154. Estimates from $W_{p, \alpha}$ to $W_{q, \beta}$ for the Solutions of the Petrovskii Well Posed Cauchy Problems

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1. Introduction and results.

In this note, we shall consider the Cauchy problem

$$
\begin{cases}\frac{\partial u(t, x)}{\partial t}=P(D) u(t, x) & (t, x) \in(0, \infty) \times R^{n}  \tag{1}\\ u(0, x)=u_{0}(x) & x \in R^{n}\end{cases}
$$

Here $P(D)$ is the pseudo-differential operator of order $d$, that is,

$$
\begin{equation*}
P(D) u=F^{-1}(S \hat{u}), \quad u \in \mathcal{S}^{\prime N}, \tag{2}
\end{equation*}
$$

where $S=\left(s_{i j}\right)_{1 \leqslant i, j \leqslant N}$ is the $N \times N$ matrix of functions $s_{i j}$ in $C^{\infty}\left(R^{n}\right)$ which satisfy, for all multi-indices $\sigma=\left(\sigma_{1}, \cdots, \sigma_{n}\right)$,

$$
\begin{equation*}
\left|D^{\sigma} s_{i j}(y)\right| \leqslant C_{\sigma}(1+|y|)^{d-|\sigma|} \tag{3}
\end{equation*}
$$

where $C_{o}$ are constants depending on $\sigma, D^{\sigma}=\left(\partial / \partial y_{1}\right)^{\sigma_{1}} \cdots\left(\partial / \partial y_{n}\right)^{\sigma_{n}}$ and $|\sigma|=\sigma_{1}+\cdots+\sigma_{n}$. The matrix $S$ will be called the symbol of $P$. In the above, $\mathcal{S}^{\prime N}, F^{-1}$ and $\hat{u}$ denote the space of all $N$-tuples of distributions in the dual space $\mathcal{S}^{\prime}$ of the Schwartz space $\mathcal{S}$, the inverse Fourier transformation and the Fourier transform of $u$, respectively. We assume that the order $d$ of $P$ is positive.

Let $\lambda_{j}(y)$ denote the eigenvalues of $S(y)$ for $j=1,2, \cdots, N$. We say that the Cauchy problem (1) is Petrovskii well posed if

$$
\begin{equation*}
\operatorname{Re} \lambda_{j}(y) \leqslant \Lambda, \quad 1 \leqslant j \leqslant N, y \in R^{n} \tag{4}
\end{equation*}
$$ are valid for some constant $\Lambda$. When the Cauchy problem (1) is Petrovskii well posed, we can solve the problem in $\mathcal{S}^{\prime N}$ and the solution can be written as

(5) $u(t)=E(t) u_{0}=F^{-1}\left(\exp (t S) \hat{u}_{0}\right) \quad$ for $u_{0} \in \mathcal{S}^{\prime N}$.

We call the operator $E(t): u_{0} \rightarrow u(t)$ the solution operator.
Let $1 \leqslant p \leqslant \infty$. For $u \in L_{p}^{N}$ (the space of all $N$-tuples of functions in $L_{p}\left(R^{n}\right)$ ), we set

$$
\|u\|_{p}= \begin{cases}\left(\int_{R^{n}}|u(x)|^{p} d x\right)^{1 / p} & \text { if } p<\infty \\ \text { ess sup }\left\{|u(x)| ; x \in R^{n}\right\} & \text { otherwise } .\end{cases}
$$

For $\alpha \geqslant 0$, let $v_{\alpha}(y)=\left(1+|y|^{2}\right)^{\alpha / 2}$ and

$$
\|u\|_{p, \alpha}=\left\|F^{-1}\left(v_{\alpha} \hat{u}\right)\right\|_{p} \quad \text { for } u \in L_{p}^{N} .
$$

We define $W_{p, \alpha}^{N}=\left\{u \in L_{p}^{N} ;\|u\|_{p, \alpha}<\infty\right\}$.
Henceforth, for given $p$ and $q$, we set $\gamma(p, q)=\max (1 / 2-1 / p$, $1 / q-1 / 2,0$ ). Our results are the following.

Theorem 1. Assume that the Cauchy problem (1) is Petrovskii well posed. Suppose that $1 \leqslant p \leqslant q \leqslant \infty$. If

$$
\begin{equation*}
\alpha-\beta>n(1 / p-1 / q)+n d \gamma(p, q)+(N-1) d \tag{6}
\end{equation*}
$$

then the inequality
$\left\|E(t) u_{0}\right\|_{q, \beta} \leqslant C(t)\left\|u_{0}\right\|_{p, \alpha}, \quad u_{0} \in W_{p, \alpha}^{N}$
holds with some function $C(t)$ which is bounded by a constant multiple of $e^{1 t}(1+t)^{N-1+n \gamma(p, q)}$. Moreover, if $1<p \leqslant 2 \leqslant q<\infty$, then the inequality (6) is valid even when $\alpha-\beta=n(1 / p-1 / q)+n d \gamma(p, q)+(N-1) d$.

