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# The Cauchy problem for nonlinear wave equations in the homogeneous Sobolev space 

by

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Abstract. - We consider the Cauchy problem for nonlinear wave equations in the homogeneous Sobolev space $\dot{H}^{\mu}\left(\mathbf{R}^{n}\right)$, where $n \geqslant 2$ and $0 \leqslant \mu<n / 2$ using the generalized Strichartz estimates given by J. Ginibre and G. Velo (1995). © Elsevier, Paris

Key words: Nonlinear wave equations, Cauchy problem, Strichartz' estimate, minimal regularity, homogeneous Besov estimates

Résumé. - Nous considérons le problème de Cauchy pour des équations d'ondes non linéaires dans l'espace de Sobolev homogène $\dot{H}^{\mu}\left(\mathbf{R}^{n}\right)$, où $n \geqslant 2$ et $0 \leqslant \mu<n / 2$, utilisant les estimées de Strichartz généralisées par J. Ginibre et G. Velo (1995). © Elsevier, Paris

## 1. INTRODUCTION

We study the Cauchy problem for nonlinear wave equations of the form

$$
\begin{equation*}
\partial_{t}^{2} u-\Delta u=f(u) \tag{1.1}
\end{equation*}
$$

[^0]in the homogeneous Sobolev space $\dot{H}^{\mu}\left(\mathbf{R}^{n}\right)$ with $n \geqslant 2$ and $0 \leqslant \mu<n / 2$, where $\Delta$ denotes the Laplacian in $\mathbf{R}^{n}$ and the typical form of $f(u)$ is the single power interaction $\lambda|u|^{p-1} u$ with $\lambda \in \mathbf{R}$ and $1<p<\infty$. As usually done, with data $u(0)=\phi, \partial_{t} u(0)=\psi$ we regard (1.1) as the following integral equation.
\[

$$
\begin{equation*}
u(t)=\Phi(u)(t) \equiv \dot{K}(t) \phi+K(t) \psi+\int_{0}^{t} K(t-\tau) f(u(\tau)) d \tau \tag{1.2}
\end{equation*}
$$

\]

where $\dot{K}(t)=\cos t \sqrt{-\Delta}, K(t)=(\sin t \sqrt{-\Delta}) / \sqrt{-\Delta}$.
There are many papers on the Cauchy problem for (1.1) and large time behavior of global solutions, see [2,4,5,7-15,17-23]. Recently, in [15] H. Lindblad and C.D. Sogge studied (1.1) in the Sobolev space with minimal regularity assumptions on the data. One of the key ingredients in [15] is generalized Strichartz estimates on the free wave equation. Those estimates are described exclusively in terms of the homogeneous Sobolev space, and accordingly, the associated estimates on the nonlinear term are required to take a form in the framework of the homogeneous Sobolev spaces.

Unfortunately, however, when it comes to the Leibniz rule for fractional derivatives, it sometimes happens that additional regularity assumptions on $f$ would be necessary more than one needed.

Meanwhile, we have recently found that the problem could be efficiently dealt in the framework of the homogeneous Besov spaces [16], see also [3,7]. Moreover, the Strichartz estimate are now available in the fully extended version, especially in the homogeneous Besov setting [10].

The purpose of this paper is to reexamine the results of [15] on the Cauchy problem for (1.1) in the homogeneous Sobolev spaces by means of a number of sharp estimates described in terms of the homogeneous Besov spaces. As a result of the homogeneous Besov technique, we have refined and generalized the previous results in some directions. To state our theorem, we make a series of definitions.

