t / \tau, x / \tau+y(t / \tau), \eta(t / \tau)) g_{\tau}^{\beta}(t, x) /|\beta|!+\mathbb{O}\left(\tau^{-\kappa-2}\right),
\end{aligned}
$$

and if we expand each term in this expression as a Taylor series at $\eta(0)$ we obtain

$$
\begin{align*}
& q_{1}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right) \\
& \quad=\sum_{|\beta| \leq \kappa+1} \sum_{|\gamma| \leq \kappa+1-|\beta|} q_{1}^{(\beta+\gamma)}(t / \tau, x / \tau+y(t / \tau), \eta(0)) g_{\tau}^{\beta}(t, x)(\eta(t / \tau)-\eta(0))^{\gamma} /(|\beta|!|\gamma|!)  \tag{4-29}\\
& +\mathbb{O}\left(\tau^{-\kappa-2}\right),
\end{align*}
$$

where we regard $\eta(t / \tau)-\eta(0)$ as a finite Taylor series $\eta^{\prime}(0) t / \tau+\eta^{\prime \prime}(0) t^{2} /\left(2 \tau^{2}\right)+\cdots$ of sufficiently high order to maintain control of the error term in (4-29). If we for each multiindex $\beta$ let $G_{\tau}^{\beta}(t, x)$ be given by

$$
G_{\tau}^{\beta}(t, x)=\sum_{\gamma_{1}+\gamma_{2}=\beta}(\eta(t / \tau)-\eta(0))^{\gamma_{1}} g_{\tau}^{\gamma_{2}}(t, x) /\left(\left|\gamma_{1}\right|!\left|\gamma_{2}\right|!\right) \quad \text { for } \gamma_{j} \in \mathbb{N}^{n-1}
$$

then the required order of the Taylor expansion $\eta(t / \tau)-\eta(0)$ will ultimately depend on $\beta$, so we can write

$$
\begin{align*}
q_{1}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau\right. & +y(t / \tau))) \\
& =\sum_{|\beta| \leq \kappa+1} q_{1}^{(\beta)}(t / \tau, x / \tau+y(t / \tau), \eta(0)) G_{\tau}^{\beta}(t, x)+\mathbb{O}\left(\tau^{-\kappa-2}\right) \tag{4-30}
\end{align*}
$$

and we can always bound $G_{\tau}^{\beta}(t, x)$ by a constant times $\tau^{-|\beta|}$. As will be evident in a moment, the value of $G_{\tau}^{\beta}(t, x)$ for $|\beta|>0$ is not important. For notational purposes, denote by $G_{0}^{\beta}(t, x)$ the limit of $\tau^{|\beta|} G_{\tau}^{\beta}(t, x)$ as $\tau \rightarrow \infty$. Since $G_{\tau}^{\beta}(t, x)=1$ when $\beta=0$ it is clear that $G_{0}^{0}(t, x)=1$.

For each $\beta$ we must now write $q_{1}^{(\beta)}(t / \tau, x / \tau+y(t / \tau), \eta(0))$ as a Taylor expansion in $t$ and $x$ at 0 and $y(0)$, respectively. As before, for given $j$ and $\alpha$, we will only have to consider Taylor expansions of $q_{-j}^{(\alpha)}$ of order $\kappa-j-|\alpha|$. By (4-23) and (4-30) we have

$$
\begin{align*}
& \lambda_{-1}(t / \tau, x / \tau+y(t / \tau)) \\
& \quad=\sum_{k+|\alpha|+|\beta| \leq \kappa+1} \phi_{0}(t / \tau, x / \tau+y(t / \tau))\left((t / \tau)^{k}(x / \tau+y(t / \tau)-y(0))^{\alpha} G_{\tau}^{\beta}(t, x)\right.  \tag{4-31}\\
& \\
& \left.\times \partial_{t}^{k} q_{1(\alpha)}^{(\beta)}(0, y(0), \eta(0)) /(k!|\alpha|!)+\mathbb{O}\left(\tau^{-\kappa-2}\right)\right)
\end{align*}
$$

where we in $(x / \tau+y(t / \tau)-y(0))^{\alpha}$ regard $y(t / \tau)-y(0)$ as a finite Taylor series of sufficiently high order to maintain control of the error terms.

In the way that we expressed the term $q_{1}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right)$ by (4-31), we can get similar expressions of appropriate order for the terms $q_{-j}^{(\gamma)}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right)$ that appear in (4-23). For each $j$ and $\gamma$ we have

$$
\begin{align*}
& q_{-j}^{(\gamma)}\left(t / \tau, x / \tau+y(t / \tau), w_{x}^{\prime}(t / \tau, x / \tau+y(t / \tau))\right) \\
& \quad=\sum_{k+|\alpha|+|\beta| \leq \kappa-j-|\gamma|}(t / \tau)^{k}(x / \tau+y(t / \tau)-y(0))^{\alpha} G_{\tau}^{\beta}(t, x) \partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}(0, y(0), \eta(0)) /(k!|\alpha|!)  \tag{4-32}\\
& \quad+\mathbb{O}\left(\tau^{-\kappa-1+j+|\gamma|}\right)
\end{align*}
$$

