# 188. On Certain Mixed Problem for Hyperbolic Equations of Higher Order. III 

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1. Introduction. Let $R_{+}^{n}$ be the open half space $\{(x, y) ; x>0$, $\left.y \in \boldsymbol{R}^{n-1}\right\}$. We consider the mixed problem ( $P, B_{j} ; j=1, \cdots, l$ ), briefly ( $P, B_{j}$ ), for hyperbolic equations of order $m$ in $(0, T) \times \boldsymbol{R}_{+}^{n}(0<T<\infty)$ :

$$
\begin{array}{ll}
\left(P\left(D_{t}, D_{x}, D_{y}\right) u\right)(t, x, y)=f(t, x, y) & \text { in }(0, T) \times \boldsymbol{R}_{+}^{n}, \\
\left(B_{j}\left(D_{t}, D_{x}, D_{y}\right) u\right)(t, 0, y)=0(j=1, \cdots, l) & \text { in }(0, T) \times \boldsymbol{R}^{n-1}, \\
\left(D_{t}^{k} u\right)(0, x, y)=0(k=0,1, \cdots, m-1) & \text { in } \boldsymbol{R}_{+}^{n},
\end{array}
$$

where $D_{t}=\frac{\partial}{\partial t}, D_{x}=-i \frac{\partial}{\partial x}, D_{y}=\left(-i \frac{\partial}{\partial y_{1}}, \cdots,-i \frac{\partial}{\partial y_{n-1}}\right)$ and $i=\sqrt{-1}$.
The purpose of this paper is to determine the necessary and sufficient conditions for $L^{2}$-well-posedness in the following sense.

Definition. The mixed problem $\left(P, B_{j}\right)$ is $L^{2}$-well-posed if and only if there exist constants $T$ and $T^{\prime}$ with $0<T^{\prime} \leq T$ which satisfy the following condition:

For every $f \in H^{1}\left((-\infty, T) \times \boldsymbol{R}_{+}^{n}\right)$ with $f=0(t<0)$ the mixed problem ( $P, B_{j}$ ) has a unique solution $u \in H^{m}\left(\left(0, T^{\prime}\right) \times \boldsymbol{R}_{+}^{n}\right)$ so that

$$
\sum_{k=0}^{m-1} \int_{0}^{T^{\prime}}\left\|\left(D_{t}^{k} u\right)(t, \cdot, \cdot)\right\|_{m-k-1}^{2} d t \leq C \int_{0}^{T}\|f(t, \cdot, \cdot)\|_{0}^{2} d t
$$

where a constant $C$ depends only on $T$.
In § 2 we give certain necessary and sufficient conditions for $L^{2}$ -well-posedness (Theorem 1) and investigate zeros of the Lopatinskii's determinant under $L^{2}$-well-posedness (Theorem 2).

In T. Shirota and K. Asano [5] it has been shown by semi-group method that the mixed problem $\left(P, D_{x}^{2 j-1} ; j=1, \cdots, l\right)(m=2 l)$ is well posed in the $L^{2}$-sense ${ }^{1)}$ if $P(D)=P\left(D_{t}, D_{x}, D_{y}\right)$ does not contain the terms of odd order relative to $D_{x}$. As one of the applications of Theorems 1 and 2 we show that, in the case of constant coefficients, the above condition for $P(D)$ is necessary to be well posed in the $L^{2}$-sense for the above mixed problem. This assertion is found in Theorem 4 in § 3.

The details and other results will be published elsewhere.
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[^0]2. Necessary and sufficient conditions for $L^{2}$-well-poseness. Let $P(D)$ and $B_{j}(D)(j=1, \cdots, l)$ be homogenuous differential operators of order $m$ and $m_{j}\left(m_{j}<m\right)$ with constant coefficients respectively. We assume that $P(D)$ is strongly hyperbolic relative to $t$-direction and the hyperplane $x=0$ is non-characteristic for $P(D)$. Then it is easily seen that the number $l(m-l)$ of the roots $\lambda_{j}^{+}(\tau, \sigma)(j=1, \cdots, l)\left(\lambda_{k}^{-}(\tau, \sigma)\right.$ $(k=1, \cdots, m-l)$ ) in $\lambda$ of the characteristic equation $P(\tau, \lambda, \sigma)=0$ located in the upper (lower) half $\lambda$-plane is constant for any ( $\tau, \sigma$ ) with $\operatorname{Re} \tau>0$ and $\sigma \in \boldsymbol{R}^{n-1}$ respectively.