Definition 1.1. - For $s \geqslant-1$ and $p \geqslant 1$, we define a class of functions $G(s, p)$ in $C(\mathbf{C}, \mathbf{C})$ as follows. We say $f \in G(s, p)$ if $f$ satisfies either of the following conditions:

1. For some nonnegative integers $a, b$ with $p=a+b$,

$$
f(z)=C_{1}+C_{2} z^{a} \bar{z}^{b},
$$

where $C_{1}$ and $C_{2}$ are constants and $C_{1}$ is disregarded if $s \leqslant 0$.
2. $[s]+1<p, f \in C^{[s]+1}(\mathbf{C}, \mathbf{C})$.

$$
f(0)=\cdots=f^{([s]+1)}(0)=0
$$

where $f(0)=0$ may be disregarded if $s>0$. Moreover, $f$ satisfies the estimates for all $z, w \in \mathbf{C}$

$$
\begin{align*}
& \left|f^{([s]+1)}(z)-f^{([s]+1)}(w)\right| \\
& \quad \leqslant\left\{\begin{array}{cl}
C\left(|z|^{p-[s]-2}\right. \\
\left.+|w|^{p-[s]-2}\right)|z-w| & \text { if }[s]+2 \leqslant p \\
C|z-w|^{p-[s]-1} & \text { if }[s]+1<p<[s]+2,
\end{array}\right. \tag{1.3}
\end{align*}
$$

where $[s]$ denotes the largest integer less than or equal to $s$, but $[0]=-1$. We call $s$ the first index of $G(s, p)$.

DEFINITION 1.2.-Let $\varepsilon>0$. Let $\Omega_{\varepsilon}$ be

$$
\begin{gathered}
\Omega_{\varepsilon} \equiv\{(1 / q, 1 / r) \mid 0 \leqslant 1 / q, 1 / r \leqslant 1 / 2, \varepsilon \leqslant 1 / r \leqslant 1 / 2-2 /((n-1) q) \\
\left.(1 / q, 1 / r) \notin B_{\varepsilon}(1 / 2,1 / 2-1 /(n-1))\right\}
\end{gathered}
$$

where $B_{\varepsilon}(1 / 2,1 / 2-1 /(n-1))$ denotes an open ball with radius $\varepsilon$ and center at $(1 / 2,1 / 2-1 /(n-1))$. Let $0 \leqslant \mu<n / 2$. Let $\Omega_{\varepsilon, \mu}$ be

$$
\begin{aligned}
\Omega_{\varepsilon, \mu} \equiv & \left\{(1 / q, 1 / r, \rho) \mid(1 / q, 1 / r) \in \Omega_{\varepsilon}, 0 \leqslant \rho \leqslant \mu\right. \\
& \mu=\rho+n(1 / 2-1 / r)-1 / q, 0 \leqslant 1 / q \leqslant n / 2-\mu-n \varepsilon\} .
\end{aligned}
$$

DEFINITION 1.3. - For any $-\infty \leqslant a \leqslant 0 \leqslant b \leqslant \infty$, we define an interval $I \equiv[a, b] \cap \mathbf{R}$ with length $|a-b|$ and for $R>0$ a function space $X_{\varepsilon}(I, R)$ with metric $d$ by

$$
\begin{aligned}
X_{\varepsilon}(I, R) \equiv & \left\{u \in \bigcap_{(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}} L^{q}\left(I, \dot{B}_{r}^{\rho}\right) \mid\right. \\
& \left.\max _{(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}}\left\|u ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant R\right\}, \\
d(u, v) \equiv & \max _{(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}}\left\|u-v ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| .
\end{aligned}
$$

In our theorem below, $\|(\phi, \psi)\|_{\mu}$ denotes $\max \left(\left\|\phi ; \dot{H}^{\mu}\right\|,\left\|\psi ; \dot{H}^{\mu-1}\right\|\right)$, $\alpha$ denotes the lower root of the quadratic equation

$$
\begin{equation*}
F(x) \equiv x^{2}-\left(\left(n^{2}-3\right) /(2 n-2)\right) x+\left(n^{2}+n+4\right) /(4 n-4)=0 \tag{1.4}
\end{equation*}
$$

It follows that $\min (1, \alpha)=1$ for $n \leqslant 6$ and $(n+1) /(2 n-2)<\alpha<1$ for $n \geqslant 7$. Finally, $\beta(\mu)$ is given by

$$
\beta(\mu) \equiv \frac{n^{2}+n+4-2(n-3) \mu}{2(n-1)(n-2 \mu)} .
$$

It follows that $\beta(\alpha)=\alpha, \beta((n-4) / 2)=1$ and that $\beta(\mu)$ is a strictly increasing function in $\mu$.