Theorem 2. If $\alpha-\beta<n(1 / p-1 / q)+n d \gamma(p, q)+(N-1) d$, then there exists a pseudo-differential operator $P(D)$ of order $d$ for which the Cauchy problem (1) is Petrovskii well posed and the solution operator $E(t)$ is not bounded from $W_{p, \alpha}^{N}$ to $W_{q, \beta}^{N}$ for each $t>0$. Further, if $d \neq 1$ and if $p=1$ or $q=\infty$, then the same conclusion as above holds for $\alpha-\beta$ $=n(1 / p-1 / q)+n d \gamma(p, q)+(N-1) d$.

Remarks. Theorem 1 is a generalization of the results obtained by Sjöstrand [8] (for the Schrödinger equation) and the author [7] (for the case that $N=1$ and $S$ is a pure imaginary polynomial function).

Considering $L_{p}-L_{q}$ estimates for pseudo-differential operators, Hörmander has obtained the essentially same result as Theorems 1 and 2 for the case $d<1$ in [5].
2. Proof of Theorem 1.

We first define

$$
M_{p, q}^{N}=M_{p, q}^{N}\left(R^{n}\right)=\left\{A=\left(a_{i j}\right)_{1 \leqslant i, j \leqslant N} ; a_{i j} \in \mathcal{S}^{\prime}, M_{p, q}^{N}(A)<\infty\right\}
$$

where

$$
M_{p, q}^{N}(A)=\sup \left\{\left\|F^{-1}(A \hat{u})\right\|_{p} ; u \in \mathcal{S}^{N} \text { with }\|u\|_{p}=1\right\} .
$$

When $N=1$, we merely write $M_{p, q}$ for $M_{p, q}^{N}$ and, in case $p=q$, we shall omit the subscript $q$ of $M_{p, q}^{N}$. We refer to Hörmander [4] and Brenner [2] for the relevant facts about $M_{p, q}^{N}$.

The following Lemma 1 is fundamental.
Lemma 1. $A=\left(a_{i j}\right)_{1 \leqslant i, j \leqslant N}$ belongs to $M_{p, q}^{N}$ if and only if $a_{i j} \in M_{p, q}$ for all $i, j, 1 \leqslant i, j \leqslant N$. Moreover, the inequality
(7) $\quad c M_{p, q}^{N}(A) \leqslant \max \left(M_{p, q}\left(a_{i j}\right) ; 1 \leqslant i, j \leqslant N\right) \leqslant C M_{p, q}^{N}(A)$ holds for some constants $c, C>0$.

The proof is easy and so we omit it.
We need two more lemmas to prove the theorem. For any $N \times N$ matrix $A, m(A)$ will denote the matrix norm, that is,

$$
m(A)=\sup \left\{|A u| ; u \in R^{N},|u|=1\right\} .
$$

Lemma 2 (Bernstein's theorem), Let $J=[n / 2]+1$. Let $A=\left(a_{i j}\right)$ be a $N \times N$ matrix satisfying $a_{i j} \in C^{J}\left(R^{n}\right)$ for all $i, j, 1 \leqslant i, j \leqslant N$. Suppose that $m\left(D^{\sigma} A\right) \in L_{2}\left(R^{n}\right)$ for all $\sigma,|\sigma| \leqslant J$. Then, $A \in M_{1}^{N}$ and the inequality

$$
\begin{equation*}
M_{1}^{N}(A) \leqslant C\|m(A)\|_{2}^{1-n /(2 J)}\left(\sum_{|\sigma|=J}\left\|m\left(D^{\circ} A\right)\right\|_{2}\right)^{n /(2 J)} \tag{8}
\end{equation*}
$$

holds for some constant $C>0$.
Proof. By the usual Bernstein's theorem (see e.g. Sjöstrand [8]), we have

$$
\begin{aligned}
M_{1}\left(a_{i j}\right) & \leqslant C\left\|a_{i j}\right\|_{2}^{1-n /(2 J)}\left(\sum_{|\sigma|=J}\left\|D^{\sigma} a_{i j}\right\|_{2}\right)^{n /(2 J)} \\
& \leqslant C\|m(A)\|_{2}^{1-n /(2 J)}\left(\sum_{|\sigma|=J}\left\|m\left(D^{\sigma} A\right)\right\|_{2}\right)^{n /(2 J)}
\end{aligned}
$$