This together with (4-23) gives

$$
\begin{array}{r}
\lambda_{J}(t / \tau, x / \tau+y(t / \tau))=\sum_{j+l+|\gamma|=J} \sum_{k+|\alpha|+|\beta| \leq \kappa-j-|\gamma|}(t / \tau)^{k}(x / \tau+y(t / \tau)-y(0))^{\alpha} G_{\tau}^{\beta}(t, x) \\
\times D_{x}^{\gamma} \phi_{l}(t / \tau, x / \tau+y(t / \tau)) \frac{\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}(0, y(0), \eta(0))}{k!|\alpha|!} \\
+\mathbb{O}\left(\tau^{-\kappa-1+j+|\gamma|)}\right)
\end{array}
$$

where $-1 \leq j \leq J$ and $l \geq 0$. Using that by assumption the Taylor coefficients $\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}(0, y(0), \eta(0))$ vanish for all $-1 \leq j+k+|\alpha|+|\beta|+|\gamma|<m$, and $\tau^{-J-k-|\alpha|}=\tau^{|\beta|} \tau^{-j-k-|\alpha|-|\beta|-|\gamma|-l}$ when $J=j+l+|\gamma|$, the equation above yields

$$
\begin{array}{r}
\sum_{J=-1}^{m} \tau^{-J} \lambda_{J}(t / \tau, x / \tau+y(t / \tau))=\sum_{j+l+|\gamma|=-1}^{m} \sum_{j+k+|\alpha|+|\beta|+|\gamma|=m} \tau^{-m-l} t^{k}\left(x+y^{\prime}(0) t\right)^{\alpha} \tau^{|\beta|} G_{\tau}^{\beta}(t, x) \\
\times D_{x}^{\gamma} \phi_{l}(t / \tau, x / \tau+y(t / \tau)) \frac{\partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}(0, y(0), \eta(0))}{k!|\alpha|!} \\
\end{array}
$$

where $\tau^{|\beta|} G_{\tau}^{\beta}(t, x) \rightarrow G_{0}^{\beta}(t, x)$ as $\tau \rightarrow \infty$. As we can see, the expression above is $\mathcal{O}\left(\tau^{-m-1}\right)$ as soon as $l>0$, so in view of (4-22) and (4-24) we obtain

$$
\begin{align*}
\lim _{\tau \rightarrow \infty} \tau^{m} I_{\tau}=\iint & H(t, x) e^{i t w_{0}^{\prime}(0)+i\langle x, \eta(0)\rangle} \\
& \times\left(\sum_{j+k+|\alpha|+|\beta|+|\gamma|=m} t^{k}\left(x+y^{\prime}(0) t\right)^{\alpha} G_{0}^{\beta}(t, x) D_{x}^{\gamma} \phi_{0}(0, y(0))\right.  \tag{4-33}\\
& \left.\times \partial_{t}^{k} q_{-j(\alpha)}^{(\beta+\gamma)}(0, y(0), \eta(0)) /(k!|\alpha|!)\right) d t d x
\end{align*}
$$

Recall (4-20) and choose $\phi_{0}$ so that $D_{x}^{\beta_{0}} \phi_{0}(0, y(0))=1$, but so that $D_{x}^{\gamma} \phi_{0}(0, y(0))=0$ for all other $\gamma$ such that $|\gamma| \leq\left|\beta_{0}\right|$ (see (4-14)). By the choice of our ordering $>_{t}$ we have $\partial_{t}^{k} q_{-j(\alpha)}^{\left(\beta+\beta_{0}\right)}(0, y(0), \eta(0))=0$ for all $\beta$ such that $|\beta|>0$ as long as $j+k+|\alpha|+|\beta|+\left|\beta_{0}\right|=m$. Hence, with this choice of $\phi_{0}$, equation (4-33) takes the form

$$
\begin{align*}
& \lim _{\tau \rightarrow \infty} \tau^{m} I_{\tau}=\iint H(t, x) e^{i t w_{0}^{\prime}(0)+i\langle x, \eta(0)\rangle} \\
& \times\left(\sum_{j+k+|\alpha|+\left|\beta_{0}\right|=m} t^{k}\left(x+y^{\prime}(0) t\right)^{\alpha} \partial_{t}^{k} q_{-j(\alpha)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0)) /(k!|\alpha|!)\right) d t d x \tag{4-34}
\end{align*}
$$

so as promised, the value of $G_{0}^{\beta}(t, x)$ for $|\beta|>0$ does not matter. (Note that $G_{0}^{0}(t, x)$ is present in (4-34) as the constant factor 1.) As in the proof of Theorem 2.20, some of the Taylor coefficients in (4-34) may be zero. In particular, the expression may well contain Taylor coefficients that preceed $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0))$ in the ordering, and those are by assumption zero. In contrast to the proof of Theorem 2.20 we shall have to exploit this fact, since the coefficients of most of the monomials in (4-34) will be linear combinations of the Taylor coefficients due to the factor $\left(x+y^{\prime}(0) t\right)^{\alpha}$. However, the ordering $>_{t}$ was chosen so that there can be no nonzero Taylor coefficient $\partial_{t}^{k} q_{-j(\alpha)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0))$ such that $k+|\alpha|>k_{0}+\left|\alpha_{0}\right|$, or $k+|\alpha|=k_{0}+\left|\alpha_{0}\right|$ and $k<k_{0}$. This follows immediately from the
choice of lexicographic order on the $n$-tuple $(k, \alpha) \in \mathbb{N}^{n}$. (Recall that in the definition of the ordering $>_{t}$, $x$ denoted all the variables in $\mathbb{R}^{n}$, while here we denote those variables by $(t, x)$.) Hence, the only coefficient of the monomial $t^{k_{0}} x^{\alpha_{0}}$ in (4-34) is $\partial_{t}^{k_{0}} q_{-j_{0}\left(\alpha_{0}\right)}^{\left(\beta_{0}\right)}(0, y(0), \eta(0))$. We may therefore, as in the proof of Theorem 2.20, choose $H$ so that the limit in (4-34) is nonzero. Since this contradicts (4-26), Theorem 2.21 follows in view of the discussion following Lemma 4.1.