Theorem 1.1. - Let $n \geqslant 2,0 \leqslant \mu<n / 2$ and $2 /(n-2 \mu) \leqslant p-1$. Let $n, f, p$ satisfy any of the following conditions:
(A1) $n=2, f \in G(0, p)$ and

$$
p-1 \leqslant \begin{cases}(2(n+1) / n+4 \mu) /(n+1-4 \mu) & \text { for } 0 \leqslant \mu \leqslant 1 / 4, \\ 4 /(n+1-4 \mu) & \text { for } 1 / 4<\mu<1 / 2, \\ 4 /(n-2 \mu) & \text { for } 1 / 2 \leqslant \mu<n / 2\end{cases}
$$

(A2) $n \geqslant 3,0 \leqslant \mu \leqslant \min (1, \alpha), f \in G(0, p)$ and

$$
\begin{aligned}
p-1 & \leqslant(3(n+1)+2(n-1) \mu) /\left(n^{2}+n-4 n \mu\right) \\
\text { for } 0 & \leqslant \mu<(n-3) /(2 n-2), \\
p-1 & <4 /(n+1-4 \mu) \text { for } \mu=(n-3) /(2 n-2), \\
p-1 & \leqslant \begin{cases}4 /(n+1-4 \mu) & \text { for }(n-3) /(2 n-2)<\mu<1 / 2, \\
4 /(n-2 \mu) & \text { for } 1 / 2 \leqslant \mu \leqslant \min (1, \alpha) .\end{cases}
\end{aligned}
$$

(A3) $n \geqslant 7, \alpha<\mu<(n-4) / 2, f \in G(\mu-\beta(\mu), p)$ and

$$
p-1 \leqslant 4 /(n-2 \mu) .
$$

(A4) $n \geqslant 3, \max (1,(n-4) / 2) \leqslant \mu<n / 2, f \in G(\mu-1, p)$ and $p-1 \leqslant 4 /(n-2 \mu)$.
Let $\varepsilon>0$ be sufficiently small. Then for any data $(\phi, \psi) \in \dot{H}^{\mu} \times \dot{H}^{\mu-1}$ there exists a unique local solution of (1.2) in $X_{\varepsilon}(I, R)$ with $|I|>0$ sufficiently small and $R$ sufficiently large. Moreover if $p-1=4 /(n-2 \mu)$ and $\|(\phi, \psi)\|_{\mu}$ is sufficiently small, there exists a unique global solution in $X_{\varepsilon}((-\infty, \infty), R)$ with $R$ sufficiently small.

On the solutions given by above, we have the following results:
(1) $\left(u, \partial_{t} u\right)$ is continuous in time with respect to the norm $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$.
(2) The solution $u$ depends on the data $(\phi, \psi)$ continuously. Namely let $v$ be the solution of (1.2) with data $\left(\phi_{0}, \psi_{0}\right)$ such that $\left\|\left(\phi-\phi_{0}, \psi-\psi_{0}\right)\right\|_{\mu}$ tends to zero, then $d(u, v) \rightarrow 0$ for $p \notin J, v \rightarrow u$ in $\mathcal{D}^{\prime}\left(\mathbf{R}^{n+1}\right)$ for $p \in J$, where $\mathcal{D}^{\prime}\left(\mathbf{R}^{n+1}\right)$ denotes the space of distribution and $J$ denotes an
interval defined only for (A3) and (A4) as

$$
J \equiv([\mu-\beta(\mu)]+1,[\mu-\beta(\mu)]+2)
$$

for (A3),

$$
J \equiv([\mu],[\mu]+1)
$$

for (A4).
(3) Let $p-1=4 /(n-2 \mu)$. There exists a pair $\left(\phi_{+}, \psi_{+}\right)$in $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$ such that

$$
\left\|u(t)-\dot{K}(t) \phi_{+}-K(t) \psi_{+} ; \dot{H}^{\mu}\right\| \rightarrow 0 \quad \text { as } t \rightarrow \infty
$$

