# PROPAGATION OF WAVE FRONT SETS OF SOLUTIONS OF THE CAUCHY PROBLEM FOR HYPERBOLIC EQUATIONS IN GEVREY CLASSES 

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Introduction. In the recent paper [18] the second author has constructed the fundamental solution of the Cauchy problem for hyperbolic equations in Gevrey classes, and investigated the propagation of wave front sets of their solutions in Gevrey classes by assuming the constant multiplicities of their characteristic roots. The purpose of the present paper is to study the propagation of wave front sets in Gevrey classes for solutions of hyperbolic equations with characteristic roots of variable multiplicities and to give a similar result to the one for the $C^{\infty}$-case obtained by Kumano-go and the second author [10]. Main results of the present paper are announced in [15] and [19].

Let $\mathcal{L}$ be an $l \times l$ hyperbolic system of the form

$$
\mathcal{L}=D_{t}-\left[\begin{array}{ccc}
\lambda_{1}\left(t, X, D_{x}\right) & 0  \tag{1}\\
& \ddots & \\
0 & & \lambda_{l}\left(t, X, D_{x}\right)
\end{array}\right]+\left(b_{j k}\left(t, X, D_{x}\right)\right)
$$

$$
\text { on } \quad[0, T] \times R_{x}^{n}
$$

with real symbols $\lambda_{j}(t, x, \xi)$ in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$ and symbols $b_{j k}(t, x, \xi)$ in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{\sigma}\right)(0 \leqq \sigma<1 / \kappa)$. Here, for $\kappa>1$ and a real $m$ we denote by $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{m}\right)$ a class of symbols $p(t, x, \xi)$ of pseudo-differential operators satisfying for any multi-indices $\alpha, \beta$ and non-negative integer $\gamma$

$$
\begin{align*}
&\left|\partial_{t}^{\gamma} \partial_{\xi}^{\alpha} \partial_{x}^{\beta} p(t, x, \xi)\right| \leqq C M^{-\left(\left|\alpha_{\mid}+\left|\beta_{1}\right| \gamma\right)\right.}(\alpha!\beta!\gamma!)^{\kappa}\langle\xi\rangle^{m-\left|\alpha_{1}\right|}  \tag{2}\\
& \text { for } \quad(t, x, \xi) \in[0, T] \times R_{x}^{n} \times R_{\xi}^{n},
\end{align*}
$$

with constants $C$ and $M(>0)$ independent of $\alpha, \beta$ and $\gamma$. Throughout the present paper we assume the symbols $\lambda_{j}$ are positively homogeneous in $\boldsymbol{\xi}$ (for $|\xi| \geqq 1$ ), that is, $\lambda_{j}$ satisfy

$$
\lambda_{j}(t, x, \theta \xi)=\theta \lambda_{j}(t, x, \xi) \quad \text { for } \quad \theta \geqq 1 \quad \text { and } \quad|\xi| \geqq 1 .
$$

[^0]Let $\mathscr{D}_{L^{2}}^{(\kappa)^{\prime}}$ denote a class of ultradistributions defined by [6], that is,

$$
\mathscr{D}_{L^{2}}^{(\kappa))^{\prime}}=\underset{\varepsilon_{\downarrow 0}}{\operatorname{proj}} \lim \mathscr{D}_{L^{2}, \varepsilon^{(\kappa)}}^{()^{\prime}} .
$$

Here, for $\varepsilon>0, \mathscr{D}_{L^{\prime}, \mathrm{s}}^{(\kappa)^{\prime}}$ is a dual space of the Hilbert space

$$
\mathscr{D}_{L^{2}, \mathrm{e}}^{(\kappa)}=\left\{u(x) \in L^{2} ; \exp \left(\varepsilon\langle\xi\rangle^{1 / \kappa}\right) \hat{u}(\xi) \in L^{2}\right\}
$$

and $\hat{u}(\xi)$ is the Fourier transform of $u(x)$ (see [20]). If $u \in \mathscr{D}_{L^{2}}^{(\kappa))^{\prime}}$ and $\kappa_{1} \geqq \kappa$ we denote by $\mathrm{WF}_{G\left(\kappa_{1}\right)}(u)$ the wave front set of $u$ in the Gevrey class of order $\kappa_{1}$ defined as follows:

Definition. Let $\left(x^{0}, \xi^{0}\right)$ be a point in $T^{*}\left(R^{n}\right) \backslash 0$ and let $u \in \mathscr{D}_{\left.L^{2}\right)^{\prime}}^{(\text {. }}$. The point $\left(x,{ }^{0} \xi^{0}\right)$ does not belong to $\mathrm{WF}_{G\left(\kappa_{1}\right)}(u)$ for $\kappa_{1} \geqq \kappa$ if there exists a symbol $a(x, \xi)$ in $S_{G(k)}^{0}$ (see Definition 1.1-ii) in Section 1) with $a\left(x^{0}, \theta \xi^{0}\right) \neq 0(\theta \geqq 1)$ such that $f(x) \equiv a\left(X, D_{x}\right) u$ belongs to $\gamma^{\left(\kappa_{1}\right)}\left(R_{x}^{n}\right)$, that is, it satisfies

$$
\left|\partial_{x}^{\infty} f(x)\right| \leqq C M^{-|\infty|} \alpha!^{\kappa_{1}} \quad \text { for all } \quad x \in R_{x}^{n}
$$

This definition is equivalent to that of Hörmander [2] if $u \in \mathcal{E}^{\prime}$ (see Theorem 3 of [20]).

Consider the Cauchy problem

$$
\begin{equation*}
\mathcal{L} U(t)=0 \quad(t \in[0, T]), \quad U(0)=G \in \mathscr{G}_{L^{\prime}}^{(k)}, \tag{7}
\end{equation*}
$$

for $\mathcal{L}$ of the form (1) with $\lambda_{j} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$ and $b_{j k} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{\sigma}\right)$ for $0 \leqq \sigma<1 / \kappa$. Then we obtain the following:

Theorem 1 (cf. Theorem 3.4 of [10]). Assume $\lambda_{j}(t, x, \xi)$ are homogeneous for $|\xi| \geqq 1$. Then, for any initial data $G \in \mathscr{D}_{L^{2}}^{(x)}{ }^{\prime}$ there exists a unique solution $U(t)$ of $(7)$ in $\mathscr{B}^{\infty}\left([0, T] ; \mathscr{D}_{L^{2}}^{(\kappa)^{\prime}}\right)$ and it satisfies

$$
\begin{equation*}
\mathrm{WF}_{G\left(\kappa_{1}\right)}(U(t)) \subset \Gamma\left(t ; \mathrm{WF}_{G\left(\kappa_{1}\right)}(G)\right) \quad(0<t \leqq T) \tag{8}
\end{equation*}
$$

for any $\kappa_{1}$ satisfying $\kappa \leqq \kappa_{1}<1 / \sigma$.
The theorem of this type in the $C^{\infty}$-case was given in [10] and [11]. In (8) the set $\Gamma\left(t_{0} ; V\right)$, for a fixed $t_{0} \in(0, T]$ and a conic set $V$ in $T^{*}\left(R_{x}^{n}\right) \backslash 0$, is defined as follows: First, we define $\Gamma_{\varepsilon}^{\nu}\left(t_{0} ; V\right)(\varepsilon>0, \nu=0,1,2, \cdots)$ as the conic hull of the set of end points (at $t=t_{0}$ ) of all $\varepsilon$-admissible trajectories of, at most, step $\nu$ issuing from $(y, \eta) \in V$ for large $|\eta|$. Then, the set $\Gamma\left(t_{0} ; V\right)$ is defined by

$$
\begin{equation*}
\Gamma\left(t_{0} ; V\right)=\bigcap_{\varepsilon>0} \bigcup_{\nu=0}^{\infty} \Gamma_{\varepsilon}^{\nu}\left(t_{0} ; V_{\varepsilon}\right), \tag{9}
\end{equation*}
$$

where $V_{\mathrm{g}}$ denotes an $\varepsilon$-conic neighborhood of $V$ defined by

$$
V_{\mathrm{z}}=\{(x, \xi) ;|x-y| \leqq \varepsilon,|\xi /|\xi|-\eta /|\eta|| \leqq \varepsilon,(y, \eta) \in V\} .
$$

Roughly speaking, the $\varepsilon$-admissible trajectory is an $\varepsilon$-appoximation of the so-called broken null-bicharacteristic flow (see its precise definition in §4). The estimate (8) seems to be loose apparently because the limiting curve of $\varepsilon$-admissible trajectories ( $\varepsilon \downarrow 0$ ) is not always broken null-bicharacteristic flow (see last sections in [4] and [5]). However, a result about the optimality of (8) was shown by the first author [14].

Next, we consider an application of Theorem 1 to the Cauchy problem

$$
\left\{\begin{array}{l}
L u=0 \quad(0 \leqq t \leqq T),  \tag{10}\\
\partial_{t}^{j} u(0)=g_{j} \in \mathscr{D}_{L^{2}}^{\left(k^{\prime}\right)}
\end{array} \quad(j=0, \cdots, m-1)\right.
$$

for a single hyperbolic operator

$$
\begin{equation*}
L=D_{t}^{m}+\sum_{j=0}^{m-1} \sum_{|\alpha| \leq m-j} a_{j, a}(t, x) D_{x}^{\infty} D_{t}^{j} \quad \text { on } \quad[0, T] \times R_{x}^{n} \tag{11}
\end{equation*}
$$

with coefficients $a_{j, \alpha}(t, x)$ in a Gevrey class $\gamma^{(k)}\left([0, T] \times R_{x}^{n}\right)$, that is, they satisfy

$$
\left|\partial_{t}^{\gamma} \partial_{x}^{\beta} a_{j, \alpha}(t, x)\right| \leqq C M^{-(\gamma+|\beta|)}(\beta!\gamma!)^{\kappa} \quad \text { on } \quad[0, T] \times R_{x}^{n} .
$$

As shown in [18] the problem (10) can be reduced to the equivalent Cauchy problem (7) with $\sigma=(r-q) / r$ and is $\gamma^{(\kappa)}$-well-posed for $1 \leqq \kappa<1 / \sigma$ (cf. [12]) if there exist regularly hyperbolic differential operators $L_{1}, L_{2}, \cdots, L_{r}$ with coefficients in $\gamma^{(\kappa)}\left([0, T] \times R_{x}^{n}\right)$ such that $L$ has a form

$$
\begin{equation*}
L=L_{1} L_{2} \cdots L_{r}+\sum_{j=0}^{m-q} a_{j}\left(t, X, D_{x}\right) D_{t}^{j} \tag{12}
\end{equation*}
$$

with $a_{j}(t, x, \xi)$ in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{m-q-j}\right)$ and $1 \leqq q \leqq r$. From this reduction it follows that for any $t \in(0, T]$

$$
\begin{equation*}
\left.\mathrm{WF}_{G\left(\kappa_{1}\right)}(u(t)) \subset \Gamma\left(t ; \bigcup_{j=0}^{m-1} \mathrm{WF}_{G\left(\kappa_{1}\right)}\right)\left(g_{j}\right)\right) \quad \text { for } \quad \kappa \leqq \kappa_{1}<r /(r-q) . \tag{13}
\end{equation*}
$$

In the present paper, we shall consider the reduction to an equivalent problem (7) from the problem (10) for a hyperbolic operator of more general form than (12), which is inspired by the work of Komatsu [7] (see (5.1) and Theorem 2 in §5). In the case that the maximal multiplicity of characteristic roots of $L$ is at most three, we also clarify the conditions of lower order terms of (11) in order that the problem (10) is reduced to an equivalent problem (7) of a hyperbolic system (1) with a given $\sigma(<1)$ (see Theorem 3 in §5). We remark that the hyperbolic operator $L$ has always the form (12) with $q=1$ if characteristic
roots of $L$ belong to $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$ and we admit $L_{j}$ to be pseudo-differential cperators with respect to $x$.

For the hyperbolic operator $L$ without any conditions on lower order terms and without assuming the smoothness of characteristic roots, Wakabayashi [21] has recently investigated the propagation of wave front sets for solutions of (10) in Gevrey class of order $\kappa_{1}$ satisfying $\kappa \leqq \kappa_{1}<r /(r-1)$, where $r$ is the maximal multiplicity of characteristic roots. The method of [21] is based on the construction of a parametrix of $L$, as in Bronshtein [1], and on the notion of "flows" $K_{z}^{+}$in $T^{*}\left(R_{t}^{1} \times R_{x}^{n}\right)$ emanating from a point $z \in T^{*}\left(R_{t}^{1} \times R_{x}^{n}\right)$. When characteristic roots of hyperbolic operator $L$ are smooth, that is, they belong to $G^{(\kappa)}\left([0, T] ; S_{G(x)}^{1}\right)$, for any $t_{0} \in(0, T]$ and any closed set $V$ in $T^{*}\left(R_{x}^{n}\right)$, it follows that

$$
\begin{equation*}
\Gamma\left(t_{0}, V\right)=\left\{\pi\left(K_{z}^{+} \cap\left\{t=t_{0}\right\}\right) ; z \in \pi^{-1}(V) \cap\{t=0\} \cap p^{-1}(0)\right\} \tag{14}
\end{equation*}
$$

Here $p=p(t, x, \tau, \xi)$ is the principal symbol of $L$ and $\pi$ is the natural projection from $T^{*}\left(R_{t}^{1} \times R_{x}^{n}\right)$ to $T^{*}\left(R_{x}^{n}\right)$ (see Theorem 4 in $\S 6$, cf. Theorem 4.4 in [22]). So, our result (8) is the same as the one in [21] in the case that characteristic roots of $L$ are smooth.

The plan of the present paper is as follows: In $\S \S 1-4$ we prove Theorem 1. $\S \S 1-3$ are devoted to preparatory lemmas and in $\S 4$ we complete the proof of Theorem 1 with the precise definition of $\Gamma\left(t_{0} ; V\right)$. In $\S 5$, we show a method of the reduction of the form (1), and give, as an application of Theorem 1, a result on the propagation of wave front sets for the Cauchy problem (10). In $\S 6$ we show the equivalence of the estimate given by the flows $K_{z}^{+}$ of Wakabayashi [21], [22] and the one given by the set $\Gamma\left(t_{0} ; V\right)$ of Kumano-go, Taniguchi and Tozaki [10], [11].

More precisely, we shall state the main idea of §1-3. In $\S 1$ we separate the symbol of the multi-product of Fourier integral operators to the sum of a main symbol and a regularizer and give the precise estimate for the part of the regularizer. To obtain this estimate we represent each factor of the multiproduct to the sum of symbols depending on a parameter $\zeta$. Then, we can use the similar discussions as in [17]. In $\S 2$ we estimate the part of the main symbol of the multi-product of Fourier integral operators which is given by the oscillatory integral of the multiple symbol. In [18], to estimate this we transform the multi-product of Fourier integral operators to the multi-product of pseudo-differential operators multiplied by a Fourier integral operator, using the decomposition $I_{\phi^{*}} R I_{\phi}$ of the identity operator. In the present paper, since we estimate main symbols represented by oscillatory integrals of multiple symbols we use the transformation of oscillatory integrals which corresponds to the one by means of the decomposition $I=I_{\phi^{*}} R I_{\phi}$. In $\S 3$, we give a method of the integration by parts for the symbol represented by an iterated integral
of Volterra type. To show the corresponding estimate for the $C^{\infty}$-case, in [10] we have estimated the iterated integral after we have simplified the multiproduct of Fourier integral operators to one Fourier integral operator with multi-phase. But in the Gevrey case we can not employ this method since we use the equations of the critical points $X_{v}^{j}=\nabla_{\xi} \phi_{j}\left(t_{j-1}, t_{j} ; X_{v}^{j-1}, \Xi_{v}^{j}\right)$, $\Xi_{v}^{j}=$ $\nabla_{x} \phi_{j+1}\left(t_{j}, t_{j+1} ; X_{\nu}^{j}, \Xi_{\nu}^{j+1}\right),(j=1, \cdots, \nu)$ to obtain the simplified symbol of a ( $\nu+1)$-multi-product and we have no uniform estimate in the Gevrey class for the solutions of the equations of the critical points. Here the uniform estimate means the estimate independent of $\nu$. So, we use the integration by parts for the iterated integral of Volterra type before simplifying the multiple symbol. It should be noted that to perform this method we must treat the oscillatory integral of the multiple symbol instead of the multi-product of Fourier integral operators and so we estimate in $\S 2$ the simplified symbol derived from the multiple symbol.

## 1. Fourier integral operators in Gevrey classes

First we recall symbol classes introduced in [18] and [20], which are subclasses of a symbol class $S^{m}$ studied in [9]. In what follows we tacitly use the notation in [9] and [20] and assume that the constant $\kappa$ is always larger than 1.

Definition 1.1. i) We say that a symbol $p(x, \xi)\left(\in S^{m}\right)$ belongs to a class $S_{G(\kappa, \infty)}^{m}$ if for any multi-index $\alpha$ there exists a constant $C_{\alpha}$ such that

$$
\begin{equation*}
\left|p_{(\beta)}^{(\alpha)}(x, \xi)\right| \leqq C_{\alpha} M^{-|\beta|} \beta!^{\aleph}\langle\xi\rangle^{m-\left|\alpha_{1}\right|} \tag{1.1}
\end{equation*}
$$

holds for a constant $M$ independent of $\alpha$ and $\beta$.
ii) We say that a symbol $p(x, \xi)\left(\in S^{m}\right)$ belongs to a class $S_{G(k)}^{m}$ if

$$
\begin{equation*}
\left|p_{(\beta)}^{(\alpha)}(x, \xi)\right| \leqq C M^{-(|\alpha|+|\beta|)}(\alpha!\beta!)^{\kappa}\langle\xi\rangle^{m-|\alpha|} \tag{1.2}
\end{equation*}
$$

holds for constants $C$ and $M$ independent of $\alpha$ and $\beta$.
iii) We say that a symbol $p(x, \xi)\left(\in S^{-\infty}\right)$ belongs to a class $\mathcal{R}_{G(x)}$ if for any $\alpha$ there exists a constant $C_{\alpha}$ such that

$$
\begin{equation*}
\left|p_{(\beta)}^{(\alpha)}(x, \xi)\right| \leqq C_{\alpha} M^{-|\beta|} \beta!^{\kappa} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa}\right) \tag{1.3}
\end{equation*}
$$

holds for positive constants $M$ and $\varepsilon$ independent of $\alpha$ and $\beta$.
Remark. The inequality (1.3) is equivalent to the condition that

$$
\left|p_{\beta}^{(\alpha)}(x, \xi)\right| \leqq C_{\alpha} M^{-(|\beta|+N)}(\beta!N!)^{\kappa}\langle\xi\rangle^{-|\alpha|-N}
$$

holds for any integer $N$ with a constant $M$ independent of $N, \alpha$ and $\beta$.
Following Definition 1.1 of [18], for a $\tau \in[0,1)$ we define a class $\mathscr{P}_{G(x)}(\tau)$
of phase functions of Fourier integral operators as follows: We say that a real valued function $\phi(x, \xi)$ is a phase function belonging to a class $\mathscr{P}_{G(x)}(\tau)$ if $J(x, \xi) \equiv \phi(x, \xi)-x \cdot \xi$ satisfieds

$$
\left\{\begin{array}{l}
\sum_{|\alpha|+|\beta| \leq 2}\left|J_{(\beta)}^{(\alpha)}(x, \xi)\right| \mid\langle\xi\rangle^{1-|\alpha|} \leqq \tau,  \tag{1.4}\\
\left|J_{(\beta)}^{(\alpha)}(x, \xi)\right| \leqq \tau M^{-\left(\left|\alpha_{1}+\left|\beta_{1}\right|\right.\right.}(\alpha!\beta!)^{\kappa}\langle\xi\rangle^{1-|\alpha|}
\end{array}\right.
$$

for a constant $M$ independent of $\alpha$ and $\beta$. We put $\mathscr{P}_{G(\kappa)}=\bigcup_{0 \leq \tau<1} \mathscr{P}_{G(k)}(\tau)$.
Let $\phi(x, \xi)$ be a phase function in $\mathscr{P}_{G(x)}$. Then, a Fourier integral operator $P_{\phi}=p_{\phi}\left(X, D_{x}\right)$ with phase function $\phi(x, \xi)$ and a symbol $\sigma\left(P_{\phi}\right)=p(x, \xi)$ in $S_{G(k, \infty)}^{m}$ is defined by

$$
\begin{equation*}
P_{\phi} u(x)=O_{s}-\iint e^{i\left(\phi(x, \xi)-x^{\prime} \cdot \xi\right)} p(x, \xi) u\left(x^{\prime}\right) d x^{\prime} d \xi \quad \text { for } \quad u \in \mathcal{S}, \tag{1.5}
\end{equation*}
$$

where $d \xi=(2 \pi)^{-n} d \xi, \mathcal{S}$ is the Schwartz space of rapidly decreasing functions on $R_{x}^{n}$ and the right hand side of (1.5) is the oscillatory integral defined in Chapter 1 of [9]. We denote the set of such Fourier integral operators by $S_{G(\kappa, \infty), \phi}^{m}$. If $\phi=x \cdot \xi$, the set $S_{G(k, \infty), \phi}^{m}$ is the one of pseudo-differential operators. In this case we write it simply by $S_{G(\kappa, \infty)}^{m}$. Since $S_{G(\kappa)}^{m}$ and $\mathcal{R}_{G(\kappa)}$ are subclasses of $S_{G(\kappa, \infty)}^{m}$, we similarly denote by $S_{G(k), \phi}^{m}$ and $\mathcal{R}_{G(\kappa), \phi}$ the corresponding classes of Fourier integral operators, which are represented by the formula (1.5). We can identify $\mathcal{R}_{G(\kappa), \phi}$ with the set $\mathcal{R}_{G(\kappa)}$ of pseudo-differential operators because it follows from (1.4) that if $p(x, \xi)$ belongs to $\mathcal{R}_{G(\kappa)}$ then $e^{i J(x, \xi)} p(x, \xi)$ so does. Here and in what follows, we use the same notation $\mathcal{R}_{G(\kappa)}$ for the class of pseudodifferential operators with symbols in $\mathcal{R}_{G(\kappa)}$ because no confusion occurs between the class of symbols and the one of operators.

As proved in [20] we have
Proposition 1.2 (see Theorem 2 and Lemma 2.1 of [20]).
i) Let $p(x, \xi)$ belong to $S_{G(x, \infty)}^{m}$ and let $\phi(x, \xi)$ belong to $\mathscr{P}_{G(k)}$. Then the Fourier integral operator $P_{\phi}=p_{\phi}\left(X, D_{x}\right)$ maps $\mathscr{D}_{L^{2}}^{(k)}$ into itself.
ii) Let $R$ be a pseudo-differential operator in $\mathcal{R}_{G(k)}$ and let $u$ belong to $\mathscr{D}_{L^{2}}^{(k)^{\prime}}$. Then we have $R u \in \gamma^{(k)}\left(R_{x}^{n}\right)$.

Following $\S 2$ of [18] we denote for $\phi \in \mathscr{P}_{G(k)}$

$$
\begin{gathered}
L_{G(x)}^{m}(\phi)=\left\{p_{\phi}\left(X, D_{x}\right)=p_{\phi}^{o}\left(X, D_{x}\right)+\tilde{p}_{\phi}\left(X, D_{x}\right) ;\right. \\
\left.p^{o}(x, \xi) \in S_{G(x)}^{m}, \tilde{p}(x, \xi) \in \mathscr{R}_{G(x)}\right\},
\end{gathered}
$$

that is, symbolically $L_{G(k)}^{m}(\phi)=S_{G(k), \phi}^{m}+\mathcal{R}_{G(\kappa), \phi}^{m}$. In what follows we often say that $p^{0}(x, \xi)$ is a main symbol of $p_{\phi}\left(X, D_{x}\right)$. If $\phi(x, \xi)=x \cdot \xi$ we denote $L_{G(\kappa)}^{m}(\phi)$ simply by $L_{G(k)}^{m}$.

For a sequence $\left\{\phi_{j}\right\}$ of phase functions $\phi_{j}(x, \xi) \in \mathscr{P}_{G(k)}\left(\tau_{j}\right)$, we consider multi-products

$$
\begin{equation*}
P_{1, \phi_{1}} P_{2, \phi_{2}} \cdots P_{\nu+1, \phi_{\nu+1}} \tag{1.6}
\end{equation*}
$$

of Fourier integral operators $P_{j, \phi_{j}}$ in $L_{G(\kappa)}^{\sigma}\left(\phi_{j}\right)$ with $\sigma \geqq 0$. As in $\S 2$ of [18] we assume the following:
(A-1) If we set $J_{j}(x, \xi)=\phi_{j}(x, \xi)-x \cdot \xi,\left\{J_{j} / \tau_{j}\right\}$ is bounded in $S_{G(k)}^{1}$ and an inequality $\sum_{j=1}^{\infty} \tau_{j} \leqq \tau^{0}$ holds for a small constant $\tau^{0}$.
(A-2) If we write $P_{j, \phi_{j}} \equiv p_{j, \phi_{j}}\left(X, D_{x}\right)=p_{j, \phi_{j}}^{o}\left(X, D_{x}\right)+\tilde{p}_{j, \phi_{j}}\left(X, D_{x}\right) \in S_{G(\kappa), \phi_{j}}^{\sigma}+$ $\mathcal{R}_{G(k), \phi_{j}}$ the set $\left\{p_{j}^{o}(x, \xi)\right\}$ is bounded in $S_{G(k)}^{\sigma}$ and the set $\left\{\tilde{p}_{j}(x, \xi)\right\}$ is bounded in $\mathcal{R}_{G(\mathrm{k})}$.

Remark. Concerning the bounded set in $S_{G(k)}^{m}$ or $\mathcal{R}_{G(k)}$, see remarks after Definition 1.1 in [18].

We assume $\tau^{\circ}$ in (A-1) small enough so that Proposition 2.4 in [18] and Lemmas 1.4-1.6 below hold. Then, a multi-product $\Phi_{\nu+1} \equiv \Phi_{\nu+1}(x, \xi)=\phi_{1} \#$ $\phi_{2} \# \cdots \# \phi_{\nu+1}(x, \xi)$ of phase functions $\phi_{1}(x, \xi), \phi_{2}(x, \xi), \cdots, \phi_{\nu+1}(x, \xi)$ is defined by

$$
\Phi_{\nu+1}(x, \xi)=\sum_{j=1}^{\nu}\left(\phi_{j}\left(X_{\nu}^{j-1}, \Xi_{\nu}^{j}\right)-X_{\nu}^{j} \cdot \Xi_{\nu}^{j}\right)+\phi_{\nu+1}\left(X_{\nu}^{\nu}, \xi\right) \quad\left(X_{\nu}^{0}=x\right)
$$

and it belongs to $\mathscr{P}_{G(k)}$, where $\left\{X_{v}^{j}, \Xi_{v}^{j}\right\}_{j=1}^{\nu} \equiv\left\{X_{v}^{j}, \Xi_{v}^{j}\right\}_{j=1}^{\nu}(x, \xi)$ is a solution of

$$
\left\{\begin{align*}
x^{j}= & \nabla_{\xi} \phi_{j}\left(x^{j-1}, \xi^{j}\right)  \tag{1.7}\\
\xi^{j}= & \nabla_{x} \phi_{j+1}\left(x^{j}, \xi^{j+1}\right) \\
& \left(j=1, \cdots, \nu ; \quad x^{0}=x, \quad \xi^{\nu+1}=\xi\right)
\end{align*}\right.
$$

Recall that the multi-product (1.6) is a Fourier integral operator in $L_{G(\kappa)}^{(+1) \sigma}\left(\Phi_{\nu+1}\right)$ with the above phase function $\Phi_{\nu+1}(x, \xi)$ and its symbol $q_{\nu+1}(x, \xi)=\sigma\left(P_{1, \phi_{1}} P_{2, \phi_{2}}\right.$ $\cdots P_{\nu+1, \phi_{\nu+1}}$ ) is written as

$$
\begin{align*}
& q_{\nu+1}(x, \xi) \\
& =O_{s}-\int \cdots \int \exp \left[i\left(\psi_{\nu+1}\left(x, \xi^{\nu}, \tilde{x}^{\nu}, \xi\right)-\Phi_{\nu+1}(x, \xi)\right)\right]  \tag{1.8}\\
& \quad \times \prod_{j=1}^{\nu+1} p_{j}\left(x^{j-1}, \xi^{j}\right) d \tilde{x}^{\nu} d \tilde{\xi}^{\nu} \quad\left(x^{0}=x, \xi^{\nu+1}=\xi\right)
\end{align*}
$$

where $\tilde{x}^{\nu}=\left(x^{1}, x^{2}, \cdots, x^{\nu}\right), \hat{\xi}^{\nu}=\left(\xi^{1}, \xi^{2}, \cdots, \xi^{\nu}\right), d \tilde{x}^{\nu}=d x^{1} \cdots d x^{\nu}, d \tilde{\xi}^{\nu}=d \xi^{1} \cdots d \xi^{\nu}$ and

$$
\begin{gather*}
\left.\Psi_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)=\sum_{j=1}^{\nu}\left(\phi_{j}\left(x^{j-1}, \xi^{j}\right)-x^{j} \cdot \xi^{j}\right)\right)+\phi_{\nu+1}\left(x^{\nu}, \xi\right)  \tag{1.9}\\
\left(x^{0}=x\right)
\end{gather*}
$$

In the above the right hand side of (1.8) is an oscillatory integral, whose well-
definedness will be proved in the next section in a more general form. The aim of this section is to find the main symbol of $q_{\nu+1}(x, \xi)$ for (1.8).