## Appendix A.

Here we prove a few results used in the main text, related to how the property that all terms in the asymptotic expansion of the total symbol have vanishing Taylor coefficients is affected by various operations.

Lemma A.1. Suppose $X$ and $Y$ are two $C^{\infty}$ manifolds of the same dimension n. Let $K \subset T^{*}(X) \backslash 0$ and $K^{\prime} \subset T^{*}(Y) \backslash 0$ be compactly based cones and let $\chi$ be a homogeneous symplectomorphism from a conic neighborhood of $K^{\prime}$ to one of $K$ such that $\chi\left(K^{\prime}\right)=K$. Let $A \in I^{m^{\prime}}\left(X \times Y, \Gamma^{\prime}\right)$ and $B \in$ $I^{m^{\prime \prime}}\left(Y \times X,\left(\Gamma^{-1}\right)^{\prime}\right)$, where $\Gamma$ is the graph of $\chi$, and assume that $A$ and $B$ are properly supported and noncharacteristic at the restriction of the graphs of $\chi$ and $\chi^{-1}$ to $K^{\prime}$ and to $K$ respectively, while $W F^{\prime}(A)$ and $W F^{\prime}(B)$ are contained in small conic neighborhoods. If $R$ is a properly supported classical pseudodifferential operator in $Y$, then each term in the asymptotic expansion of the total (left) symbol of $R$ has vanishing Taylor coefficients at a point $(y, \eta) \in K^{\prime}$ if and only if each term in the asymptotic expansion of the total (left) symbol of the pseudodifferential operator ARB in $X$ has vanishing Taylor coefficients at $\chi(y, \eta) \in K$.

Proof. We may assume that we have a homogeneous generating function $\varphi \in C^{\infty}$ for the symplectomorphism $\chi$; see [Grigis and Sjöstrand 1994, pages 101-103]. Then $\chi$ is locally of the form

$$
(\partial \varphi(x, \eta) / \partial \eta, \eta) \mapsto(x, \partial \varphi(x, \eta) / \partial x)
$$

and $A$ and $B$ are given by

$$
\begin{aligned}
& A u(x)=\frac{1}{(2 \pi)^{n}} \iint e^{i(\varphi(x, \zeta)-z \cdot \zeta)} a(x, z, \zeta) u(z) d z d \zeta \\
& B v(y)=\frac{1}{(2 \pi)^{n}} \iint e^{i(y \cdot \theta-\varphi(s, \theta))} b(y, s, \theta) v(s) d s d \theta
\end{aligned}
$$

Since $R$ is properly supported we may assume that

$$
\begin{equation*}
R u(z)=\frac{1}{(2 \pi)^{n}} \int e^{i z \cdot \eta} r(z, \eta) \hat{u}(\eta) d \eta \quad \text { for } u \in C_{0}^{\infty}(Y) \tag{A-1}
\end{equation*}
$$

where $r(z, \eta)=\sigma_{R}$ is the total symbol of $R$. Hence
$A R B u(x)=\frac{1}{(2 \pi)^{3 n}} \int e^{i(\varphi(x, \zeta)-z \cdot \zeta+(z-y) \cdot \sigma+y \cdot \theta-\varphi(s, \theta))} a(x, z, \zeta) r(z, \sigma) b(y, s, \theta) u(s) d s d \theta d y d \sigma d z d \zeta$,
since $B$ being properly supported implies that $B u \in C_{0}^{\infty}(Y)$ when $u \in C_{0}^{\infty}(Y)$. Using integration by parts in $z$, we see that we can insert a cutoff $\phi((\zeta-\sigma) /|\sigma|)$ in the last integral without changing the operator
$A R B \bmod \Psi^{-\infty}$. If we make the change of variables $\tau=\zeta-\sigma$, then (A-2) takes the form

$$
\begin{aligned}
A R B u(x)= & \frac{1}{(2 \pi)^{3 n}} \int \phi(\tau /|\sigma|) e^{i(\varphi(x, \tau+\sigma)-z \cdot(\tau+\sigma)+(z-y) \cdot \sigma+y \cdot \theta-\varphi(s, \theta))} \\
& \times a(x, z, \tau+\sigma) r(z, \sigma) b(y, s, \theta) u(s) d s d \theta d y d \sigma d z d \tau+L u
\end{aligned}
$$

with $L \in \Psi^{-\infty}$. If $\Omega \subset \mathbb{R}^{2 n}$ is open and $\tilde{\varphi} \in C^{\infty}(\Omega, \mathbb{R})$ is a phase function with a nondegenerate critical point $x_{0} \in \Omega$ such that $d \tilde{\varphi} \neq 0$ everywhere else, then [Grigis and Sjöstrand 1994, Proposition 2.3] states, in particular, that for every compact $M \subset \Omega$ and every $u \in C^{\infty}(\Omega) \cap \mathscr{E}^{\prime}(M)$ we have

$$
\begin{equation*}
\left|\int e^{i \lambda \tilde{\varphi}(x)} u(x) d x-e^{i \lambda \tilde{\varphi}\left(x_{0}\right)} A_{0} u\left(x_{0}\right) \lambda^{-n}\right| \leq C_{M} \lambda^{-n-1} \sum_{|\alpha| \leq 2 n+3} \sup \left|\partial^{\alpha} u(x)\right| \quad \text { for } \lambda \geq 1 \tag{A-3}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{0}=\frac{(2 \pi)^{n} \cdot e^{i \pi \operatorname{sgn} \tilde{\varphi}^{\prime \prime}\left(x_{0}\right) / 4}}{\left|\operatorname{det} \tilde{\varphi}^{\prime \prime}\left(x_{0}\right)\right|^{1 / 2}} \tag{A-4}
\end{equation*}
$$