(4) Let $p-1=4 /(n-2 \mu)$. Let $\gamma>0$ be sufficiently small. Then for ally data $\left(\phi_{-}, \psi_{-}\right)$which satisfies $\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}<\gamma$, there exists a global solution u and a pair $\left(\phi_{+}, \psi_{+}\right)$in $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$ such that

$$
\left\|u(t)-\dot{K}(t) \phi_{ \pm}-K(t) \psi_{ \pm} ; \dot{H}^{\mu}\right\| \rightarrow 0 \quad \text { as } t \rightarrow \pm \infty .
$$

Moreover if $p \notin J$, then the map $\left(\phi_{-}, \psi_{-}\right) \mapsto\left(\phi_{+}, \psi_{+}\right)$is continuous in $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$.

Remark 1. - By dilation argument, it is natural to call $p=1+$ $4 /(n-2 \mu)$ the critical exponent for the well-posedness of the Cauchy problem for (1.2) in $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$. On the other hand, H. Lindblad and C.D. Sogge $[14,15]$ showed the ill-posedness in the following three cases:
(a) $p>1+4 /(n+1-4 \mu)$ with $n=2$ and $1 / 4<\mu \leqslant 1 / 2$;
(b) $p>1+4 /(n+1-4 \mu)$ with $n \geqslant 3$ and $(n-3) /(2 n-2)<\mu \leqslant$ $1 / 2$;
(c) $p=2$ with $n=3$ and $\mu=0$.

Remark 2. - We use the homogeneous Besov space for the linear and nonlinear estimates for (1.2), by which it becomes easy to deal with the fractional derivative of nonlinear term (see Propositions 2.1 and 2.2). For the definition of the homogeneous Besov space and its properties, we refer to $[1,8,10,24]$. Our results for the local and global solvability of (1.2) in the homogeneous Sobolev space $\dot{H}^{\mu}$ with $3 / 2<\mu<n / 2$ and the corresponding results on scattering are new.

## 2. ESTIMATES FOR NONLINEAR TERMS

Proposition 2.1. - Let $s>0,1 \leqslant p$ and $f \in G(s, p)$. Let $1 \leqslant \ell$ $<\infty, 2 \leqslant q<\infty, 2 \leqslant r \leqslant \infty$ with $1 / \ell=(p-1) / q+1 / r$. Then

$$
\begin{equation*}
\left\|f(u) ; \dot{B}_{\ell}^{s}\right\| \leqslant C\left\|u ; \dot{B}_{q}^{0}\right\|^{p-1}\left\|u ; \dot{B}_{r}^{s}\right\| \tag{2.5}
\end{equation*}
$$

$$
\begin{aligned}
& \left\|f(u)-f(v) ; \dot{B}_{\ell}^{s}\right\| \\
& \qquad \begin{array}{l}
\leqslant \\
\quad \max \left(\left\|u ; \dot{B}_{q}^{0}\right\|,\left\|v ; \dot{B}_{q}^{0}\right\|\right)^{p-1}\left\|u-v ; \dot{B}_{r}^{s}\right\| \\
\quad+ \\
\quad C \max \left(\left\|u ; \dot{B}_{q}^{0}\right\|,\left\|v ; \dot{B}_{q}^{0}\right\|\right)^{p-2}\left\|u-v ; \dot{B}_{q}^{0}\right\| \max \left(\left\|u ; \dot{B}_{r}^{s}\right\|,\left\|v ; \dot{B}_{r}^{s}\right\|\right) \\
\quad+ \\
\quad C \max \left(\left\|u ; \dot{B}_{q}^{0}\right\|,\left\|v ; \dot{B}_{q}^{0}\right\|\right)^{[s]}\left\|u-v ; \dot{B}_{q}^{0}\right\|^{p-[s]-1} \\
\quad \times \max \left(\left\|u ; \dot{B}_{r}^{s}\right\|,\left\|v ; \dot{B}_{r}^{s}\right\|\right)
\end{array}
\end{aligned}
$$