Let $\chi(\xi)$ be a function in $\gamma^{(\kappa)}\left(R_{\xi}^{n}\right)$ satisfying

$$
\begin{equation*}
0 \leqq \chi(\xi) \leqq 1, \quad \chi(\xi)=1(|\xi| \leqq 2 / 5), \quad \chi(\xi)=0 \quad(|\xi| \geqq 1 / 2) . \tag{1.10}
\end{equation*}
$$

For a main symbol $p_{j}^{o}(x, \xi) \in S_{G(k)}^{o}$ of $P_{j, \phi_{j}}$ and a parameter $\zeta \in R^{n}$, we set

$$
\left\{\begin{array}{l}
p_{j}^{+}(x, \xi ; \zeta)=\chi(4(\xi-\zeta) \mid\langle\zeta\rangle) p_{j}^{o}(x, \xi),  \tag{1.11}\\
p_{j}^{-}(x, \xi ; \zeta)=p_{j}^{o}(x, \xi)-p_{j}^{\dagger}(x, \xi ; \zeta)
\end{array}\right.
$$

and consider $P_{j, \phi_{j}}^{ \pm}(\zeta)=p_{j, \phi_{j}}^{ \pm}\left(X, D_{x} ; \zeta\right)$ Fourier integral operators with a parameter $\zeta$. Set for $k=1, \cdots, \nu+1$
and set

$$
q_{\nu+1}^{o}(x, \xi ; \zeta)=\sigma\left(P_{1, \phi_{1}}^{+}(\zeta) \cdots P_{\nu+1, \phi_{\nu+1}}^{+}(\zeta)\right) .
$$

Then, for any fixed $\zeta$ we get the decomposition

$$
\begin{equation*}
q_{\nu+1}(x, \xi)=q_{\nu+1}^{o}(x, \xi ; \zeta)+\sum_{k=1}^{\nu+1} q_{\nu+1}^{k}(x, \xi ; \zeta)+\sum_{k=1}^{\nu+1} \tilde{q}_{\nu+1}^{k}(x, \xi) . \tag{1.13}
\end{equation*}
$$

In (1.13) we set $\zeta=\xi$, where $\xi$ is the fiber variable of the simplified symbol of (1.6). Then, we have

Lemma 1.3. The symbol $r_{\nu+1}(x, \xi) \equiv q_{\nu+1}(x, \xi)-q_{\nu+1}^{o}(x, \xi ; \xi)$ belongs to $\mathcal{R}_{G(\mathrm{k})}$ and it satisfies for any $\alpha$ and $\beta$

$$
\begin{equation*}
\left|r_{\nu+1(\beta)}^{(\alpha)}(x, \xi)\right| \leqq C_{\alpha} A^{\nu} M^{-|\beta|} \beta!^{\kappa} \nu!^{\sigma_{\kappa}} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa}\right) \tag{1.14}
\end{equation*}
$$

with positive constants $\varepsilon, A$ and $M$ independent of $\alpha, \beta$ and $\nu$. Here, $C_{\alpha}$ is a conttant independent of $\beta$ and $\nu$.

Together with Lemma 2.1 in the following section $q_{v+1}^{o}(x, \xi ; \xi)$ is a main symbol of (1.6).

Remark 1. The constants $\varepsilon, A, M$ and $C_{\alpha}$ in (1.14) are determined only by the dimension $n, \sigma, \tau^{0}$ and constants $C, M, C_{\alpha}, \varepsilon$ in (1.2), (1.3) and (1.4) for $p_{j}^{o}(x, \xi), \tilde{p}_{j}(x, \xi)$ and $\phi_{j}(x, \xi)$.

Remark 2. The estimate (1.14) still holds even if we replace $\Phi_{\nu+1}(x, \xi)$
by $x \cdot \xi$ in the corresponding formula (1.8) for the multi-products (1.12). This follows from (1.4) (see the discussion above Proposition 1.2).

We begin the proof of Lemma 1.3 with the estimation of the third term of (1.13). By means of Theorem 2.1 and Proposition 2.2 (and its Remark 2) in [18] it is clear that $\tilde{q}_{v+1}^{k}(x, \xi)$ satisfies the same inequalities as (1.14) because $\langle\xi\rangle^{(\nu+1) \sigma} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa}\right) \leqq A_{1}^{v} \nu!^{\sigma_{\kappa}} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa} / 2\right)$ for a suitable constant $A_{1}$. So, for the proof of the lemma it suffices to show

$$
\begin{equation*}
\left|\partial_{\xi}^{\alpha} \partial_{x}^{\beta} q_{\nu+1}^{k}(x, \xi ; \xi)\right| \leqq C_{\alpha} A^{\nu} M^{-|\beta|} \beta!^{\kappa} \nu!^{\sigma_{\kappa}} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa}\right) . \tag{1.15}
\end{equation*}
$$

To this end, we prepare the following three lemmas which are versions of propositions in $\S 2$ of [18]. Let $p(x, \xi ; \zeta)$ be a symbol in $S_{G(k)}^{m}$ with a parameter $\zeta \in R^{n}$ such that for any $\alpha, \alpha^{\prime}$ and $\beta$ we have

$$
\begin{equation*}
\left|\partial_{\xi}^{\alpha} \partial_{\xi}^{\alpha^{\prime}} \partial_{x}^{\beta} p\right| \leqq C M^{-\left(|\alpha|+\left|\alpha^{\prime}\right|+|\beta|\right)}\left(\alpha!\alpha^{\prime}!\beta!\right)^{\kappa}\langle\xi\rangle^{m-\left|\alpha^{\alpha}\right|-\left|\alpha^{\prime}\right|} . \tag{1.16}
\end{equation*}
$$

Lemma 1.4 (cf. Proposition 2.2 in [18]). Let $P_{1}(\zeta) \equiv p_{1}\left(X, D_{x} ; \zeta\right)$ be a pseudo-differential operator with a symbol $p_{1}(x, \xi ; \zeta) \in S_{G(k)}^{m}$ satisfying (1.16) and set $P_{2, \phi} \equiv p_{2, \phi}\left(X, D_{x}\right)$ be a Fourier integral operator with phase function $\phi(x, \xi)$ in $\mathcal{P}_{G(\mathrm{k})}$ and symbol $p_{2}(x, \xi)$ in $S_{G(\mathrm{k})}^{m^{\prime}}$. Then, we have the following:
i) The product $P_{1}(\zeta) P_{2, \phi}$ belongs to $L_{G(\kappa)}^{m+m^{\prime}}(\phi)$ and has the form

$$
\begin{equation*}
P_{1}(\zeta) P_{2, \phi}=q_{\phi}^{o}\left(X, D_{x} ; \zeta\right)+\widetilde{q}_{\phi}\left(X, D_{x} ; \zeta\right) \tag{1.17}
\end{equation*}
$$

with symbols $q^{\circ}(x, \xi ; \zeta)$ and $\tilde{q}(x, \xi ; \zeta)$ satisfying

$$
\begin{array}{r}
\left|\partial_{\xi}^{\alpha} \partial_{\xi}^{\alpha^{\prime}} \partial_{x}^{\beta} q^{o}(x, \xi ; \zeta)\right| \leqq C_{1} M_{1}^{-\left(|\alpha|+\left|\alpha^{\prime}\right|+|\beta| \mid\right.}\left(\alpha!\alpha^{\prime}!\beta!\right)^{\kappa} \\
\times\langle\xi\rangle^{m+m^{\prime}|\alpha|+\left|\alpha^{\prime}\right|}, \\
\left|\partial_{\xi}^{\alpha} \partial_{\zeta}^{\alpha^{\prime}} \partial_{x}^{\beta} \widetilde{q}(x, \xi ; \zeta)\right| \leqq C_{\alpha, \alpha^{\prime}} M_{1}^{-|\beta|} \beta!^{\kappa} \exp \left(-\varepsilon\langle\xi\rangle^{1 / k}\right) \tag{1.19}
\end{array}
$$

for constants $C_{1}, M_{1}, C_{\alpha, \alpha^{\prime}}$ and a positive constant $\varepsilon$.
ii) Let the symbol $p_{1}(x, \xi ; \zeta)$ satisfy for $a \delta>0$

$$
\begin{equation*}
p_{1}(x, \xi ; \zeta)=0 \quad \text { if } \quad|\xi-\zeta| \leqq \delta\langle\zeta\rangle . \tag{1.20}
\end{equation*}
$$

Then, there exists a positive constant $\tau^{0} \equiv \tau^{0}(\delta)$ such that for $\phi(x, \xi) \in \mathscr{P}_{G(k)}\left(\tau^{0}\right)$ the decomposition (1.17) still holds with (1.18)-(1.19) and the main symbol $q^{\circ}(x, \xi ; \zeta)$ of (1.17) satisfies (1.20) with $\delta$ replaced by $\delta^{\prime}>0$ depending only on $\delta$.
iii) For the product $P_{2, \phi} P_{1}(\zeta)$ we have the same statements as i) and ii).

Remark. The constants $C_{1}, M_{1}$ and $C_{\alpha, \alpha^{\prime}}$ are determined by the constants $\tau$ and $M$ in (1.4) for $\phi(x, \xi)$, the constants $C$ and $M$ in (1.16) for $p_{1}(x, \xi ; \zeta)$ and those in (1.2) for $p_{2}(x, \xi)$.

Proof. The first statement easily follows from the proof of Proposition
2.2 in [18]. In order to show ii) we recall that the main symbol of the product $P_{1}(\zeta) P_{2, \phi}$ is defined by the first term of the right hand side of (4.4) in [18], that is,

$$
\begin{array}{r}
q^{o}\left(x, \xi^{\prime} ; \zeta\right)=O_{s}-\iint e^{i \psi} p_{1}(x, \xi ; \zeta) \chi\left(\left(\xi-\xi^{\prime}\right)\left|<\xi^{\prime}\right\rangle\right)  \tag{1.21}\\
\times p_{2}\left(x^{\prime}, \xi^{\prime}\right) d x^{\prime} d \xi
\end{array}
$$

where $\psi=x \cdot \xi-x^{\prime} \cdot \xi+\phi\left(x^{\prime}, \xi^{\prime}\right)-\phi\left(x, \xi^{\prime}\right)$ and $\chi$ is a function satisfying (1.10). In (1.21) we replace $\chi(\xi)$ by $\chi(\xi / \theta)$ for a small $\theta>0$, that is,

$$
\begin{align*}
q^{o}\left(x, \xi^{\prime} ; \zeta\right)=O_{s}-\iint e^{i \psi} p_{1}(x, \xi ; \zeta) & \chi\left(\left(\xi-\xi^{\prime}\right) /\left(\theta\left\langle\xi^{\prime}\right\rangle\right)\right)  \tag{1.21}\\
& \times p_{2}\left(x^{\prime}, \xi^{\prime}\right) d x^{\prime} d \xi
\end{align*}
$$

Then, we can prove (1.20) for $q^{o}(x, \xi ; \zeta)$ if we take $\theta$ small enough. This exchange is harmless for the proof of (1.18) and (1.19) if we use $\theta>5 \tau / 2$ when we prove (1.19). The proof of iii) is similar to those of i) and ii). Q.E.D.

Let $I_{\phi}$ (resp. $I_{\phi^{*}}$ ) denote the Fourier (resp. the conjugate Fourier) integral operator with symbol 1 .

Lemma 1.5. Assume that $P_{\phi}(\zeta) \equiv p_{\phi}\left(X, D_{x} ; \zeta\right)$ is a Fourier integral operator with a symbol $p(x, \xi ; \zeta) \in S_{G(k)}^{m}$ satisfying (1.16). Then, we have

$$
\begin{equation*}
P_{\phi}(\zeta) I_{\phi^{*}} \in L_{G(\kappa)}^{m}, \quad I_{\phi^{*}} P_{\phi}(\zeta) \in L_{G(k)}^{m} \tag{1.22}
\end{equation*}
$$

and about the symbols of the product operators we have the corresponding results to (1.17)-(1.19) with $m+m^{\prime}$ replaced by $m$ in (1.18). If $p(x, \xi ; \zeta)$ satisfies (1.20), then there exists a positive $\tau^{0} \equiv \tau^{0}(\delta)$ satisfying the following property: If $\phi(x, \xi)$ $\in \mathscr{P}_{G(k)}\left(\tau^{0}\right)$ then, adding to (1.17)-(1.19), the main symbols of $P_{\phi}(\zeta) I_{\phi^{*}}$ and $I_{\phi^{*}} P_{\phi}$ $(\zeta)$ satisfy (1.10) with $\delta$ replaced by $\delta^{\prime}\left(0<\delta^{\prime}<\delta\right)$.

Remark. We have the similar statement as in the remark of Lemma 1.4.
Proof. The formula (1.22) is the same as (2.10) in [18]. For the proof of the last statement we replace $\chi(\xi) \in \gamma^{(\kappa)}\left(R_{\xi}^{n}\right)$ by $\chi(\xi / \theta)$ with a sufficiently small $\theta>0$ when we proceed the proof of Proposition 2.3 in [18]. Then, by means of the inequality (2.3)-a) in [17] we obtain the desired main symbol keeping the properties (1.17)-(1.19).
Q.E.D.

Lemma 1.6. Let $\phi_{j}(x, \xi)$ belong to $\mathscr{P}_{G(k)}\left(\tau_{j}\right), j=1,2, \tau_{1}+\tau_{2} \leqq \tau^{0}$ for a sufficiently small $\tau^{0}>0$, and let $P_{\phi_{2}}(\zeta) \in S_{G(\kappa), \phi_{2}}^{m}$ be the same as in Lemma 1.5. Then there exists a pseudo-differential operator $P^{\prime}(\zeta) \equiv p^{\prime}\left(X, D_{x} ; \zeta\right)=p^{\prime 0}\left(X, D_{x} ; \zeta\right)$ $+\tilde{p}^{\prime}\left(X, D_{x} ; \zeta\right)$ in $L_{G(\kappa)}^{m}$ such that

$$
\begin{equation*}
I_{\phi_{1}} P_{\phi_{2}}(\zeta)=P^{\prime}(\zeta) I_{\phi_{1} \phi_{2}} \tag{1.23}
\end{equation*}
$$

and $p^{\prime o}(x, \xi ; \zeta)$ and $\tilde{p}^{\prime}(x, \xi ; \zeta)$ satisfy (1.18) (with $m+m^{\prime}$ replaced by $m$ ) and (1.19), respectively. Furthermore, if $p(x, \xi ; \zeta)$ satisfies $(1.20)$ the main symbol $p^{\prime 0}$ $(x, \xi ; \zeta)$ also satisfies (1.20) with $\delta$ replaced by $\delta^{\prime}\left(0<\delta^{\prime}<\delta\right)$, provided that we take $\tau^{0}$ sufficiently small corresponding to $\delta$.

Remark. We have the similar assertion as in the remark of Lemma 1.4.
Proof. The formula (1.23) follows from Propositions 2.2, 2.3, 2.5 and Corollary 2.8 in [18], as in the proof of Lemma 2.10 in [17]. For the proof of (1.20) for $p^{\prime o}(x, \xi ; \zeta)$ we use Lemmas 1.4 and 1.5 repeatedly. Then, we get the lemma.
Q.E.D.

Now, we are prepared to prove (1.15).
Proof of (1.15). It follows from Lemma 1.6 that there exist pseudodifferential operators $P_{j}^{\prime}(\zeta)=p_{j}^{\prime}\left(X, D_{x} ; \zeta\right)(j=1, \cdots, k)$ and $P_{j}^{\prime}=p_{j}^{\prime}\left(X, D_{x}\right)(j=$ $k+1, \cdots, \nu+1)$ in $L_{G(k)}^{\sigma}$ such that

$$
\begin{cases}I_{\Phi_{j-1}} P_{j, \phi_{j}}^{+}(\zeta)=P_{j}^{\prime}(\zeta) I_{\Phi_{j}}, & j=1, \cdots, k-1 \quad\left(\Phi_{0}=x \cdot \xi\right),  \tag{1.24}\\ I_{\Phi_{k-1}} P_{\bar{k}, \phi_{k}}(\zeta)=P_{k}^{\prime}(\zeta) I_{\Phi_{k}}, & \\ I_{\Phi_{j-1}} P_{j, \phi_{j}}=P_{j}^{\prime} I_{\Phi_{j}}, & j=k+1, \cdots, \nu+1\end{cases}
$$

As in the last paragraph of $\S 2$ in [18] it follows that

$$
\begin{aligned}
& P_{1, \phi_{1}}^{+}(\zeta) \cdots P_{k-1, \phi_{k-1}}^{+}(\zeta) P_{k, \phi_{k}}^{-}(\zeta) P_{k+1, \phi_{k+1}} \cdots P_{\nu+1, \phi_{\nu+1}} \\
& \quad=P_{1}^{\prime}(\zeta) \cdots P_{k-1}^{\prime}(\zeta) P_{k}^{\prime}(\zeta) P_{k+1}^{\prime} \cdots P_{\nu+1}^{\prime} I_{\Phi_{\nu+1}} .
\end{aligned}
$$

If we apply Theorem 2.6 of [18] to the multi-product $Q_{v+1, k}^{\prime}(\zeta)=P_{1}^{\prime}(\zeta) \cdots P_{k-1}^{\prime}(\zeta)$ - $P_{k}^{\prime}(\zeta) P_{k+1}^{\prime} \cdots P_{v+1}^{\prime}$ of pseudo-differential operators we have

$$
Q_{\nu+1, k}^{\prime}(\zeta)=q_{\nu+1}^{k, o}\left(X, D_{x} ; \zeta\right)+q_{\nu+1}^{\nu, \infty}\left(X, D_{x} ; \zeta\right) \in L_{G(k)}^{(\nu+1) \sigma} .
$$

As in the proof of Lemmas $1.4-1.6$ we exchange $\chi(\xi)$ by $\chi(\xi / \theta)$ for a sufficiently small $\theta>0$ in the proof of Proposition 5.1 in [18]. Then, in view of (1.16), the symbols $q_{v+1}^{k, o}(x, \xi ; \zeta)$ and $q_{v+1}^{k, \infty}(x, \xi ; \zeta)$ satisfy

$$
\left\{\begin{align*}
&\left|\partial_{\xi}^{\alpha} \partial_{\xi}^{\alpha^{\prime}} \partial_{x}^{\beta} q_{\nu+1}^{k, o}(x, \xi ; \zeta)\right| \leqq A^{\nu} M^{-\left(\left|\alpha^{\alpha}\right|+\left|\alpha^{\prime}\right|+\left|\beta^{\beta}\right|\right.}\left(\alpha!\alpha^{\prime}!\beta!\right)^{\kappa}  \tag{1.25}\\
& \times\langle\xi\rangle^{(\nu+1)^{\sigma-|\alpha|-\left|\alpha^{\prime}\right|}}, \\
&\left|\partial_{\xi}^{\alpha} \partial_{\xi}^{\alpha^{\prime}} \partial_{x}^{\beta} q_{\nu+1}^{k, \infty}(x, \xi ; \zeta)\right| \leqq C_{\alpha, \alpha^{\prime}} \mathcal{A}^{\nu} M^{-|\beta|} \beta!^{\kappa} \nu!^{\sigma \kappa} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa}\right)
\end{align*}\right.
$$

and, in addition, the main symbol $q_{\nu+1}^{k, o}(x, \xi ; \zeta)$ satisfies (1.20) for a $\delta^{\prime}>0$ independent of $\nu$. Here, we used the main symbol of $P_{k}^{\prime}(\zeta)$ satisfies (1.20) for a small $\delta>0$. Finally we use Lemma 1.4. Then the main symbol of $q_{v+1}^{k, o}\left(X, D_{x} ; \zeta\right) I_{\Phi_{\nu+1}}$ vanishes when $\zeta=\xi$. Noting (2.7) of [18] and the remarks of Proposition 2.2 of [18] we obtain (1.15) by (1.25). This concludes the proof of Lemma 1.3.
Q.E.D.

## 2. Multiple symbols and lemmas

For $x^{j}, \xi^{j} \in R^{n}(j=1, \cdots, \nu)$ we write $\tilde{x}^{\nu}=\left(x^{1}, \cdots, x^{\nu}\right)$ and $\xi^{\nu}=\left(\xi^{1}, \cdots, \xi^{\nu}\right)$. We consider a multiple symbol $p_{\nu+1}\left(x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) \in C^{\infty}$ satisfying for any $\alpha, \beta$, $\widetilde{\alpha}^{\nu}=\left(\alpha^{1}, \cdots, \alpha^{\nu}\right)$ and $\tilde{\beta}^{\nu}=\left(\beta^{1}, \cdots, \beta^{\nu}\right)$

$$
\begin{align*}
& \left|\partial_{x}^{\beta} \partial_{\tilde{\xi}^{\nu}}^{\tilde{\sim}^{\nu}} \partial_{\tilde{x}^{\nu}}^{\tilde{\nu}^{\nu}} \partial_{\xi}^{\alpha} p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)\right| \\
& \quad \leqq C_{o} M^{-\left(|\alpha|+\left|\beta_{|+|}^{\alpha} \tilde{\alpha}^{\nu}\right|+\left|\tilde{\beta}^{\nu}\right|\right)}\left(\alpha!\beta!\tilde{\alpha}^{\nu}!\tilde{\beta}^{\nu}!\right)^{\kappa} \prod_{j=1}^{\nu+1}\left\langle\xi^{j}\right\rangle^{m-\mid \alpha j_{\mid}}  \tag{2.1}\\
& \quad\left(\alpha^{\nu+1}=\alpha, \quad \xi^{\nu+1}=\xi\right)
\end{align*}
$$

for positive constants $C_{o}$ and $M$ independent of $\nu$. We consider a simplified symbol $q_{v+1}(x, \xi)$ defined by

$$
\left.\left.\begin{array}{rl}
q_{\nu+1}(x, \xi)=O_{s}-\int \ldots \int & \exp [
\end{array} \quad\left(\Psi_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)-\Phi_{\nu+1}(x, \xi)\right)\right]\right)
$$

where $\Psi_{\nu+1}\left(x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ is defined by (1.9) for a sequence $\left\{\phi_{j}\right\}$ of phase functions $\phi_{j}(x, \xi)$ satisfying (A-1) in the preceding section. The integral of the right hand side of (2.2) means the oscillatory integral, that is,

$$
\begin{align*}
\lim _{\varepsilon \neq 0} \int \ldots \int \exp & {\left[i\left(\Psi_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)-\Phi_{\nu+1}(x, \xi)\right)\right] } \\
& \times \prod_{j=1}^{\nu} \chi\left(\varepsilon x^{j}\right) \chi\left(\varepsilon \xi^{j}\right) p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{x}^{\nu} d \tilde{\xi}^{\nu} \tag{2.3}
\end{align*}
$$

with $\chi \in C^{\infty}$ satisfying (1.10). We shall show this limit is well-defined. Set $\tilde{p}_{\nu+1}\left(x, \quad \hat{\xi}^{\nu}, \quad \tilde{x}^{\nu}, \quad \xi\right)=p_{\nu+1}\left(x, \quad \tilde{\xi}^{\nu}, \quad \tilde{x}^{\nu}, \quad \xi\right) \prod_{j=1}^{\nu+1} \exp \left(i J_{j}\left(x^{j-1}, \xi^{j}\right)\right) \cdot \exp \left(-i \Phi_{\nu+1}(x, \xi)\right)$ $\left(x^{0}=x, \xi^{\nu+1}=\xi\right)$ with $J_{j}(x, \xi)=\phi_{j}(x, \xi)-x \cdot \xi$. Then, we have

$$
\begin{equation*}
\left|\partial_{\tilde{\xi}^{\nu}}^{\tilde{\alpha}^{\nu}} \partial_{\tilde{x}^{\nu}}^{\tilde{\beta}^{\nu}} \tilde{p}_{v+1}\right| \leqq C_{\tilde{\alpha}^{\nu}, \tilde{\beta}^{\nu}}\left\langle\xi^{1}\right\rangle^{m_{1}} \prod_{j=1}^{\nu}\left\langle\xi^{j+1}\right\rangle^{m_{j+1}+\left|\beta^{\beta j}\right|} . \tag{2.4}
\end{equation*}
$$

Set $n_{o}=[n / 2]+1$ and define integers $l_{j}(j=1,2, \cdots)$ inductively by $l_{1}=\left[\left(m_{1}+n\right) /\right.$ $2]+1, l_{j}=\left[\left(m_{1}+\cdots+m_{j}+n\right) / 2\right]+l_{1}+\cdots+l_{j-1}$. Then, it follows from the integration by parts that the limit (2.3) equals

$$
\begin{aligned}
\int \cdots \int \exp & {\left[i\left(\sum_{j=1}^{v}\left(x^{j-1}-x^{j}\right) \cdot \xi^{j}+x^{\nu} \cdot \xi\right)\right] } \\
\times & \prod_{j=1}^{\nu}\left\{\left(1+\left|x^{j-1}-x^{j}\right|^{2}\right)^{-n_{o}}\left(1-\Delta_{\xi^{j}}\right)^{n^{\circ}}\right\} \\
\cdot & {\left[\prod_{j=1}^{\nu}\left(1+\left|\xi^{j}-\xi^{j+1}\right|^{2}\right)^{-l_{j}\left(1-\Delta_{x^{j}}\right)^{l_{j}}}\right.} \\
\cdot & p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{x}^{\nu} d \tilde{\xi}^{\nu} \quad\left(x^{0}=x, \quad \xi^{\nu+1}=\xi\right),
\end{aligned}
$$

which is well-defined by means of (2.4). Therefore, we can exchange the order of integration and differentiate the right hand side of (2.2) under the integral sign (that is, we obtain the Fubini theorem and Lebesgue's convergence theorem for the oscillatory integral, see $\S 6$ of Chapter 1 of [9]).

Set $p_{\nu+1}\left(x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)=\prod_{j=1}^{\nu+1} p_{j}\left(x^{j-1}, \xi^{j}\right)\left(x^{0}=x, \xi^{\nu+1}=\xi\right)$ for a sequence $\left\{p_{j}\right\}$ of symbols $p_{j}(x, \xi)$ satisfying (A-2) in the preceding section. Then, $p_{v+1}\left(x, \tilde{\xi}^{\nu}\right.$, $\tilde{x}^{\nu}, \xi$ ) satisfies (2.1). Taking the decomposition (1.13) with $\zeta=\xi$ and Lemma 1.3 into account, we may investigate only $q_{\nu+1}^{o}(x, \xi ; \xi)$ in (1.13). So, in what follows we may assume

$$
\begin{equation*}
\left|\xi^{j}-\xi\right| \leqq\langle\xi\rangle / 8 \quad \text { on } \quad \operatorname{supp} p_{\nu+1} \tag{2.5}
\end{equation*}
$$

Then, the condition (2.1) is reduced to

$$
\begin{align*}
& \left|\partial_{x}^{\beta} \partial_{\tilde{\xi}^{\nu}}^{\tilde{\alpha}^{\nu}} \partial_{\tilde{x}_{\nu}^{\nu}}^{\tilde{\beta}_{\nu}^{\nu}} \partial_{\xi}^{\alpha} p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)\right|  \tag{2.6}\\
& \quad \leqq C_{0} M^{-\left(\left|\left.\right|^{\alpha}\right|+\left|\beta_{1}\right|+\left|\tilde{\alpha}^{\nu}\right|+\left|\tilde{\beta}^{\nu}\right|\right)}\left(\alpha!\beta!\tilde{\alpha}^{\nu}!\tilde{\beta}^{\nu}!\right)^{\kappa}\langle\xi\rangle^{m-\left|\alpha_{1-1}\right| \tilde{\alpha}^{\nu} \mid} .
\end{align*}
$$

with $m=\sum_{j=1}^{\nu+1} m_{j}$. So, till the end of Section 3 we always assume (2.5) and (2.6).
Lemma 2.1. Assume (2.5), (2.6) and (A-1) in §1. Then, for $q_{\nu+1}(x, \xi)$ defined by (2.2) we have

$$
\begin{equation*}
\left|q_{\nu+1(\beta)}^{(\alpha)}(x, \xi)\right| \leqq C_{o} A^{\nu} M_{1}^{-(|\alpha|+|\beta|}(\alpha!\beta!)^{\kappa}\langle\xi\rangle^{m-|\alpha|} \tag{2.7}
\end{equation*}
$$

for any $\alpha$ and $\beta$, where $A$ and $M_{1}$ are independent of $\nu$.
Remark. For the multiple symbol satisfying only (2.1) we can obtain the conclusion of the lemma modulo the regularizer satisfying (1.14). The proof, however, is fairly long. So, in this paper we restrict ourselves to prove the lemma in the above form.

For any $k \in\{1, \cdots, \nu\}$ we write $\tilde{x}^{\nu, k}=\left(x^{1}, \cdots, x^{k-1}, x^{k+1}, \cdots, x^{\nu}\right)$ and set

$$
\begin{align*}
\Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)= & \left(\phi_{1} \# \cdots \# \phi_{k}\right)\left(x, \xi^{k}\right)-x^{k} \cdot \xi^{k} \\
& +\left(\phi_{k+1} \# \cdots \# \phi_{\nu+1}\right)\left(x^{k}, \xi\right) . \tag{2.8}
\end{align*}
$$

Then, we have
Lemma 2.2. Assume (2.5), (2.6) and (A-1). Then, the symbol $q_{\nu+1, k}\left(x, \xi^{k}\right.$, $\left.x^{k}, \xi\right), k=1, \cdots, \nu$, defined by

$$
\begin{align*}
& q_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)  \tag{2.9}\\
& =O_{s}-\int \cdots \int \exp \left[i\left(\Psi_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)-\Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)\right)\right] \\
& \quad \times p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{x}^{\nu, k} d \xi^{\nu, k}
\end{align*}
$$

satisfy

$$
\begin{align*}
& \left|\partial_{\xi}^{\alpha} \partial_{\xi^{k}}^{\alpha k} \partial_{x}^{\beta} \partial_{x^{k}}^{\rho_{k}^{k}} q_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)\right| \\
& \leqq C_{o} A^{\nu-1} M_{1}^{-\left(\left|\alpha^{\alpha}\right|+\left|\alpha^{k}\right|+\left|\beta^{\beta}\right|+\left|\beta^{k}\right|\right)}\left(\alpha!\alpha^{k}!\beta!\beta^{k}!\right)^{\kappa}  \tag{2.10}\\
& \times\langle\xi\rangle^{m-\left|\alpha^{\alpha}\right|-\left|\alpha^{k}\right|}
\end{align*}
$$

for constants $A$ and $M_{1}$ independent of $\nu$ and $k$.
Remark. From the definition (2.9) it follows immediately that if the support of $p_{v+1}$ with respect to $\left(x^{k}, \xi^{k}\right)$ is contained in a subset $\Omega$ of $R_{\left(x^{k}, \xi^{k}\right)}^{2 n}$ then the support of $q_{\nu+1, k}$ with respect to $\left(x^{k}, \xi^{k}\right)$ is also contained in the same subset $\Omega$.

In the following we prove Lemma 2.1 and Lemma 2.2. For the proof we employ the following, which is the case $\nu=1$ in Lemma 2.1 if we set $\zeta=\xi$.