It is clear that the result extends to the setting $\Omega=T^{*}(\mathcal{N}) \backslash 0$, where $\mathcal{N}$ is a $C^{\infty}$ manifold of dimension $n$. In order to apply the result, we put $\sigma=\lambda \omega$, and make the change of variables $\tau=\lambda \tilde{\tau}$. After dropping the tilde we obtain

$$
\begin{aligned}
A R B u(x)= & \frac{\lambda^{2 n}}{(2 \pi)^{3 n}} \int \phi(\tau /|\omega|) e^{i \lambda(\varphi(x, \tau+\omega)-z \cdot(\tau+\omega)+y \cdot \theta / \lambda+(z-y) \cdot \omega-\varphi(s, \theta) / \lambda)} \\
& \times a(x, z, \lambda(\tau+\omega)) r(z, \lambda \omega) b(y, s, \theta) u(s) d s d \theta d y d \omega d z d \tau+L u
\end{aligned}
$$

where we have used the fact that $\varphi$ is homogeneous of degree 1 in the fiber. For the $z, \tau$-integration we have the nondegenerate critical point given by $\tau=0, z=\varphi_{\zeta}^{\prime}(x, \tau+\omega)$. Note that since $\varphi_{\zeta}^{\prime}$ is homogeneous of degree 0 in the fiber we have $\varphi_{\zeta}^{\prime}(x, \sigma / \lambda)=\varphi_{\zeta}^{\prime}(x, \sigma)$, so this critical point corresponds to the critical point for the $z, \zeta$-integration given by $\zeta=\sigma, z=\varphi_{\zeta}^{\prime}(x, \sigma)$. Hence the expression above together with (A-3) imply that

$$
A R B u(x)=C \lambda^{2 n} \int e^{i(\varphi(x, \lambda \omega)+y \cdot \theta-y \cdot \lambda \omega-\varphi(s, \theta))} w(x, y, s, \omega, \theta) u(s) d s d \theta d y d \omega+L u
$$

where

$$
\begin{aligned}
w(x, y, s, \omega, \theta) & =\left.\frac{A_{0}}{\lambda^{n}} a(x, z, \lambda(\tau+\omega)) r(z, \lambda \omega) b(y, s, \theta) \phi(\tau /|\omega|)\right|_{\substack{\tau=0 \\
z=\varphi_{\zeta}^{\prime}(x, \omega)}} \\
& =\frac{A_{0}}{\lambda^{n}} a\left(x, \varphi_{\zeta}^{\prime}(x, \omega), \lambda \omega\right) r\left(\varphi_{\zeta}^{\prime}(x, \omega), \lambda \omega\right) b(y, s, \theta)
\end{aligned}
$$

with an error of order $\mathbb{O}\left(\lambda^{-n-1}\right)$. Note that $A_{0}$ is now a function of $x$ and $\omega$, since the matrix corresponding to $\tilde{\varphi}^{\prime \prime}\left(x_{0}\right)$ in (A-4) is given by the block matrix

$$
F=\left(\begin{array}{cc}
0 & -\mathrm{Id}_{n}  \tag{A-5}\\
-\mathrm{Id}_{n} & \varphi_{\zeta \zeta}^{\prime \prime}(x, \omega)
\end{array}\right),
$$

where $\mathrm{Id}_{n}$ is the identity matrix on $\mathbb{R}^{n}$. Clearly the determinant of $F$ is either 1 or -1 , so $F$ is nonsingular. Furthermore, $F$ depends smoothly on the parameters $x$ and $\omega$ since $\varphi \in C^{\infty}$, so the eigenvalues of $F$
are continuous in $x$ and $\omega$. Hence it follows that the signature of $F$ is constant, for if not there has to exist an eigenvalue vanishing at some point $(x, \omega)$, contradicting the nonsingularity of $F$. Reverting to the variable $\sigma=\lambda \omega$ we thus obtain

$$
A R B u(x)=C \int e^{i(\varphi(x, \sigma)+y \cdot(\theta-\sigma)-\varphi(s, \theta))} \tilde{w}(x, y, s, \sigma, \theta) u(s) d s d \theta d y d \sigma+L u
$$

where

$$
\tilde{w}(x, y, s, \sigma, \theta)=a\left(x, \varphi_{\zeta}^{\prime}(x, \sigma), \sigma\right) r\left(\varphi_{\zeta}^{\prime}(x, \sigma), \sigma\right) b(y, s, \theta)
$$