Hence, by Lemma 1, we get

$$
M_{1}^{N}(A) \leqslant C^{\prime}\|m(A)\|_{2}^{1-n /(2 J)}\left(\sum_{\mid \sigma=J}\left\|m\left(D^{\sigma} A\right)\right\|_{2}\right)^{n /(2 J)}
$$

This proves Lemma 2.
Lemma 3. Let $B$ be a $N \times N$ matrix and $\lambda_{j}, 1 \leqslant j \leqslant N$, be eigenvalues of $B$. Set $\Lambda=\max \left(\operatorname{Re} \lambda_{j} ; 1 \leqslant j \leqslant N\right)$. The following estimate holds:

$$
\begin{equation*}
m\left(e^{B}\right) \leqslant e^{A_{j=0}^{N-1}}(2 m(B))^{j} \tag{9}
\end{equation*}
$$

For the proof of this lemma, we refer to Gelfand-Shilov [3].
Proof of Theorem 1. Without any loss of generality, we may assume $\beta=0$. Let us set $A(t, y)=\left(1+|y|^{2}\right)^{-\alpha / 2} e^{t S(y)}$ for $(t, y) \in(0, \infty) \times R^{n}$. We shall show below that

$$
\begin{equation*}
M_{p, q}^{N}(A(t)) \leqslant C(t) \tag{10}
\end{equation*}
$$

which proves Theorem 1 when $p<\infty$. When $p=\infty$, Theorem 1 is a stronger assertion than (10) and we need a slight modification. For such a modification, see the author [7].

We divide our consideration into three cases. We first consider the case $1 \leqslant p \leqslant 2 \leqslant q \leqslant \infty$. When $p \neq 1$ and $q \neq \infty$, we set $\zeta=n(1 / p-1 / 2)$ and $\eta=n(1 / 2-1 / q)$. Putting $A^{\prime}(t, y)=v_{\xi}(y) A(t, y) v_{\eta}(y)$ for $y \in R^{n}$, we have
(11)

$$
m\left(A^{\prime}(t, y)\right) \leqslant C_{1} C(t)
$$

by Lemma 3 and the assumption on $\alpha$, and hence $A^{\prime}(t) \in M_{2}^{N}$.
By the Hardy-Littlewood-Sobolev theorem, we see that $v_{-\zeta} \in M_{p, 2}^{N}$ and $v_{-\eta} \in M_{2, q}^{N}$. Therefore, $A(t)=v_{-\zeta} A^{\prime}(t) v_{-\eta} \in M_{p, q}^{N}$ and (12)

$$
M_{p, q}^{N}(A(t)) \leqslant C_{2} C(t)
$$

By the assumption on $\alpha$, it is possible to take $\zeta>n(1 / p-1 / 2)$ (when $p=1$ ) and $\eta>n(1 / 2-1 / q)$ (when $q=\infty$ ), so that the inequality (11) holds. So we can show (12) in the same manner as above in case $p$ $=1$ or $q=\infty$.

Next we turn to the case $1 \leqslant p \leqslant q<2$. We divide $\alpha$ into $\alpha=\alpha^{\prime}+\alpha^{\prime \prime}$ where $\alpha^{\prime}>n d(1 / q-1 / 2)+(N-1) d$ and $\alpha^{\prime \prime}>n(1 / p-1 / q)$. Define $A^{\prime \prime}(t, y)$ $=v_{-\alpha^{\prime}}(y) e^{t S(y)}$ for $(t, y) \in(0, \infty) \times R^{n}$. Let us choose a function $\phi_{0}(r)$ in $C^{\infty}(R)$ which equals to 1 for $r<1$ and vanishes for $r>2$, and we put $\phi_{k}(r)=\phi_{0}\left(2^{-k} r\right)-\phi_{0}\left(2^{-k+1} r\right)$ for $k=1,2, \cdots$. We decompose $A^{\prime \prime}$ as $A^{\prime \prime}$ $=\sum_{k=0}^{\infty} A_{k}^{\prime \prime}$, where $A_{k}^{\prime \prime}(t, y)=\phi_{k}(|y|) A^{\prime \prime}(t, y)$. By Lemma 3, we have

$$
m\left(e^{t S(y)}\right) \leqslant C e^{\Lambda t}(1+t)^{N-1}(1+|y|)^{(N-1) d}
$$

Using this estimate and (3), we easily get

$$
m\left(D^{\sigma} A_{k}^{\prime \prime}(t, y)\right) \leqslant C e^{\Lambda t}(1+t)^{|\sigma|+N-1} 2^{(d-1) k|\sigma|-\alpha^{\prime} k+(N-1) d k}
$$