Lemma 2.3. Let $\phi_{j} \in \mathcal{P}_{G(k)}\left(\tau_{j}\right), j=1,2\left(\tau_{1}+\tau_{2} \leqq 1 / 4\right)$ and set $\psi\left(x, \xi^{\prime}, x^{\prime}, \zeta\right)$ $=\phi_{1}\left(x, \xi^{\prime}\right)-x^{\prime} \cdot \xi^{\prime}+\phi_{2}\left(x^{\prime}, \zeta\right)-\Phi_{2}(x, \zeta)$, where $\Phi_{2}(x, \xi)=\phi_{1} \# \phi_{2}(x, \xi)$. Let $p\left(x, \xi^{\prime}\right.$, $\left.x^{\prime}, \zeta ; \xi\right) \in C^{\infty}$ satisfy

$$
\begin{align*}
\mid \partial_{\xi}^{\alpha} \partial_{\xi^{\prime}}^{\alpha^{\prime}} \partial_{\xi}^{\alpha^{\prime \prime}} \partial_{x}^{\beta} \partial_{x^{\prime}}^{\beta^{\prime}} & p \mid \leqq C_{o} M^{-\left(\left|\alpha_{1}\right|+\left|\alpha^{\prime}\right|+\left|\alpha^{\prime \prime}\right|+\left|\beta^{\beta}\right|+\left|\beta^{\prime}\right|\right)} \\
& \times\left(\alpha!\alpha^{\prime}!\alpha^{\prime \prime}!\beta!\beta^{\prime}!\right)^{\kappa}\langle\xi\rangle^{m-|\kappa|-\left|\alpha^{\prime}\right|-\left|\alpha^{\prime \prime}\right|} \tag{2.11}
\end{align*}
$$

and assume

$$
\begin{equation*}
\left|\xi^{\prime}-\xi\right| \leqq\langle\xi\rangle / 8, \quad|\zeta-\xi| \leqq\langle\xi\rangle / 4 \quad \text { on } \quad \operatorname{supp} p \tag{2.12}
\end{equation*}
$$

Set

$$
\begin{equation*}
q(x, \zeta ; \xi)=O_{s}-\iint e^{i \psi\left(x, \zeta^{\prime}, x^{\prime}, \zeta\right)} p\left(x, \xi^{\prime}, x^{\prime}, \zeta ; \xi\right) d x^{\prime} d \xi^{\prime} \tag{2.13}
\end{equation*}
$$

Then, there exist constants $A$ and $M_{1}$ depending only on $M$ such that

$$
\begin{array}{r}
\left|\partial_{\xi}^{\alpha} \partial_{\xi}^{\alpha^{\prime}} \partial_{x}^{\beta} q(x, \zeta ; \xi)\right| \leqq C_{0} A M_{1}^{-\left(\left|\alpha_{1}\right|+\left|\alpha^{\prime}\right|+\left|\beta^{\beta}\right|\right)}\left(\alpha!\alpha^{\prime}!\beta!\right)^{\kappa}  \tag{2.14}\\
\times\langle\xi\rangle^{m-\left|\alpha^{\prime}\right|-\left|\alpha^{\prime}\right|}
\end{array}
$$

for any $\alpha, \alpha^{\prime}$ and $\beta$.
Proof. Consider (2.13) instead of (4.11) in [18]. Then, we can prove (2.14) in the almost same way as in the proof of Proposition 2.5 in [18]. We need not consider the part corresponding to the estimation of $p_{\infty}\left(x, \xi^{\prime}\right)$ there because of (2.12). We can begin with the step (II) of the proof of Proposition 2.5 in [18] by replacing $\chi_{o}\left(\xi, \xi^{\prime}\right)$ by $p\left(x, \xi^{\prime}, x^{\prime}, \zeta ; \xi\right)$. Hence, we get (2.14) in view of (4.25) and (4.33) in [18].
Q.E.D.

Proof of Lemma 2.1. We divide the proof into two steps.
(I) In this step we shall show

$$
\begin{align*}
& q_{\nu+1}(x, \xi) e^{i \Phi_{\nu+1}(x, \xi)} \\
& \quad=O_{s}-\int \ldots \int \exp \left(i\left(\sum_{j=1}^{\nu}\left(x^{j-1}-x^{j}\right) \cdot \xi^{j}+\Phi_{\nu+1}\left(x^{\nu}, \xi\right)\right)\right)  \tag{2.15}\\
& \quad \times p_{\nu+1}^{\prime}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{x}^{\nu} d \tilde{\xi}^{\nu} \quad\left(x^{0}=x\right),
\end{align*}
$$

where the symbol $p_{v+1}^{\prime}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ satisfies

$$
\begin{equation*}
c^{-1}\langle\xi\rangle \leqq\left\langle\xi^{j}\right\rangle \leqq c\langle\xi\rangle \quad \text { on } \quad \operatorname{supp} p_{v+1}^{\prime} \tag{2.16}
\end{equation*}
$$

and (2.6) with $C_{o}$ and $M$ replaced by $C_{o} A^{\nu}$ and $M_{1}$. Here, the constants $A$ and $M_{1}$ are independent of $\nu$. For simplicity we consider the case $\nu=2$ for a while. Since it follows that

$$
e^{i \phi_{1}\left(x, \xi^{\prime}\right)}=O_{s}-\iint e^{i z \cdot\left(\xi^{\prime}-\zeta\right)+i \phi_{1}(x, \zeta)} \chi\left(4\left(\zeta-\xi^{\prime}\right) \mid\langle\xi\rangle\right) d z d \zeta
$$

for $\chi(\xi) \in \gamma^{(k)}\left(R_{\xi}^{n}\right)$ satisfying (1.10), we get

$$
\begin{aligned}
& I_{1}\left(x, \xi^{2}, x^{2}, \xi\right) \\
& \equiv O_{s}-\iint e^{i\left(\phi_{1}\left(x, \xi^{1}\right)-x^{1} \cdot \xi^{1}+\phi_{2}\left(x^{1}, \xi^{2}\right)\right)} p_{3}\left(x, \tilde{\xi}^{2}, \tilde{x}^{2}, \xi\right) d x^{1} d \xi^{1} \\
&= O_{s}-\iiint \int e^{i\left(\phi_{1}\left(x, \zeta^{1}\right)-z^{1} \cdot\left(\zeta^{1}-\xi^{1}\right)-x^{1} \cdot \xi^{1}+\phi_{2}\left(x^{1}, \xi^{2}\right)\right)} \\
& \times \chi\left(4\left(\zeta^{1}-\xi^{1}\right)\langle\xi\rangle\right) p_{3}\left(x, \tilde{\xi}^{2}, \tilde{x}^{2}, \xi\right) d x^{1} d \xi^{1} d z^{1} d \zeta^{1} .
\end{aligned}
$$

Next, we take changes of variables as follows: First we change the variable $z^{1}$ to $y^{1}$ as

$$
\begin{equation*}
z^{1}=\tilde{\nabla}_{\xi} \phi_{1}\left(y^{1} ; \zeta^{1}, \xi^{1}\right) \equiv \int_{0}^{1} \nabla_{\xi} \phi_{1}\left(y^{1}, \xi^{1}+\theta\left(\zeta^{1}-\xi^{1}\right)\right) d \theta \tag{2.17}
\end{equation*}
$$

and then the variable $\zeta^{1}$ to $\eta^{1}$ as

$$
\begin{equation*}
\eta^{1}=\tilde{\nabla}_{x} \phi_{1}\left(x, y^{1} ; \zeta^{1}\right) \equiv \int_{0}^{1} \nabla_{x} \phi_{1}\left(y^{1}+\theta\left(x-y^{1}\right), \zeta^{1}\right) d \theta . \tag{2.18}
\end{equation*}
$$

Then, we have

$$
\begin{aligned}
& \phi_{1}\left(x, \zeta^{1}\right)-z^{1} \cdot\left(\zeta^{1}-\xi^{1}\right)=\phi_{1}\left(x, \zeta^{1}\right)-\tilde{\nabla}_{\xi} \phi_{1}\left(y^{1} ; \zeta^{1}, \xi^{1}\right) \cdot\left(\zeta^{1}-\xi^{1}\right) \\
& \quad=\phi_{1}\left(x, \zeta^{1}\right)-\phi_{1}\left(y^{1}, \zeta^{1}\right)+\phi_{1}\left(y^{1}, \xi^{1}\right) \\
& \quad=\left(x-y^{1}\right) \cdot \tilde{\nabla}_{x} \phi_{1}\left(x, y^{1} ; \zeta^{1}\right)+\phi_{1}\left(y^{1}, \xi^{1}\right) \\
& \quad=\left(x-y^{1}\right) \cdot \eta^{1}+\phi_{1}\left(y^{1}, \xi^{1}\right)
\end{aligned}
$$

and

$$
d z^{1} d \zeta^{1}=r_{1}\left(x, \xi^{1}, y^{1}, \tilde{\nabla}_{x} \phi_{1}^{-1}\left(x, y^{1} ; \eta^{1}\right)\right) d y^{1} d \eta^{1}
$$

where $\zeta=\tilde{\nabla}_{x} \phi_{1}^{-1}(x, y ; \eta)$ is the inverse function of $\eta=\tilde{\nabla}_{x} \phi_{1}(x, y ; \zeta)$ and

$$
r_{1}(x, \xi, y, \zeta)=\operatorname{det} \frac{\partial}{\partial x} \tilde{\nabla}_{\xi} \phi_{1}(y ; \zeta, \xi) \cdot\left(\operatorname{det} \frac{\partial}{\partial \xi} \tilde{\nabla}_{x} \phi_{1}(x, y ; \zeta)\right)^{-1}
$$

In the above we denote $\frac{\partial}{\partial x} f(x)$, for a vector $f={ }^{t}\left(f_{1}, \cdots, f_{n}\right)$ of functions $f_{j}(x)$, as $\frac{\partial}{\partial x} f(x)=\left(\partial f_{j} / \partial x_{k} \underset{k \rightarrow 1, \cdots, n}{j \downarrow 1, \cdots, n}\right)$.
Now, we set

$$
\begin{equation*}
\tilde{r}_{1}\left(x, \xi^{\prime}, y, \zeta ; \xi\right)=\chi\left(4\left(\zeta-\xi^{\prime}\right) /\langle\xi\rangle\right) r_{1}\left(x, \xi^{\prime}, y, \zeta\right) \tag{2.19}
\end{equation*}
$$

and set

$$
\begin{align*}
& q_{3}^{\prime}\left(x, \zeta^{1}, y^{1}, \xi^{2}, x^{2}, \xi ; \zeta^{\prime}\right) \\
& =O_{s}-\iint e^{i\left(\phi_{1}\left(y^{1}, \xi^{1}\right)-x^{1} \cdot \xi^{1}+\phi_{2}\left(x^{1}, \zeta^{\prime}\right)-\Phi_{2}\left(y^{1}, \zeta^{\prime}\right)\right)}  \tag{2.20}\\
& \quad \times \widetilde{r}_{1}\left(x, \xi^{1}, y^{1}, \zeta^{1} ; \xi\right) p_{3}\left(x, \tilde{\xi}^{2}, \tilde{x}^{2}, \xi\right) d x^{1} d \xi^{1}
\end{align*}
$$

with $\Phi_{2}=\phi_{1} \# \phi_{2}$. Then, by the change of variables (2.17) and (2.18), we have

$$
\begin{aligned}
& I_{1}\left(x, \xi^{2}, x^{2}, \xi\right)=O_{s}-\iint e^{i\left(\left(x-y^{1}\right) \cdot \eta^{1}+\Phi_{2}\left(y^{1}, \xi^{2}\right)\right)} \\
& \quad \times q_{3}^{\prime}\left(x, \tilde{\nabla}_{x} \phi_{1}^{-1}\left(x, y^{1} ; \eta^{1}\right), y^{1}, \xi^{2}, x^{2}, \xi ; \xi^{2}\right) d y^{1} d \eta^{1}
\end{aligned}
$$

By the same way we can obtain

$$
\begin{aligned}
& I_{2}\left(x, \zeta^{1}, y^{1}, \xi\right) \equiv O_{s}-\iint e^{i\left(\Phi_{2}\left(y^{1}, \xi^{2}\right)-x^{2} \cdot \xi^{2}+\phi_{3}\left(x^{2}, \xi\right)\right)} \\
& \quad \times q_{3}^{\prime}\left(x, \zeta^{1}, y^{1}, \xi^{2}, x^{2}, \xi ; \xi^{2}\right) d x^{2} d \xi^{2} \\
& = \\
& O_{s}-\iiint \iint^{\left.i\left(\Phi_{2}\left(y^{1}, \zeta^{2}\right)-z^{2}\right) \cdot\left(\zeta^{2}-\xi^{2}\right)-x^{2} \cdot \xi^{2}+\phi_{3}\left(x^{2}, \xi\right)\right)} \\
& \left.\quad \times x\left(4\left(\zeta^{2}-\xi^{2}\right) /\langle\xi\rangle\right)\right) q_{3}^{\prime}\left(x, \zeta^{1}, y^{1}, \xi^{2}, x^{2}, \xi ; \zeta^{2}\right) d x^{2} d \xi^{2} d z^{2} d \zeta^{2} \\
& = \\
& \quad O_{s}-\iint e^{i\left(\left(y^{1}-y^{2}\right) \cdot \eta^{2}+\Phi_{3}\left(y^{2}, \xi\right)\right)} \\
& \quad \times \tilde{q}_{3}\left(x, \zeta^{1}, y^{1}, \tilde{\nabla}_{x} \Phi_{2}^{-1}\left(y^{1}, y^{2} ; \eta^{2}\right), y^{2}, \xi ; \xi\right) d y^{2} d \eta^{2}
\end{aligned}
$$

where

$$
\begin{align*}
& \widetilde{q}_{3}\left(x, \zeta^{1}, y^{1}, \zeta^{2}, y^{2}, \xi ; \zeta^{\prime}\right) \\
& =O_{s}-\iint e^{i\left(\Phi_{2}\left(y^{2}, \xi^{2}\right)-x^{2} \cdot \xi^{2}+\Phi_{3}\left(x^{2}, \zeta^{\prime}\right)-\Phi_{3}\left(y^{2}, \zeta^{\prime}\right)\right)}  \tag{2.21}\\
& \quad \times \widetilde{r}_{2}\left(y^{1}, \xi^{2}, y^{2}, \zeta^{2} ; \xi\right) q_{3}^{\prime}\left(x, \zeta^{1}, y^{1}, \xi^{2}, x^{2}, \xi ; \zeta^{2}\right) d x^{2} d \xi^{2} .
\end{align*}
$$

Here, $\Phi_{3}=\Phi_{2} \# \phi_{3}\left(=\phi_{1} \# \phi_{2} \# \phi_{3}\right)$ and $\tilde{r}_{2}\left(x, \xi^{\prime}, y, \zeta ; \xi\right)$ is defined by the same way as $\widetilde{r}_{1}$ with $\phi_{1}$ replaced by $\Phi_{2}$. Hence, we obtain

$$
\begin{aligned}
& q_{3}(x, \xi) e^{i \Phi_{3}(x, \xi)} \\
& \quad=O_{s}-\iint I_{1}\left(x, \xi^{2}, x^{2}, \xi\right) e^{i\left(-x^{2} \cdot \xi^{2}+\phi_{3}\left(x^{2}, \xi\right)\right)} d x^{2} d \xi^{2}
\end{aligned}
$$

$$
\begin{gathered}
=O_{s}-\iint e^{i\left(x-y^{1}\right) \cdot \eta^{1}} I_{2}\left(x, \tilde{\nabla}_{x} \phi_{1}^{-1}\left(x, y^{1} ; \eta^{1}\right), y^{1}, \xi\right) d y^{1} d \eta^{1} \\
=O_{s}-\iiint \int \exp \left[i\left(\sum_{j=1}^{2}\left(y^{j-1}-y^{j}\right) \cdot \eta^{j}+\Phi_{3}\left(y^{2}, \xi\right)\right)\right] \\
\times p_{3}^{\prime}\left(x, \tilde{\eta}^{2}, \tilde{y}^{2}, \xi\right) d \tilde{y}^{2} d \tilde{\eta}^{2},
\end{gathered}
$$

where $p_{3}^{\prime}\left(x, \tilde{\eta}^{2}, \tilde{y}^{2}, \xi\right)=\tilde{q}_{3}\left(x, \tilde{\nabla}_{x} \phi_{1}^{-1}\left(x, y^{1}, \eta^{1}\right), y^{1}, \tilde{\nabla}_{x} \Phi_{2}^{-1}\left(y^{1}, y^{2} ; \eta^{2}\right), y^{2}, \xi ; \xi\right)$. It follows from (2.20) and (2.21) that

$$
\begin{aligned}
& p_{3}^{\prime}\left(x, \tilde{\eta}^{2}, \tilde{y}^{2}, \xi\right) \\
& \quad=p_{3}^{\prime \prime}\left(x, \tilde{\nabla}_{x} \phi_{1}^{-1}\left(x, y^{1}, \eta^{1}\right), y^{1}, \tilde{\nabla}_{x} \Phi_{2}^{-1}\left(y^{1}, y^{2} ; \eta^{2}\right), y^{2}, \xi\right)
\end{aligned}
$$

if we set

$$
\begin{aligned}
& p_{3}^{\prime \prime}\left(x, \tilde{\zeta}^{2}, \tilde{y}^{2}, \xi\right) \\
& =O_{s}-\iint e^{i \psi_{2}\left(y^{2}, \xi^{2}, x^{2}, \xi\right)} \widetilde{r}_{2}\left(y^{1}, \xi^{2}, y^{2}, \zeta^{2} ; \xi\right) d x^{2} d \xi^{2} \\
& \cdot O_{s}-\iint e^{i \psi_{1}\left(y^{1}, \xi^{1}, x^{1}, \zeta^{2}\right)} \tilde{r}_{1}\left(x, \xi^{1}, y^{1}, \zeta^{1} ; \xi\right) d x^{1} d \xi^{1} \\
& \cdot p_{3}\left(x, \tilde{\xi}^{2}, \tilde{x}^{2}, \xi\right) .
\end{aligned}
$$

Here $\psi_{j}(y, \xi, x, \zeta)=\Phi_{j}(y, \xi)-x \cdot \xi+\phi_{j+1}(x, \zeta)-\Phi_{j+1}(y, \zeta)\left(j=1, \cdots, \nu, \Phi_{1}=\phi_{1}\right.$, $\left.\Phi_{j}=\phi_{1} \# \cdots \# \phi_{j}\right)$.

Now, we consider the case for a general $\nu$. Then, repeating the above method we can prove

$$
\begin{align*}
& q_{\nu+1}(x, \xi) e^{i \Phi_{\nu+1}(x, \xi)} \\
& =O_{s}-\int \cdots \int \exp \left(i\left(\sum_{j=1}^{\nu}\left(y^{j-1}-y^{j}\right) \cdot \eta^{j}+\Phi_{\nu+1}\left(y^{\nu}, \xi\right)\right)\right)  \tag{2.22}\\
& \quad \times p_{\nu+1}^{\prime}\left(x, \tilde{\eta}^{\nu}, \tilde{y}^{\nu}, \xi\right) d \tilde{y}^{\nu} d \tilde{\eta}^{\nu}
\end{align*}
$$

for

$$
\begin{align*}
& p_{\nu+1}^{\prime}\left(x, \tilde{\eta}^{\nu}, \tilde{y}^{\nu}, \xi\right)=p_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{y}^{\nu}, \xi\right)  \tag{2.23}\\
& \quad \text { with } \quad \zeta^{j}=\tilde{\nabla}_{x} \Phi_{j}^{-1}\left(y^{j-1}, y^{j} ; \eta^{j}\right) \quad\left(j=1, \cdots, \nu, \quad \Phi_{1}=\phi_{1}\right) .
\end{align*}
$$

Here, $p_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{y}^{\nu}, \xi\right)$ is defined by

$$
\begin{align*}
& p_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{y}^{\nu}, \xi\right) \\
& =O_{s}-\iint e^{i \psi_{\nu}\left(\nu^{\nu}, \xi^{\nu}, x^{\nu}, \xi\right)} \tilde{r}_{\nu}\left(y^{\nu-1}, \xi^{\nu}, y^{\nu}, \zeta^{\nu} ; \xi\right) d x^{\nu} d \xi^{\nu} \\
& \quad \cdot O_{s}-\iint e^{i \psi_{\nu-1}\left(y^{\nu-1}, \xi^{\nu-1}, x^{\nu-1}, \zeta^{\nu}\right)}  \tag{2.24}\\
& \quad \times \widetilde{r}_{\nu-1}\left(y^{\nu-2}, \xi^{\nu-1}, y^{\nu-1}, \zeta^{\nu-1} ; \xi\right) d x^{\nu-1} d \xi^{\nu-1}
\end{align*}
$$

$$
\begin{gathered}
\cdot O_{s}-\iint e^{i \psi_{1}\left(y^{1}, \xi^{1}, x^{1}, \zeta^{2}\right)} \tilde{r}_{1}\left(x, \xi^{1}, y^{1}, \zeta^{1} ; \xi\right) d x^{1} d \xi^{1} \\
\cdot p_{v+1}\left(x, \tilde{\xi}^{v}, \tilde{x}^{\nu}, \xi\right) \\
=\prod_{j=1}^{v} O_{s}-\iint e^{i \psi_{j}\left(y^{j}, \xi^{j}, \tilde{x}^{j}, \xi^{j+1}\right)} \\
\times \tilde{r}_{j}\left(y^{j-1}, \xi^{j}, y^{j}, \zeta^{j} ; \xi\right) d x^{j} d \xi^{j} p_{v+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) \\
\left(y^{0}=x, \quad \zeta^{v+1}=\xi\right) .
\end{gathered}
$$

Since (2.22) is nothing but (2.15), it remains to prove (2.6) for $p_{\nu+1}^{\prime}$ and (2.16). From (2.19), (2.5) for $p_{v+1}$ and (2.24) we can prove

$$
\begin{align*}
& \left|\zeta^{j}-\xi\right| \leqq\langle\xi\rangle / 4 \quad j=1, \cdots, \nu  \tag{2.25}\\
& \text { on the support of the integrand of } p_{\nu+1}^{\prime \prime} .
\end{align*}
$$

This property and (2.3)-a) in [17] implies (2.16). For the proof of (2.6) for $p_{\nu+1}^{\prime}$ we set for a while

$$
\begin{aligned}
& \tilde{p}_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{\zeta}^{\prime \nu}, \tilde{y}^{\nu}, \tilde{y}^{\prime \nu}, \xi\right)=\prod_{j=1}^{\nu} O_{s}-\iint e^{i \psi_{j}\left(y^{j}, \xi^{j}, x^{j}, \xi^{j+1}\right)} \\
& \quad \times \widetilde{r}_{j}\left(y^{\prime j-1}, \xi^{j}, y^{\prime j}, \zeta^{\prime j} ; \xi\right) d x^{j} d \xi^{j} p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) \quad\left(y^{\prime 0}=x\right)
\end{aligned}
$$

and we estimate this under

$$
\begin{equation*}
\left|\zeta^{j}-\xi\right| \leqq\langle\xi\rangle / 4, \quad\left|\zeta^{\prime j}-\xi\right| \leqq\langle\xi\rangle / 4 \tag{2.26}
\end{equation*}
$$

noting (2.25). Then, applying Lemma 2.3 to each oscillatory integral

$$
\begin{gathered}
O_{s}-\iint e^{i \psi_{j}\left(y^{j}, \xi^{j}, x^{j}, \zeta^{j+1}\right)} \widetilde{r}_{j}\left(y^{\prime j-1}, \xi^{j}, y^{\prime j}, \zeta^{\prime j}\right) \\
\times p_{v+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d x^{j} d \xi^{j}
\end{gathered}
$$

regarding $y^{\prime j-1}, y^{\prime j}, \zeta^{\prime j}, \tilde{x}^{\nu, j}$ and $\tilde{\xi}^{\nu, j}$ as parameters, we get for $\zeta^{j}$ and $\zeta^{\prime j}$ satisfying (2.26)

$$
\begin{aligned}
& \left|\partial_{\xi}^{\alpha} \partial_{\tilde{\zeta}^{\nu}}^{\tilde{a}^{\nu}} \partial_{\tilde{\zeta}^{\nu}}^{\tilde{a}^{\prime \nu}} \partial_{x}^{\beta} \partial_{\tilde{y}_{\nu}^{\sim}}^{\tilde{\beta}_{\nu}^{\nu}} \partial_{\tilde{\tilde{\beta}}^{\prime \prime} \nu}^{\nu \nu} \tilde{p}_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{\zeta}^{\prime \nu}, \tilde{y}^{\nu}, \tilde{y}^{\prime \nu}, \xi\right)\right| \\
& \left.\leqq C_{o} A_{1}^{\nu} M_{1}^{-(|\alpha|+\mid} \tilde{\alpha}^{\nu} \nu_{|+|} \tilde{\alpha}^{\prime} \nu_{|+|} \mid \beta_{|+|} \tilde{\beta}^{\nu_{|+|}} \tilde{\beta}^{\prime} \nu_{1}\right) \\
& \times\left(\alpha!\widetilde{\alpha}^{\nu}!\widetilde{\alpha}^{\prime \nu}!\beta!\tilde{\beta}^{\nu}!\tilde{\beta}^{\prime \nu}!\right)^{\kappa} \\
& \times\langle\xi\rangle^{m-\left|\alpha_{1-1}\right|} \tilde{\nu_{\nu}}|-| \tilde{\alpha} \tilde{\nu}_{1}
\end{aligned}
$$

with constants $A_{1}$ and $M_{1}$ independent of $\nu$. This implies (2.6) for $p_{\nu+1}^{\prime}$ since we have (2.23) and $p_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{y}^{\nu}, \xi\right)=\widetilde{p}_{\nu+1}^{\prime \prime}\left(x, \tilde{\zeta}^{\nu}, \tilde{\zeta}^{\nu}, \tilde{y}^{\nu}, \tilde{y}^{\nu}, \xi\right)$.
(II) Take the change of variables

$$
\left\{\begin{array}{l}
x^{j}=y^{j}+x,  \tag{2.27}\\
\xi^{j}=\eta^{j}+\tilde{\nabla}_{x} \Phi_{\nu+1}\left(x, x^{\nu} ; \xi\right)
\end{array}\right.
$$

in the integral of the right hand side of (2.15). Then, since we have

$$
\begin{aligned}
& \sum_{j=1}^{\nu}\left(x^{j-1}-x^{j}\right) \cdot \xi^{j}+\Phi_{\nu+1}\left(x^{\nu}, \xi\right)-\Phi_{\nu+1}(x, \xi) \\
& \quad=\sum_{j=1}^{\nu}\left(x^{j-1}-x^{j}\right) \cdot \xi^{j}-\left(x-x^{\nu}\right) \cdot \tilde{\nabla}_{x} \Phi_{\nu+1}\left(x, x^{\nu}, \xi\right) \\
& \quad=\sum_{j=1}^{\nu}\left(x^{j-1}-x^{j}\right) \cdot\left(\xi^{j}-\tilde{\nabla}_{x} \Phi_{\nu+1}\left(x, x^{\nu} ; \xi\right)\right) \\
& \quad=-\sum_{j=1}^{\nu}\left(y^{j}-y^{j-1}\right) \cdot \eta^{j} \quad\left(x^{0}=x, y^{0}=0\right),
\end{aligned}
$$

$q_{\nu+1}(x, \xi)$ in (2.15) is written as

$$
\begin{align*}
q_{\nu+1}(x, \xi)= & O_{s}-\int \cdots \int \exp \left(-i \sum_{j=1}^{\nu}\left(y^{j}-y^{j-1}\right) \cdot \eta^{j}\right)  \tag{2.28}\\
& \times \tilde{p}_{v+1}\left(x, \tilde{\eta}^{\nu}, \tilde{y}^{\nu}, \xi\right) d \tilde{y}^{\nu} d \tilde{\eta}^{\nu},
\end{align*}
$$

where $\tilde{p}_{v+1}\left(x, \tilde{\eta}^{\nu}, \tilde{y}^{\nu}, \xi\right)$ is defined from $p_{v+1}^{\prime}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ by the change of variables (2.27). Since $p_{\nu+1}^{\prime}$ satisfies (2.6) (with $C_{o}$ replaced by $C_{o} A^{\nu}$ ) and (2.16) we may use only the step (II) in the proof of Proposition 5.1 in [18] if we consider (2.28) instead of (5.2) in [18]. In fact, if $\tau^{0}$ in (A-1) is small enough, we have $c^{-1}\langle\xi\rangle \leqq\left\langle\xi+\eta^{j}\right\rangle \leqq c^{\prime}\langle\xi\rangle$ on supp $\tilde{p}_{\nu+1}$ (with $c^{\prime}$ independent of $\nu$ ) on account of (2.16) above and (2.3)-a) in [17]. Thus the proof is completed.
Q.E.D.

The proof of Lemma 2.2 is carried out by the same way as in that of Lemma 2.1 if we note $c^{-1}\langle\xi\rangle \leqq\left\langle\xi^{k}\right\rangle \leqq c\langle\xi\rangle$ on supp $p_{\nu+1}$.