with an error of order $\mathcal{O}\left(\lambda^{-1}\right)$. Taking the limit as $\lambda \rightarrow \infty$ yields
$A R B u(x)=C \int e^{i(\varphi(x, \sigma)+y \cdot(\theta-\sigma)-\varphi(s, \theta))} a\left(x, \varphi_{\zeta}^{\prime}(x, \sigma), \sigma\right) r\left(\varphi_{\zeta}^{\prime}(x, \sigma), \sigma\right) b(y, s, \theta) u(s) d s d \theta d y d \sigma+L u$.
We can now repeat the procedure. Indeed, we can insert a cutoff $\phi((\sigma-\theta) /|\theta|)$ without changing the operator mod $\Psi^{-\infty}$, and after making the corresponding changes of variables in order to apply [Grigis and Sjöstrand 1994, Proposition 2.3] we find that for the $y, \sigma$-integration we have the nondegenerate critical point given in the original variables by $\sigma=\theta, y=\varphi_{\sigma}^{\prime}(x, \sigma)$. After taking the limit as $\lambda \rightarrow \infty$ we obtain

$$
A R B u(x)=C \int e^{i(\varphi(x, \theta)-\varphi(s, \theta))} w_{1}(x, s, \theta) u(s) d s d \theta+L_{1} u
$$

where $L_{1} \in \Psi^{-\infty}$ and

$$
\begin{equation*}
w_{1}(x, s, \theta)=a\left(x, \varphi_{\theta}^{\prime}(x, \theta), \theta\right) r\left(\varphi_{\theta}^{\prime}(x, \theta), \theta\right) b\left(\varphi_{\theta}^{\prime}(x, \theta), s, \theta\right) \tag{A-6}
\end{equation*}
$$

As before we let the factor $A_{0}$ from (A-4) be included in the constant $C$. In a conic neighborhood of $\operatorname{supp} w_{1}$ we can write

$$
\varphi(x, \theta)-\varphi(s, \theta)=(x-s) \Xi(x, s, \theta)
$$

Then $\Xi(x, x, \theta)=\varphi_{x}^{\prime}(x, \theta)$ so $\partial \Xi(x, x, \theta) / \partial \theta=\varphi_{x \theta}^{\prime \prime}(x, \theta)$ is invertible, since $\varphi_{x \theta}^{\prime \prime}(x, \theta) \neq 0$ is equivalent to the fact that the graph of $\chi$ is (locally) the graph of a smooth map. Hence $\theta \mapsto \Xi(x, s, \theta)$ is $C^{\infty}$, homogeneous of degree 1 and with an inverse having the same properties. For $s$ close to $x$, the equation $\Xi(x, s, \theta)=\xi$ then defines $\theta=\Theta(x, s, \xi)$. After a change of variables, the last integral therefore takes the form

$$
\begin{equation*}
A R B u(x)=C \int e^{i(x-s) \cdot \xi} \tilde{w}_{1}(x, s, \xi) u(s) d s d \xi+L_{1} u \tag{A-7}
\end{equation*}
$$

where $\tilde{w}_{1}(x, s, \xi)$ is just $w_{1}(x, s, \Theta(x, s, \xi))$ multiplied by a Jacobian. We note in passing that evaluating $\tilde{w}_{1}$ at a point $(x, x, \xi)$ where $\xi$ is of the form $\xi=\varphi_{x}^{\prime}(x, \eta)$ therefore involves evaluating $w_{1}$ at the point $(x, x, \eta)$. The integral (A-7) defines a pseudodifferential operator with total symbol $\rho(x, \xi)$ satisfying

$$
\begin{equation*}
\left.\rho(x, \xi) \sim \sum \frac{i^{-|\alpha|}}{\alpha!}\left(\partial_{\xi}^{\alpha} \partial_{y}^{\alpha} \tilde{w}_{1}(x, y, \xi)\right)\right|_{y=x} \tag{A-8}
\end{equation*}
$$

If the total symbol $r=\sigma_{R}$ of $R$ has vanishing Taylor coefficients at a point $(y, \eta)=\left(\varphi_{\eta}^{\prime}(x, \eta), \eta\right)$, then by examining (A-8) in decreasing order of homogeneity we find that each term of $\rho$ must have vanishing

Taylor coefficients at $(x, \xi)=\left(x, \varphi_{x}^{\prime}(x, \eta)\right)$, since by what we have shown this would involve evaluating $r(z, \sigma)$ and its derivatives at $\left(\varphi_{\eta}^{\prime}(x, \eta), \eta\right)$.

To prove the converse, choose $A_{1} \in I^{-m^{\prime \prime}}\left(X \times Y, \Gamma^{\prime}\right)$ and $B_{1} \in I^{-m^{\prime}}\left(Y \times X,\left(\Gamma^{-1}\right)^{\prime}\right)$ properly supported such that

$$
\begin{array}{ll}
K^{\prime} \cap W F\left(B A_{1}-I\right)=\varnothing, & K \cap W F\left(A_{1} B-I\right)=\varnothing, \\
K^{\prime} \cap W F\left(B_{1} A-I\right)=\varnothing, & K \cap W F\left(A B_{1}-I\right)=\varnothing .
\end{array}
$$

Then a repetition of the arguments above shows that all the terms in the asymptotic expansion of the total symbol of $B_{1} A R B A_{1}$ has vanishing Taylor coefficients at a point $(y, \eta)=\left(\varphi_{\eta}^{\prime}(x, \eta), \eta\right)$ if all the terms in the asymptotic expansion of the total symbol of $A R B$ has vanishing Taylor coefficients at $(x, \xi)=\left(x, \varphi_{x}^{\prime}(x, \eta)\right)$. Since $R$ and $B_{1} A R B A_{1}$ have the same total symbol in $K^{\prime} \bmod \Psi^{-\infty}$, the same must hold for the total symbol of $R$.