where the second and third terms on the right hand side of the last inequality are disregarded for $p<2$ and $p \notin([s]+1,[s]+2)$, respectively.

Proof. - We have already shown the first inequality in [16]. The second inequality would be proved analogously and we omit the proof.

For the proof of the next proposition, we describe fundamental relations between $1 / q$ and $1 / r$ with $(1 / q, 1 / r) \in \bigcup_{\varepsilon>0} \Omega_{\varepsilon}$.

Lemma 2.1. - Let $\mu, \rho \in \mathbf{R}$. Let $1 / q, 1 / r$ satisfy

$$
\mu=\rho-n(1 / 2-1 / r)-1 / q
$$

If $\rho, q$ satisfy any of the following conditions, then the above $1 / q, 1 / r$ satisfy $(1 / q, 1 / r) \in \bigcup_{\varepsilon>0} \Omega_{\varepsilon}$.
(1) $n=2$,

$$
\begin{aligned}
& 0 \leqslant 1 / q<n / 2-(\mu-\rho) \quad \text { for } \mu-1<p \leqslant \mu-3 / 4 \\
& 0 \leqslant 1 / q \leqslant(n-1)(\mu-\rho) /(n+1) \quad \text { for } \mu-3 / 4<p \leqslant \mu
\end{aligned}
$$

(2) $n \geqslant 3$,

$$
\begin{aligned}
& 0 \leqslant 1 / q<1 / 2 \quad \text { for } \rho=\mu-(n-1) / 2 \\
& 0 \leqslant 1 / q \leqslant 1 / 2 \quad \text { for } \mu-(n-1) / 2<\rho<\mu-(n+1) /(2 n-2) \\
& 0 \leqslant 1 / q<1 / 2 \text { for } \rho=\mu-(n+1) /(2 n-2) \\
& 0 \leqslant 1 / q \leqslant(n-1)(\mu-\rho) /(n+1) \quad \text { for } \mu-(n+1) /(2 n-2)<\rho \leqslant \mu
\end{aligned}
$$

PROPOSITION 2.2. - Let $n, \mu, p, f$ satisfy any of the assumptions in Theorem 1.1. Let $-\rho_{0}$ be the first index of $G$. Then for sufficiently small $\varepsilon>0$, there exists a pair $\left(1 / q_{0}, 1 / r_{0}\right) \in \Omega_{\varepsilon}$ with

$$
\mu=1-\left(\rho_{0}+n\left(1 / 2-1 / r_{0}\right)-1 / q_{0}\right)
$$

and two triplets

$$
\left(1 / q_{i}, 1 / r_{i}, \rho_{i}\right) \in \Omega_{\varepsilon, \mu}, \quad i=1,2
$$

such that

$$
\begin{equation*}
\left\|f(u) ; L^{q_{0}^{\prime}}\left(I, \dot{B}_{r_{0}^{\prime}}^{-\rho_{0}}\right)\right\| \leqslant C|I|^{\sigma}\|u\|_{1}^{p-1}\|u\|_{2} \tag{2.6}
\end{equation*}
$$