Throughout this paper we use the following Fourier-Laplace transforms and norms.

$$
\begin{aligned}
& \hat{u}(\tau, \lambda, \sigma)=\int_{0}^{\infty} d t \int_{0}^{\infty} d x \int_{R^{n-1}} e^{-\tau t-i \lambda x-i \sigma y} u(t, x, y) d y \\
& \hat{u}(\tau, x, \sigma)=\int_{0}^{\infty} d t \int_{R^{n-1}} e^{-\tau t-i \sigma y} u(t, x, y) d y \\
& \|u(t, \cdot, \cdot)\|_{k}^{2}=\sum_{j=0}^{k}\left\|\left(D_{t}^{j} u\right)(t, \cdot, \cdot)\right\|_{k-j}^{2} \\
& \|\hat{u}(\tau, \cdot, \cdot)\|_{k}^{2}=\sum_{j=0}^{k} \int_{R^{n-1}}\left(|\tau|^{2}+|\sigma|^{2}\right)^{k-j} d \sigma \int_{0}^{\infty}\left|\left(D_{x}^{j} \hat{u}\right)(\tau, x, \sigma)\right|^{2} d x(\operatorname{Re} \tau \geq \gamma>0),
\end{aligned}
$$

where $\sigma y=\sigma_{1} y_{1}+\cdots+\sigma_{n-1} y_{n-1},|\sigma|^{2}=\sigma_{1}^{2}+\cdots+\sigma_{n-1}^{2}, \gamma$ is arbitrarily fixed and $\left\|\|_{h}\right.$ is the norm in Sobolev space $H^{h}\left(\boldsymbol{R}_{+}^{n}\right)(h=0,1, \cdots)$.

We define the Lopatinskii's determinant $R(\tau, \sigma)$ as follows:

$$
\begin{aligned}
& B(\tau, \sigma)=\operatorname{det}\left(B_{1}\left(\tau, \lambda_{k}^{+}(\tau, \sigma), \sigma\right), \cdots, B_{l}\left(\tau, \lambda_{k}^{+}(\tau, \sigma), \sigma\right) ; k \downarrow 1, \cdots, l\right), \\
& R(\tau, \sigma)=B(\tau, \sigma) / \prod_{1 \leq j<k \leq l}\left(\lambda_{k}^{+}(\tau, \sigma)-\lambda_{j}^{+}(\tau, \sigma)\right) .
\end{aligned}
$$

Note that $R(\tau, \sigma)$ is analytic in $\operatorname{Re} \tau>0$ and real analytic in $\boldsymbol{R}^{n-1}$. Let $V$ be the $\operatorname{set}\left\{(\tau, \sigma) ; R(\tau, \sigma)=0, \operatorname{Re} \tau>0, \sigma \in R^{n-1}\right\}$ and $S(\tau)$ the analytic variety $V \cap\left\{(\tau, \sigma) ; \sigma \in \boldsymbol{R}^{n-1}\right\}$. Then we have $\alpha V=V$ and $\alpha S(\tau)=S(\alpha \tau)$ for every $\alpha>0$.

Applying now the Fourier-Laplace transform to the equations in the problem $\left(P, B_{j}\right)$ we obtain the boundary value problem ( $\hat{P}, \hat{B}_{j}$ ) of the ordinary differential equations depending parameters ( $\tau, \sigma$ ) with $\operatorname{Re} \tau>0$ and $\sigma \in \boldsymbol{R}^{n-1}$ :

$$
\begin{aligned}
& \left(P\left(\tau, D_{x}, \sigma\right) \hat{u}\right)(\tau, x, \sigma)=\hat{f}(\tau, x, \sigma) \quad \text { in } \boldsymbol{R}_{+}^{1}, \\
& \left(B_{j}\left(\tau, D_{x}, \sigma\right) \hat{u}\right)(\tau, 0, \sigma)=0 \quad(j=1, \cdots, l) .
\end{aligned}
$$