In the rest of this section we shall give another fundamental lemma by means of Lemma 2.2. Let $\left\{X_{v}^{j}, \Xi_{v}^{j}\right\}_{v=1}^{j}(x, \xi)$ denote the solution of

$$
\left\{\begin{array}{l}
x^{j}=\nabla_{\xi} \phi_{j}\left(x^{j-1}, \xi^{j}\right), \\
\xi^{j}=\nabla_{x} \phi_{j+1}\left(x^{j}, \xi^{j+1}\right), \quad j=1, \cdots, \nu \quad\left(x^{0}=x, \xi^{\nu+1}=\xi\right) .
\end{array}\right.
$$

We remark that $\left\{X_{\nu}^{j}, \Xi_{\nu}^{j} /\langle\xi\rangle\right\}_{j, \nu}$ are bounded in $S_{G(k)}^{0}$ (see Proposition 2.4 in [18] and its proof). For a $\delta>0$ and $k=1, \cdots, \nu$ we set

$$
\begin{aligned}
\chi_{\delta, k} & \equiv \chi_{\delta, k}\left(x^{k}, \xi^{k} ; x, \xi\right) \\
& =\chi\left(\left(X_{\nu}^{k}(x, \xi)-x^{k}\right) / \delta\right) \chi\left(\left(\Xi_{\nu}^{k}(x, \xi)-\xi^{k}\right) /(\delta\langle\xi\rangle)\right)
\end{aligned}
$$

for $\chi \in \gamma^{(k)}\left(R_{\xi}^{n}\right)$ satisfying (1.10). Set

$$
\begin{align*}
& p_{\nu+1, \delta}^{o}\left(x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)  \tag{2.29}\\
& \quad=\prod_{j=1}^{v} \chi_{\delta, j}\left(x^{j}, \xi^{j} ; x, \xi\right) p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)
\end{align*}
$$

and set for $k=1, \cdots, \nu$

$$
\begin{equation*}
p_{\nu+1, \delta}^{k}\left(x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)=\chi_{\delta, 1} \cdots \chi_{\delta, k-1}\left(1-\chi_{\delta, k}\right) p_{\nu+1} . \tag{2.30}
\end{equation*}
$$

Then, we get $p_{\nu+1}=p_{\nu+1, \delta}^{o}+\sum_{k=1}^{\nu} p_{\nu+1, \delta}^{k}$.
Lemma 2.4. Let $p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ satisfy (2.5) and (2.6). Let $r_{\nu+1, \delta}(x, \xi)$ denote the symbol defined by (2.2) with $p_{\nu+1}$ replaced by $p_{\nu+1}-p_{\nu+1, \delta}^{o}$. Then, for any $\delta$ there exist positive constants $\varepsilon, A_{1}$ and $M_{1}$ independent of $\nu$ such that for any $\alpha$ and $\beta$ we have

$$
\begin{equation*}
\left.\left|\partial_{\xi}^{\alpha} \partial_{x}^{\beta} r_{\nu+1, \delta}(x, \xi)\right| \leqq C_{o} A_{1}^{\nu} M_{1}^{-(|\alpha|+|\beta|)} \alpha!\beta!\right)^{\kappa}\langle\xi\rangle^{m-|\alpha|} \exp \left(-\varepsilon\langle\xi\rangle^{1 / \kappa}\right) . \tag{2.31}
\end{equation*}
$$

Remark. As in the remark of Lemma 1.3 the estimate (2.31) holds even if we replace $\Phi_{\nu+1}(x, \xi)$ by $x \cdot \xi$ in (2.2).

Proof. Let $q_{\nu+1, \delta, k}\left(x, \xi^{k}, x^{k}, \xi\right)$ denote a symbol defined by (2.9) with $p_{\nu+1}$ replaced by $p_{\nu+1, \delta}^{k}$ and set

$$
\begin{align*}
q_{\nu+1, \delta}^{k}(x, \xi)= & O_{s}-\iint e^{i\left(\Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)-\Phi_{\nu+1}(x, \xi)\right)}  \tag{2.32}\\
& \times q_{\nu+1, \delta, k}\left(x, \xi^{k}, x^{k}, \xi\right) d x^{k} d \xi^{k}
\end{align*}
$$

Then we have $r_{\nu+1, \delta}(x, \xi)=\sum_{j=1}^{\nu+1} q_{\nu+1, \delta}^{k}(x, \xi)$. So, it suffices to estimate each $q_{v+1}^{k}(x, \xi)$. Since $\left(X_{v}^{k}, \Xi_{v}^{k}\right)(x, \xi)$ is the solution of the equation

$$
\left\{\begin{array}{l}
x^{k}=\nabla_{\xi}\left(\phi_{1} \# \cdots \# \phi_{k}\right)\left(x, \xi^{k}\right), \\
\xi^{k}=\nabla_{x}\left(\phi_{k+1} \# \cdots \# \phi_{\nu+1}\right)\left(x^{k}, \xi\right),
\end{array}\right.
$$

it follows from (2.8), (2.30) and Remark of Lemma 2.2 that

$$
\begin{array}{r}
\langle\xi\rangle\left|\nabla_{\xi^{k}} \Phi_{\nu+1 k}\left(x, \xi^{k}, x^{k}, \xi\right)\right|+\left|\nabla_{x^{k}} \Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)\right|  \tag{2.33}\\
\geqq c_{1}\langle\xi\rangle \quad \text { on } \operatorname{supp} q_{\nu+1, \delta, k}
\end{array}
$$

for a constant $c_{1}$ determined by $\delta$ and $\tau^{0}$ in (A-1). Set $L=-i\left(\langle\xi\rangle^{2}\left|\nabla_{\xi^{k}} \Phi_{\nu+1, k}\right|^{2}\right.$ $\left.+\left|\nabla_{x^{k}} \Phi_{\nu+1, k}\right|^{2}\right)^{-1}\left(\langle\xi\rangle^{2} \nabla_{\xi}{ }^{k} \Phi_{\nu+1, k} \cdot \nabla_{\xi^{k}}+\nabla_{x^{k}} \Phi_{\nu+1 k} \cdot \nabla_{x^{k}}\right)$. Then, we get $L[\exp (i$ $\left.\left.\times\left(\Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)-\Phi_{\nu+1}(x, \xi)\right)\right)\right]=\exp \left(i\left(\Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)-\Phi_{\nu+1}(x, \xi)\right)\right)$. Now, we integrate (2.32) by parts. Then, we have

$$
q_{\nu+1, \delta}^{k}(\xi, x)=O_{s}-\iint e^{i\left(\Phi_{\nu+1, k}-\Phi_{\nu+1}\right)}\left(L^{t}\right)^{N} q_{\nu+1, \delta, k} d x^{k} d \xi^{k}
$$

where $L^{t}$ is the transposed operator of $L$. It follows from (2.10) and (2.33) that there exist constants $A_{2}$ and $M_{2}$ independent of $\nu$ such that we have for any $N$

$$
\begin{aligned}
&\left|\partial_{\xi}^{\alpha} \partial_{\xi^{\prime}}^{\alpha^{\prime}} \partial_{x}^{\beta} \partial_{x^{\prime}}^{\beta^{\prime}}\left(\left(L^{t}\right)^{N} q_{v+1, \delta, k}\right)\right| \leqq\left(C_{o} A_{2}^{\nu-1}\right) M_{2}^{-\left(\left|\alpha^{\alpha}\right|+\left|\alpha^{\prime}\right|+\left|\beta^{\prime}\right|+\left|\beta^{\prime}\right| \mid-N\right.} \\
& \times\left(\alpha!\alpha^{\prime}!\beta!\beta^{\prime}!N!\right)^{\kappa}\langle\xi\rangle^{m-\left|\alpha^{\alpha}\right|-\left|\alpha^{\prime}\right|-N} \quad\left(x^{\prime}=x^{k}, \xi^{\prime}=\xi^{k}\right) .
\end{aligned}
$$

Now, we apply Lemma 2.3 for the case $\zeta=\xi$ by setting $p_{v+1}=\left(L^{t}\right)^{N} q_{v+1, \delta, k}$ and $C_{o}$ and $m$ in (2.6) as $C_{o} A_{2}^{\nu-1} N!^{\kappa}$ and $m-N$, respectively. Then, for any $N$ the estimate

$$
\begin{equation*}
\left|\partial_{\xi}^{\alpha} \partial_{x}^{\beta} q_{v+1, \delta}^{k}\right| \leqq C_{o} A_{3}^{v} M_{3}^{-(|\alpha|+|\beta|+N)}(\alpha!\beta!N!)^{\kappa}\langle\xi\rangle^{m-|\alpha|-N} \tag{2.34}
\end{equation*}
$$

holds with constants $A_{3}$ and $M_{3}$ independent of $N$ and $\nu$. Consequently, by means of the remark after Definition 1.1 we obtain (2.31) from (2.34). Q.E.D.

We end this section by the following remark: In the following sections we may assume that there exists a constant $M_{0}$ independent of $\nu$ such that for a multiple symbol $p_{\nu+1}\left(x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ satisfying (2.5) and (2.6) we have $\left|x-x^{\nu}\right| \leqq M_{o}$ on $\operatorname{supp} p_{v+1}$. Indeed, it follows from (1.5) in [11] that $\left|x-X_{\nu}^{\nu}(x, \xi)\right| \leqq c_{0}$ for a constant $c_{o}$ independent of $\nu$ and the symbol $q_{\nu+1}(x, \xi)$ defined by (2.2) satisfies (2.31) if $p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ vanishes when $\left|x-x^{\nu}\right| \leqq M_{0}$. This result also follows from Lemma 2.4.

## 3. Integration by parts with respect to time variables

Let $\phi(t, s ; x, \xi)$ be a solution of an eiconal equation

$$
\left\{\begin{array}{l}
\partial_{t} \phi=\lambda\left(t, x, \nabla_{x} \phi\right),  \tag{3.1}\\
\phi_{\mid t=s}=x \cdot \xi
\end{array}\right.
$$

Proposition 3.1. Assume that $\lambda(t, x, \xi)$ is a real symbol in $G^{(\kappa)}([0, T] ;$ $\left.S_{G(k)}^{1}\right)$. Then there exists a $T_{o}>0$ such that the solution $\phi(t, s ; x, \xi)$ of (3.1) exists uniquely in $\left\{(t, s) ; 0 \leqq t, s \leqq T_{o}\right\} \equiv\left[0, T_{o}\right]^{2}$ and belongs to $\mathscr{P}_{G(\kappa)}\left(\tilde{c}_{o}|t-s|\right)$ for a constant $\tilde{c}_{o}$ independent of $t$ and $s$. Furthermore, there exist constants $C$, $M_{1}$ and $M_{2}$ such that $\phi(t, s ; x, \xi)$ satisfies for any $\alpha, \beta, \gamma$ and $\gamma^{\prime}$

$$
\begin{equation*}
\left|\partial_{t}^{\gamma} \partial_{s}^{\gamma^{\prime}} \partial_{x}^{\beta} \partial_{\xi}^{\alpha} \phi\right| \leqq C M_{1}^{-(|\alpha|+|\beta|)} M_{2}^{-\left(\gamma+\gamma^{\prime}\right)}\left(\alpha!\beta!\gamma!\gamma^{\prime}!\right)^{\kappa}\langle\xi\rangle^{1-|\alpha|} \tag{3.2}
\end{equation*}
$$

for $t, s \in\left[0, T_{o}\right]$.
Proof. Assertions except the last one are the same as those of Proposition 3.1 in [18]. Since it follows that $\partial_{t} \phi=\lambda\left(t, x, \nabla_{x} \phi\right)$ and $\partial_{s} \phi=-\lambda\left(s, \nabla_{\xi} \phi, \xi\right)$ we obtain (3.2) by the inductive method with respect to $\gamma$ and $\gamma^{\prime}$. Q.E.D.

Let $\left\{\lambda_{j}(t, x, \xi)\right\}_{j=1}^{\infty}$ be a bounded set of real symbols in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$ and let $\psi_{j}(t, s ; x, \xi)$ denote the solution of the eiconal equation (3.1) with $\lambda=\lambda_{j}$. Let $\nu$ and $\mu$ be non-negative integers such that $\nu \geqq 3$ and $0 \leqq \mu \leqq \nu-2$. For a fixed positive $t_{0} \leqq T_{o}$ we set

$$
\begin{equation*}
\Delta_{\mu}=\left\{\tilde{t}^{\mu}=\left(t_{1}, \cdots, t_{\mu}\right) \in R^{\mu} ; 0 \leqq t_{\mu} \leqq t_{\mu_{-1}} \leqq \cdots \leqq t_{1} \leqq t_{0}\right\} \tag{3.3}
\end{equation*}
$$

when $\mu \geqq 1$. Let $\Sigma_{\mu}$ be a subset of $\{2, \cdots, \nu\}$ and denote it as

$$
\begin{equation*}
\Sigma_{\mu}=\left\{j_{1}, \cdots, j_{\mu_{+1}}\right\} \tag{3.4}
\end{equation*}
$$

with $2 \leqq j_{1}<j_{2}<\cdots<j_{\mu_{+1}} \leqq \nu$. For convenience we use, in the following, the notation $\left(t_{0}, \tilde{t}^{\mu}\right)=t_{0}$ when $\mu=0$. For a $\tilde{t}^{\mu} \in \Delta_{\mu}$ (when $\mu \geqq 1$ ) and $\Sigma_{\mu}$ we define a set of phase functions $\left\{\phi_{j}\left(t_{0}, \tilde{t}^{\mu} ; x, \xi\right)\right\}_{j=1}^{\nu+1} \equiv\left\{\phi_{j}\left(t_{0}, \tilde{t}^{\mu}\right)\right\}_{j+1}^{\nu-1}$ by

$$
\left\{\begin{array}{lll}
\phi_{j}\left(t_{0}, \tilde{t}^{\mu} ; x, \xi\right)=x \cdot \xi & \text { for } j \notin \Sigma_{\mu} \\
\phi_{j_{k}}\left(t_{0}, \tilde{t}^{\mu} ; x, \xi\right)=\psi_{j_{k}}\left(t_{k-1}, t_{k} ; x, \xi\right) & \text { for } & j_{k} \in \Sigma_{\mu} \\
\left(k=1, \cdots, \mu+1, t_{\mu_{+1}}=0\right) &
\end{array}\right.
$$

Then it follows from Proposition 3.1 that the set $\left\{\phi_{j}\left(t_{0}, \tilde{t}^{\mu}\right)\right\}_{j=1}^{\nu+1}$ satisfies the assumption (A-1) in § 1 if $T_{o}$ is small enough. Define $\Psi_{\nu+1, \Sigma_{\mu}}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ by the formula (1.9) for the set $\left\{\phi_{j}\left(t_{0}, \tilde{t}^{\mu}\right)\right\}_{j=1}^{\nu+1}$. In the following we shorten $\Psi_{\nu+1, \Sigma_{\mu}}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ to $\Psi_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu}\right)$ or $\Psi_{\nu+1}^{\mu}$. Then, we have

$$
\begin{align*}
\partial_{t_{k}} \Psi_{\nu+1}^{\mu}= & \left(\partial_{s} \psi_{j_{k}}\right)\left(t_{k-1}, t_{k} ; x^{j_{k}-1}, \xi^{j_{k}}\right) \\
& +\left(\partial_{t} \psi_{j_{k+1}}\right)\left(t_{k}, t_{k+1} ; x^{j_{k+1}-1}, \xi^{j_{k+1}}\right) \\
=- & -\lambda_{j_{k}}\left(t_{k}, \nabla_{\xi} \psi_{j_{k}}\left(t_{k-1}, t_{k} ; x^{j_{k}-1}, \xi^{j_{k}}\right), \xi^{j_{k}}\right)  \tag{3.5}\\
& +\lambda_{j_{k+1}}\left(t_{k}, x^{j_{k+1}-1}, \nabla_{x} \psi_{j_{k+1}}\left(t_{k}, t_{k+1} ; x^{j_{k+1}-1}, \xi^{j_{k+1}}\right)\right) \\
& \left(k=1, \cdots, \mu, t_{\mu_{+1}}=0\right) .
\end{align*}
$$

We note that (3.5) depends only on (3+4n) variables ( $t_{k-1}, t_{k}, t_{k+1}, x^{j_{k}-1}, \xi^{j_{k}}$, $x^{j_{k+1}}, \xi^{j_{k+1}}$ ), which is the key point of our discussions. Let $\mu \geqq 1$. For $k$ $k \in\{1, \cdots, \mu+1\}$ and $\Sigma_{\mu}$ in (3.4) we denote $\Sigma_{\mu, k}=\left\{j_{1}, \cdots, j_{k-1}, j_{k+1}, \cdots, j_{\mu_{+1}}\right\}$. Noting $\psi_{j_{k} \mid t=s}=x \cdot \xi$, we have

$$
\Psi_{\nu+1, \Sigma_{\mu, k}}^{\mu-1}\left(t_{0}, \tilde{s}^{\mu-1} ; x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)=\Psi_{\nu+1, \Sigma_{\mu}}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)_{t_{k}=t_{k-1}}
$$

by setting $\tilde{s}^{\mu-1}=\left(t_{1}, \cdots, t_{k-1}, t_{k+1}, \cdots, t_{\mu}\right)$, where $\Psi_{\nu+1, \Sigma_{\mu} \mid t_{k}=t_{k-1}}^{\mu}$ for $k=\mu+1$
 $\Psi_{\nu+1}^{\mu-1, k}\left(t_{0}, \tilde{s}^{\mu-1}\right)$ or $\Psi_{\nu+1}^{\mu-1, k}$ in what follows. From (3.5) it is easy to see

$$
\begin{align*}
\partial_{s_{k-1}} \Psi_{\nu+1}^{\mu-1, k}\left(t_{0}, \tilde{s}^{\mu-1}\right)= & \left(\partial_{s} \psi_{j_{k-1}}\right)\left(s_{k-2}, s_{k-1} ; x^{j_{k-1}-1}, \xi^{j_{k-1}}\right) \\
& +\left(\partial_{i} \psi_{j_{k+1}}\right)\left(s_{k-1}, s_{k} ; x^{j_{k+1}-1}, \xi^{j_{k+1}}\right),  \tag{3.6}\\
& \left(k=2, \cdots, \mu, s_{0}=t_{0}, s_{\mu}=0\right),
\end{align*}
$$

and

$$
\begin{gather*}
\partial_{s_{j}} \Psi_{\nu+1}^{\mu-1, k}\left(t_{0}, \tilde{s}^{\mu-1}\right)= \begin{cases}\partial_{t_{j}} \Psi_{\nu+1}\left(t_{0}, \tilde{t}^{\mu}\right) & \text { for } j<k-1, \\
\partial_{t_{j+1}} \Psi_{\nu+1}\left(t_{0}, \tilde{t}^{\mu}\right) & \text { for } j \geqq k,\end{cases}  \tag{3.7}\\
(k=1, \cdots, \mu+1),
\end{gather*}
$$

with

$$
\left(t_{0}, \tilde{t}^{\mu}\right)= \begin{cases}\left(s_{0}, s_{1}, \cdots, s_{k-2}, s_{k-1}, s_{k-1}, s_{k}, \cdots, s_{\mu-1}\right),\left(s_{0}=t_{0}\right) \\ & \text { for } k \in\{1, \cdots, \mu\} \\ \left(t_{0}, s_{1}, s_{2}, \cdots, s_{\mu_{-1}}, 0\right) & \text { for } k=\mu+1\end{cases}
$$

Let $p_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\nu} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ be a multiple symbol with parameters $t_{0}$ and $\tilde{t}^{\mu} \in \Delta_{\mu}$ satisfying (2.5). We say that $p_{\nu+1}^{\mu}$ satisfies the condition $\mathscr{I}(C, h, M, m)$ for an integer $h \geqq 0$, a real $m$ and constants $C$ and $M$ if we have for any $\alpha, \beta$, $\tilde{\alpha}^{\nu}, \tilde{\beta}^{\nu}$ and $k$

$$
\begin{align*}
& \leqq C M^{-\left(\left|\tilde{\alpha}^{\nu}\right|+\left|\tilde{\beta}^{\nu}\right|_{|+|\alpha|+|\beta|+k)}\left(\alpha!\beta!\tilde{\alpha}^{\nu}!\tilde{\beta}^{\nu}!(k+h)!\right)^{\kappa}\langle\xi\rangle^{m-|\alpha|-\left|\tilde{\alpha}^{\nu}\right|_{1}} . ~ . ~ . ~ . ~\right.} \tag{3.8}
\end{align*}
$$

We consider the symbol $q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)$ defined by

$$
\begin{align*}
& q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)= \int_{\Delta_{\mu}} d \tilde{t}^{\mu} O_{s}-\iint^{i\left(\Psi_{v+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu}\right)-x \cdot \xi\right)} \\
& \times p_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{x}^{\nu} d \xi^{\nu}  \tag{3.9}\\
& \equiv O_{s}-\iint^{2}\left(\int_{\Delta_{\mu}} e^{i\left(\Psi_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu}\right)-x \cdot \xi\right)}\right. \\
&\left.\times p_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{t}^{\mu}\right) d \tilde{x}^{\nu} d \tilde{\xi}^{\nu}
\end{align*}
$$

when $\mu \geqq 1$ and

$$
\begin{equation*}
q_{\nu+1}^{0}\left(t_{0} ; x, \xi\right)=O_{s}-\iint e^{i\left(\phi_{\nu+1}^{0}\left(t_{0}\right)-x \cdot \xi\right)} p_{\nu+1}^{0}\left(t_{0} ; x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d \tilde{x}^{\nu} d \xi^{\nu} \tag{3.9}
\end{equation*}
$$

for $\mu=0$. Let $\left\{X_{v}^{j}, \Xi_{v}^{j}\right\}\left(t_{0}, \tilde{t}^{\mu} ; x, \xi\right)$ be a solution of (1.7) with $\left\{\phi_{j}\left(t_{0}, \tilde{t}^{\mu}\right)\right\}_{j=1}^{\nu+1}$. First, we assume that there exist positive constants $\varepsilon$ and $\delta$ such that

$$
\begin{align*}
& \operatorname{supp} p_{\nu+1}^{\mu} \cap\left(\bigcap_{k=1}^{\mu}\left\{\left|\partial_{t_{k}} \Psi_{\nu+1}^{r}\right| \leqq \varepsilon\langle\xi\rangle \mid 2\right\}\right)  \tag{3.10}\\
& \cap\left(\bigcap_{j=1}^{\nu}\left\{X_{\nu}^{j}-x^{j}\left|\leqq \delta,\left|\Xi_{\nu}^{j}-\xi^{j}\right| \leqq \delta\langle\xi\rangle\right\}\right)=\phi\right.
\end{align*}
$$

In (3.10) the second factor is the whole space if $\mu=0$. Let $\mu \geqq 1$ and let $p_{\nu+1}^{\mu-1, k}\left(t_{0}, \tilde{s}^{\mu-1} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ denote $p_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)_{t_{t_{k}=t_{k-1}}}$ for some $k \in$ $\{1,2, \cdots, \mu+1\}$ with $\tilde{s}^{\mu-1}=\left(t_{1}, \cdots, t_{k-1}, t_{k+1}, \cdots, t_{\mu}\right)$, where $p_{\nu+1 t_{k}=t_{k-1}}^{\mu}$ for $k=\mu+1$ means $p_{\nu+11^{t}{ }_{\mu}=0^{\circ}}^{\mu}$. Then $p_{\nu+1}^{\mu-1, k}$ satisfies the condition $\mathcal{I}(C, h, M, m)$ if $p_{\nu+1}^{\mu}$ so does.
By setting $t_{k}=t_{k-1}$ for a $k$ with $1 \leqq k \leqq \mu$ or $t_{\mu}=0$ for $k=\mu+1$ we can define a "child" $\left\{p_{\nu+1}^{\mu-1}, \Psi_{\nu+1}^{\mu-1}\right\} \equiv\left\{p_{\nu+1}^{\mu-1, k}, \Psi_{\nu+1}^{\mu-1, k}\right\}(k=1, \cdots, \mu+1)$ of $\left\{p_{\nu+1}^{\mu}, \Psi_{\nu+1}^{\mu}\right\}$, and moreover can define descendent sequences of $\left\{p_{\nu+1}^{\mu}, \Psi_{\nu+1}^{\mu}\right\}$ successively until $\mu=0$. We finally assume that (3.10) holds for $\left\{p_{\nu+1}^{\nu}, \Psi_{\nu+1}^{\mu}\right\}$ and all children of all its descendent sequences. In what follows we denote this assumption by $B(\varepsilon, \delta)$.

Lemma 3.2. Let $p_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ be a multiple symbol satisfying (2.5) and the condition $\mathcal{G}\left(C_{o}, 0, M, m\right)$. Assume that $\left\{p_{\nu+1}^{\mu}, \Psi_{\nu+1}^{\mu}\right\}$ satisfies $B(\varepsilon, \delta)$. Then there exist constants $C, \varepsilon^{\prime}>0, A_{1}$ and $M_{1}$ independent of $\nu, \mu$
such that for any $\alpha$ and $\beta$

$$
\begin{equation*}
\left|\partial_{\xi}^{\alpha} \partial_{x}^{\beta} q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)\right| \leqq C C_{o}\left(A_{1}^{\nu+\mu} / \mu!\right) M_{1}^{-\left(\left|\alpha_{1}+\right| \beta^{\beta}\right)}(\alpha!\beta!)^{\kappa} e^{-\varepsilon^{\prime}\langle\xi\rangle^{1 / k}} \tag{3.11}
\end{equation*}
$$

where $q_{\nu+1}^{\mu}$ is the symbol defined by (3.9) or (3.9)'. Furthermore, if we resolve the constant $C$ in (3.11) into $\tilde{C} C_{1}^{m}[m]!^{\kappa}$ when $m>0$, we can take constants $\widetilde{C}, C_{1}$, $A_{1}$ and $M_{1}$ independent also of $m$ and have

$$
(3.11)^{\prime} \quad\left|\partial_{\xi}^{\alpha} \partial_{x}^{\beta} q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)\right| \leqq \tilde{C} C_{o} C_{1}^{m}[m]!^{\kappa}\left(A_{1}^{\nu+\mu} / \mu!\right) M_{1}^{-\left(\left|\alpha^{\alpha}+|\beta|\right)\right.}(\alpha!\beta!)^{\kappa} e^{-\varepsilon^{\prime}\left\langle\langle \rangle^{1 / \kappa}\right.}
$$

For the proof of this lemma it suffices to show

$$
\begin{equation*}
\left|q_{\gamma+1}^{\mu}\left(t_{0} ; x, \xi\right)\right| \leqq C C_{0}\left(A_{1}^{\nu+\mu} / \mu!\right)\langle\xi\rangle^{m} e^{-\varepsilon^{\prime}\langle\xi\rangle^{1 / \kappa}} \tag{3.12}
\end{equation*}
$$

for constants $C$ and $A_{1}$ independent of $\nu, \mu$ and $m$. In fact, differentiating with respect to $x$ and $\xi$, we get in view of $\phi_{1}=\phi_{\nu+1}=x \cdot \xi$

$$
\begin{aligned}
& \partial_{\xi}^{\alpha} \partial_{x}^{\beta} q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)=\int_{\Delta_{\mu}} d \tilde{t}^{\mu} O_{s}-\iint e^{i\left(\Psi_{\nu+1}-x \cdot \xi\right)} \\
& \quad \times \sum_{\substack{\alpha^{\prime}+\alpha^{\prime \prime}=\alpha \\
\beta^{\prime}+\beta^{\prime \prime}=\beta}} \frac{\alpha!}{\alpha^{\prime}!\alpha^{\prime \prime}!} \frac{\beta!}{\beta^{\prime}!\beta^{\prime \prime}!}\left(i\left(\xi^{1}-\xi\right)\right)^{\beta^{\prime}} \partial_{x}^{\beta^{\prime \prime}}\left(\left(i\left(x^{\nu}-x\right)\right)^{\alpha^{\prime}} \partial_{\xi^{\prime \prime}}^{\alpha^{\prime \prime}} p_{\nu+1}^{\mu}\right) \tilde{x}^{\nu} d \tilde{\xi}^{\nu},
\end{aligned}
$$

when $\mu \geqq 1$. It follows from the last remark in Section 2 that we may assume $\left|x^{\nu}-x\right| \leqq M_{o}$ on $\operatorname{supp} p_{\nu+1}^{\mu}$. Noting that $\left|\xi^{j}-\xi\right| \leqq\langle\xi\rangle / 2$ on $\operatorname{supp} p_{\nu+1}^{\mu}$ and for $N=\left|\beta^{\prime}\right|+[\max (m, 0)]$ an estimate $\langle\xi\rangle^{N} \exp \left(-\varepsilon^{\prime}\langle\xi\rangle^{1 / \kappa}\right) \leqq M^{-N} N!^{\kappa} \times$ $\exp \left(-\varepsilon^{\prime}\langle\xi\rangle^{1 / k} / 2\right)$ holds for a sufficiently small $M>0$, we get the assertion by means of the Leibniz formula. Similarly, we get the assertion for the case of $\mu=0$.

For the proof of Lemma 3.2 we prepare the following lemma, which is the direct consequence of Lemma 2.4.

Lemma 3.3. Let $p_{\nu+1}^{\mu}$ satisfy the same condition as in Lemma 3.2. Assume that we have (3.10) and

$$
\begin{equation*}
\operatorname{supp} p_{\nu+1}^{\mu} \subset \bigcap_{k=1}^{\mu}\left\{\left|\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right| \leqq \varepsilon\langle\xi\rangle / 2\right\} \tag{3.13}
\end{equation*}
$$

when $\mu \geqq 1$. Then there exist constants $\varepsilon^{\prime}>0, A_{1}$ and $M_{1}$ independent of $\nu$ such that

$$
\begin{equation*}
\left|q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)\right| \leqq C_{o}\left(A_{1}^{\nu+\mu} / \mu!\right)\langle\xi\rangle^{m} e^{-\varepsilon^{\prime}\langle\xi\rangle^{1 / \kappa}} . \tag{3.14}
\end{equation*}
$$

Proof. Set

$$
\tilde{q}_{\nu+1}\left(t_{0}, \tilde{t}^{\mu} ; x, \xi\right)=O_{s}-\iint e^{i\left(\Psi_{\nu+1}^{\mu}-x \cdot \xi\right)} p_{\nu+1}^{\mu} d \tilde{x}^{\nu} d \hat{\xi}^{\nu}
$$

It follows from Lemma 2.4 and its remark that $\tilde{\boldsymbol{q}}_{\nu+1}$ satisfies the estimate (2.31)
with $\alpha=\beta=0$ uniformly with respect to $\tilde{t}^{\mu} \in \Delta_{\mu}$, since the term defined by (2.29) for $p_{\nu+1}^{\mu}$ vanishes from (3.10) and (3.13). Noting that the volume of $\Delta_{\mu}$ is equal to $t_{0}^{\mu} / \mu$ ! we get the desired estimate.
Q.E.D.

