# 133. A Characterization of $L^{2}$-well Posedness for Iterations of Hyperbolic Mixed Problems of Second Order 

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§ 1. Introduction and theorem. We are concerned with an iterated mixed problem as follows:

$$
\left(\tilde{P}, \tilde{B}_{1}, \cdots, \tilde{B}_{m}\right) \begin{cases}\tilde{P}(x, D) u=f & \text { in } \Omega \\ \tilde{B}_{j}\left(x^{\prime}, D\right) u=g_{j} & \text { on } \Gamma, j=1, \cdots, m\end{cases}
$$

Here $\Omega$ and $\Gamma$ are the open half space $\left\{x=\left(x^{\prime}, x_{n}\right)=\left(x_{0}, x^{\prime \prime}, x_{n}\right) ; x_{0} \in R^{1}\right.$, $\left.x^{\prime \prime} \in R^{n-1}, x_{n}>0\right\}(n \geqq 2)$ and its boundary respectively, and for covariable ( $\tau, \sigma, \lambda$ ) of ( $x_{0}, x^{\prime \prime}, x_{n}$ ) the principal symbols $\tilde{P}^{0}(x, \tau, \sigma, \lambda), \tilde{B}_{j}^{0}\left(x^{\prime}, \tau, \sigma, \lambda\right)$ of $\tilde{P}, \tilde{B}_{j}$ have the following forms:
$\tilde{P}^{0}=P_{1}^{0} \cdots P_{m}^{0}, \tilde{B}_{1}^{0}=B_{1}^{0}, \tilde{B}_{2}^{0}=B_{2}^{0} P_{1}^{0}, \tilde{B}_{3}^{0}=B_{3}^{0} P_{2}^{0} P_{1}^{0}, \cdots, \tilde{B}_{m}^{0}=B_{m}^{0} P_{m-1}^{0} \cdots P_{1}^{0}$, where $P_{j}^{0}, j=1, \cdots, m$ are $x_{0}$-hyperbolic homogeneous operators of second order whose normal cones cut by $\tau=1$ don't intersect each other and are bounded surfaces in the ( $\sigma, \lambda$ ) space for every fixed $x \in \Gamma$. Furthermore $B_{j}^{0}$ is a homogeneous boundary differential operator at most of first order such that $\Gamma$ is noncharacteristic for $B_{j}^{0}$. All the coefficients are assumed to be real and smooth in $\bar{\Omega}$ and to be constant near the infinity (see [2], [3], [8]).

Definition. The problem $\left(\tilde{P}, \tilde{B}_{1}, \cdots, \tilde{B}_{m}\right)$ is said to be $L^{2}$-well posed if and only if there exist positive constants $C$ and $\gamma_{0}$ such that for every $\gamma \geqq \gamma_{0}$ and $f \in H_{1, r}(\Omega)$ the problem $\left(\tilde{P}, \tilde{B}_{1}, \cdots, \tilde{B}_{m}\right)$ with $g_{j}=0, j=1, \cdots$, $m$ has a unique solution $u$ in $H_{2 m, r}(\Omega)$ satisfying

$$
\begin{equation*}
r\|u\|_{2 m-1, r} \leqq C\|f\|_{0, r} . \tag{1.1}
\end{equation*}
$$

(For function spaces see, e.g., [7]).
Now we have
Theorem. The problem $\left(\tilde{P}, \tilde{B}_{1}, \cdots, \tilde{B}_{m}\right)$ is $L^{2}$-well posed if and only if all the frozen constant coefficients problems $\left(\tilde{P}^{0}, \tilde{B}_{1}^{0}, \cdots, \tilde{B}_{m}^{0}\right)_{x^{\prime}}$ at boundary points $x^{\prime} \in \Gamma$ are "uniformly $L^{2}$-well posed", that is, ( $\tilde{P}^{0}, \tilde{B}_{1}^{0}$, $\left.\ldots, \tilde{B}_{m}^{0}\right)_{x^{\prime}}$ is $L^{2}$-well posed for every $x^{\prime} \in \Gamma$ and the constants $C$ in (1.1) with respect to these problems are independent of the parameter $x^{\prime}$.
§ 2. Outline of the proof. It is enough to prove the "if" part, because of Theorem 1 and Lemma 2.2 in [1]. Let $\tilde{L}\left(x^{\prime}, \tau, \sigma\right)$ and $L_{j}\left(x^{\prime}\right.$, $\tau, \sigma), j=1, \cdots, m$ be the Lopatinskii determinants of $\left(\tilde{P}^{0}, \tilde{B}_{1}^{0}, \cdots, \tilde{B}_{m}^{0}\right)$ and ( $P_{j}^{0}, B_{j}^{0}$ ) respectively. Then it follows from (3.2) and Theorem 1 in [2] respectively that

$$
\begin{equation*}
\tilde{L}=L_{1} \cdots L_{m} \cdot(\text { nonzero factor }) \tag{2.1}
\end{equation*}
$$