Let $\left\{e_{k}: k=1, \ldots, n\right\}$ be a basis for $\mathbb{R}^{n}$, let $(U, x)$ be local coordinates on a smooth manifold $M$ of dimension $n$, and let

$$
\left\{\frac{\partial}{\partial x_{k}}: k=1, \ldots, n\right\}
$$

be the induced local frame for the tangent bundle $T M$. Since the local frame fields commute, we can use standard multiindex notation to express the partial derivatives $\partial_{x}^{\alpha} f$ of $f \in C^{\infty}(U)$.

Lemma A.2. Let $M$ be a smooth manifold of dimension $n$, and for $j \geq 1$ let $p, q_{j}, g_{j} \in C^{\infty}(M)$. Let $\left\{\gamma_{j}\right\}_{j=1}^{\infty}$ be a sequence in $M$ such that $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$, and assume that $p(\gamma)=p\left(\gamma_{j}\right)=0$ for all $j$, and that $d p(\gamma) \neq 0$. Let $(U, x)$ be local coordinates on $M$ near $\gamma$, and suppose that there exists a smooth function $q \in C^{\infty}(M)$ such that

$$
\partial_{x}^{\alpha} q(\gamma)=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} q_{j}\left(\gamma_{j}\right) \quad \text { for all } \alpha \in \mathbb{N}^{n}
$$

If $q_{j}-p g_{j}$ vanishes of infinite order at $\gamma_{j}$ for all $j$, then there exists a smooth function $g \in C^{\infty}(M)$ such that $q-p g$ vanishes of infinite order at $\gamma$. Furthermore,

$$
\begin{equation*}
\partial_{x}^{\alpha} g(\gamma)=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} g_{j}\left(\gamma_{j}\right) \quad \text { for all } \alpha \in \mathbb{N}^{n} \tag{A-9}
\end{equation*}
$$

Proof. We have stated the result for a manifold, but since the result is purely local we may assume that $M \subset \mathbb{R}^{n}$ in the proof. It is also clear that we may assume that there exists an open neighborhood $थ$ of $\gamma$ such that $\gamma_{j} \in U$ for $j \geq 1$, and that $d p \neq 0$ in $U$. By shrinking $U$ if necessary, we can then find a unit vector $v \in \mathbb{R}^{n}$ such that $\partial_{\nu} p(w)=\langle v, d p(w)\rangle \neq 0$ for $w \in U$. (We will identify a tangent vector $v \in \mathbb{R}^{n}$ at $\gamma$ with $\partial_{\nu} \in T_{\gamma} \mathbb{R}^{n}$ through the usual vector space isomorphism.) Hence $\partial_{\nu} p(w)$ is invertible in $\because$, and we let $\left(\partial_{\nu} p(w)\right)^{-1} \in C^{\infty}(\vartheta)$ denote its inverse. By an orthonormal change of coordinates we may even assume that $\partial_{\nu} p(w)=\partial_{e_{1}} p(w)$. In accordance with the notation used in the statement of the lemma, we shall write $\partial_{x_{k}} p(w)$ for the partial derivatives $\partial_{e_{k}} p(w)$ and denote by $\left(\partial_{x_{1}} p(w)\right)^{-1}$ the inverse of $\partial_{\nu} p(w)=\partial_{x_{1}} p(w)$ in $थ$.

Now

$$
\begin{equation*}
0=\partial_{x_{1}}\left(q_{j}-p g_{j}\right)\left(\gamma_{j}\right)=\partial_{x_{1}} q_{j}\left(\gamma_{j}\right)-\partial_{x_{1}} p\left(\gamma_{j}\right) g_{j}\left(\gamma_{j}\right) \tag{A-10}
\end{equation*}
$$

for all $j$ since $p\left(\gamma_{j}\right)=0$. Since $\lim _{j} \partial_{x_{1}} q_{j}\left(\gamma_{j}\right)=\partial_{x_{1}} q(\gamma)$ by assumption, equation (A-10) yields

$$
\begin{equation*}
\lim _{j \rightarrow \infty} g_{j}\left(\gamma_{j}\right)=\left(\partial_{x_{1}} p(\gamma)\right)^{-1} \partial_{x_{1}} q(\gamma)=a \in \mathbb{C} \tag{A-11}
\end{equation*}
$$

We claim that we can in the same way determine

$$
\lim _{j \rightarrow \infty}\left(\partial_{x}^{\alpha} g_{j}\right)\left(\gamma_{j}\right)=a_{(\alpha)} \in \mathbb{C} \quad \text { for any } \alpha \in \mathbb{N}^{n}
$$

We start by determining $\lim _{j \rightarrow \infty} \partial g_{j}\left(\gamma_{j}\right) / \partial x_{k}=a_{(k)}$ for $1 \leq k \leq n$. By the hypotheses of the lemma we have

$$
\begin{align*}
0 & =\partial_{x_{k}} \partial_{x_{l}}\left(q_{j}-p g_{j}\right)\left(\gamma_{j}\right) \\
& =\partial_{x_{k}} \partial_{x_{l}} q_{j}\left(\gamma_{j}\right)-\partial_{x_{k}} \partial_{x_{l}} p\left(\gamma_{j}\right) g_{j}\left(\gamma_{j}\right)-\partial_{x_{k}} p\left(\gamma_{j}\right) \partial_{x_{l}} g_{j}\left(\gamma_{j}\right)-\partial_{x_{l}} p\left(\gamma_{j}\right) \partial_{x_{k}} g_{j}\left(\gamma_{j}\right) \tag{A-12}
\end{align*}
$$

since $p\left(\gamma_{j}\right)=0$. For $k=l=1$ we obtain from (A-11) and (A-12)