Let $\tilde{\chi}(t)$ be a function in $\gamma^{(\kappa)}\left(R_{t}^{1}\right)$ satisfying

$$
0 \leqq \tilde{\chi} \leqq 1, \quad \tilde{\chi}=1 \quad(|t| \leqq 1 / 4), \quad \tilde{\chi}=0 \quad(|t| \geqq 1 / 2) .
$$

We set

$$
\left\{\begin{array}{l}
\chi_{k}^{0}\left(t_{0}, \tilde{t}^{\mu} ; x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)=\tilde{\chi}\left(\partial_{t_{k}} \Psi_{\nu+1}^{\mu} /(\varepsilon\langle\xi\rangle)\right), \\
\chi_{k}^{1}=1-\chi_{k}^{0}, \\
k \in\{1, \cdots, \mu\}
\end{array}\right.
$$

Setting

$$
H_{\mu}=\left\{\tilde{h}^{\mu}=\left(h_{1}, \cdots, h_{\mu}\right) ; h_{j}=0,1\right\}
$$

we divide $p_{\nu+1}^{\mu}$ into $2^{\mu}$ terms:

$$
p_{\nu+1}^{\mu}=\sum_{\tilde{h}^{\mu} \in H_{\mu}} p_{\nu+1}^{\mu}, \tilde{h}^{\mu},
$$

where $p_{\nu+1,}^{\mu}, \tilde{h}^{\mu}=\prod_{k=1}^{\mu} \chi_{k}^{k_{k}} p_{\nu+1}^{\mu}$. In view of this division, for the proof of Lemma 3.2, that is, for the proof of (3.12), it suffices to show the following:

Lemma 3.4. Let $p_{\nu+1}^{\mu}\left(t_{0}, \tilde{t}^{\mu} ; x, \hat{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ be a multiple symbol satisfying (2.5) and $\mathcal{I}\left(C_{\mu}, h, M, m\right)$. Assume that $\left\{p_{v .1}^{\mu}, \Psi_{\nu+1}^{\mu}\right\}$ satisfies the condition $B(\varepsilon, \delta)$. We assume furthermore that

$$
\begin{cases}\text { for each } k \in\{1, \cdots, \mu\} & \text { it follows that either }  \tag{*}\\ \left|\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right| \leqq \varepsilon\langle\xi\rangle / 2 & \text { for all } \quad\left(t_{0}, \tilde{t}^{\mu}, x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) \text { on } \operatorname{supp} p_{\nu+1}^{\mu} \\ \text { or } & \\ \left|\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right| \geqq \varepsilon\langle\xi\rangle / 4 & \text { for all } \quad\left(t_{0}, \tilde{t}^{\mu}, x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) \text { on } \operatorname{supp} p_{\nu+1}^{\mu}\end{cases}
$$

when $\mu \geqq 1$. Then there exist constants $A, A_{1}$ and $M_{1}$ independent of $\nu$ and $\mu$ such that for any $N$ we have

$$
\begin{equation*}
\left|q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)\right| \leqq 5 C_{\mu} A^{\nu} M_{1}^{-N}\langle\xi\rangle^{m-N} \sum_{j=0}^{N \wedge \mu} \frac{(N-j+h)!^{\kappa}}{(\mu-j)!} A_{1}^{\mu-j} \tag{3.15}
\end{equation*}
$$

where $N \wedge \mu=\min (N, \mu)$.
Proof. We shall prove the lemma by induction on $\mu$. When $\mu=0$ the estimate (3.15) for any integer $N \geqq 0$ follows from Lemma 3.3. Suppose $\mu \geqq 1$ and that the lemma for any $h$ is valid until $\mu-1$. If $\left|\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right| \leqq \varepsilon\langle\xi\rangle / 2$ on $\operatorname{supp} p_{\nu+1}^{\mu}$ for all $k$ then the estimate (3.15) follows also from Lemma 3.3. Suppose that there exists a $k \in\{1, \cdots, \mu\}$ such that $\left|\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right| \geqq \varepsilon\langle\xi\rangle / 4$ on
$\operatorname{supp} p_{\nu+1}^{\nu}$. Note that

$$
\int_{\Delta_{\mu}} d \tilde{t}^{\mu}=\int_{\Delta_{\mu-1}^{\prime}} d t_{1} \cdots d t_{k-1} d t_{k+1} \cdots d t_{\mu} \int_{t_{k-1}}^{t_{k-1}} d t_{k}
$$

where $\Delta_{\mu-1}^{\prime}$ is defined by (3.3) with respect to $t_{1}, \cdots, t_{k-1}, t_{k+1}, \cdots, t_{\mu}$. Integrating by parts with respect to $t_{k}$ we have

$$
\begin{align*}
\int_{t_{k+1}}^{t_{k-1}} e^{i \Psi_{\nu+1}^{\mu}} p_{\nu+1}^{\mu} d t_{k}= & -\int_{t_{k+1}}^{t_{k-1}} e^{i \Phi_{\nu+1}^{\mu}} \partial_{t_{k}}\left(\left(i \partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right)^{-1} p_{\nu+1}^{\mu}\right) d t_{k} \\
& +e^{i \Psi_{\nu+1}^{\mu}, 1, k}\left[\left(i \partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right)^{-1} p_{\nu+1}^{\mu}\right]_{t_{k}=t_{k-1}}  \tag{3.16}\\
& -e^{i \Psi_{\nu+1}^{\mu}, 1, k}\left[\left(i \partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right)^{-1} p_{\nu+1}^{\mu}\right]_{t_{k}=t_{k+1}}
\end{align*}
$$

Hence we get

$$
q_{v+1}^{\mu}\left(t_{0} ; x, \xi\right)=q_{v+1}^{\mu} 1_{0}\left(t_{0} ; x, \xi\right)+q_{v+1,0}^{\mu-1,1}\left(t_{0} ; x, \xi\right)+q_{v+1,1}^{\mu-1,1}\left(t_{0} ; x, \xi\right),
$$

where the multiple symbol of each term corresponds to the one of (3.16). In order to apply the hypothesis of the induction, in view of (3.6) and (3.7) we divide the multiple symbol of $q_{\nu+1, j}^{\mu-1,1}\left(t_{0} ; x, \xi\right), j=0,1$, into two terms by multiplying the partition of unity $\left\{\tilde{\chi}\left(\partial_{s_{k^{\prime}}} \Psi_{\nu+1}^{\mu-1, k+j} /(\varepsilon\langle\xi\rangle)\right), 1-\chi\left(\partial_{s_{k^{\prime}}} \Psi_{\nu+1}^{\mu-1, k+j} /(\varepsilon\langle\xi\rangle)\right)\right\}$ with $k^{\prime}=k-1+j$. Then we get the division

$$
q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)=q_{\nu+1}^{\mu}\left(t_{0} ; x, \xi\right)+\sum_{j=1}^{4} \tilde{q}_{\nu+1, j}^{\mu-1,1}\left(t_{0} ; x, \xi\right)
$$

where each term satisfies the assumptions of the lemma. Repeating the same procedure as above for $q_{\nu+1}^{\mu}{ }_{1}$ again and moreover repeating $N$ times, we finally obtain

$$
\begin{equation*}
q_{\nu+1}^{\mu}=q_{\nu+1}^{\mu, N}+\sum_{K=1}^{N} \sum_{j=1}^{4} \tilde{q}_{\nu+1 ; j}^{\mu-1, K} \tag{3.17}
\end{equation*}
$$

Each term satisfies the assumptions of the lemma. More precisely, multiple symbols of $q_{\nu+1}^{\mu, N}$ and $\tilde{q}_{\nu+1}^{\nu-1, K}$ satisfy conditions $\mathscr{I}\left(C_{\mu} M_{2}^{-N}, h+N, M, m-N\right)$ and $\mathscr{I}\left(C_{\mu} M_{2}^{-K}, h+K-1, M, m-K\right)$, respectively, if we take another constant $M_{2}$ independent of $\nu, h$ and $K$. To prove this, taking another small $M$ if necessary, we may assume that $r\left(\theta_{1}, \theta_{2}, \theta_{3}, y^{1}, \eta^{1}, y^{2}, \eta^{2}\right) \equiv\left(\left(\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\right)\left(\theta_{1}, \theta_{2}, \theta_{3}, y^{1}, \eta^{1}, y^{2}, \eta^{2}\right)\right)^{-1}$ satisfies

$$
\left|\partial_{\tilde{\boldsymbol{\gamma}}^{3}}^{\tilde{\gamma}^{3}} \partial_{\tilde{y}_{2}^{2}}^{\tilde{\beta}^{2}} \partial_{\tilde{\eta}^{2}}^{\tilde{\alpha}^{2}} r\right| \leqq C(2 M)^{-\mid \tilde{\alpha}^{2}+\tilde{\beta}^{2}+\tilde{\gamma}^{3}}\left(\tilde{\alpha}^{2}!\tilde{\beta}^{2}!\tilde{\gamma}^{3}!\right)^{\kappa}\langle\xi\rangle^{-\left(1+\left|\tilde{\alpha}^{2}\right|\right)}
$$

for a constant $C$. Then, noting that $r$ depends only on $3+4 n$ variables, we get the desired properties for multiple symbols of $q_{\nu+1}^{\mu, N}$ and $q_{\nu+1}^{\mu-1, K}$.

Now, we use Lemma 2.1 by setting $C_{o}=C_{\mu} M_{2}^{-N}(h+N)!^{\kappa}$ and replacing $m$ by $m-N$. Then, from the fact that the multiple symbol of $q_{v+1}^{u, N}$ satisfies
$\mathcal{I}\left(C_{\mu} M_{2}^{-N}, h+N, M, m-N\right)$ we have

$$
\begin{equation*}
\left|q_{v+1}^{v_{i}^{N}}\right| \leqq C_{\mu} M_{2}^{-N} A^{\nu} A_{1}^{\mu}(h+N)!^{\kappa}\langle\xi\rangle^{m-N} / \mu! \tag{3.18}
\end{equation*}
$$

for a constant $A_{1}$. Similarly we have

$$
\begin{equation*}
\sum_{j=1}^{4}\left|\tilde{q}_{v+1, j}^{\mu-1, N}\right| \leqq 4 C_{\mu} M_{2}^{-N} A^{\nu} A_{1}^{\mu-1}(h+N-1)!^{\kappa}\langle\xi\rangle^{m-N} /(\mu-1)! \tag{3.19}
\end{equation*}
$$

For $\tilde{q}_{\nu+1}^{\mu-1, K}, K=1, \cdots, N-1$ we use the hypothesis of the induction with $N$ replaced by $N-K$. Then we have

$$
\begin{align*}
& \sum_{K=1}^{N-1} \sum_{j=1}^{4}\left|\tilde{q}_{\nu+1,1, j}^{\mu}\right| \leqq \sum_{K=1}^{N}\left(5 C_{\mu} M_{2}^{-K} A^{\nu} M_{1}^{-(N-K)}\langle\xi\rangle^{(m-K)-(N-K)}\right.  \tag{3.20}\\
& \quad \times \frac{(N-K) \wedge(\mu-1)}{\sum_{j=0}^{(\mu)}} \frac{(N-K-j+h+K-1)!^{K}}{(\mu-1-j)!} A_{1}^{\mu-1-j} \\
& \leqq 20 C_{\mu} A^{\nu} M_{1}^{-N} \sum_{K=1}^{N}\left(M_{1} / M_{2}\right)^{K}\langle\xi\rangle^{m-N} \sum_{j=0}^{(N \lambda)-1} \frac{(N-(j+1)+h)!}{(\mu-(j+1))!} A_{1}^{\mu-(j+1)}
\end{align*}
$$

Summing up (3.17)-(3.20) we get (3.15) if $M_{1}$ is sufficiently smaller than $M_{2}$. This concludes the proof.
Q.E.D.

Finally, we give a simple proposition for the argument of the next section.

## Proposition 3.5. Assume that

$$
\begin{equation*}
\left|\partial_{t_{k}} \Psi_{\nu+1}^{\mu}\left(t_{0}, \tilde{z}^{\mu}, x, \xi^{\nu}, \tilde{x}^{\nu}, \xi\right)\right| \leqq \varepsilon\langle\xi\rangle / 2 \tag{3.21}
\end{equation*}
$$

for a $k \in\{1, \cdots, \mu\}$ and that

$$
\begin{gather*}
\left|X_{\nu}^{j}\left(t_{0}, \tilde{t}^{\mu}, x, \xi\right)-x^{j}\right| \leqq \delta,\left|\Xi_{\nu}^{j}\left(t_{0}, \tilde{Z}_{\mu}, x, \xi\right)-\xi^{j}\right| \leq \delta\langle\xi\rangle  \tag{3.22}\\
\text { for any } j \in\{1, \cdots, \nu\} .
\end{gather*}
$$

If $\delta$ is sufficiently small, we have

$$
\begin{equation*}
\left|\lambda_{j_{k}}\left(t_{k}, Y^{k}, H^{k}\right)-\lambda_{j_{k+1}}\left(t_{k}, Y^{k}, H^{k}\right)\right| \leqq \varepsilon\langle\xi\rangle, \tag{3.23}
\end{equation*}
$$

where $\left\{Y_{\mu}^{k}, H_{\mu}^{k}\right\}$ is the solution of (1.7) with $\nu=\mu$ and with $\left\{\phi_{j}\right\}_{j=1}^{\nu+1}$ replaced by $\left\{\psi_{j}\right\}_{j=1}^{p_{j 1}}$.

Proof. The estimate (3.23) follows easily from (3.5) and (1.7) because we have $\left(X_{v}^{j}, \Xi_{j}^{j}\right)=\left(X_{\nu}^{j}, \Xi_{\nu}^{j}\right)=\left(Y_{\mu,}^{k}, H_{k}^{k}\right)$ for $j_{k} \leqq j<j_{k+1}, k \in\{0, \cdots, \mu+1\}$, where $j_{0}=1$ and $j_{\mu+2}=\nu+1$.
Q.E.D.

## 4. Proof of Theorem 1

Before the proof we state the definition of $\varepsilon$-admissible trajectories, following [4] and [5]. Let $\lambda_{j}(t, x, \xi), j \in\{1, \cdots, l\}$ be characteristic roots of $\mathcal{L}$,
given in Introduction. Namely, $\lambda_{j}(t, x, \xi)$ belong to $G^{(\kappa)}\left([0, T] ; S_{G(k)}^{1}\right)$ and satisfy $\lambda_{j}(t, x, \theta \xi)=\theta \lambda_{j}(t, x, \xi)$ for $\theta \geqq 1$ and $|\xi| \geqq 1$. We say that a curve $\{(t, x(t), \xi(t))\} \subset[0, T] \times T^{*}\left(R_{x}^{n}\right)$ is the bicharacteristic curve with respect to $\lambda_{j}$ through $(s, y, \eta)$ if $\{x(t), \xi(t)\}$ satisfies the equation

$$
\left\{\begin{array}{l}
d x / d t=-\nabla_{\xi} \lambda_{j}(t, x, \xi), \quad d \xi / d t=\nabla_{x} \lambda_{j}(t, x, \xi) \\
(x, \xi)_{\mid t=s}=(y, \eta)
\end{array}\right.
$$

We denote by $\chi_{j}(t, s)$ a transformation

$$
T^{*}\left(R^{n}\right) \backslash 0 \ni \rho \equiv(y, \eta) \rightarrow \chi_{j}(t, s) \rho=(x(t), \xi(t)) \in T^{*}\left(R^{n}\right) \backslash 0 .
$$

For an integer $\nu \geqq 0$, let $\Pi_{\nu+1}$ denote a set of ( $\nu+1$ )-repeated permutations $\left(j_{1}, j_{2}, \cdots, j_{\nu+1}\right)$ with $j_{k} \in\{1, \cdots, l\}$ and let $\Pi_{\nu+1}^{o}$ denote a subset of $\Pi_{\nu+1}$ whose elements $J_{\nu} \equiv\left(j_{1} \cdots, j_{\nu+1}\right)$ satisfy $j_{k} \neq j_{k+1}$ for any $k$. Let $t_{0}$ be a fixed point in ( $0, T]$ and let $\Delta_{\nu}^{o}$ denote the interior of $\Delta_{\nu}$ defined by (3.3). A continuous curve $\left\{(t, x(t), \xi(t)) ; t \in\left[0, t_{0}\right]\right\}$ is called a trajectory of step $\nu$, issuing from $\rho$, if for some $J_{\nu}=\left(j_{1}, \cdots, j_{\nu+1}\right) \in \Pi_{\nu+1}^{o}$ and some $\tilde{t}^{\nu} \in \Delta_{\nu}^{o}$ it is the bicharacteristic curve with respect to $\lambda_{j_{k}}$ when $t \in\left[t_{k}, t_{k-1}\right]\left(k=1, \cdots, \nu+1, t_{\nu+1}=0\right)$ and $(x(0), \xi(0))=\rho$. We often denote the trajectory by $C\left(J_{v}, \tilde{t}^{\nu}, \rho\right)$. A point

$$
\chi_{j_{1}}\left(t_{0}, t_{1}\right) \chi_{j_{2}}\left(t_{1}, t_{2}\right) \cdots \chi_{j_{v+1}}\left(t_{v}, 0\right) \rho
$$

is called the end point (at $t=t_{0}$ ) of the trajectory. For an $\varepsilon \geqq 0$ we say that the trajectory is $\varepsilon$-admissible if

$$
\left|\lambda_{j_{k}}\left(t_{k}, \rho_{k}\right)-\lambda_{j_{k+1}}\left(t_{k}, \rho_{k}\right)\right| \leqq \varepsilon\langle\eta\rangle, \quad k=1, \cdots, \nu
$$

where $\rho_{k} \equiv\left(x^{k}, \xi^{k}\right)=\chi_{j_{k+1}}\left(t_{k}, t_{k+1}\right) \cdots \chi_{j v+1}\left(t_{v}, 0\right) \rho$. We remark that the bicharacteristic curve is also a trajectory of step 0 and it is always 0 -admissible.

Since $\left\{\lambda_{j}\right\}$ is bounded in the symbol class $S^{1}$ and each $\lambda_{j}(t, x, \xi)$ is homogeneous for $|\xi| \geqq 1$, it is easy to see

Proposition 4.1. Let $\left\{\left(t, \rho\left(t ; \rho_{o}\right)\right) ; t \in\left[0, t_{0}\right]\right\}$ denote a trajectory $C\left(J_{\nu}, \tilde{t}^{\nu}, \rho_{o}\right)$ for $J_{\nu} \in \Pi_{\nu+1}^{o}, \tilde{t}^{\nu} \in \Delta_{\nu}^{o}$ and $\rho_{o} \in T^{*}\left(R^{n}\right) \backslash 0$. Then, there exists a positive constant $c>0$ independent of $J_{\nu}$ and $\tilde{t}^{\nu}$ such that for $\rho_{o} \equiv\left(x_{o}, \xi_{o}\right), \rho_{o}^{\prime} \equiv\left(x_{o}^{\prime}, \xi_{o}^{\prime}\right)$

$$
\left|\rho\left(t ; \rho_{o}\right)-\rho\left(t ; \rho_{o}^{\prime}\right)\right| \leqq e^{c t}\left|\rho_{o}-\rho_{o}^{\prime}\right|,
$$

if $\left|\xi_{0}\right|$ and $\left|\xi_{o}^{\prime}\right|$ are large enough. Here $\left|\rho-\rho^{\prime}\right|=\left|x-x^{\prime}\right|+\left|\xi /|\xi|-\xi^{\prime} /\left|\xi^{\prime}\right|\right|$ for $\rho \equiv(x, \xi), \rho^{\prime} \equiv\left(x^{\prime}, \xi^{\prime}\right) \in T^{*}\left(R^{n}\right) \backslash 0$.

By means of Proposition 4.1 we may replace the definition of $\Gamma\left(t_{0} ; V\right)$ in Introduction by the following

$$
\left\{\begin{array}{l}
\Gamma_{\mathrm{e}}\left(t_{0} ; V\right)=\text { the closure of } \bigcup_{\nu=0}^{\infty} \Gamma_{\mathrm{e}}^{\nu}\left(t_{0} ; V_{\mathrm{z}}\right)  \tag{4.1}\\
\Gamma\left(t_{0} ; V\right)=\bigcap_{\mathrm{e}>0} \Gamma_{\mathrm{e}}\left(t_{0} ; V\right)
\end{array}\right.
$$

In fact, for any $\varepsilon>0$ there exists an $\varepsilon^{\prime}\left(0<\varepsilon^{\prime}<\varepsilon\right)$ such that $\Gamma_{\mathrm{e}^{\prime}}\left(t_{0} ; V\right) \subset$ $\bigcup_{\nu=0}^{\infty} \Gamma_{\varepsilon}^{\nu}\left(t_{0} ; V_{\mathrm{z}}\right)$. It follows from (4.1) that $\Gamma\left(t_{0} ; V\right)$ is closed in $T^{*}\left(R^{n}\right) \backslash 0$.

Let $\phi_{j}(t, s ; x, \xi)$ be the solution of (3.1) with $\lambda=\lambda_{j}(j \in\{1, \cdots, l\})$, where $\lambda_{j}$ are characteristic roots of $\mathcal{L}$ of (1). From Proposition 3.1 we can find a small constant $T_{o}$ such that the following property holds: Let $t_{0}$ be a positive constant smaller than $T_{o}$. Then, for any fixed $\tilde{t}^{\nu} \in \Delta_{\nu}$ and $J_{\nu}=\left(j_{1}, \cdots, j_{\nu+1}\right) \in$ $\Pi_{\nu+1}$ the set of phase functions $\left\{\phi_{j_{k}}\left(t_{k-1}, t_{k} ; x, \xi\right)\right\}_{k=1}^{\nu+1}\left(t_{v+1}=0\right)$ satisfies the assumption (A-1) in Section 1. As in the $C^{\infty}$-case ([10], pp. 185-186) the fundamental solution $E(t, s)$ of (1) is constructed in the form

$$
\begin{align*}
& E(t, s)= \sum_{j=1}^{t} I_{j, \phi_{j}}(t, s)+\sum_{\nu=1} \sum_{\substack{j_{k} \in \in, \cdots, \cdots, l_{1} \\
k=1, \cdots, \cdots+1}}^{\infty} \int_{s}^{t} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-1}} I_{j_{1}, \phi_{j_{1}}}\left(t, t_{1}\right)  \tag{4.2}\\
& \times W_{j_{2}, \phi_{j_{2}}}\left(t_{1}, t_{2}\right) \cdots W_{j_{\nu+1}, \phi_{j+1}}\left(t_{v}, s\right) d t_{\nu} \cdots d t_{1} \\
&\left(t_{0}=t\right) \quad \text { for } 0 \leqq t, s \leqq T_{o},
\end{align*}
$$

where $I_{j, \phi_{j}}(t, s)$ is a matrix of Fourier integral operators with phase function $\phi_{j}(t, s ; x, \xi)$ and with symbol $1\left((j, j)\right.$ element) or 0 (others), and $W_{j, \phi_{j}}(t, s)$ is the one with symbol $w_{j}(t, s ; x, \xi)=w_{j}^{o}(t, s ; x, \xi)+\widetilde{w}_{j}(t, s ; x, \xi)$. Here $w_{j}^{o}(t, s ; x, \xi) \in G^{(\kappa)}\left(\left[0, T_{o}\right] \times\left[0, T_{o}\right] ; S_{G(k)}^{\sigma}\right)$ and $\widetilde{w}_{j}(t, s ; x, \xi) \in G^{(\kappa)}\left(\left[0, T_{o}\right] \times\right.$ $\left.\left[0, T_{o}\right] ; \mathcal{R}_{G(x)}\right)$, that is, $w_{j}^{o}$ satisfies

$$
\begin{equation*}
\left|\partial_{t}^{\gamma} \partial_{s}^{\gamma^{\prime}} \partial_{\xi}^{\alpha} \partial_{x}^{\beta} w_{j}^{o}(t, s ; x, \xi)\right| \leqq C M^{-\left(\gamma+\gamma^{\prime}+\left|\left.\right|_{|+|}\right| \beta \mid\right)}\left(\alpha!\beta!\gamma!\gamma^{\prime}!\right)^{\kappa}\langle\xi\rangle^{\sigma-\left|\alpha^{\alpha}\right|} \tag{4.3}
\end{equation*}
$$

for constants $C$ and $M$ independent of $\alpha, \beta, \gamma$ and $\gamma^{\prime}$, and $\widetilde{w}_{j}$ satisfies

$$
\begin{equation*}
\left|\partial_{t}^{\gamma} \partial_{s}^{\gamma^{\prime}} \partial_{\xi}^{\alpha} \partial_{x}^{\beta} \widetilde{w}_{j}(t, s ; x, \xi)\right| \leqq C_{\alpha} M^{-\left(\gamma+\gamma^{\prime}+|\beta|\right)}\left(\beta!\gamma!\gamma^{\prime}!\right)^{\kappa} e^{-\varepsilon\langle\xi\rangle^{1 / \kappa}} \tag{4.4}
\end{equation*}
$$

for constants $M$ and $\varepsilon>0$ independent of $\alpha, \beta, \gamma$ and $\gamma^{\prime}$, and for a constant $C_{\alpha}$ depending only on $\alpha$. It follows from Theorem 2.1 in [18] that the right hand side of (4.2) is transformed to that of (1.12) in [20] (cf. Theorem 3.2 in [18]). Therefore, from Proposition 1.1 in [20] we can find a solution $U(t)$ in $\mathscr{B}^{\circ}\left(\left[0, T_{o}\right]\right.$; $\left.\mathscr{D}_{L}^{(\mathbb{K})^{\prime}}\right)$ for the problem (7) as $U(t)=E(t, 0) G$. For the proof of the existence of the solution $U(t)$ for $t \in\left(T_{o}, T\right]$ it suffices to consider the product $E(t, 0)$ $\equiv E\left(t, k T_{o}\right) E\left(k T_{o},(k-1) T_{o}\right) \cdots E\left(T_{o}, 0\right)$ of the fundamental solutions if $t \in\left[k T_{o}\right.$, $\left.(k+1) T_{o}\right]$. Finally, we note that $E(t, s)$ of (1.12) in [20] maps $\mathscr{B}^{\infty}\left(\left[0, T_{o}\right] ; \mathcal{B}_{L^{2}}^{(k)}\right)$ to itself. So, the uniqueness of the solution also follows from the usual duality method.

For the proof of the inclusion (8) we prepare

Proposition 4.2. Let $V$ be a closed conic set in $T^{*}\left(R_{x}^{n}\right)$ and let $\Gamma_{e}\left(t_{0} ; V\right)$ be a set defined by (4.1) for an $\varepsilon>0$ and $0<t_{0} \leqq T$. Let $a(x, \xi)$ and $b(x, \xi)$ be symbols in $S_{G(k)}^{0}$ satisfying

$$
\begin{align*}
& \left\{\begin{array}{l}
\operatorname{supp} b \subset V_{\mathrm{\varepsilon} / 2}, \\
|x-y| \geqq \varepsilon / 2 \text { or } \quad|\xi /|\xi|-\eta /|\eta|| \geqq \varepsilon / 2
\end{array}\right.  \tag{4.5}\\
& \quad \text { if }(x, \xi) \in \operatorname{supp} a \text { and }(y, \eta) \in \Gamma_{\mathrm{e}}\left(t_{0} ; V\right) .
\end{align*}
$$

Then, for the fundamental solution $E(t, s)$ of (4.2) the operator $a\left(X, D_{x}\right) E\left(t_{0}, 0\right)$ $\times b\left(X, D_{x}\right)$ is a pseudo-differential operator with symbol in $\mathcal{R}_{G(k)}$.

Admitting this proposition for a moment, we first give the proof of (8) by using this. Let $U\left(t_{0}\right)$ be a solution of (7) and set $V=\mathrm{WF}_{G\left(\kappa_{1}\right)}(G)$ for $\kappa \leqq$ $\kappa_{1}<1 / \sigma$. Assume that ( $x^{0}, \xi^{0}$ ) does not belong to $\Gamma\left(t_{0} ; V\right)$. Then there exists an $\varepsilon>0$ such that $\left(x^{0}, \xi^{0}\right) \notin \Gamma_{0}\left(t_{\varepsilon} ; V\right)$. Since $\Gamma_{e}\left(t_{0} ; V\right)$ and $V$ are closed conic sets, taking another small $\varepsilon>0$ if necessary, we can find symbols $a(x, \xi)$ and $b(x, \xi)$ in $S_{G(x)}^{0}$ satisfying (4.5), $b(x, \xi)=1$ in a conic neighborhood of $V$ and $a\left(x^{0}, \theta \xi^{0}\right) \neq 0$ for $\theta \geqq 1$. Then we have

$$
\begin{aligned}
A U\left(t_{0}\right) & =A E\left(t_{0}, 0\right) G \\
& =A E\left(t_{0}, 0\right) B G+A E\left(t_{0}, 0\right)(I-B) G \in \gamma^{\left(\kappa_{1}\right)} .
\end{aligned}
$$

Here, we used Proposition 4.2 above and Lemma 2.1 in [20] for the proof of $A E\left(t_{0}, 0\right) B G \in \gamma^{\left(\kappa_{1}\right)}$, and for the proof of $A E\left(t_{0}, 0\right)(I-B) G \in \gamma^{\left(\kappa_{1}\right)}$ we used the similar discussions as in the proof of Proposition 1.1 and Theorem 4 in [20]. Then, in view of Definition , in Introduction we have $\left(x^{0}, \xi^{0}\right) \notin \mathrm{WF}_{G\left(\kappa_{1}\right)}\left(U\left(t_{0}\right)\right)$. This proves (8).