Let $R_{j}(\tau, x, \sigma)$ be the determinant replacing the $j$-column in $R(\tau, \sigma)$ by the transposed vector of ( $e^{i x \lambda_{1}^{+}(\tau, \sigma)}, \cdots, e^{i x \lambda_{i}^{+}(\tau, \sigma)}$ ) and $\Gamma=\Gamma(\tau, \sigma)$ a closed Jordan curve in the lower half $\lambda$-plane enclosing all the roots $\lambda_{k}^{-}(\tau, \sigma)$ ( $k=1, \cdots, m-l$ ). If $R(\tau, \sigma)$ is not zero for some $(\tau, \sigma)$ with $\operatorname{Re} \tau>0$ and $\sigma \in \boldsymbol{R}^{n-1}$, then it is well known that for every $\hat{f}(\tau, \cdot, \sigma) \in C_{0}^{\infty}\left(\boldsymbol{R}_{+}^{1}\right)$ the boundary value problem $\left(\hat{P}, \hat{B}_{j}\right)$ has a unique solution $\hat{u}(\tau, \cdot, \sigma) \in C^{\infty}\left(\overline{\boldsymbol{R}}_{+}^{1}\right)$, which is written in the form:
$\hat{u}(\tau, x, \sigma)=\frac{1}{2 \pi} \int_{0}^{\infty} G_{1}(x, s, \tau, \sigma) \hat{f}(\tau, s, \sigma) d s+\frac{1}{2 \pi} \int_{0}^{\infty} G_{2}(x, s, \tau, \sigma) \hat{f}(\tau, s, \sigma) d s$,
where $G_{1}(x, s, \sigma)=\int_{\Gamma} \frac{e^{i(x-s) \lambda}}{P(\tau, \lambda, \sigma)} d \lambda$,

$$
G_{2}(x, s, \tau, \sigma)=-\sum_{j=1}^{i} \frac{R_{j}(\tau, x, \sigma)}{R(\tau, \sigma)} \int_{\Gamma} \frac{B_{j}(\tau, \lambda, \sigma)}{P(\tau, \lambda, \sigma)} e^{-i s \lambda} d \lambda .
$$

Let $\Sigma_{+}$be the set $\left\{\left(\tau^{\prime}, \sigma^{\prime}\right) ;\left|\tau^{\prime}\right|^{2}+\left|\sigma^{\prime}\right|^{2}=1, \operatorname{Re} \tau^{\prime}>0, \sigma^{\prime} \in R^{n-1}\right\}$ and $\bar{\Sigma}_{+}$ its closure. Set $V^{\prime}=V \cap \Sigma_{+}$. When $A$ is a subset of $\Sigma_{+}$we denote the complement of $A$ in $\Sigma_{+}$by $A^{c}$. Then we have the following

Theorem 1. If $S(\tau)$ is not the whole space $R^{n-1}$ for every $\tau$ with $\operatorname{Re} \tau>0$, then the mixed problem $\left(P, B_{j}\right)$ is $L^{2}$-well-posed if and only if the following condition (I) is satisfied:

For every $\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \in\left(\bar{\Sigma}_{+}-\Sigma_{+}\right) \cup V^{\prime}$ there exist a neighbourhood $U\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)$ and a constant $C\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)$ such that for any $\left(\tau^{\prime}, \sigma^{\prime}\right)$ $\in U\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \cap \Sigma_{+} \cap V^{\prime c}$

$$
\left\|\left(D_{x}^{k} G_{2}\right)(x, s, \tau, \sigma)\right\|_{\mathcal{L}_{\left(L^{2}(s>0), L^{2}(x>0)\right)} \leq} \frac{C\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)}{\operatorname{Re} \tau^{\prime}}
$$

$$
(k=0,1, \cdots, m-1)
$$

where $\left\|\|_{\mathcal{L}_{\left(L^{2}(s>0), L^{2}(x>0)\right)}}\right.$ is the operator norm from $L^{2}(s>0)$ to $L^{2}(x>0)$.
To prove Theorem 1 we need the following lemmas. Hereafter we denote various positive constants by $C$.

Lemma 1. If a polynomial $P(\tau, \xi)$ of degree $m$ is strongly hyperbolic relative to $\tau$, then we have for any $(\tau, \xi)$ with $\operatorname{Re} \tau>0$ and $\xi \in \boldsymbol{R}^{n}$

$$
|P(\tau, \xi)|^{2} \geq C(\operatorname{Re} \tau)^{2}\left(|\tau|^{2}+|\xi|^{2}\right)^{m-1}
$$

Lemma 2. If the assumption in Theorem 1 and the condition (I) are satisfied, then for every $(\tau, \sigma) \notin V$ and $f \in H^{k+1}\left((-\infty, \infty) \times \boldsymbol{R}_{+}^{n}\right)$ $(k=0,1, \cdots)$ with $f=0(t<0)$ the boundary value problem $\left(\hat{P}, \hat{B}_{j}\right)$ has a unique solution $\hat{u}(\tau, \cdot, \sigma) \in H^{m+k}\left(\boldsymbol{R}_{+}^{1}\right)$ so that
$(\operatorname{Re} \tau)^{2}\|\hat{u}(\tau, \cdot, \cdot)\|_{m-1+k}^{2} \leq C\|\hat{f}(\tau, \cdot, \cdot)\|_{k}^{2}$ for any $\tau$ with $\operatorname{Re} \tau \geq \gamma>0$, where $\gamma$ is arbitrarily fixed and a constant $C$ depends only on $\gamma$.