and that every constant coefficients problem $\left(P_{j}^{0}, B_{j}^{0}\right)_{x^{\prime}}\left(x^{\prime} \in \Gamma, j=1, \cdots\right.$, $m$ ) is $L^{2}$-well posed. Hence we find by virtue of Lemma 4.1 in [2] that $\tilde{L}\left(x^{\prime}, \tau, \sigma\right)$ vanishes at a point $\left(x^{0}, \tau^{0}, \sigma^{0}\right) \in \Gamma \times\left(R^{n} \backslash 0\right)$ if and only if there is an index $j$ such that $L_{j}\left(x^{0}, \tau^{0}, \sigma^{0}\right)=0$ and $P_{j}^{0}\left(x^{0}, \tau^{0}, \lambda\right)$ has a double real zero $\lambda$. Furthermore by our hypothesis with respect to the normal cones we see that such $j$ is uniquely determind for given ( $x^{0}, \tau^{0}, \sigma^{0}$ ). Now we shall reduce ( $\tilde{P}, \tilde{B}_{1}, \cdots, \tilde{B}_{m}$ ) to a system of first order which involves tangential pseudo-differential operators, by a usual transformation (see for instance [3]). Then the uniform $L^{2}$-well posedness of the problems $\left(\tilde{P}^{0}, \tilde{B}_{1}^{0}, \cdots, \tilde{B}_{m}^{0}\right)_{x^{\prime}}, x^{\prime} \in \Gamma$ implies the following inequalities:

$$
\begin{align*}
& \left|b_{i j}\left(x^{\prime}, \tau, \sigma\right)\right|  \tag{2.2}\\
& \quad \leqq C|\operatorname{Im} \tau|^{-1}\left|\operatorname{Im} \lambda_{i}^{+}\left(x^{\prime}, \tau, \sigma\right) \operatorname{Im} \lambda_{j}^{-}\left(x^{\prime}, \tau, \sigma\right)\right|^{\sharp}, i, j=1, \cdots, l,
\end{align*}
$$

with a constant $C=C\left(x^{0}, \tau^{0}, \sigma^{0}\right)$ independent of not only $(\tau, \sigma)$ but $x^{\prime}$, where $\operatorname{Im} \tau<0$ and $\left(x^{\prime}, \tau, \sigma\right)$ varies near such a point $\left(x^{0}, \tau^{0}, \sigma^{0}\right) \in \Gamma$ $\times\left(R^{n} \backslash 0\right)$ as $\tilde{L}\left(x^{0}, \tau^{0}, \sigma^{0}\right)=0$ (see Theorem 4.1, $\alpha$ ), (i) in [9]). Here $b_{i j}$ is the function defined in Definition 4.2 of [9] and $\lambda_{j}^{ \pm}\left(x^{\prime}, \tau, \sigma\right), j=1, \cdots, m$ are zeros of $\tilde{P}^{0}\left(x^{\prime}, \tau, \sigma, \lambda\right)$ with positive imaginary part and negative one when $\operatorname{Im} \tau<0$ respectively such that $\lambda_{j}^{ \pm}\left(x^{0}, \tau^{0}, \sigma^{0}\right), j=1, \cdots, m$ are simple real, double real or nonreal if $j<l, j=l$ or $j>l$ respectively. Furthermore it follows from (2.1) and [10], Lemma 2.2 that $\tilde{L}$ is decomposed as follows:

$$
\begin{equation*}
\tilde{L}\left(x^{\prime}, \tau, \sigma\right)=\left(\sqrt{\tau-\theta\left(x^{\prime}, \sigma\right)}-D\left(x^{\prime}, \sigma\right)\right) \tilde{L}^{(1)}\left(x^{\prime}, \tau, \sigma\right) \tag{2.3}
\end{equation*}
$$

if $\operatorname{Im} \tau \leqq 0$, where $D\left(x^{0}, \sigma^{0}\right)=0, \tilde{L}^{(1)}\left(x^{0}, \tau^{0}, \sigma^{0}\right) \neq 0, \sqrt{ }$ stands for the branch with $\sqrt{1}=1$ and $\theta$ is the real-valued function with $\theta\left(x^{0}, \sigma^{0}\right)=\tau^{0}$ defined in Lemma 3.1 of [9].

Now from (2.2) and (2.3) we find that

$$
\begin{equation*}
\sum_{j=1}^{l-1}\left(\left|\left(b_{j l} \tilde{L}\right)\left(x^{\prime}, \tau, \sigma\right)\right|^{2}+\left|\left(b_{l j} \tilde{L}\right)\left(x^{\prime}, \tau, \sigma\right)\right|^{2}\right) \leqq C\left|D\left(x^{\prime}, \sigma\right)\right|, \tag{2.4}
\end{equation*}
$$

if $\tau=\theta\left(x^{\prime}, \sigma\right)$,
where $C=C\left(x^{0}, \tau^{0}, \sigma^{0}\right)$. Furthermore by means of (2.20) and (2.21) in [10] we observe that (2.4) gives

$$
\begin{equation*}
\sum_{j=1}^{l-1}\left(\left|\left(b_{j l} \tilde{L}\right)\left(x^{\prime}, \tau, \sigma\right)\right|^{2}+\left|\left(b_{l j} \tilde{L}\right)\left(x^{\prime}, \tau, \sigma\right)\right|^{2} \leqq-C Q\left(x^{\prime}, \tau, \sigma\right),\right. \tag{2.5}
\end{equation*}
$$