Hence, by Lemma 2,

$$
M_{1}^{N}\left(A_{k}^{\prime \prime}(t)\right) \leqslant C e^{\Lambda t}(1+t)^{n / 2+(N-1)} 2^{n k d / 2-\alpha^{\prime} k+(N-1) d k}
$$

It is easy to see from Lemmas 1 and 3 that

$$
M_{2}^{N}\left(A_{k}^{\prime \prime}(t)\right) \leqslant C\left\|m\left(A_{k}^{\prime \prime}(t)\right)\right\|_{\infty} \leqslant C^{\prime} e^{\Lambda t}(1+t)^{N-1} 2^{(N-1) d k-\alpha^{\prime} k}
$$

Applying the Riesz-Thorin's convexity theorem, we obtain

$$
M_{q}^{N}\left(A_{k}^{\prime \prime}(t)\right) \leqslant C e^{\Lambda t}(1+t)^{N-1+n(1 / q-1 / 2)} 2^{(N-1) d k-\alpha^{\prime} k+n(1 / q-1 / 2) k}
$$

Summing over all $k$, we have

$$
M_{q}^{N}\left(A^{\prime \prime}(t)\right) \leqslant \sum_{k=0}^{\infty} M_{q}^{N}\left(A_{k}^{\prime \prime}(t)\right) \leqslant C e^{4 t}(1+t)^{N-1+n(1 / q-1 / 2)}
$$

On the other hand, by the Hardy-Littlewood-Sobolev theorem, we get $v_{-\alpha^{\prime \prime}} \in M_{p, q}^{N}$. Therefore, we have

$$
M_{p, q}^{N}(A(t)) \leqslant C(t)
$$

In case $2<p \leqslant q \leqslant \infty$, our theorem is easily shown by the standard duality argument. This finishes the proof.
3. Proof of Theorem 2.

We begin with the well-known lemma.
Lemma 4. Set $v_{\delta}(y)=\left(1+|y|^{2}\right)^{\delta / 2}$ for $y \in R^{n}$. If $\delta>-n(1 / p-1 / q)$ and $1 \leqslant p \leqslant q \leqslant \infty$, then $v_{\delta} \notin M_{p, q}$. Moreover, in case $p=1$ or $q=\infty$, $v_{\delta} \notin M_{p, q}$ for $\delta=-n(1 / p-1 / q)$.

For the proof, see Stein [9].
The next lemma was proved by Wainger [10] $(0<d<1$ and $p=q)$, Hörmander [5] $(0<d<1)$ and Sjöstrand [8] $(d>1$ and $p=q)$. Here and later the letter $\psi$ denotes a function in $C^{\infty}(R)$ satisfying $\psi(r)=1$ for $r>2$ and $\psi(r)=0$ for $r<1$, and let $w_{\delta}(y)=\psi(|y|)|y|^{-\delta} \exp \left(i|y|^{d}\right)$ for $y \in R^{n}$ and $d>0$.

Lemma 5. If $d \neq 1,1 \leqslant p \leqslant q<2$ and $\delta<n(1 / p-1 / q)+n d(1 / q-1 / 2)$, then $w_{\delta} \notin M_{p, q}\left(R^{n}\right)$. Especially, if $p=1$, then $w_{\delta} \notin M_{p, q}$ for $\delta=n(1 / p-1 / q)$ $+n d(1 / q-1 / 2)$.

Proof. First we assume $\delta<n(1 / p-1 / q)+n d(1 / q-1 / 2)$. Let $p^{\prime}$ $=p /(p-1)$ and $\hat{g}(y)=\psi(|y|)|y|^{-\theta}$ with $\theta=n / p^{\prime}+n(1 / p-1 / q)+n d(1 / q$ $-1 / 2)-\delta$. We know that $g \in L_{p}\left(R^{n}\right)$ since $\theta>n / p^{\prime}$ (see Sjöstrand [8]). Putting $\hat{f}(y)=w_{\delta}(y) \hat{g}(y)$ for $y \in R^{n}$, the asymptotic behavior of $f$ is as follows: (i) If $d>1$, then

$$
\begin{equation*}
|f(x)|=C_{d, \delta+\theta}|x|^{(n-\delta-\theta-n d / 2) /(d-1)}+O\left(|x|^{\omega}\right) \tag{13}
\end{equation*}
$$

as $|x| \rightarrow \infty$, where $\omega<(n-\delta-\theta-n d / 2) /(d-1)$ and where $C_{d, \delta+\theta}$ is a positive constant.
(ii) If $d<1$, then

$$
\begin{equation*}
|f(x)|=C_{d, \delta+\theta}|x|^{(n-\delta-\theta-n d / 2) /(d-1)}+O\left(|x|^{\omega}\right) \tag{14}
\end{equation*}
$$

as $|x| \rightarrow 0$, where $\omega>(n-\delta-\theta-n d / 2) /(d-1)$ and where $C_{d, \delta+\theta}$ is a posi-
tive constant.
Since $q(n-\delta-\theta-n d / 2) /(d-1)=-n, f$ does not belong to $L_{p}\left(R^{n}\right)$. This means $w_{\delta} \notin M_{p, q}\left(R^{n}\right)$.