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \partial_{x_{1}} g_{j}\left(\gamma_{j}\right)=\left(\partial_{x_{1}} p(\gamma)\right)^{-1}\left(\partial_{x_{1}}^{2} q(\gamma)-\partial_{x_{1}}^{2} p(\gamma) a\right) / 2 \tag{A-13}
\end{equation*}
$$

This allows us to solve for $\partial_{x_{k}} g_{j}\left(\gamma_{j}\right)$ in (A-12) by choosing $l=1$. If $b \in \mathbb{C}$ denotes the limit in (A-13) and $a \in \mathbb{C}$ is given by (A-11) we thus obtain

$$
\lim _{j \rightarrow \infty} \partial_{x_{k}} g_{j}\left(\gamma_{j}\right)=\left(\partial_{x_{1}} p(\gamma)\right)^{-1}\left(\partial_{x_{1}} \partial_{x_{k}} q(\gamma)-\partial_{x_{1}} \partial_{x_{k}} p(\gamma) a-\partial_{x_{k}} p(\gamma) b\right) \quad \text { for } 2 \leq k \leq n
$$

Now assume that for some $m \geq 3$ we have in this way determined

$$
\lim _{j \rightarrow \infty} \partial_{x_{k_{1}}} \ldots \partial_{x_{k_{m-2}}} g_{j}\left(\gamma_{j}\right), \quad \text { for } k_{i} \in[1, n], \text { with } i \in[1, m-2] .
$$

To shorten notation, we will use the (nonstandard) multiindex notation introduced on page 465; to every $\alpha \in \mathbb{N}^{n}$ with $|\alpha|=m$ corresponds precisely one $m$-tuple $\beta=\left(k_{1}, \ldots, k_{m}\right)$ of nondecreasing numbers $1 \leq k_{1} \leq \cdots \leq k_{m} \leq n$ such that $\partial_{x}^{\beta}$ equals $\partial_{x}^{\alpha}$. Throughout the rest of this proof we shall let $\beta$ represent such an $m$-tuple, and we let

$$
\hat{\beta}_{i}=\left(k_{1}, \ldots, k_{i-1}, k_{i+1}, \ldots, k_{m}\right)
$$

As before we have

$$
\begin{equation*}
0=\partial_{x}^{\beta}\left(q_{j}-p g_{j}\right)\left(\gamma_{j}\right)=\partial_{x}^{\beta} q_{j}\left(\gamma_{j}\right)-\partial_{x}^{\beta} p\left(\gamma_{j}\right) g_{j}\left(\gamma_{j}\right)-\cdots-\sum_{i=1}^{m} \partial_{x_{k_{i}}} p\left(\gamma_{j}\right) \partial_{x}^{\hat{\beta}_{i}} g_{j}\left(\gamma_{j}\right) \tag{A-14}
\end{equation*}
$$

by assumption. If we choose $k_{i}=1$ for all $1 \leq i \leq m$, the last sum is just $m \partial_{x_{1}} p\left(\gamma_{j}\right) \partial_{x_{1}}^{m-1} g_{j}\left(\gamma_{j}\right)$, and since the limit of all other terms on the right side are known by the induction hypothesis, we thus obtain the value of the limit of $\partial_{x_{1}}^{m-1} g_{j}\left(\gamma_{j}\right)$ from (A-14) by first multiplying by $m^{-1}\left(\partial_{x_{1}} p\left(\gamma_{j}\right)\right)^{-1}$ and then letting $j \rightarrow \infty$. Denote this limit by $c \in \mathbb{C}$. If we choose $k_{i} \neq 1$ for precisely one $i \in[1, m]$, say $k_{m}=k$, then the last sum in (A-14) satisfies

$$
\sum_{i=1}^{m} \partial_{x_{k_{i}}} p\left(\gamma_{j}\right) \partial_{x}^{\hat{\beta}_{i}} g_{j}\left(\gamma_{j}\right)=\partial_{x_{k}} p\left(\gamma_{j}\right) \partial_{x_{1}}^{m-1} g_{j}\left(\gamma_{j}\right)+(m-1) \partial_{x_{1}} p\left(\gamma_{j}\right) \partial_{x_{1}}^{m-2} \partial_{x_{k}} g_{j}\left(\gamma_{j}\right)
$$

so by the same argument as before we can obtain the value of $\lim _{j \rightarrow \infty} \partial_{x_{1}}^{m-2} \partial_{x_{k}} g_{j}\left(\gamma_{j}\right)$ for $2 \leq k \leq n$ by multiplying by $(m-1)^{-1}\left(\partial_{x_{1}} p\left(\gamma_{j}\right)\right)^{-1}$ and using $\partial_{x_{1}}^{m-1} g_{j}\left(\gamma_{j}\right) \rightarrow c$ when taking the limit as $j \rightarrow \infty$ in (A-14). Continuing this way it is clear that we can successively determine

$$
\lim _{j \rightarrow \infty} \partial_{x_{k_{1}}} \ldots \partial_{x_{k_{m-1}}} g_{j}\left(\gamma_{j}\right) \quad \text { for any } 1 \leq k_{1} \leq \ldots \leq k_{m-1} \leq n
$$

which completely determines $\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} g_{j}\left(\gamma_{j}\right)=a_{(\alpha)}$ for $\alpha \in \mathbb{N}^{n}$ with $|\alpha|=m-1$, proving the claim.
By Borel's theorem there exists a smooth function $g \in C^{\infty}(M)$ such that

$$
\partial_{x}^{\alpha} g(\gamma)=a_{(\alpha)}=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} g_{j}\left(\gamma_{j}\right) \quad \text { for all } \alpha \in \mathbb{N}^{n}
$$

Since $q-p g$ vanishes of infinite order at $\gamma$ by construction, this completes the proof.
The lemma will be used to prove the following result for homogeneous smooth functions on the cotangent bundle.