$$
\text { if } \tau=\theta\left(x^{\prime}, \sigma\right),
$$

since $Q\left(x^{\prime}, \tau, \sigma\right) \leqq 0$ for such $\tau$ by the realness of the coefficients and $L^{2}-$ well posedness of ( $\left.\tilde{P}^{0}, \tilde{B}_{1}^{0}, \cdots, \tilde{B}_{m}^{0}\right)_{x^{\prime}}$. Here $C=C\left(x^{0}, \tau^{0}, \sigma^{0}\right)$ is a positive constant, $Q$ is the function defined in Lemma 6.1 of [9] and we restrict ourselves to the case ( $a$ ) in Lemma 3.1 of [9]. Notice that (2.5) is equivalent to (6.5) in [9], because of Definitions 4.1, 4.2 and (6.3) in [9]. Thus we can obtain a priori estimate for ( $\tilde{P}, \tilde{B}_{1}, \cdots, \tilde{B}_{m}$ ), as in [9] or [3] (see (2.32) in [3]). The same argument may be applied to an adjoint prob-
lem and therefore we complete our proof (see Remark 1) in § 5 of [7]). § 3. Remark. In particular let $m=2$,

$$
\begin{array}{ll}
P_{j}^{0}(\tau, \sigma, \lambda)=\tau^{2}-a_{j}^{2}\left(|\sigma|^{2}+\lambda^{2}\right), & j=1,2 \text { and } \\
B_{j}^{0}\left(x^{\prime}, \tau, \lambda\right)=\lambda-C_{j}\left(x^{\prime}\right) \tau, & j=1,2,
\end{array}
$$

where $a_{j}, j=1,2$ are constants such that $0<a_{2}<a_{1}$ (see [2]). Then we have

Corollary. The problem $\left(\tilde{P}, \tilde{B}_{1}, \tilde{B}_{2}\right)$ is $L^{2}$-well posed if and only if (3.1)

$$
C_{j}\left(x^{\prime}\right) \geqq 0, \quad x^{\prime} \in \Gamma, \quad j=1,2,
$$

and for every $x^{0} \in \Gamma$ with $C_{1}\left(x^{0}\right)=0$ there is a positive constant $K$ such that

$$
\begin{equation*}
C_{2}\left(x^{\prime}\right)^{2} \leqq K C_{1}\left(x^{\prime}\right) \quad \text { for } x^{\prime} \text { near } x^{0} \tag{3.2}
\end{equation*}
$$

To prove the above fact it is enough to show that (3.2) is equivalent to (2.5) in this case. Let ( $\left.x^{0}, \tau^{0}, \sigma^{0}\right) \in \Gamma \times\left(R^{n} \backslash 0\right)$ be a point such that $L_{1}\left(x^{0}, \tau^{0}, \sigma^{0}\right)=0$ and $P_{1}^{0}\left(\tau^{0}, \sigma^{0}, \lambda\right)$ has a double real zero, say,

$$
\tau^{0}=a_{1}\left|\sigma^{0}\right| \quad \text { and } \quad C_{1}\left(x^{0}\right)=0
$$

Then we find that $\theta(\sigma)=a_{1}|\sigma|$,

$$
\begin{aligned}
& Q\left(x^{\prime}, \tau, \sigma\right)=-C_{1}\left(x^{\prime}\right) \tau\left(|\tau|^{2}+|\sigma|^{2}\right)^{-1 / 2}, \quad b_{12}\left(x^{\prime}, \tau, \sigma\right)=0 \text { and } \\
& \left(b_{21} \tilde{L}\right)\left(x^{\prime}, \tau, \sigma\right)=\left(C_{1}\left(x^{\prime}\right)-C_{2}\left(x^{\prime}\right)\right) \cdot(\text { nonzero factor }),
\end{aligned}
$$

which implies our assertion.
It is known that for every $x^{\prime} \in \Gamma$ the constant coefficients problem ( $\left.\tilde{P}^{0}, \tilde{B}_{1}^{0}, \tilde{B}_{2}^{0}\right)_{x^{\prime}}$ is $L^{2}$-well posed if and only if (3.1) is valid and $C_{1}\left(x^{\prime}\right)=0$ implies $C_{2}\left(x^{\prime}\right)=0$ (see Theorem 1 of [3] and Lemma 4.1 of [2]). Thus the inequality (3.2) shows that the $L^{2}$-well posedness of the variable coefficients problem ( $\tilde{P}, \tilde{B}_{1}, \tilde{B}_{2}$ ) need not follow from that of the constant coefficients problems ( $\left.\tilde{P}^{0}, \tilde{B}_{1}^{0}, \tilde{B}_{2}^{0}\right)_{x^{\prime}}$ for all $x^{\prime} \in \Gamma$, in contrast with the case of second order or $2 \times 2$ systems of first order (see [4], [10]).

The method of considerations used in proving Theorem is applicable to more general cases. The details will be published in Hokkaido Math. J.

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