Now, we return to the proof of Proposition 4.2. First, we consider the case $T=T_{o}$. Regard $A E\left(t_{0}, 0\right) B$ as a pseudo-differential operator. Then its symbol is a sum of $\sigma\left(a\left(X, D_{x}\right) I_{j, \phi_{j}}\left(t_{0}, 0\right) b\left(X, D_{x}\right)\right)$ and the terms of the form

$$
\begin{align*}
& q_{\nu+1}\left(t_{0} ; x, \xi\right)  \tag{4.6}\\
& \equiv \int_{\Delta_{\nu-2}} d \tilde{t}^{\nu-2} U_{s}
\end{aligned} \quad-\iint^{i\left(\Psi_{\nu+1}\left(t_{0}, \tilde{t}^{\nu-2} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)-x \cdot \xi\right)} \begin{aligned}
& \times p_{\nu+1}\left(t_{0}, \tilde{t}^{\nu-2} ; x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right) d x d \tilde{\xi}^{\nu},
\end{align*}
$$

with $\left(j_{1}, \cdots, j_{\nu-1}\right) \in \Pi_{\nu-1}, \nu \geqq 3$. Here $\Psi_{\nu+1}$ is defined by (1.9) with $\phi_{1}=\phi_{\nu+1}$ $=x \cdot \xi$ and $\phi_{k+1}$ replaced by $\phi_{j_{k}}(k=1, \cdots, \nu-1)$, and $p_{\nu+1} \equiv p_{\nu+1}\left(t_{0}, \tilde{t}^{\nu-1} ; x, \tilde{\xi}^{\nu}\right.$, $\tilde{x}^{\nu}, \boldsymbol{\xi}$ ) is a multiple symbol defined by

$$
p_{\nu+1}=a\left(x, \xi^{1}\right)\left(\prod_{k=1}^{\nu-1} w_{j_{k}}\left(t_{k-1}, t_{k} ; x^{k}, \xi^{k+1}\right)\right) b\left(x^{\nu}, \xi^{\nu}\right) \quad\left(t_{\nu-1}=0\right)
$$

So, for the proof of Proposition 4.2 with $T=T_{o}$ it suffices to show

Proposition 4.3. Let $p_{v+1}$ be as above. Then the symbol $q_{v+1}\left(t_{0} ; x, \xi\right)$ ( $\nu \geqq 2$ ) defined by (4.6) belongs to $\mathcal{R}_{G(\kappa)}$ and it satisfies for positive constants $\varepsilon^{\prime}$, $A$ and $M$ independent of $\nu$

$$
\left|\partial_{\xi}^{\alpha} \partial_{x}^{\beta} q_{\nu+1}\left(t_{0} ; x, \xi\right)\right| \leqq C_{\alpha} A^{\nu} M^{-|\beta|} \beta!_{\nu}^{\kappa}!^{\sigma_{\alpha-1}-1} e^{-\varepsilon^{\prime}\langle\xi\rangle^{1 / \kappa}}
$$

for $\alpha$ and $\beta$. Here $C_{\alpha}$ is a constant independent of $\beta$ and $\nu$.
Proof. In view of Lemma 1.3 and its remark we may assume that $p_{\nu+1}$ satisfies (2.5). Hence it suffices to check the conditions $\mathcal{I}\left(A_{0}^{\nu}, 0, M,(\nu-2) \sigma\right)$ (for a constant $A_{0}$ independent of $\nu$ ) and $B(\varepsilon, \delta)$ in Lemma 3.2 by means of (3.11)' with $m=(\nu-2) \sigma$ and the fact that we have $[(\nu-2) \sigma]!^{\kappa} \leqq A^{\nu} \nu!^{\sigma^{\kappa}}$ for a constant $A$ independent of $\nu$. The condition $\mathscr{I}\left(A_{o}^{v}, 0, M,(\nu-2) \sigma\right)$ follows from (4.3) clearly. If $\left|\partial_{t_{k}} \Psi_{\nu+1}\left(t_{0}, \tilde{t}^{\nu-2} ; x, \xi\right)\right| \leqq \varepsilon\langle\xi\rangle / 2$ for all $k \in\{1, \cdots, \nu-2\}$ and if ( $x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi$ ) satisfies (3.22) for a sufficiently small $\delta>0$ then it follows from Proposition 3.5 that there exists an $\varepsilon$-admissible trajectory of step $\mu+1$, from $\left(X_{\nu}^{\nu}, \xi\right)$, whose end point is ( $x, \Xi_{v}^{1}$ ). Here $\mu$ is a number of elements in $\left\{k ; j_{k} \neq j_{k+1}\right\}$ for a permutation $\left(j_{1}, \cdots, j_{\nu-1}\right)$ which determines $\Psi_{\nu+1}$. Since $\left(x^{\nu}, \xi\right)$ and $\left(x, \xi^{1}\right)$ are contained in $\delta$-conic neighborhoods of ( $X_{\nu}^{\nu}, \xi$ ) and ( $x, \Xi_{1}^{\nu}$ ), respectively, we see by means of the choice of $a(x, \xi)$ and $b(x, \xi)$ that $\left\{p_{\nu+1}, \Psi_{\nu+1}\right\}$ satisfies (3.10) if $\delta$ also satisfies $\delta \leqq \varepsilon / 2$. Next, we consider a child $\left\{p_{\nu+1}^{\nu-3}, \Psi_{\nu+1}^{\nu-3}\right\}$ of $\left\{p_{\nu+1}, \Psi_{\nu+1}\right\}$. Then under $\left|\partial_{t_{j}} \Psi_{\nu+1}^{\nu-3}\right| \leqq \varepsilon\langle\xi\rangle / 2, j=1, \cdots, \nu-3$, we find the existence of an $\varepsilon$-admissible trajectory of step, at most, $\mu+1$ from ( $X_{\nu}^{\nu}, \xi$ ), whose end point is ( $x, \Xi_{v}^{1}$ ). Noting the choice of $a(x, \xi)$ and $b(x, \xi)$ we also see that this child satisfies (3.10). Repeating this procedure, we finally see that $\left\{p_{\nu+1}, \Psi_{\nu+1}\right\}$ satisfies the condition $B(\varepsilon, \delta)$.
Q.E.D.

We proceed to the proof of Proposition 4.2 for the general case. For simplicity we assume $T_{o} \leqq T \leqq 3 T_{o} / 2$. Then, from the uniqueness of the problem (7) it follows that for any $s$ with $T_{o} / 2 \leqq s \leqq T_{o}$

$$
\begin{equation*}
E\left(t_{0}, 0\right)=E\left(t_{0}, s\right) E(s, 0) \quad\left(T_{o} \leqq t_{0} \leqq T\right) \tag{4.7}
\end{equation*}
$$

Let $\omega(s)$ be a function in $\gamma^{(\kappa)}\left(R_{t}^{1}\right)$ satisfying $\operatorname{supp} \omega \subset\left(T_{o} / 2, T_{o}\right)$ and $\int_{T_{0} / 2}^{T_{o}} \omega(s) d s=1$. Then, from (4.7) we have

$$
E\left(t_{0}, 0\right)=\int_{T_{0} / 2}^{T_{o}} E\left(t_{0}, 0\right) \omega(s) d s=\int_{T_{o} / 2}^{T_{o}} E\left(t_{0}, s\right) \omega(s) E(s, 0) d s
$$

So, in order to show Proposition 4.2 it suffices to show

$$
\begin{equation*}
\int_{T_{o / 2}}^{T_{o}} A E\left(t_{0}, s\right) \omega(s) E(s, 0) B d s \in \mathcal{R}_{G(x)} \tag{4.8}
\end{equation*}
$$

with $A=a\left(X, D_{x}\right)$ and $B=b\left(X, D_{x}\right)$ whose symbols $a(x, \xi)$ and $b(x, \xi)$ satisfy
(4.5). Since $t_{0}-s \leqq T_{o}$ and $s \leqq T_{o}$, we can apply the discussions of Section 1 to each term in $E\left(t_{0}, s\right)$ and $E(s, 0)$, and we obtain

$$
\begin{align*}
& \sigma\left(\int_{T_{0} / 2}^{T o} A E\left(t_{0}, s\right) \omega(s) E(s, 0) B d s\right)(x, \xi) \\
& =\sum_{\nu=0}^{\infty} \sum_{\nu^{\prime}=0}^{\infty} \sum_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}} \int_{T_{0} / 2}^{T_{o}}\left\{\int_{s}^{t_{0}} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-1}}\right. \\
& \quad \cdot \int_{0}^{s} \int_{0}^{t_{1}^{\prime}} \cdots \int_{0}^{t_{\nu}^{\prime}-1}\left\{O_{s}-\int \cdots \int \exp \left(i\left(\Psi_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}-x \cdot \xi\right)\right)\right.  \tag{4.9}\\
& \quad \times p_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}\left(t_{0}, \tilde{t}^{\nu}, s, \tilde{t}^{\prime \nu^{\prime}} ; x, \xi^{0}, x^{0}, \tilde{\xi}^{\nu+v^{\prime}+1}, \tilde{x}^{\nu+v^{\prime}+1}, \xi\right) \\
& \left.\left.\quad \cdot d x^{0} d \xi^{0} d \tilde{x}^{\nu+\nu^{\prime}+1} d \tilde{\xi}^{\nu+v^{\prime}+1}\right\} d \tilde{t}^{\nu} d \tilde{t}^{\prime \nu^{\prime}}\right\} d s+r\left(t_{0}, s\right)
\end{align*}
$$

with $r\left(t_{0}, s\right) \in \mathcal{R}_{G(\mathrm{k})}$, where

$$
\begin{aligned}
\Psi_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}= & x \cdot \xi^{0}-x^{0} \cdot \xi^{0} \\
& +\sum_{k=1}^{\nu}\left(\phi_{j_{k}}\left(t_{k-1}, t_{k} ; x^{k-1}, \xi^{k}\right)-x^{k} \cdot \xi^{k}\right) \\
& +\phi_{j_{\nu+1}}\left(t_{v}, s ; x^{\nu}, \xi^{\nu+1}\right)-x^{\nu+1} \cdot \xi^{\nu+1} \\
& +\sum_{k=1}^{\nu^{\prime}}\left(\phi_{j_{k}^{\prime}}\left(t_{k-1}^{\prime}, t_{k}^{\prime} ; x^{\nu+k}, \xi^{\nu+k+1}\right)-x^{\nu+k+1} \cdot \xi^{\nu+k+1}\right) \\
& +\phi_{j^{\prime} \prime \prime+1}\left(t_{\nu^{\prime}}^{\prime}, 0 ; x^{\nu+\nu^{\prime}}, \xi^{\nu+\nu^{\prime}+1}\right) \\
& -x^{\nu+\nu^{\prime}+1} \cdot \xi^{\nu+\nu^{\prime}+1}+x^{\nu+v^{\prime}+1} \cdot \xi \quad\left(t_{0}^{\prime}=s\right) .
\end{aligned}
$$

Here, $\sum_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}$ means the summation which is taken over all $J_{\nu+1}=\left(j_{1}, \cdots, j_{\nu+1}\right)$ and $J_{\nu^{\prime}+1}^{\prime}=\left(j_{1}^{\prime}, \cdots, j_{\nu^{\prime}+1}^{\prime}\right)$ with $j_{k}, j_{k}^{\prime}=1, \cdots, l$; and in (4.9), $\int_{s}^{t_{0}} \cdots \int_{s}^{t_{\nu-1}}\left(\right.$ resp. $\int_{0}^{s} \cdots$ $\int_{0}^{t_{\nu^{\prime}-1}^{\prime}}$ ) for the case $\nu=0$ (resp. $\nu^{\prime}=0$ ) means that we do not integrate the integrand with respect to $\tilde{t}^{\nu}$-variables (resp. $\tilde{t}^{\nu^{\prime}}$-variables). In (4.9) the symbol $p_{J_{\nu+1}, J_{v^{\prime}+1}^{\prime}}$ satisfies

$$
\begin{aligned}
& \left\{\begin{array}{l}
\left|\xi^{j}-\xi^{\nu+1}\right| \leqq\left\langle\xi^{\nu+1}\right\rangle / 8 \quad(j=0, \cdots, \nu), \\
\left|\xi^{j}-\xi\right| \leqq\langle\xi\rangle / 8 \quad\left(j=\nu+2, \cdots, \nu+\nu^{\prime}+1\right)
\end{array}\right. \\
& \text { on } \quad \operatorname{supp} p_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}} \text {, } \\
& \left.p_{J_{v+1}, J_{v^{\prime}+1}^{\prime}}\right|_{s=T_{o} / 2}=0, p_{J_{v+1}, J_{v^{\prime}+1}^{\prime}} l_{s=T_{O}}=0
\end{aligned}
$$

and the pair $\left\{p_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}, \Psi_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}\right\}$ with $\Psi_{J_{\nu+1}, J_{\nu^{\prime}+1}^{\prime}}$ satisfies condition similar to $B(\varepsilon, \delta)$ (we note that in this case we pose the condition (3.10) with $X_{\nu}^{j}$ and $\Xi_{\nu}^{j}$ replaced by the points $\chi_{j_{k}^{\prime}}\left(t_{k-1}^{\prime}, t_{k}^{\prime}\right) \cdots \chi_{j_{\nu^{\prime}+1}}\left(t_{\nu^{\prime}}^{\prime}, 0\right)\left(x^{\nu+\nu^{\prime}+1}, \xi\right)\left(1 \leqq k \leqq \nu^{\prime}+1\right.$; $\left.t_{0}^{\prime}=s, t_{\nu^{\prime}+1}^{\prime}=0\right)$ or $\chi_{j_{k}}\left(t_{k-1}, t_{k}\right) \cdots \chi_{j_{\nu+1}}\left(t_{\nu}, s\right) \chi_{j_{1}^{\prime}}\left(s, t_{1}^{\prime}\right) \cdots \chi_{j_{\nu^{\prime}+1}^{\prime}}\left(t_{\nu^{\prime}}^{\prime}, 0\right)\left(x^{\nu+\nu^{\prime}+1}, \xi\right)(2 \leqq$ $k \leqq \nu+1 ; t_{\nu+1}=s$ ) in the trajectories). Replacing $r\left(t_{0}, s\right)$ in (4.9) by another
symbol in $\mathcal{R}_{G(\kappa)}$ we may assume moreover that

$$
\left|\xi^{\nu+1}-\xi^{\nu+2}\right| \leqq\left\langle\xi^{\nu+2}\right\rangle / 2 \quad \text { on } \quad \operatorname{supp} p_{J_{\nu+1},}, J_{\nu^{\prime}+1}^{\prime}
$$

Then, we can prove (4.8) by the similar discussions as the case of $T \leqq T_{o}$ if we use Lemma $2.1^{\prime}$ below instead of Lemma 2.1 and the fact that with $\Psi=$ $\Psi_{J_{v+1},} J_{\nu^{\prime}+1}^{\prime}$ and $p=p_{J_{\nu+1},}, J_{\nu^{\prime}+1}^{\prime}$

$$
\begin{align*}
& \int_{T_{o} / 2}^{T_{o}}\left\{\int_{s}^{t_{0}} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-1}} \int_{0}^{s} \int_{0}^{t_{1}^{\prime}} \cdots \int_{0}^{t_{\nu}^{\prime}-1} \partial_{s} e^{i(\Psi-x \cdot \xi)}\right.  \tag{4.10}\\
& \left.\times\left(i \partial_{s} \Psi\right)^{-1} p d \tilde{t}^{\nu} d \tilde{t}^{\prime \nu^{\prime}}\right\} d s \\
& =\int_{T_{0} / 2}^{T_{o}}\left\{\frac { \partial } { \partial s } \left\{\int_{s}^{t_{0}} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-1}} \int_{0}^{s} \int_{0}^{t_{1}^{\prime}} \cdots \int_{0}^{t_{\nu_{\prime}^{\prime}-1}} e^{i(\Psi-x \cdot \xi)}\right.\right. \\
& \left.\left.\times\left(i \partial_{s} \Psi\right)^{-1} p d \tilde{t}^{\nu} d \tilde{t}^{\prime \nu^{\prime}}\right\}\right\} d s \\
& -\left[-\int_{T_{o} / 2}^{T_{o}}\left\{\int_{s}^{t_{0}} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-2}} \int_{0}^{s} \int_{0}^{t_{1}^{\prime}} \cdots \int_{0}^{t_{\nu_{\nu}^{\prime}-1}} e^{i(\Psi-x \cdot \xi)}\right.\right. \\
& \left.\times\left(\left(i \partial_{s} \Psi\right)^{-1} p\right)_{t_{\nu}=s} d t_{1} \cdots d t_{\nu-1} d \tilde{t}^{\prime \prime}\right\} d s \\
& +\int_{T_{0} / 2}^{T_{O}}\left\{\int_{s}^{t_{0}} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-1}} \int_{0}^{s} \int_{0}^{t_{2}^{\prime}} \cdots \int_{0}^{t_{\nu^{\prime}-1}} \epsilon^{i(\Psi-x \cdot \xi)}\right. \\
& \left.\times\left(\left(i \partial_{s} \Psi\right)^{-1} p\right)_{\mid t_{1}^{\prime}=s} d \tilde{t}^{\nu} d t_{2}^{\prime} \cdots d t_{\nu}^{\prime}\right\} d s \\
& +\int_{T_{0} / 2}^{T_{o}}\left\{\int_{s}^{t_{0}} \int_{s}^{t_{1}} \cdots \int_{s}^{t_{\nu-1}} \int_{0}^{s} \int_{0}^{t_{1}^{\prime}} \cdots \int_{0}^{t_{\nu^{\prime}-1}^{\prime}} e^{i(\Psi-x \cdot \xi)}\right. \\
& \left.\left.\times \partial_{s}\left\{\left(i \partial_{s} \Psi\right)^{-1} p\right\} d \tilde{t}^{\nu} d \tilde{t}^{\prime \nu^{\prime}}\right\} d s\right]
\end{align*}
$$

holds and that the first term in the right member of (4.10) is zero. Hence, the proof of Proposition 4.2 for the case of $T_{o} \leqq T \leqq 3 T_{o} / 2$ is reduced to the proof of Lemma 2.1' below.

Lemma 2.1'. Let $p_{\nu+1}\left(x, \tilde{\xi}^{\nu}, \tilde{x}^{\nu}, \xi\right)$ satisfy (2.6) and consider $q_{\nu+1}(x, \xi)$ defined by (2.2) with $\Phi_{\nu+1}(x, \xi)$ replaced by $x \cdot \xi$. Suppose that for a $k$ the variables $\left(\xi^{1}, \cdots, \xi^{\nu}\right)$ are divided into two groups $\left(\xi^{1}, \cdots, \xi^{k}\right)$ and $\left(\xi^{k+1}, \cdots, \xi^{\nu}\right)$ and they satisfy

$$
\left\{\begin{array}{lr}
\left|\xi^{j}-\xi^{k}\right| \leqq\left\langle\xi^{k}\right\rangle / 8 & (j=1, \cdots, k-1) \\
\left|\xi^{j}-\xi\right| \leqq\langle\xi\rangle / 8 & (j=k+1, \cdots, \nu) \\
\left|\xi^{k}-\xi^{k+1}\right| \leqq\left\langle\xi^{k+1}\right\rangle / 2 & \text { on } \operatorname{supp} p_{\nu+1}
\end{array}\right.
$$

Moreover, let $\phi_{j} \in \mathcal{P}_{G(k)}\left(\tau_{j}\right)$ and assume $\sum_{j=1}^{k} \tau_{j} \leqq \tau^{0}$ and $\sum_{j=k+1}^{\nu+1} \tau_{j} \leqq \tau^{0}$ with $\tau^{0}$ in (A-1). Then, there exist constants $A$ and $C_{\alpha}$, and for any $\varepsilon(>0)$ there exists a constant $M=M_{\mathrm{e}}$ such that

$$
\begin{equation*}
\left|q_{\nu+1(\beta)}^{(\alpha)}(x, \xi)\right| \leqq C_{o} C_{\alpha} A^{\nu} M^{-|\beta|} \beta!^{\kappa}\langle\xi\rangle^{m} \exp \left(\varepsilon\langle\xi\rangle^{1 / \kappa}\right) \tag{4.11}
\end{equation*}
$$

## hold.

Remark. In the above lemma we need not assume $\sum_{j=1}^{\nu+1} \tau_{j} \leqq \tau^{0}$.
Proof. Let $q_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)$ be the symbol defined by (2.9). Then, since $\sum_{j=1}^{k} \tau_{j} \leqq \tau^{0}$ and $\sum_{j=k+1}^{\nu+1} \tau_{j} \leqq \tau^{0}$, we can apply the discussions in Section 2 to the integrals with respect to ( $x^{1}, \cdots, x^{k-1}, \xi^{1}, \cdots, \xi^{k-1}$ ) and ( $x^{k+1}, \cdots, x^{\nu}, \xi^{k+1}, \cdots, \xi^{\nu}$ ) individually and obtain (2.10) with

$$
\begin{equation*}
q_{\nu+1}(x, \xi)=O_{s}-\iint e^{i\left(\Phi_{\nu+1, k}^{\left.\left(x, \xi^{k}, x^{k}, \xi\right)-x \cdot \xi\right)}\right.} q_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right) d x^{k} d \xi^{k} \tag{4.12}
\end{equation*}
$$

where $\Phi_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)$ is defined by (2.8). Write (4.12) as

$$
q_{v+1}(x, \xi)=O_{s}-\iint e^{i\left(x-x^{k}\right) \cdot\left(\xi-\xi^{k}\right)} \tilde{q}_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right) d x^{k} d \xi^{k},
$$

where

$$
\tilde{q}_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)=e^{i \tilde{J}_{\nu+1}\left(x, \xi^{k}, x^{k}, \xi\right)} q_{\nu+1, k}\left(x, \xi^{k}, x^{k}, \xi\right)
$$

and

$$
\begin{aligned}
& \tilde{J}_{\nu+1}\left(x, \xi^{k}, x^{k}, \xi\right)=\left\{\left(\phi_{1} \# \cdots \# \phi_{k}\right)\left(x, \xi^{k}\right)-x \cdot \xi^{k}\right\} \\
& \\
& \quad+\left\{\left(\phi_{k+1} \# \cdots \# \phi_{v+1}\right)\left(x^{k}, \xi\right)-x^{k} \cdot \xi\right\} .
\end{aligned}
$$

Then, together with the discussion (1.6)-(1.7) in [18] for $\tilde{J}_{\nu+1}\left(x, \xi^{k}, x^{k}, \xi\right)$, we can easily prove (4.11) by using the factor $\exp \left(\varepsilon\langle\xi\rangle^{1 / k}\right)$ in the right member of (4.11).
Q.E.D.

## 5. Hyperbolic differential operators

Let $L$ be a single hyperbolic operator of order $m$ which has a form

$$
\begin{equation*}
L=L_{1} L_{2} \cdots L_{r}+A_{1} L_{2} \cdots L_{r}+A_{2} L_{3} \cdots L_{r}+\cdots+A_{r-1} L_{r}+A_{r} \tag{5.1}
\end{equation*}
$$

where $L_{k}(k=1, \cdots, r)$ are regularly hyperbolic operators with coefficients in $\gamma^{(k)}\left([0, T] \times R_{x}^{n}\right)$ and $A_{k}(k=1, \cdots, r)$ are differential operators with coefficients in $\gamma^{(\kappa)}\left([0, T] \times R_{x}^{n}\right)$ satisfying

$$
\begin{equation*}
\operatorname{Ord} A_{k} \leqq \operatorname{Ord}\left(L_{1} \cdots L_{k}\right)-k / \mu \tag{5.2}
\end{equation*}
$$

for a constant $\mu \geqq 1$. We assume $\kappa<\mu /(\mu-1)$. The form (5.1) is a generalization of (12) in Introduction. In fact, (12) is derived from (5.1) by setting $A_{1}=A_{2}=\cdots=A_{r-1}=0$. We remark that any hyperbolic operator with characteristic roots of constant multiplicity can be written in the form (5.1) if the constant $\mu$ is defined as the irregularity of the hyperbolic operator (cf. Theorem 3.1 of [7], see also Lemma 4.1 of [3]).

Let $m_{k}$ denote the order of $L_{k}$ and let $\lambda_{k, j}(t, x, \xi), j=1, \cdots, m_{k}$, be characteristic roots of $L_{k}$. We may assume $\lambda_{k, j} \in G^{(k)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$ by multiplying a cut function with respect to $\xi$ if necessary. Since $L_{k}$ is a regularly hyperbolic operator, there exist $\lambda_{k, j}^{\prime}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(k)}^{0}\right)$ such that

$$
\begin{align*}
L_{k}= & \left(D_{t}-\lambda_{k, 1}\left(t, X, D_{x}\right)-\lambda_{k, 1}^{\prime}\left(t, X, D_{x}\right)\right) \cdots  \tag{5.3}\\
& \times\left(D_{t}-\lambda_{k, m_{k}}\left(t, X, D_{x}\right)-\lambda_{k, m_{k}}^{\prime}\left(t, X, D_{x}\right)\right) \\
& +\sum_{j=0}^{m_{k}-1} b_{k, j}\left(t, X, D_{x}\right) D_{t}^{j}
\end{align*}
$$

with $b_{k, j}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(x)}^{-2 m}\right)$. Here the equality (5.3) means that it holds modulo regularizers of the form $\sum_{j=0} r_{j}\left(t, X, D_{x}\right) D_{t}^{j}$ with $r_{j} \in \mathcal{R}_{G(\kappa)}$ for any fixed $t$. Since we may disregard the contribution of such regularizers in our discussion, till the end of this section the equality means that it holds modulo regularizers. Set $\bar{m}_{0}=0, \bar{m}_{k}=m_{1}+\cdots+m_{k}\left(\bar{m}_{r}=m\right)$ and

$$
\begin{gather*}
\partial_{j}=D_{t}-\lambda_{k, j-\bar{m}_{k-1}}\left(t, X, D_{x}\right)-\lambda_{k, j-\bar{m}_{k-1}}^{\prime}\left(t, X, D_{x}\right)  \tag{5.4}\\
\text { if } \bar{m}_{k-1}<j \leqq \bar{m}_{k} .
\end{gather*}
$$

Proposition 5.1. Let L be a hyperbolic operator of order $m$ which has the form (5.1). Assume (5.2). Then, $L$ can be written as

$$
\begin{align*}
L= & \partial_{1} \cdots \partial_{m} \\
& +\sum_{p=1}^{m_{1}-1} \sum_{1 \leqq j_{1}<\cdots<j_{p} \leq m_{1}} b_{j_{1} \cdots j_{p}}^{1} \partial_{j_{1}} \cdots \partial_{j_{p}} \partial_{m_{1}+1} \cdots \partial_{m} \\
& +\sum_{p=1}^{\bar{m}_{2}-2} \sum_{1 \leqq j_{1}<\cdots<j_{p} \leq m_{2}} b_{j_{1} \cdots j_{p}}^{2} \partial_{j_{1}} \cdots \partial_{j_{p}} \partial_{\bar{m}_{2}+1} \cdots \partial_{m}  \tag{5.5}\\
& +\cdots \\
& +\sum_{p=0}^{m-r} \sum_{1 \leqq j_{1}<\cdots<j_{p} \leq m} b_{j_{1} \cdots j_{p}}^{r} \partial_{j_{1}} \cdots \partial_{j_{p}},
\end{align*}
$$

where $\partial_{j}$ are defined by (5.4) and $b_{j_{1} \ldots j_{p}}^{k}$ is a pseudo-differential operator $b_{j_{1} \ldots j_{p}}^{k}$ $\left(t, X, D_{x}\right)$ with symbol $b_{j_{1} \cdots j_{p}}^{k}(t, x, \xi) \in G^{(k)}\left([0, T] ; S_{G(k)}^{k-k / \mu}\right)$.

For the proof we prepare
Lemma 5.2. Let $s$ be a positive integer and let $\partial_{j}(j=1, \cdots, s)$ denote $D_{t}-\lambda_{j}\left(t, X, D_{x}\right)$ for some $\lambda_{j}(t, x, \xi) \in G^{(k)}\left([0, T] ; S_{G(k)}^{1}\right)$. Assume $\mid \lambda_{j}(t, x, \xi)-$ $\lambda_{k}(t, x, \xi) \mid \geqq c_{o}\langle\xi\rangle$ for a constant $c_{o}>0$ if $j \neq k$ and $|\xi|$ is large. Let $A$ be an operator of the form

$$
\begin{align*}
& A=\sum_{k=1}^{s-1} b_{k}\left(t, X, D_{x}\right) D_{t}^{k}  \tag{5.6}\\
& \quad \text { for } b_{k}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{s-1-k}\right) .
\end{align*}
$$

Then, we can worite $A$ as

$$
A=\sum_{1 \leq j_{1}<\cdots<j_{k} \leq s}^{k<s} a_{j_{1} \cdots j_{k}}\left(t, X, D_{x}\right) \partial_{j_{1}} \cdots \partial_{j_{k}}+a\left(t, X, D_{x}\right)
$$

for some $a(t, x, \xi)$ and $a_{j_{1} \ldots j_{k}}(t, x, \xi)$ in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)$.
The proof of this lemma easily follows from the induction on $s$.
Proof of Proposition 5.1. Let $b(t, x, \xi)$ be a symbol in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{-m}\right)$. Then for any integer $0 \leqq k \leqq m$ there exist $a_{j}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)(j=0$, $\cdots, k)$ such that $b D_{t}^{k}=a_{0}+\sum_{j=2}^{k} a_{j} \partial_{m-j+1} \cdots \partial_{m-1} \partial_{m}$. On the other hand, from (5.3) we can write $L_{1} L_{2} \cdots L_{r}$ as

$$
\begin{aligned}
L_{1} L_{2} \cdots L_{r}= & \left(\partial_{1} \cdots \partial_{m}\right)\left(\partial_{m_{1}+1} \cdots \partial_{\bar{m}_{2}}\right) \cdots\left(\partial_{\bar{m}_{r-1}+1} \cdots \partial_{\bar{m}_{r}}\right) \\
& +\sum_{k=1}^{r} L_{1} \cdots L_{k-1}\left(\sum_{j=0} b_{k, j}\left(t, X, D_{x}\right) D_{t}^{j}\right) \partial_{\bar{m}_{k}+1} \cdots \partial_{\bar{m}_{r}}
\end{aligned}
$$

and the second term of the right hand side can be rewritten as

$$
\sum_{j=0}^{m-1} b_{j}\left(t, X, D_{x}\right) D_{t}^{j}
$$

for some $b_{j}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{-m}\right)$. Hence we may assume $L_{k}=\partial_{\bar{m}_{k-1}+1} \cdots$ $\partial_{\bar{m}_{k}}$, and so $L_{k} \cdots L_{r}=\partial_{\bar{m}_{k-1}+1} \cdots \partial_{m}$. Since $D_{t}^{m_{k}}$ can be written as $L_{k}-\sum_{j=1}^{m_{k}} b_{k}^{j} D_{t}^{m_{k}-j}$ for some $b_{k}^{j} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{j}\right)$, we may also assume that the order of $A_{k}$ with respect to $D_{i}$ is smaller than or equal to $\bar{m}_{k}-k$. Consequently, $A_{k}$ can be written as the finite sum of operators of the form $\Lambda^{k-k / \mu} A_{k, 1} A_{k, 2} \cdots A_{k, k}$, where $\Lambda=\left\langle D_{x}\right\rangle$ and $A_{k, j}(j=1, \cdots, k)$ is the operator of the form (5.6) with $s=m_{j}$. Applying Lemma 5.2 to each $A_{k, j}$ and $\left\{\partial_{\bar{m}_{j-1}+1}, \cdots, \partial_{\bar{m}_{j}}\right\}$, we have

$$
\begin{aligned}
& A_{k}=\sum_{p=1}^{\bar{m}_{k}-k} \sum_{1 \leqq j_{1}<j_{2}<\cdots<j_{p} \leq m_{k}} b_{j_{1}}^{k} \cdots j_{p} \partial_{j_{1}} \cdots \partial_{j_{p}}, \\
& b_{j_{1} \cdots j_{p}}^{k} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{k-k / \mu}\right),
\end{aligned}
$$

and this gives (5.5).
Q.E.D.