Proposition A.3. For $j \geq 1$ let $p, q_{j}, g_{j} \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$, where $p$ and $q_{j}$ are homogeneous of degree $m$ and the $g_{j}$ are homogeneous of degree 0 . Let $\left\{\gamma_{j}\right\}_{j=1}^{\infty}$ be a sequence in $T^{*}\left(\mathbb{R}^{n}\right) \backslash 0$ such that $\gamma_{j} \rightarrow \gamma$ as $j \rightarrow \infty$, and assume that $p(\gamma)=p\left(\gamma_{j}\right)=0$ for all $j$, and that $d p(\gamma) \neq 0$. If there exists a smooth function $q \in C^{\infty}\left(T^{*}(\mathbb{R})^{n} \backslash 0\right)$, homogeneous of degree $m$, such that

$$
\partial_{x}^{\alpha} \partial_{\xi}^{\beta} q(\gamma)=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} q_{j}\left(\gamma_{j}\right) \quad \text { for all }(\alpha, \beta) \in \mathbb{N}^{n} \times \mathbb{N}^{n}
$$

and if $q_{j}-p g_{j}$ vanishes of infinite order at $\gamma_{j}$ for all $j$, then there exists a $g \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$, homogeneous of degree 0 , such that $q-p g$ vanishes of infinite order at $\gamma$. Furthermore,

$$
\begin{equation*}
\partial_{x}^{\alpha} \partial_{\xi}^{\beta} g(\gamma)=\lim _{j \rightarrow \infty} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} g_{j}\left(\gamma_{j}\right) \quad \text { for all }(\alpha, \beta) \in \mathbb{N}^{n} \times \mathbb{N}^{n} \tag{A-15}
\end{equation*}
$$

Proof. Let $\pi: T^{*}\left(\mathbb{R}^{n}\right) \backslash 0 \rightarrow S^{*}\left(\mathbb{R}^{n}\right)$ be the projection. Since $d p(\gamma) \neq 0$ it follows from homogeneity that $d p(\pi(\gamma)) \neq 0$. By using the homogeneity of $q, q_{j}$ and $g_{j}$ we may even assume that $\gamma$ and $\gamma_{j}$ belong to $S^{*}\left(\mathbb{R}^{n}\right)$ for $j \geq 1$ to begin with.

Now, the radial vector field $\xi \partial_{\xi}$ applied $k$ times to $a \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right) \backslash 0\right)$ equals $l^{k} a$ if $a$ is homogeneous of degree $l$. For any point $w \in S^{*}\left(\mathbb{R}^{n}\right)$ with $w=\left(w_{x}, w_{\xi}\right)$ in local coordinates on $T^{*}\left(\mathbb{R}^{n}\right)$ it is easy to see that

$$
T_{w} S^{*}\left(\mathbb{R}^{n}\right)=\left\{(u, v) \in \mathbb{R}^{n} \times \mathbb{R}^{n}:\left\langle w_{\xi}, v\right\rangle=0\right\}
$$

Therefore a basis for $T_{w} S^{*}\left(\mathbb{R}^{n}\right)$ together with the radial vector field $\left(\xi \partial_{\xi}\right)_{w}$ at $w$ constitutes a basis for $T_{w} T^{*}\left(\mathbb{R}^{n}\right)$. This implies that if we can find a homogeneous function $g$ such that $q-p g$ vanishes of infinite order in the directions $T_{\gamma} S^{*}\left(\mathbb{R}^{n}\right)$, then $q-p g$ vanishes of infinite order at $\gamma$, for the derivatives involving the radial direction are determined by lower order derivatives in the directions $T_{\gamma} S^{*}\left(\mathbb{R}^{n}\right)$.

By the hypotheses of the proposition together with an application of Lemma A.2, we find that there exists a function $\tilde{g} \in C^{\infty}\left(T^{*}\left(\mathbb{R}^{n}\right)\right)$, not necessarily homogeneous, such that $q-p \tilde{g}$ vanishes of infinite order at $\gamma$ and (A-15) holds for $\tilde{g}$. The function $g(x, \xi)=\tilde{g}(x, \xi / \| \xi \mid)$ coincides with $\tilde{g}$ on $S^{*}\left(\mathbb{R}^{n}\right)$. In particular, all derivatives of $g$ and $\tilde{g}$ in the directions $T_{\gamma} S^{*}\left(\mathbb{R}^{n}\right)$ are equal at $\gamma$. Thus, by the arguments
above we conclude that $q-p g$ vanishes of infinite order at $\gamma$. Since $g$ and $g_{j}$ are homogeneous of degree 0 , the same arguments also imply that (A-15) holds for $g$, which completes the proof.

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## ANALYSIS \& PDE

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[^1]:    ${ }^{1}$ If the factor $x_{2}$ in $\operatorname{Im} p_{1}$ is raised to the power 3 for example, then it turns out that $\left\{\gamma\left(x_{1}\right): 0<x_{1}<2\right\}$ is a one-dimensional bicharacteristic interval of $p_{1}$, and not only a bicharacteristic of the real part. It is obviously not minimal though, nor does it contain any minimal characteristic points.