Lemma 3. If the assumption in Theorem 1 and the condition (I) are satisfied, then for every $a>0$ and $f \in H^{k+1}\left((-\infty, \infty) \times \boldsymbol{R}_{+}^{n}\right)$ $(k=0,1, \cdots)$ with $f=0(t<0)$ the mixed problem $\left(P, B_{j}\right)(t a k i n g T=\infty)$ has a unique solution $u$ which satisfies $e^{-a t} u \in H^{m+k}\left((0, \infty) \times \boldsymbol{R}_{+}^{n}\right)$ and the following estimate

$$
\left.\int_{0}^{\infty} e^{-2 a t}| | u(t, \cdot, \cdot)\| \|_{m-1+k}^{2} d t \leq \frac{C}{a^{2}} \int_{0}^{\infty} e^{-2 a t}| | \right\rvert\, f(t, \cdot, \cdot)\| \|_{k}^{2} d t,
$$

where a constant $C$ does not depend on $u, f$ and $a$.
The following lemma is used in proof of necessity of Theorem 1.
Lemma 4. Let $f$ be a function in $L^{2}\left((-\infty, \infty) \times \boldsymbol{R}_{+}^{n}\right)$ whose support is contained in $(0, T) \times \boldsymbol{R}_{+}^{n}$ and $u$ a function satisfying $e^{-a t} u$ $\in H^{m}\left((0, \infty) \times \boldsymbol{R}_{+}^{n}\right)$ for some $a>0$ and $\left(D_{t}^{k} u\right)(0, x, y)=0(k=0,1, \cdots, m-1)$ in $\boldsymbol{R}_{+}^{n}$. If

$$
\int_{0}^{\infty} e^{-2 a t}\| \| u(t, \cdot, \cdot)\left\|_{m-1}^{2} d t \leq C \int_{0}^{\infty} \mid\right\| f(t, \cdot, \cdot) \|_{0}^{2} d t
$$

then we have
$\|\hat{u}(\tau, \cdot, \cdot)\|\left\|_{m-1}^{2} \leq C_{0} C \int_{-\infty}^{\infty}\right\| \hat{f}(a+i \eta, \cdot, \cdot) \|_{0}^{2} d \eta \quad$ for any $\tau$ with $\operatorname{Re} \tau=a$, where the constant $C_{0}$ depends on a and the support of $f$.

Next we state the following theorem which shows that $S(\tau)$ must be the cone surface with its vertex at the origin in $\boldsymbol{R}^{n-1}$.

Theorem 2. Suppose that the hyperplane $x=0$ is non-characteristic for $B_{j}(D)(j=1, \cdots, l)$ and $m_{1}<\cdots<m_{l}$. If the mixed problem $\left(P, B_{j}\right)$ is $L^{2}$-well-posed, then the varieties $S(\tau)$ don't depend on $\tau$ with $\operatorname{Re} \tau>0$.

By Theorem 2 and the theory of characters of unitary group [6] we obtain

Corollary. Under the same assumptions in Theorem 2, if $B_{j}(D)$ does not contain the terms relative to $D_{t}$ and the mixed problem $\left(P, B_{j}\right)$ is $L^{2}$-well-posed, the $S(\tau)$ is empty for any $\tau$ with $\operatorname{Re} \tau>0$.
3. Applications. First we describe necessary and sufficient conditions for $L^{2}$-well-posedness by the terms of reflection coefficients. To define reflection coefficients, for every $\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \in \bar{\Sigma}_{+}-\bar{\Sigma}_{+}$we rearrange the roots $\lambda_{j}^{+}\left(\tau^{\prime}, \sigma^{\prime}\right)(j=1, \cdots, l)$ in a sufficiently small neighbour$\operatorname{hood} U\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \cap \bar{\Sigma}_{+}$such that $\lambda_{j_{1}}^{+}\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)=\cdots=\lambda_{j_{2}-1}^{+}\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)\left(j_{1}=1\right), \cdots$, $\lambda_{j_{q}}^{+}\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)=\cdots=\lambda_{j_{+1}-1}^{+}\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)\left(j_{q+1}-1=l\right)$. Then we define reflection coefficients $C_{k, j}\left(\tau^{\prime}, \lambda, \sigma^{\prime}\right)\left(k=1, \cdots, q ; j=j_{k}, \cdots, j_{k+1}-1\right)$ by the equality