We turn to the case $p=1$ and $\delta=n(1 / p-1 / q)+n d(1 / q-1 / 2)$. It is well-known (see Hörmander [4]) that

$$
M_{1, q}=F L_{q} \text { for } q>1 \text { and } M_{1}=F M
$$

Here $F L_{q}$ denotes the space of all Fourier transforms of functions in $L_{q}$ and $F M$ denotes the space of all Fourier-Stieltjes transforms of bounded measures.

On the other hand, the inverse Fourier transform of the function $w_{\delta}$ is asymptotically described by the right hand side of (13) or (14) with $\theta=0$. So we have $w_{\delta} \notin M_{1, q}$ since $q(n-\delta-n d / 2) /(d-1)=-n$. Thus we have proved Lemma 5.

Proof of Theorem 2. Seeing that $M_{p, q}^{N}=\{0\}$ for $p>q$ (see Hörmander [4]), we assume below that $p \leqslant q$. We may also assume $\beta=0$. The proof will be divided into three cases.

We first treat the case $p \leqslant 2 \leqslant q$. Let us define the $N \times N$ matrix $S$ by

$$
S(y)=\left(1+|y|^{2}\right)^{d / 2}\left(\begin{array}{lll}
0 & 1 & 0  \tag{15}\\
& \ddots & 0 \\
0 & & 1 \\
0 & & 0
\end{array}\right) \quad \text { for } y \in R^{n}
$$

The $(1, N)$ element of $e^{t S(y)}$ is given by

$$
\frac{1}{(N-1)!}\left(1+|y|^{2}\right)^{(N-1) d / 2} t^{N-1} .
$$

Therefore, in view of Lemmas 1 and 4, we see that

$$
\left(1+|y|^{2}\right)^{-\alpha / 2} e^{t S(y)} \notin M_{p, q}^{N} .
$$

It is now easily checked that the pseudo-differential operator $P(D)$ defined by (2) satisfies the desired properties.

We turn to the case $d \neq 1$ and $1 \leqslant p \leqslant q<2$ (or $2<p \leqslant q \leqslant \infty$ ). Let us set

$$
S(y)=i \psi(|y|)|y|^{\alpha}\left(\begin{array}{lll}
1 & 1 & 0  \tag{16}\\
& \ddots & \\
& \ddots & 1 \\
0 & & 1
\end{array}\right), y \in R^{n}
$$

where the matrix is $N \times N$. Then, the ( $1, N$ ) element of the matrix function $e^{t S(y)}$ is

$$
\frac{i^{N-1}}{(N-1)!} \psi(|y|)^{N-1}|y|^{N-1} \exp \left(i \psi(|y|)|y|^{\mid d}\right) t^{N-1} .
$$

When $1 \leqslant p \leqslant q<2$, by Lemma 5 , we easily see that the $(1, N)$ element of $\left(1+|y|^{2}\right)^{-\alpha / 2} e^{t S(y)}$ does not belong to $M_{p, q}^{N}$. When $2<p \leqslant q \leqslant \infty$, it is shown by the duality argument that the same is also true. It follows from Lemma 1 that $\left(1+|y|^{2}\right)^{-\alpha / 2} e^{t S(y)} \notin M_{p, q}^{N}$.

In case when $d=1$ and $1 \leqslant p \leqslant q<2$ (or $2<p \leqslant q \leqslant \infty$ ), we construct the pseudo-differential operator having the required properties for each $p, q$ and $\alpha$. For fixed $p, q$ and $\alpha$ satisfying $\alpha<n(1 / p-1 / q)+n \gamma(p, q)$ $+N-1$, we choose a number $d^{\prime}$ which is smaller than 1 and satisfies $\alpha<n(1 / p-1 / q)+n d^{\prime} \gamma(p, q)+(N-1) d^{\prime}$. Then, the symbol $S$ givin by (16) replaced $d$ by $d^{\prime}$ defines the pseudo-differential operator having the required properties. The proof is completed.

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