Theorem 2 (cf. Proposition 3.3 in [18]). Let L be a hyperbolic operator of the form (5.1). Assume (5.2). Set $\sigma=1-1 / \mu$. Then there exists a hyperbolic system $\mathcal{L}$ of the form (1) with $b_{j, k}(t, x, \xi)$ in $G^{(\kappa)}\left([0, T] ; S_{G\left(\kappa_{k}\right)}^{\sigma}\right)$ such that the Cauchy problem (10) for $L$ can be reduced to the equivalent Cauchy problem (7) for $\mathcal{L}$.

Consequently we have (13) for $\kappa_{1}$ satisfying $\kappa \leqq \kappa_{1}<1 / \sigma=\mu /(\mu-1)$ concerning the propagation of wave front sets in Gevrey classes of solutions of (10).

Proof. For $1 \leqq p<m$ we set

$$
\Pi_{p}^{o}=\left\{J=\left(j_{1}, \cdots, j_{p}\right) ; 1 \leqq j_{1}<\cdots<j_{p} \leqq m\right\}
$$

and for $J \in \Pi_{p}^{o}$ we denote the length $p$ of $J$ by $|J|$. Let $\Pi_{p}$ be a subset of $\Pi_{o}^{p}$ whose element $J=\left(j_{1}, \cdots, j_{p}\right)$ satisfies the following: Set $S_{J}=\left\{j_{1}, \cdots, j_{p}\right\}$. If there exists a $k \in\{1, \cdots, r\}$ such that $\left\{\bar{m}_{k-1}+1, \cdots, \bar{m}_{k}\right\} \supset S_{J}$ then the set $\left\{\bar{m}_{k}+1, \cdots, m\right\}$ is also contained in $S_{J} . \quad$ Set $\Pi=\{0\} \cup\left(\bigcup_{p=1}^{m-1} \Pi_{p}\right)$ and denote the number of elements of $\Pi$ by $l$. Let $u$ be a solution of (10) and set

$$
\left\{\begin{array}{l}
u_{0}=\Lambda^{(m-1) \sigma} u  \tag{5.7}\\
u_{J}=\Lambda^{(m-1-p) \sigma} \partial_{J} u, \quad J=\left(j_{1}, \cdots, j_{p}\right) \in \Pi_{q}, 1 \leqq p \leqq m,
\end{array}\right.
$$

where $\partial_{J}=\partial_{j_{1}} \cdots \partial_{j_{p}}$. Then, from (5.5) we can write $L u=0$ as

$$
\begin{equation*}
L u=\partial_{1} u_{(2, \cdots, m)}+\sum_{J \in \mathrm{I}} b_{J}\left(t, X, D_{x}\right) u_{J}=0 \tag{5.8}
\end{equation*}
$$

for some $b_{J}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{\sigma}\right)$. For $J \in \Pi$ we set $j_{0}=\max \{j ; 1 \leqq j \leqq$ $\left.m, j \notin S_{J}\right\}$, where we denote $S_{J}=\phi$ for $J=0$. We shall show

$$
\begin{gather*}
\partial_{j_{0}} \Lambda^{(m-1-|J|) \sigma} \partial_{J}=\left\{\begin{array}{l}
\sum_{\tilde{T} \in \mathbb{I}} b_{\widetilde{J}}^{J} \Lambda^{(m-1-\mid \tilde{J})^{\sigma}} \partial_{\widetilde{J}} \quad \text { if } \quad|J| \leqq m-2, \\
\partial_{1} \cdots \partial_{m}+\sum_{\tilde{J} \in \mathbb{I}} b_{\widetilde{J}}^{J} \Lambda^{(m-1-|I|)^{\prime} \sigma} \partial_{\tilde{J}}
\end{array}\right.  \tag{5.9}\\
\text { for } b_{J}^{J}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(k)}^{\sigma}\right) .
\end{gather*}
$$

Then, together with (5.8) we have for any $J \in \Pi$

$$
\partial_{j_{0}} u_{J}=\left\{\begin{array}{rc}
\sum_{\tilde{J} \in \mathrm{I}} b_{\tilde{J}}^{J}\left(t, X, D_{x}\right) u_{\tilde{J}} \quad \text { if }|J| \leqq m-2 \\
-\sum_{\tilde{\tau} \in \mathrm{II}} b_{\tilde{J}}\left(t, X, D_{x}\right) u_{\tilde{J}}+\sum_{\tilde{J} \in \mathbb{I}} b_{\tilde{J}}^{J}\left(t, X, D_{x}\right) u_{\tilde{J}} \\
\text { if }|J|=m-1
\end{array}\right.
$$

This shows that the $l$-dimensional vector $U=\left(u_{J}\right)_{J \in I I}$ satisfies $\mathcal{L} U=0$ for a system $\mathcal{L}$ of the form (1). In this way we reduce the problem (10) for $L$ to a problem (7) for $\mathcal{L}$. The fact that (10) and (7) are equivalent is verified by the same way as in [13] and [8].

So, it remains to prove (5.9). To prove this it suffices to show

$$
\begin{align*}
& \partial_{j_{0}} \partial_{J}=\left\{\begin{array}{l}
\sum_{\tilde{J} \in \mathrm{I}} a_{\tilde{J}}\left(t, X, D_{x}\right) \partial_{\tilde{J}} \quad \text { if }|J| \leqq m-2 \\
\partial_{1} \cdots \partial_{m}+\sum_{\tilde{J} \in \mathbb{I}} a_{\widetilde{J}}\left(t, X, D_{x}\right) \partial_{\tilde{J}} \quad \text { if } \quad|J|=m-1
\end{array}\right.  \tag{5.10}\\
& \text { for } \quad a_{\widetilde{J}}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right) \text {. }
\end{align*}
$$

Since (5.10) follows immediately for the case when $J=0$ or $j_{0}<j$ for any $j \in S_{J}$, we take $J \in \Pi$ such that $j_{0}>j$ holds for some $j \in S_{J}$. Note that $J$ can be written as $\left(j_{1}, \cdots, j_{\tilde{p}}, \bar{m}_{k-1}+1, \cdots, m\right)$ for some $\tilde{p}$ and $k$. Let $p^{\prime} \in\{1, \cdots, \tilde{p}\}$ be a maximal integer such that $j_{p^{\prime}}<j_{0}$ and write $\partial_{j_{0}} \partial_{J}$ as

$$
\begin{align*}
\partial_{j_{0}} \partial_{J}= & \partial_{j_{1}} \cdots \partial_{j_{p^{\prime}}} \partial_{j_{0}} \partial_{j_{p^{\prime}+1}} \cdots \partial_{j_{p}} \partial_{\bar{m}_{k-1}+1} \cdots \partial_{m}  \tag{5.11}\\
& +\sum_{q=1}^{p^{\prime}} \partial_{j_{1}} \cdots \partial_{j_{q-1}}\left[\partial_{j_{0}}, \partial_{j_{q}}\right] \partial_{j_{q+1}} \cdots \partial_{j_{p}} \partial_{\bar{m}_{k-1}+1} \cdots \partial_{m} .
\end{align*}
$$

The definition of $\Pi_{p}$ implies that for any $j_{q}$ there exists a $j^{\prime} \notin S_{J}$ such that $\bar{m}_{k^{\prime}-1}+1 \leqq j^{\prime} \leqq \bar{m}_{k^{\prime}}$, holds with $k^{\prime}$ satisfying $\bar{m}_{k^{\prime}-1}+1 \leqq j_{q} \leqq \bar{m}_{k^{\prime}}$. Hence, by means of the regular hyperbolicity of $L_{k^{\prime}}$ we have $\left|\lambda_{j_{q}}-\lambda_{j^{\prime}}\right| \geqq c_{o}\langle\xi\rangle$ for some $c_{o}>0$ if $|\xi|$ is large enough. Consequently we have

$$
\begin{aligned}
{\left[\partial_{j_{0}}, \partial_{j_{q}}\right]=} & a_{1} \partial_{j_{q}}+a_{2} \partial_{j^{\prime}}+a_{3} \\
& \text { for } a_{j}(t, x, \xi) \in G^{(\kappa)}\left([0, T] ; S_{G(k)}^{0}\right) .
\end{aligned}
$$

In view of this, each term with commutator [ $\left.\partial_{j_{0}}, \partial_{j_{q}}\right]$ in (5.11) can be written as the linear combination of $\partial_{J}, \partial_{J^{\prime}}, \partial_{J_{q}}$ and their minor operators $\partial_{J^{\prime \prime}}$ with coefficients in $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)$, where $J_{q} \in \Pi_{p-1}$ is defined by $S_{J_{q}}=S_{J} \backslash\left\{j_{q}\right\}$, $J^{\prime}$ is the permutation $\left(j_{1}, \cdots, j_{q-1}, j^{\prime}, j_{q+1}, \cdots, j_{p}, \bar{m}_{k-1}+1, \cdots, m\right.$ ) and the permutation $J^{\prime \prime}$ of the minor operator $\partial_{J^{\prime \prime}}$ is defined by $S_{J^{\prime \prime}}=S_{J} \backslash \tilde{S}$ or $S_{J^{\prime}} \backslash \tilde{S}$ or $S_{J_{q}} \backslash \tilde{S}$ for a subset $\tilde{S}$ of $\left\{j_{1}, \cdots, j_{q-1}\right\}$. For the operator $\partial_{J^{\prime}}$ or its minor operators $\partial_{J^{\prime \prime}}$, if $J^{\prime} \notin \Pi$ or $J^{\prime \prime} \notin \Pi$, we repeat the above discussions until $\partial_{J^{\prime}}$ and $\partial_{J^{\prime \prime}}$ are represented as the linear combination of $\partial_{J}$ with $J$ in $\Pi$. Then, we get (5.10) and we can complete the proof of the theorem.
Q.E.D.

As another application of Theorem 1 we consider an operator $L$ of the form

$$
\begin{equation*}
L=L_{1} L_{2} L_{3}+P_{1} L_{1}+P_{2} L_{2}+P_{3} L_{3}+P_{4} . \tag{5.12}
\end{equation*}
$$

Here, $L_{j}, j=1,2,3$, are regularly hyperbolic differential operators of order $m_{j}\left(m_{1}+m_{2}+m_{3}=m\right)$ and $P_{1}, P_{2}, P_{3}$ and $P_{4}$ are differential operators of order, at most, $m-m_{1}-1, m-m_{2}-1, m-m_{3}-1$ and $m-1$, respectively, with coefficients in $\gamma^{(k)}\left([0, T] \times R_{x}^{n}\right)$. If we admit $L_{j}$ and $P_{j}$ in (5.12) to be pseudodifferential operators with respect to $x$, a hyperbolic operator with characteristic roots of the maximal multiplicity at most three always has the form (5.12), provided that its characteristic roots belong to $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$. The assumption of differential operators with respect to $x$ is not necessary for the argument in what follows.

Theorem 3. Let $L$ be a hyperbolic operator of the form (5.12). Then, the problem (10) for $L$ can be reduced to the equivalent problem (7) for an operator of the form (1) with $\mathcal{L}$ satisfying the following:


Remark. When the operator $L$ is a differential operator whose maximal multiplicity is at most three, it seems that the cases i)-iv) cover all the cases which we can consider as the conditions on lower order terms for any given constant $\sigma<1$. In the above we make a convention: the terms of the forms $A_{1} L_{2} L_{3}, A_{2} L_{1} L_{3}$ and $A_{3} L_{1} L_{2}$ are absorbed in $L_{1} L_{2} L_{3}$ by modifying the lower order terms of $L_{j}$.

As in the proof of Proposition 5.1 we may assume $L_{1}=\partial_{1} \cdots \partial_{m_{1}}, L_{2}=$ $\partial_{m_{1}+1} \cdots \partial_{m_{1}+m_{2}}, L_{3}=\partial_{m_{1}+m_{2}+1} \cdots \partial_{m}$, where $\partial_{j}(j=1, \cdots, m)$ are defined by (5.4). Let $\Pi_{p}^{o}$ and $\Pi_{p}$ denote the same sets as in the proof of Theorem 2.

Proof of Case i) (cf. [16]). Since $\operatorname{Ord}\left(P_{2} L_{2}\right)$ and $\operatorname{Ord}\left[P_{1}, L_{1}\right]$ are smaller than or equal to $m-3$ we can write

$$
\begin{equation*}
L=L_{1}\left(L_{2} L_{3}+P_{1}\right)+P_{3} L_{3}+P_{4} . \tag{5.13}
\end{equation*}
$$

Let $u$ be a solution of (10) and set $J_{0}=\left(m_{1}+1, \cdots, m\right)$. For $J \in \Pi$ we set

$$
\left\{\begin{array}{l}
\left\{\begin{array}{l}
u_{J}=\partial_{J} u \quad \text { if } \quad S_{J} \perp S_{J_{0}} \\
u_{J}=\partial_{j_{1}} \cdots \partial_{j_{k}}\left(L_{2} L_{3}+P_{1}\right) u
\end{array}\right.  \tag{5.14}\\
\qquad \text { if } J=\left(j_{1}, \cdots, j_{k}, m_{1}+1, \cdots, m\right) \in \Pi,
\end{array}\right.
$$

where $\partial_{J} u=u$ if $J=0$. Using Lemma 5.2 as in the proof of Proposition 5.1 we have from (5.13)

$$
\begin{align*}
\partial_{1} u_{(2, \cdots, m)} & =-P_{3} L_{3} u-P_{4} u \\
= & \sum_{J \ngtr J_{0}} a_{J} \partial_{J} u=\sum_{J \ngtr J_{0}} a_{J} u_{J},  \tag{5.15}\\
\text { with } & a_{J} \in G^{(k)}\left([0, T] ; S_{G(\kappa)}^{0}\right)
\end{align*}
$$

where we denote $J \supset J^{\prime}$ if $J, J^{\prime} \in \Pi$ satisfy $S_{J} \supset S_{J^{\prime}}$. Let $J$ be an element of $\Pi$ for which the set $S_{J}$ contains just $\tilde{p}_{0} \equiv m_{2}+m_{3}-1$ elements of $\left\{m_{1}+1, \cdots, m\right\}$.

Then, denoting $J=\left(j_{1}, \cdots, j_{p-\tilde{p}_{0}}, j_{p-\tilde{p}_{0}+1}, \cdots, j_{p}\right)$ with $j_{1}, \cdots, j_{p-\tilde{p}_{0}} \in\left\{1, \cdots, m_{1}\right\}$ and $j_{k} \geqq m_{1}+1$ for $k>p-\widetilde{p}_{0}$, we have from (5.10)

$$
\begin{align*}
\partial_{j_{0}} \partial_{J} & =\partial_{j_{1}} \cdots \partial_{j_{p}-\tilde{p}_{0}} L_{2} L_{3}+\sum_{\tilde{J} \in I} a_{\widetilde{J}} \partial_{\tilde{J}}  \tag{5.16}\\
& =\partial_{j_{1}} \cdots \partial_{j_{p-\tilde{p}_{0}}}\left(L_{2} L_{3}+P_{1}\right)+\sum_{\tilde{J} \in I I} a_{\widetilde{J}}^{\prime} \partial_{\widetilde{J}}
\end{align*}
$$

with $a_{\tilde{J}}$ and $a_{\tilde{J}}^{\prime} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)$. Here we used the fact that the order of $\partial_{j_{1}} \cdots \partial_{j_{p-\tilde{p}_{0}}} P_{1}$ is smaller than or equal to $m-3$. Hence, we can reduce the problem (10) to (7) by (5.15), (5.10) and (5.16). This concludes the proof of Case i).

Proof of Case i)'. We add to the set $\Pi$ the set $\Pi^{\prime}=\left\{J=\left(j_{1}, \cdots, j_{p}\right.\right.$, $\left.m_{1}+1, \cdots, \bar{m}_{2}\right) ; j_{1}<j_{2}<\cdots<j_{p^{\prime}}, j_{1}, \cdots, j_{p^{\prime}} \in\left\{1, \cdots, m_{1}, \bar{m}_{2}+1, \cdots, m\right\},\left\{1, \cdots, m_{1}\right\} \nsubseteq$ $\left.S_{J},\left\{\bar{m}_{2}+1, \cdots, m\right\} \nsubseteq S_{J}\right\}$ and we define $u_{J}$ as in (5.14). Then we have

$$
\begin{align*}
\partial_{1} u_{(2, \cdots, m)}= & -P_{2} L_{2} u-P_{3} L_{3} u-P_{4} u  \tag{5.15}\\
= & \sum_{J>J_{0}} a_{J} u_{J} \\
& \quad \text { for } a_{J} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right),
\end{align*}
$$

where $J$ belong to $\Pi \cup \Pi^{\prime}$. Consequently, in view of the proof of Case i) it suffices to derive equations for $u_{J}$ with $J \in \Pi^{\prime}$. Let $J=\left(j_{1}, \cdots, j_{k}, m_{1}+1, \cdots\right.$, $\left.\bar{m}_{2}\right) \in \Pi^{\prime}$ and $j_{0}=\max \left\{j ; 1 \leqq j \leqq m, j \notin S_{J}\right\}$. For the case where the number of elements in $\left\{j_{1}, \cdots, j_{k}\right\} \cap\left\{\bar{m}_{2}+1, \cdots, m\right\}$ is smaller than $m_{3}-1$, from the same discussion to prove (5.10) we have

$$
\begin{align*}
\partial_{j_{0}} \partial_{J}= & \sum_{\tilde{J} \in \mathbb{I}^{\prime}} a_{\widetilde{J}} \partial_{\tilde{J}} u+a L_{2} u  \tag{5.17}\\
& \text { for } \quad a_{\tilde{J}}, a \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right) .
\end{align*}
$$

So, we assume the number of elements in $\left\{j_{1}, \cdots, j_{k}\right\} \cap\left\{\bar{m}_{2}+1, \cdots, m\right\}$ is equal to $m_{3}-1$. Then, we have from $\left[L_{2}, L_{3}\right]=B_{2} L_{2}+B_{3} L_{3}+B_{4}$

$$
\begin{align*}
\partial_{j_{0}} \partial_{J} u= & \partial_{J_{1}} L_{3} L_{2} u+\sum_{\tilde{J} \in \Pi^{\prime}} a_{\widetilde{J}} \partial_{\widetilde{J}} u+a L_{2} u \\
= & \partial_{J_{1}}\left(L_{2} L_{3}+P_{1}\right) u+\partial_{J_{1}} B_{2} L_{2} u+\partial_{J_{1}} B_{3} L_{3} u \\
& +\partial_{J_{1}}\left(B_{4}-P_{1}\right) u+\sum_{\tilde{J} \in \Pi^{\prime}} a_{\widetilde{J}} \partial_{\widetilde{J}} u+a L_{2} u  \tag{5.18}\\
= & \partial_{J_{1}}\left(L_{2} L_{3}+P_{1}\right) u+\sum_{\tilde{J} \in \mathbb{I} \cup \Pi^{\prime}} a_{\widetilde{J}}^{\prime} \partial_{\widetilde{J}} u,
\end{align*}
$$

where $J_{1} \in \Pi_{k+1-m_{3}}$ with $S_{J_{1}}=\left\{j_{1}, \cdots, j_{k}\right\} \backslash\left\{\bar{m}_{2}+1, \cdots, m\right\}$ and $a_{\tilde{J}}, a, a_{\tilde{J}}^{\prime} \in$ $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)$. Here we used Lemma 5.2 to represent $\partial_{J_{1}} B_{2} L_{2}+\partial_{J_{1}} B_{3} L_{3}+$ $\partial_{J_{1}}\left(B_{4}-P_{1}\right)$ as a linear combination of $\partial_{\tilde{J}}\left(\tilde{J} \in \Pi \cup \Pi^{\prime}\right)$ with coefficients in $G^{(\kappa)}\left([0, T] ; S_{G(k)}^{0}\right)$. Combining (5.15)', (5.17), (5.18) and the results in the
proof of Case i), we obtain the systemization for Case i)' and we can concludes the proof of this case.

Remark. The condition that for any $j \in\left\{m_{1}+1, \cdots, \bar{m}_{2}\right\}$ and $j^{\prime} \in\left\{\bar{m}_{2}+1\right.$, $\cdots, m\}$ the equation

$$
\left[\partial_{j}, \partial_{j^{\prime}}\right]=a_{j j^{\prime}} \partial_{j}+b_{j j^{\prime}} \partial_{j^{\prime}}+c_{j j^{\prime}}
$$

holds with symbols $a_{j j^{\prime}}, b_{j j^{\prime}}$ and $c_{j j^{\prime}}$ in $G^{(k)}\left([0, T] ; S_{G(k)}^{0}\right)$ implies that $\left[L_{2}, L_{3}\right]=$ $B_{2} L_{2}+B_{3} L_{3}+B_{4}$ in the condition for the case i)' if we admit $B_{j}$ to be pseudodifferential operators with respect to $x$.

Proofs of Case ii) and iv). Proofs are the direct consequence of Theorem 2. Indeed, in the case ii) (resp. iv)) the operator $L$ can be written as $L_{1} L_{2} L_{3}+P_{4}, P_{4}=\sum_{j=0}^{m-q} a_{j} D_{t^{\prime}}^{j}, a_{j} \in G^{(\kappa)}\left([0, T] ; S_{G}^{m-q-j}\right)^{\prime-j}$, with $q=2$ (resp. $q=1$ ).

Proof of Case iii). Let $\left\{j_{1}, \cdots, j_{p}\right\}$ be a subset of $\{1, \cdots, m\}$. By induction on $p$ it is easy to see

$$
\begin{align*}
& \partial_{j_{1}} \cdots \partial_{j_{q-1}} \partial_{j_{q}} \partial_{j_{q+1}} \cdots \partial_{j_{q}} \\
& =\partial_{j_{1}} \cdots \partial_{j_{q-1}} \partial_{j_{q+1}} \partial_{j_{q}} \partial_{j_{q+2}} \cdots \partial_{j_{p}}+\sum_{\substack{\begin{subarray}{c}{p^{\prime} \leqslant \text { Ip }^{\prime}-2} }}\end{subarray}} \Lambda^{\left[\left(p^{\prime}-p^{\prime}\right) / 2\right]} a_{J} \partial_{J}  \tag{5.19}\\
& \text { with } \quad a_{J} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)
\end{align*}
$$

because $\left[\partial_{j}, \partial_{k}\right],\left[\partial_{j^{\prime}},\left[\partial_{j}, \partial_{k}\right]\right], \cdots$ belong to $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right)$. Set $\tilde{\Pi}=\bigcup_{p=0}^{m-1} \Pi_{p}^{o}$ $\left(\Pi_{0}^{o}=\{0\}\right)$. We shall prove

$$
\begin{align*}
L=\partial_{1} \partial_{2} \cdots \partial_{m}+\sum_{J \in \mathbb{K}_{m-1}} a_{J} \partial_{J}+\sum_{\substack{J \in \mathbb{I}_{p} \\
p \leqq m-2}} \Lambda^{[(m-p) / 2]} a_{J} \partial_{J}  \tag{5.20}\\
\text { with } \quad a_{J} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right) .
\end{align*}
$$

As in the beginning of the proof of Proposition 5.1 we may assume that the order of $P_{4}\left(\operatorname{resp} . P_{j}, j=1,2,3\right)$ with respect to $D_{t}$ is smaller than or equal to $m-3$ (resp. $m-m_{j}-2$ ) by adding the second term of the right hand side of (5.20). Next, we apply Lemma 5.2 for $\Lambda^{-1} P_{j}(j=1,2,3,4)$ and use (5.20) for the terms of the form $\partial_{j_{1}} \cdots \partial_{j_{\tilde{p}}} L_{k}$ with $k=1,2$ and $\tilde{p} \leqq m-m_{k}-2$. Then, $P_{1} L_{1}+P_{2} L_{2}+P_{3} L_{3}+P_{4}$ can be written as the linear combination of $\Lambda^{(m-p) \sigma} \partial_{J}$ $\left(J \in \Pi_{p}, p \leqq m-2\right)$ with coefficients in $G^{(\kappa)}\left([0, T] ; S_{G(k)}^{0}\right)$. Let $J \in \Pi_{p}^{o}$ $(0 \leqq p \leqq m-1)$ and set $j_{0}=\max \left\{j ; 1 \leqq j \leqq m, j \notin S_{J}\right\}$. By means of (5.19) we have

$$
\begin{array}{r}
\partial_{j_{0}} \Lambda^{(m-1-p) \sigma} \partial_{J}=\Lambda^{\sigma}\left(\Lambda^{(m-1-(p+1) \sigma} \partial_{J^{\prime}}\right)+\sum_{\tilde{J} \in \mathbb{I}_{p} \prime} \Lambda^{\sigma} a_{\tilde{J}} \Lambda^{\left(m-1-p^{\prime}\right) \sigma} \partial_{\tilde{J}}  \tag{5.21}\\
\text { with } \quad a_{\widetilde{J}} \in G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{0}\right) .
\end{array}
$$

Here $J^{\prime} \in \Pi_{p+1}$ satisfies $S_{J^{\prime}}=S_{J} \cup\left\{j_{0}\right\}$. The conjunction of (5.20) and (5.21)
shows that the vector $U=\left(u_{J}\right)_{J \in \tilde{\mathrm{I}}}$ defined by (5.7) with $\Pi_{p}$ replaced by $\Pi_{p}^{o}$ satisfies $\mathcal{L} U=0$ for a system $\mathcal{L}$ of the form (1). This proves the reduction of (10) to (7) for Case iii).
Q.E.D.

As shown in Theorem 3 it seems to be very difficult to find the conditions on lower order terms of a hyperbolic operator which the problem (10) is reduced to an equivalent problem (7) of a hyperbolic system (1) with a given $\sigma(<1)$.

## 6. Equivalence of two estimates

In this section we assume that characteristic roots $\lambda_{j}(t, x, \xi)$ of $\mathcal{L}$ belong to $\mathscr{B}^{\infty}\left([0, T] ; S^{1}\right)$ instead of $G^{(\kappa)}\left([0, T] ; S_{G(\kappa)}^{1}\right)$ and are homogeneous for $|\xi| \geqq 1$. Set

$$
\left\{\begin{array}{l}
p(\tilde{x}, \tilde{\xi})=\sum_{j=1}^{l} p_{j}(\tilde{x}, \tilde{\xi}),  \tag{6.1}\\
p_{j}(\tilde{x}, \tilde{\xi})=\tau-\lambda_{j}(t, \tilde{x}, \tilde{\xi}), \quad j=1, \cdots, l,
\end{array}\right.
$$

where $\tilde{x}=(t, x)$ and $\tilde{\xi}=(\tau, \xi)$. In what follows we write $z=(\tilde{x}, \tilde{\xi}) \in T^{*}\left(R^{n+1}\right) \simeq$ $R^{2 n+2}$ and $\delta z=(\delta \tilde{x}, \delta \tilde{\xi}) \in T_{z}\left(T^{*}\left(R^{n+1}\right)\right) \simeq R^{2 n+2}$.