$$
\sum_{j=1}^{l} \frac{R_{j}\left(\tau^{\prime}, x, \sigma^{\prime}\right)}{R\left(\tau^{\prime}, \sigma^{\prime}\right)} B_{j}\left(\tau^{\prime}, \lambda, \sigma^{\prime}\right)=\sum_{k=1}^{q} \sum_{j=j_{k}}^{j_{k+1}^{+1-1}} C_{k, j}\left(\tau^{\prime}, \lambda, \sigma^{\prime}\right) \gamma_{k, j}\left(\tau^{\prime}, x, \sigma^{\prime}\right),
$$

where $\gamma_{k, j_{k}}\left(\tau^{\prime}, x, \sigma^{\prime}\right)=e^{i x x j_{k}\left(\tau^{\prime}, \sigma^{\prime}\right)}$,

$$
\begin{aligned}
\gamma_{k, j}\left(\tau^{\prime}, x, \sigma^{\prime}\right)= & x^{j-j_{k}} \int_{0}^{1} d \theta_{1} \cdots \int_{0}^{1} \theta_{1}^{j-j_{k}-1} \cdots \theta_{j-j_{k}}^{-}{ }^{-} e^{i x g} j^{\left(\tau^{\prime}, \sigma^{\prime} ; \theta\right)} d \theta_{j-j_{k}}, \\
g_{j}\left(\tau^{\prime}, \sigma^{\prime} ; \theta\right)= & \lambda_{j_{k}^{+}}^{+}\left(\tau^{\prime}, \sigma^{\prime}\right)+\left(\lambda_{j_{k}+1}^{+1}\left(\tau^{\prime}, \sigma^{\prime}\right)-\lambda_{j_{k}^{\prime}}^{+}\left(\tau^{\prime}, \sigma^{\prime}\right)\right) \theta_{1}+\cdots \\
& +\left(\lambda_{j}^{+}\left(\tau^{\prime}, \sigma^{\prime}\right)-\lambda_{j-1}^{+}\left(\tau^{\prime}, \sigma^{\prime}\right)\right) \theta_{1} \cdots \theta_{j-j_{k}}\left(j_{k}<j<j_{k+1}\right) .
\end{aligned}
$$

The following condition is introduced by S . Agmon [1].
Condition (\#). The multiplicity of a real root $\lambda(\tau, \sigma)$ in $\lambda$ of the characteristic equation $P(\tau, \lambda, \sigma)=0$ is at most double for every $(\tau, \sigma)$ with $\operatorname{Re} \tau=0$ and $\sigma \in \boldsymbol{R}^{n-1}$.

Then we have the following
Theorem 3. Suppose the condition (\#). If $S(\tau)$ is not the whole space $\boldsymbol{R}^{n-1}$ for every $\tau$ with $\operatorname{Re} \tau>0$, then the mixed problem $\left(P, B_{j}\right)$ is $L^{2}$-well-posed if and only if the following condition (II) is satisfied:

For every $\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \in\left(\bar{\Sigma}_{+}-\Sigma_{+}\right) \cup V^{\prime}$ there exist a neighbourhood
(II) $U\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)$ and a constant $C\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right)$ such that for any $\left(\tau^{\prime}, \sigma^{\prime}\right)$ $\in U\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \cap \Sigma_{+} \cap V^{\prime c}$

$$
\begin{gathered}
\left\|\int_{\Gamma} \frac{C_{k, j}\left(\tau^{\prime}, \lambda, \sigma^{\prime}\right)}{P\left(\tau^{\prime}, \lambda, \sigma^{\prime}\right)} e^{-i s \lambda} d \lambda\right\|_{L^{2}(s>0)} \leq C\left(\tau_{0}^{\prime}, \sigma_{0}^{\prime}\right) \frac{\left|\operatorname{Im} \lambda_{j}^{+}\left(\tau^{\prime}, \sigma^{\prime}\right)\right|^{\frac{1}{2}}}{\operatorname{Re} \tau^{\prime}}, \\
\left(k=1, \cdots, q ; j=j_{k}, \cdots, j_{k+1}-1\right)
\end{gathered}
$$

From Theorem 3 we obtain the following
Theorem 4. Let $P(D)$ and $Q(D)$ be homogenuuous differential operators, which don't contain the terms of odd order relative to $D_{x}$, of order $2 l$ and $2 l-1$ with constant coefficients respectively. If $P(D)$ satisfies the condition (\#), then the mixed problem $\left(P(D)+\varepsilon D_{x} Q(D)\right.$, $D_{x}^{2 j-1} ; j=1, \cdots, l$ ) is not well posed in the $L^{2}$-sense for a sufficiently small $\varepsilon$ with certain fixed sign.

## References

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[^0]:    1) This term means that in our definition one changes the inequality into the energy inequality.