For the case where $p(z)$ has the form (6.1), we shall define the "flows" $K_{z}^{+}\left(z \in T^{*}\left(R^{n+1}\right) \backslash 0\right.$ following [22] and [21]: We first define the localization $p_{z}(\delta z)$ at $z \in T^{*}\left(R^{n+1}\right) \backslash 0$ by

$$
p(z+s \delta z)=s^{\mu}\left(p_{z}(\delta z)+o(1)\right) \quad \text { as } \quad s \rightarrow 0,
$$

where $p_{z}(\delta z) \equiv 0$ (in $\delta z$ ) is a homogeneous polynomial of $\delta z \in T_{z}\left(T^{*}\left(R^{n+1}\right)\right.$ ). Since $p$ has the form (6.1) the localization $p_{z}(\delta z)$ is simply given by

$$
\begin{equation*}
p_{z}(\delta)=\left(\prod_{j \notin \Sigma_{z}} p_{j}(z)\right) \prod_{j \in \Sigma_{z}}\left(\nabla_{\tilde{x}} p_{j}(z) \cdot \delta \tilde{x}+\nabla_{\tilde{\xi}} p_{j}(z) \cdot \delta \tilde{\xi}\right), \tag{6.2}
\end{equation*}
$$

where $\Sigma_{z}$ is a maximal subset of $\{1, \cdots, l\}$ satisfying $z \in \bigcap_{j \in \Sigma_{z}} p_{j}^{-1}(0)$. Here $p_{j}^{-1}(0)=\left\{z \in T^{*}\left(R^{n+1}\right) \backslash 0 ; p_{j}(z)=0\right\}$. Let $\Gamma_{z}$ denote the connected component of $\left\{\delta z \in T_{z}\left(T^{*}\left(R^{n+1}\right)\right) ; p_{z}(\delta z) \neq 0\right\}$ which contains $(0 ; 1,0, \cdots, 0)$. Then it follows from (6.2) that

$$
\begin{equation*}
\Gamma_{z}=\bigcap_{j \in \Sigma_{z}}\left\{\delta z ; \sigma\left(H_{p_{j}}(z), \delta z\right)>0\right\} \tag{6.3}
\end{equation*}
$$

where $\sigma\left(\delta z^{\prime}, \delta z\right)=\delta \tilde{x}^{\prime} \cdot \delta \tilde{\xi}-\delta \tilde{\xi}^{\prime} \cdot \delta \tilde{x}$ and $H_{p_{j}}(z)$ denotes $\left(\nabla_{\tilde{\xi}} p_{j}(z),-\nabla_{\tilde{x}} p_{j}(z)\right)$. Set

$$
\Gamma_{z}^{\sigma}=\left\{\delta z ; \sigma\left(\delta z, \delta z^{\prime}\right) \geqq 0 \quad \text { for any } \quad \delta z^{\prime} \in \Gamma_{z}\right\}
$$

Then by means of (6.3) we have

$$
\begin{equation*}
\Gamma_{z}^{\sigma}=\left\{\sum_{j \in \Sigma_{z}} \alpha_{j} H_{p_{j}}(z) ; \alpha_{j} \geqq 0\right\} \tag{6.4}
\end{equation*}
$$

Now, we define $K_{z}^{+}$as

$$
\begin{align*}
K_{z}^{+}= & \left\{z(s) \in T^{*}\left(R^{n+1}\right) ;\{z(t)\}\right. \text { is Lipschitz continuous curve }  \tag{6.5}\\
& \text { satisfying } \left.(d / d s) z(s) \in \Gamma_{z(s)}^{\sigma} \text { (a.e. s) and } z(0)=z, s \geqq 0\right\}
\end{align*}
$$

for $z \in T^{*}\left(R^{n+1}\right) \backslash 0$.
Theorem 4. Let $V$ be a closed conic set in $T^{*}\left(R^{n}\right) \backslash 0$ and let $t_{0} \in(0, T]$. Then we have

$$
\begin{equation*}
\Gamma\left(t_{0} ; V\right)=\left\{\pi\left(K_{z}^{+} \cap\left\{t=t_{0}\right\}\right) ; z \in \pi^{-1}(V) \cap\{t=0\} \cap p^{-1}(0)\right\}, \tag{6.6}
\end{equation*}
$$

where $\pi$ is the natural projection from $T^{*}\left(R_{\tilde{x}}^{n+1}\right)$ to $T^{*}\left(R_{x}^{n}\right)$ and $p^{-1}(0)=\bigcup_{j=1}^{l} p_{j}^{-1}(0)$.
An inclusion relation $\Gamma\left(t_{0} ; V\right) \supset\{\cdot\}$ was proved in Theorem 4.4 of [22]. We remark that the assumption (L.2) of [21] is verified because $\lambda_{j}(t, x, \xi) \in$ $\mathscr{B}^{\infty}\left([0, T] ; S^{1}\right)$. So, in what follows we shall show another inclusion relation. Suppose that $\delta_{0} \in \Gamma\left(t_{0} ; V\right)$. Then for any $\varepsilon>0$ there exists an $\varepsilon$-admissible trajectory linking $\delta_{0}$ and $\rho_{\mathrm{e}} \in V$. Taking a subsequence of $\left\{\rho_{\mathrm{e}}\right\}_{\mathrm{e}>0}$, if necessary, we may assume that $\rho_{\mathrm{g}}$ converges to a point $\rho_{0} \in V$ because $V$ is closed. It follows from Proposition 4.1 that for any $\varepsilon>0$ there exists an $\varepsilon$-admissible trajectory issuing from $\rho_{0}$ whose end point $\delta_{\varepsilon}$ converges to $\delta_{0}$. From the $\varepsilon$-admissible trajectory $\left\{(t, x(t), \xi(t)) ; t \in\left[0, t_{0}\right]\right\} \subset R_{t} \times T^{*}\left(R_{x}^{n}\right)$ we make a lift $\tilde{C_{8}} \equiv\left\{(t, x(t), \tau(t), \xi(t)) ; t \in\left[0, t_{0}\right]\right\} \subset T^{*}\left(R_{x}^{n+1}\right)$ by setting $\tau(t)=\lambda_{j}(t, x(t), \xi(t))$ if $\{(t, x(t), \xi(t))\}$ is the bicharacteristic curve with respect to $\lambda_{j}$. It is clear that $\widetilde{C}_{\mathrm{e}} \subset p^{-1}(0)$. Taking a subsequence $\left\{\tilde{C}_{\mathrm{\varepsilon}}\right\}_{\mathrm{g}>0}$, if necessary, we may assume that the initial point of $\tilde{C}_{\varepsilon}$ for any $\varepsilon>0$ equals a point $z^{0} \in T^{*}\left(R^{n+1}\right) \backslash 0$ with $\pi\left(z^{0}\right)=\rho_{0}$. Similarly we may assume that the end point $z_{\varepsilon}$ of $\tilde{C}_{\varepsilon}$ converges to a point $z_{0} \in T^{*}\left(R^{n+1}\right) \backslash 0$ with $\pi\left(z_{0}\right)=\delta_{0}$. Summing up, for the proof of (6.6) it suffices to show

$$
\begin{equation*}
z_{0} \in K_{z}^{+} . \tag{6.7}
\end{equation*}
$$

In order to show this we need to define the set $K_{z^{0} 0, t_{0}}^{+}(h)$ for $h>0$ which approximates $K_{2}^{+}$, following [22] and [21]. Let $K$ be a compact neighborhood of $z^{0}$ in $T^{*}\left(R^{n+1}\right)$. We assume that $K$ is large enough to contain all lifts of $\varepsilon$-admissible trajectories for $\varepsilon<1$ from $z^{0}$. For $h>0$ and $z \in K$, there exists a compact set $M(z, h)$ in $\Gamma_{z}$ such that $(0 ; 1,0, \cdots, 0) \in \dot{M}(z, h)$ and

$$
\Gamma_{z}^{\sigma} \subset M(z, h)^{\sigma} \subset\left(\Gamma_{z}^{\sigma}\right)_{h},
$$

where $\mathscr{M}$ denotes the interior of $M$ and $(\Gamma)_{h}$ is defined by

$$
(\Gamma)_{h} \equiv\left\{\delta z ; \delta z=0 \text { or }\left.| | \delta z\right|^{-1} \delta z-\left|\delta z^{\prime}\right|^{-1} \delta z^{\prime} \mid<h \text { for some } \delta z^{\prime} \in \Gamma\right\} .
$$

Here we take $M(z, h)=\{0\}$ if $\Gamma_{z}^{\sigma}=\{0\}$, that is, if $z \notin p^{-1}(0)$. By Theorem 2.3
of [22], for each $h>0$ and $z \in K$ there exists $r(z, h)>0$ such that $r(z, h)<h$ and

$$
\begin{equation*}
M(z, h) \subset \Gamma_{z^{1}} \quad \text { for } \quad z^{1} \in U(z, h) \equiv\left\{z^{1} ;\left|z^{1}-z\right|<r(z, h)\right\} . \tag{6.8}
\end{equation*}
$$

(In our special case, this fact follows easily from (6.3)). Since $K$ is compact, there exists a finite number of $z^{h, j} \in K(1 \leqq j \leqq N(h))$ such that $K \subset \bigcup_{j=1}^{N(h)} U^{\prime}\left(z^{k, j}, h\right)$, where $U^{\prime}(z, h)=\left\{z^{\prime} ;\left|z^{\prime}-z\right| \leqq r(z, h) / 2\right\}$. We remark that $z \notin p^{-1}(0)$ for any $z \in U^{\prime}\left(z^{h, j}, h\right)$ if $z^{h, j} \notin p^{-1}(0)$. In fact, if $z^{h, j} \notin p^{-1}(0)$ then we have $M\left(z^{h, j}, h\right)^{\sigma}$ $=\{0\}$ and hence, it follows from (6.8) that for $z \in U^{\prime}\left(z^{h, j}, h\right)$ we have $\Gamma_{z}^{\sigma}=\{0\}$, that is, $z \notin p^{-1}(0)$. Now, we define $K_{z^{0}, t_{0}}^{+}(h)$ as follows: A point $z \in K \cap$ $\left\{0 \leqq t \leqq t_{0}\right\}$ belongs to $K_{z^{0}, t_{0}}^{+}(h)$ if there exist $j_{0}, \cdots, j_{\nu}$ and $\tilde{z}^{1}, \cdots, \tilde{z}^{\nu-1}$ such that $\tilde{z}^{k} \in U^{\prime}\left(z^{h, j_{k}}, h\right)(0 \leqq k \leqq \nu)$ and

$$
\begin{equation*}
\tilde{z}^{k+1}-\tilde{z}^{k} \in M\left(z^{h} \cdot j_{k}, h\right)^{\sigma} \cap\{\delta z ;|\delta z|<\rho(h)\} \quad(0 \leqq k \leqq \nu-1), \tag{6.9}
\end{equation*}
$$

where $\tilde{z}^{0}=z^{0}, z^{\nu}=z$ and $\rho(h)=\min _{1 \leqq j \leqq N(h)} r\left(z^{h, j}, h\right) / 2$. We remark that $K_{z^{0}, t_{0}}^{+}(h)$ is well-defined because the assumption (L-2) of [22] is valid (see pp. 1160 in [22]).

Proposition 6.1 (see Theorem 2.4 of [22] and Theorem 3.3 of [21]). It follows that

$$
\begin{equation*}
\bigcap_{h>0} \overline{K_{2^{0}, t_{0}}^{+}(h)}=K_{2^{0}}^{+} \cap\left\{0 \leqq t \leqq t_{0}\right\}, \tag{6.10}
\end{equation*}
$$

where $\bar{K}$ denotes the closure of $K$.
By means of this proposition, for the proof of (6.7) it suffices to show

$$
\left\{\begin{array}{l}
\text { for auy } h>0 \text { there exists an } \varepsilon_{0}>0 \text { such that the end }  \tag{6.11}\\
\text { poiut } z_{\mathrm{\varepsilon}} \text { of the lift } \tilde{C_{\varepsilon}} \text { belongs to } K_{z^{0}, t_{0}}^{+}(h) \text { if } \varepsilon \leqq \varepsilon_{0} .
\end{array}\right.
$$

Lemma 6.2. Let $\left\{z(t)=(t, x(t), \tau(t), \xi(t)) ; t \in\left[0, t_{0}\right]\right\}$ be a lift of an $\varepsilon$ admissible trajectory. Then for any two continuous points $z\left(s_{1}\right)$ and $z\left(s_{2}\right)$ on the lift we have

$$
\begin{equation*}
\left|z\left(s_{1}\right)-z\left(s_{2}\right)\right| \leqq C\left|s_{1}-s_{2}\right|+(l-1) \varepsilon, \tag{6.12}
\end{equation*}
$$

where $C$ it a positive constant and $l$ is the size of the system $\mathcal{L}$.
Proof. Let $\pi_{0}$ denote the natural projection from $T^{*}\left(R^{n+1}\right)$ to $R_{t} \times T^{*}\left(R_{x}^{n}\right)$. It is clear that

$$
\begin{equation*}
\left|\pi_{0}\left(z\left(s_{1}\right)\right)-\pi_{0}\left(z\left(s_{2}\right)\right)\right| \leqq C\left|s_{1}-s_{2}\right| . \tag{6.13}
\end{equation*}
$$

Assume $\tau\left(s_{k}\right)=\lambda_{j_{k}}\left(\pi_{0}\left(z\left(s_{k}\right)\right), k=1,2\right.$. If $j_{1}=j_{2}$ then (6.12) follows from the continuity of $\lambda_{j}$. Assume that $j_{1} \neq j_{2}$. For simplicity we consider the case for $l=2$. By taking a discontinuous point $z\left(t^{\prime}\right)$ between $z\left(s_{1}\right)$ and $z\left(s_{2}\right)$ we estimate

$$
\begin{aligned}
& \left|\tau\left(s_{2}\right)-\tau\left(s_{1}\right)\right| \\
& \quad \leqq\left|\tilde{\lambda}_{j_{1}}\left(s_{1}\right)-\tilde{\lambda}_{j_{1}}\left(t^{\prime}\right)\right|+\left|\tilde{\lambda}_{j_{1}}\left(t^{\prime}\right)-\tilde{\lambda}_{j_{2}}\left(t^{\prime}\right)\right|+\left|\tilde{\lambda}_{j_{2}}\left(t^{\prime}\right)-\tilde{\lambda}_{j_{2}}\left(s_{2}\right)\right| \\
& \quad \leqq C\left(\left|s_{1}-t^{\prime}\right|+\left|t^{\prime}-s_{2}\right|\right)+\varepsilon
\end{aligned}
$$

by using (6.13), where $\tilde{\lambda}_{j}(s)=\lambda_{j}\left(\pi_{0}(z(s))\right.$. In the general case for $l \geqq 3$ we can also estimate the difference between $\tau\left(s_{1}\right)$ and $\tau\left(s_{2}\right)$ by taking $l-1$ discontinuous points, at most, between $z\left(s_{1}\right)$ and $z\left(s_{2}\right)$. Then, we get (6.12). Q.E.D.

Let $\pi_{1}$ denote the natural projection from $T^{*}\left(R^{n+1}\right)$ to $R_{t}$. Set

$$
\mathscr{H}_{z}=\left\{\sum_{j=1}^{l} \alpha_{j} H_{p_{j}}(z) ; \sum_{j}^{l} \alpha_{j}=1, \alpha_{j} \geqq 0\right\} \quad \text { (cf. (6.4)) } .
$$

Since the Hamilton field $H_{p_{j}}(z)$ depends only on $\pi_{0}(z)$ it follows that if $z^{k}$ ( $k=1,2$ ) are two points of the lift $C_{\varepsilon}$ of an $\varepsilon$-admissible trajectory then for any $v_{2}{ }^{1} \in \mathscr{H}_{2^{1}}$ there exists a $v_{2}{ }^{2} \in \mathscr{H}_{2}{ }^{2}$ such that

$$
\begin{equation*}
\left|v_{2}{ }^{2}-v_{z^{1}}\right| \leqq C\left|\pi_{1}\left(z^{2}-z^{1}\right)\right|, \tag{6.14}
\end{equation*}
$$

where $C$ is a constant independent of the lifts $C_{\mathrm{e}}$ of $\varepsilon$-admissible trajectories.
Lemma 6.3. For any $h>0$ there is an $\varepsilon(h, l)>0$ satisfying the following property: Assume that $z^{k}(k=1,2)$ are two continuous points on a lift $\widetilde{C_{\varepsilon}}$ of an $\varepsilon$-admissible trajectory such that $\varepsilon \leqq \varepsilon(h, l)$ and

$$
h \leqq \pi_{1}\left(z^{2}-z^{1}\right) \leqq 2 h .
$$

Then there exists a $v_{2^{1}} \in \mathscr{H}_{2^{1}}$ such that

$$
\begin{equation*}
\left|\left(z^{2}-z^{1}\right) / \pi_{1}\left(z^{2}-z^{1}\right)-v_{z^{1}}\right| \leqq C h, \tag{6.15}
\end{equation*}
$$

where $C$ is a constant independent of $h$ and $\tilde{C}_{8}$.
Remark. Let $\Sigma$ be a subset of $\{1, \cdots, l\}$ such that the part of $\widetilde{C_{\mathrm{g}}}$ between $z^{1}$ and $z^{2}$ is composed of bicharacteristic curves with respect to $\lambda_{j}, j \in \Sigma$. Then we can replace $\mathscr{H}_{z^{1}}$ in the lemma by $\mathscr{H}_{z^{1}}^{\prime} \equiv\left\{\sum_{j \in \Sigma} \alpha_{j} H_{p_{j}}\left(z^{1}\right) ; \alpha_{j} \geqq 0, \sum_{j \in \Sigma} \alpha_{j}=1\right\}$.

Proof. We shall prove the lemma by the inductive method on $l$. The case for $l=1$ is trivial. So, we assume $l \geqq 2$ and suppose that the conclusion holds until $l-1$. Take continuous points $\tilde{z}^{1}, \cdots, \widetilde{z}^{\nu-1}$ on $\widetilde{C}_{\varepsilon}$ between $z^{1}$ and $z^{2}$ such that

$$
\begin{gathered}
h^{3} \leqq \pi_{1}\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) \leqq 2 h^{3}, \quad k=1, \cdots, \nu-1 \\
\left(\tilde{z}^{1}=z^{1} \text { and } \tilde{z}^{\nu}=z^{2}\right) .
\end{gathered}
$$

Setting $h_{o}=\pi_{1}\left(z^{2}-z^{1}\right)$ and $h_{k}=\pi_{1}\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right)(k=1, \cdots, \nu-1)$ we write

$$
\begin{equation*}
\left(z^{2}-z^{1}\right) / h_{o}=\sum_{k=1}^{\nu}\left(h_{k} / h_{o}\right)\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) / h_{k} . \tag{6.16}
\end{equation*}
$$

Suppose that, for any $k$, the part of $\tilde{C}_{\varepsilon}$ between $\tilde{z}^{k}$ and $\tilde{z}^{k+1}$ is composed of bicharacteristic curves with respect to, at most, $l-1$ elements of $\left\{\lambda_{j}\right\}$. Then it follows from the hypothesis of the induction that if $\varepsilon \leqq \varepsilon\left(h^{3}, l-1\right)$ we have

$$
\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) / h_{k}-v_{\tilde{z}^{k}}=O\left(h^{3}\right) \quad \text { for } \quad v_{\tilde{z}^{k}} \in \mathcal{H}_{\tilde{z}^{k}}
$$

Here $v=O\left(h^{3}\right)$ means $|v| \leqq C h^{3}$ with a constant $C$ independent of $h$ and the choice of the lift $\widetilde{C}_{\varepsilon}$ of $\varepsilon$-admissible trajectory. By using (6.14) we have for some $v_{z^{1}}^{k} \in \mathcal{H}_{2^{1}}$

$$
\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) / h_{k}-v_{z^{1}}^{k}=O(h) .
$$

Consequently, using $\sum h_{k} / h_{o}=1$ we obtain (6.15) in this case.
Consider the case that the part of $\tilde{C}_{\varepsilon}$ between $\tilde{z}^{k}$ and $\tilde{z}^{k+1}$ for some $k$ is composed of bicharacteristic curves with respect to full elements of $\left\{\lambda_{j}\right\}_{j=1}^{l}$. We denote by $k_{1}$ the minimum of such $k$ and by $k_{2}$ the integer $k_{2}$ such that $k_{2}-1$ is the maximum of the $k$ 's stated above. Now, we write (6.16) as

$$
\begin{align*}
\left(z^{2}-z^{1}\right) / h_{o}= & \sum_{k=1}^{k_{1}}\left(h_{k} / h_{o}\right)\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) / h_{k}  \tag{6.16}\\
& +\left(\tilde{z}^{k_{2}}-\tilde{z}^{k_{1}}\right) / h_{o}+\sum_{k=k_{2}}^{v}\left(h_{k} / h_{o}\right)\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) / h_{k}
\end{align*}
$$

Suppose that $\pi_{1}\left(\tilde{z}^{k_{2}}-\tilde{\boldsymbol{z}}^{k_{1}}\right) \leqq h^{2}$. Then it follows from Lemma 6.2 that we have $\left|z^{k_{2}}-z^{k_{1}}\right|=O\left(h^{2}\right)$ if $\varepsilon \leqq \min \left(\varepsilon\left(h^{3}, l-1\right), h^{2}\right)$. Hence, the second term of the right hand side of (6.16)' is estimated by the constant times of $h$. So, we get (6.15) by using the discussions of the preceding paragraph and $\sum_{k=k_{1}+1}^{k_{2}-1}\left(h_{k} / h_{o}\right) v=$ $O(h)$ for any $v \in \mathscr{H}_{2}{ }^{1}$.

Assume that $\pi_{1}\left(\tilde{z}^{k_{2}}-\tilde{z}^{k_{1}}\right) \geqq h^{2}$ and let $\varepsilon \leqq \min \left(\varepsilon\left(h^{3}, l-1\right), h^{3}\right)$. To complete the proof it suffices to show

$$
\begin{equation*}
\left(\tilde{z}^{k_{2}}-\tilde{z}^{k_{1}}\right) \mid \tilde{h}-v_{\tilde{z}} \tilde{k}_{1}=O(h) \quad \text { for some } \quad v_{\tilde{z}} \tilde{k}_{1} \in \mathscr{H}_{\tilde{z}^{k_{1}}}, \tag{6.17}
\end{equation*}
$$

where $\tilde{h}=\pi_{1}\left(\tilde{z}^{k_{2}}-\tilde{z}^{k_{1}}\right)$. Since the bicharacteristic curves with respect to full elements of $\left\{\lambda_{j}\right\}_{j=1}^{l}$ appear on the part of $\widetilde{C}_{\varepsilon}$ between $\tilde{z}^{k_{1}}$ and $\tilde{z}^{k_{1}+1}$ (also between $z^{k_{2}-1}$ and $z^{k_{2}}$ ), it follows from the continuity of $\lambda_{j}$ that we have

$$
\begin{gather*}
\lambda_{j}\left(\pi_{0}\left(z^{k_{i}}\right)\right)-\lambda_{j^{\prime}}\left(\pi_{0}\left(z^{k^{k}}\right)\right)=O\left(h^{3}\right), \quad i=1,2,  \tag{6.18}\\
\text { for any } j, j^{\prime} \in\{1, \cdots, l\}
\end{gather*}
$$

In order to simplify the notation below we denote $\tilde{z}^{k_{1}}$ and $\tilde{z}^{k_{2}}$ by $\tilde{z}^{1}$ and $\tilde{z}^{2}$, respectively, in what follows. Since $\pi_{0} \tilde{C}_{\mathrm{z}}$ is continuous in $R_{t} \times T^{*}\left(R_{x}^{n}\right)$ we have

$$
\begin{equation*}
\pi_{0}\left(z^{2}-\tilde{z}^{1}\right) / \tilde{h}-\sum_{j=1}^{i} \alpha_{j} \pi_{0}\left(H_{p_{j}}\left(z^{1}\right)\right)=O(\tilde{h}) \tag{6.19}
\end{equation*}
$$

for some $\alpha_{j} \geqq 0$ with $\sum \alpha_{j}=1$. Let $\pi_{1}^{\prime}$ be the natural projection from $T^{*}\left(R^{n+1}\right)$ to $R_{\tau}$, where $\tau$ is the dual variable of $t$. We shall show

$$
\begin{equation*}
\pi_{1}^{\prime}\left(\tilde{z}^{2}-\tilde{z}^{1}\right) \mid \tilde{h}-\sum_{j=1}^{l} \alpha_{j} \pi_{1}^{\prime}\left(H_{p_{j}}\left(\tilde{z}^{1}\right)\right)=\boldsymbol{O}(h) \tag{6.20}
\end{equation*}
$$

In view of $\tilde{h} \geqq h^{2}$ it follows from (6.18) that

$$
\begin{equation*}
\pi_{1}^{\prime}\left(\tilde{z}^{2}-\tilde{z}^{1}\right) / \tilde{h}=\sum_{j=1}^{l} \alpha_{j}\left(\lambda_{j}\left(\pi_{0}\left(z^{2}\right)\right)-\lambda_{j}\left(\pi_{0}\left(\tilde{z}^{1}\right)\right) / \tilde{h}+O(h)\right. \tag{6.21}
\end{equation*}
$$

Set $s_{k}=\pi_{1}\left(\tilde{z}^{k}\right)(k=1,2)$ and denote the $\varepsilon$-admissible trajectory $\pi_{0}\left(\check{C}_{\mathfrak{z}}\right)$ by $\{(t, x(t), \xi(t))\}$. Then we have

$$
\begin{align*}
& \sum_{j=1}^{l} \alpha_{j}\left(\lambda_{j}\left(\pi_{0}\left(\tilde{z}^{2}\right)\right)-\lambda_{j}\left(\pi_{0}\left(\tilde{z}^{1}\right)\right)\right) / \tilde{h} \\
& \quad=\sum_{j=1}^{l}\left(\alpha_{j} / \tilde{h}\right) \int_{s_{1}}^{s_{2}} \partial_{t} \lambda_{j}(t, x(t), \xi(t)) d t  \tag{6.22}\\
& \quad+(1 / \tilde{h}) \int_{s_{1}}^{s_{2}} \sum_{j=1}^{l} \alpha_{j}\left(\nabla_{x} \lambda_{j}(t, x(t),(t)) d x / d t\right. \\
& \left.\quad+\nabla_{\xi} \lambda_{j}(t, x(t), \xi(t)) d \xi / d t\right) d t
\end{align*}
$$

The second term of the right hand side is equal to

$$
\left(\sum_{j=1}^{l} \alpha_{j} \nabla_{x} \lambda_{j}\left(\tilde{z}^{1}\right), \sum_{j=1}^{l} \alpha_{j} \nabla_{\xi} \lambda_{j}\left(\tilde{z}^{1}\right)\right) \cdot\left(\frac{1}{\tilde{h}} \int_{s_{1}}^{s_{2}}(d x(t), d \xi(t))\right)+O(h) .
$$

It follows from (6.19) that

$$
\begin{aligned}
\frac{1}{\tilde{h}} \int_{s_{1}}^{s_{2}}(d x(t), d \xi(t)) & =\frac{1}{\tilde{h}}\left(\pi\left(\tilde{z}^{2}\right)-\pi\left(\tilde{z}^{1}\right)\right) \\
& =\sum_{j=1}^{l} \alpha_{j} \pi\left(H_{p_{j}}\left(\tilde{z}^{1}\right)\right)+O(h)
\end{aligned}
$$

Hence, we can estimate the second term of (6.22) by

$$
\left\{\sum_{j=1}^{i} \alpha_{j} \lambda_{j}\left(z^{1}\right)+O(h), \sum_{j=1}^{i} \alpha_{j} \lambda_{j}\left(z^{1}\right)+O(h)\right\}=O(h)
$$

where $\{$,$\} denotes the Poisson bracket in T^{*}\left(R_{x}^{n}\right)$. Note that the first term of the right hand side of (6.22) equals $\sum \alpha_{j} \pi_{1}^{\prime}\left(H_{p_{j}}\left(z^{1}\right)\right)+O(h)$. Hence we get (6.20) from (6.21). This concludes the proof of (6.17), and hence, the proof of (6.15).
Q.E.D.

Proof of (6.11). Let $h$ be a fixed positive number. As in the proof of Lemma 3.2 of [21], for any $j \in\{1, \cdots, N(h)\}$ we can find an $h(j)>0$ such that

$$
\begin{align*}
& \left(\Gamma_{z}^{\sigma}\right)_{h(j)} \subset  \tag{6.23}\\
& \quad M\left(z^{h, j}, h\right)^{\sigma} \\
& \text { for } z \in U^{\prime \prime}\left(z^{h, j}, h\right) \equiv\left\{z ;\left|z-z^{h, j}\right|<2 r\left(z^{h, j}, h\right) / 3\right\},
\end{align*}
$$

if $z^{h, j} \in p^{-1}(0)$ (see (3.4) of [21]). Set $h^{\prime}=\min \{\rho(h), h(1), \cdots, h(N(h))\}$ with a convention $h(j)=\infty$ if $z^{h, j} \notin p^{-1}(0)$. Let $h^{\prime \prime}$ be another positive number sufficiently smaller than $h^{\prime}$, which is determined later on. For the moment we take $h^{\prime \prime}$ as the number for which we can find a constant $\varepsilon\left(h^{\prime \prime}\right)>0$ satisfying the following;

$$
\left\{\begin{array}{l}
\text { if two points } z^{k}(k=1,2) \text { on the lift } \widetilde{C}_{z} \text { with } \varepsilon \leqq \varepsilon\left(h^{\prime \prime}\right)  \tag{6.24}\\
\text { satisfy } \pi_{1}\left(z^{1}-z^{2}\right) \leqq 2 h^{\prime \prime} \text { then }\left|z^{1}-z^{2}\right|<\rho(h) / 3 \text { holds } \\
\left(\text { and hence } z^{1} \in U^{\prime}\left(z^{h, j}, h\right) \text { implies } z^{2} \in U^{\prime \prime}\left(z^{h, j}, h\right)\right) .
\end{array}\right.
$$

Here we used Lemma 6.2.
We shall apply Lemma 6.3 by setting $h=h^{\prime \prime}$. Assume that $\varepsilon \leqq \min \left(\varepsilon\left(h^{\prime \prime}, l\right)\right.$, $\left.\varepsilon\left(h^{\prime \prime}\right)\right)$. Take continuous points $\tilde{z}^{1}, \cdots, z^{\nu-1}$ on $\tilde{C}_{\mathrm{g}}$ linking $z^{0}$ and $z_{\mathrm{g}}$ such that

$$
\begin{gathered}
h^{\prime \prime} \leqq \pi_{1}\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) \leqq 2 h^{\prime \prime}, \quad k=0, \cdots, \nu-1 \\
\left(\tilde{z}^{0}=z^{0}, \tilde{z}^{\nu}=z_{\mathrm{q}}\right) .
\end{gathered}
$$

Then it follows from Lemma 6.3 and its remark that for any $k \in\{0, \cdots, \nu-1\}$ there exist $\left\{\alpha_{j}\right\}$ with $\alpha_{j} \geqq 0$ and $\sum \alpha_{j}=1$ such that

$$
v \equiv\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) / \pi_{1}\left(\tilde{z}^{k+1}-\tilde{z}^{k}\right)-\sum_{j \in \Sigma_{k}} \alpha_{j} H_{p_{j}}\left(\tilde{z}^{k}\right)=O\left(h^{\prime \prime}\right)
$$

Here $\Sigma_{k}$ is a subset of $\{1, \cdots, l\}$ satisfying the following: The part of the lift $\tilde{C}_{\mathrm{e}}$ between $\tilde{z}^{k}$ and $\tilde{z}^{k+1}$ is composed of bicharacteristic curves with respect to $\lambda_{j}$ for $j \in \Sigma_{k}$. By means of (6.13), for any $j \in \Sigma_{k}$ there exists a point $z^{j}$ on $\widetilde{C_{\mathrm{z}}}$ between $\tilde{z}^{k}$ and $\tilde{z}^{k+1}$ such that $z^{j} \in p_{j}^{-1}(0)$ and

$$
v_{j} \equiv H_{p_{j}}\left(\tilde{z}^{k}\right)-H_{p_{j}}\left(z^{j}\right)=O\left(h^{\prime \prime}\right) .
$$

Note that $H_{p_{j}}\left(z^{j}\right) \in \Gamma_{z^{j}}^{\sigma}$. It follows from (6.24) that $z^{j} \in U^{\prime \prime}\left(z^{h, j_{k}}, h\right)$ holds if $z^{k} \in U^{\prime}\left(z^{h, j_{k}}, h\right)$. Hence, using (6.24) and the convexity of $M\left(z^{h, j_{k}}, h\right)^{\sigma}$ we have

$$
\begin{align*}
& \left|z^{k+1}-\tilde{z}^{k}\right|<\rho(h),  \tag{6.25}\\
& \left(\tilde{z}^{k+1}-\tilde{z}^{k}\right) /\left|\tilde{z}^{k+1}-\tilde{z}^{k}\right|  \tag{6.26}\\
& \quad=\sum_{j \in \Sigma_{k}} \alpha_{j}\left(\pi_{1}\left(\tilde{z}^{k+1}-z^{k}\right) /\left|\tilde{z}^{k+1}-\tilde{z}^{k}\right|\right)\left(H_{p_{j}}\left(z^{j}\right)+v_{j}+v\right) \in M\left(z^{h, j_{k}}, h\right)^{\sigma}
\end{align*}
$$

because $H_{p_{j}}\left(z^{j}\right)+v_{j}+v \in\left(\Gamma_{z^{j}}^{\sigma}\right)_{h^{\prime}} \subset M\left(z^{h, j_{k}}, h\right)^{\sigma}$ if $h^{\prime \prime}$ is sufficiently smaller than $h^{\prime}$. Then, (6.25) and (6.26) show that the end point $z_{\mathrm{z}}$ of $\tilde{C}_{\mathrm{z}}$ belongs to $K_{z^{0}, t_{0}}^{+}(h)$ if $\varepsilon \leqq \min \left(\varepsilon\left(h^{\prime \prime} ; l\right), \varepsilon\left(h^{\prime \prime}\right)\right)$, that is, we have proved the property (6.11). Q.E.D.

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