$$
\begin{aligned}
& \left\|f(u)-f(v) ; L^{q_{0}^{\prime}}\left(I, \dot{B}_{r_{0}^{\prime}}^{-\rho_{0}}\right)\right\| \\
& \leqslant \\
& \quad C|I|^{\sigma} \max \left(\|u\|_{1},\|v\|_{1}\right)^{p-1}\|u-v\|_{2} \\
& \quad+C|I|^{\sigma} \max \left(\|u\|_{1},\|v\|_{1}\right)^{p-2}\|u-v\|_{1} \max \left(\|u\|_{2},\|v\|_{2}\right) \\
& \quad+C|I|^{\sigma} \max \left(\|u\|_{1},\|v\|_{1}\right)^{\left[-\rho_{0}\right]}\|u-v\|_{1}^{p-\left[-\rho_{0}\right]-1} \max \left(\|u\|_{2},\|v\|_{2}\right)
\end{aligned}
$$

where $\|\cdot\|_{i}=\left\|\cdot ; L^{q_{i}}\left(I, \dot{B}_{r_{i}}^{\rho_{i}}\right)\right\|$ and $\sigma=2-(p-1)(n-2 \mu) / 2$ and the constant $C$ is independent of $I$. On the right hand side of the last inequality, the second and third terms are disregarded for $p<2$ and $p \notin\left(\left[-\rho_{0}\right]+1,\left[-\rho_{0}\right]+2\right)$, respectively.

Proof. - Let $1 / r^{*}=1 / r_{1}-\rho_{1} / n$ and $1 / r^{* *}=1 / r_{2}-\left(\rho_{0}+\rho_{2}\right) / n$. If

$$
\begin{aligned}
& \rho_{1} \geqslant 0, \quad 0 \leqslant-\rho_{0} \leqslant \rho_{2}, \quad 0<1 / r^{*} \leqslant 1 / 2 \\
& 0 \leqslant 1 / r^{* *} \leqslant 1 / 2, \quad 1 / r_{0}^{\prime}=(p-1) / r^{*}+1 / r^{* *}
\end{aligned}
$$

and

$$
\sigma=1 / q_{0}^{\prime}-(p-1) / q_{1}-1 / q_{2} \geqslant 0
$$

then by Proposition 2.1 and the embeddings $\dot{B}_{r_{1}}^{\rho_{1}} \subset \dot{B}_{r^{*}}^{0}, \dot{B}_{r_{2}}^{\rho_{2}} \subset \dot{B}_{r^{* *}}^{-\rho_{0}}$ and the Hölder inequality in time, we obtain the required inequality, where we use the embedding $L^{q} \subset \dot{B}_{q}^{0}$ with $1<q \leqslant 2$ for $\rho_{0}=0$.

By a simple calculation, we see that the above assumptions are satisfied by a pair $\left(1 / q_{0}, 1 / r_{0}\right) \in \Omega_{\varepsilon}$ and $\rho_{0}$ with

$$
1-\mu=\rho_{0}+n\left(1 / 2-1 / r_{0}\right)-1 / q_{0}
$$

and two triplets $\left(1 / q_{i}, 1 / r_{i}, \rho_{i}\right) \in \Omega_{\varepsilon, \mu}, i=1,2$, which satisfy the following conditions:
(1) $\rho_{1} \geqslant 0, \quad 0 \leqslant-\rho_{0} \leqslant \rho_{2}$.
(2) $1 / q_{i}<n / 2-\mu, \quad i=1,2$.
(3) $1 / q_{0}+1 / q_{2}=\ell\left(1 / q_{1}\right) \equiv(p-1)\left(n / 2-\mu-1 / q_{1}\right)-1$.
(4) $\sigma=2-(p-1)(n-2 \mu) / 2 \geqslant 0$.

We show the existence of the above triplets $\left(1 / q_{i}, 1 / r_{i}, \rho_{i}\right), i=0,1,2$, using Lemma 2.1. We make some comments here. By the condition (3), we must assume

$$
\mu<n / 2 \quad \text { and } \quad p-1 \geqslant 2 /(n-2 \mu)
$$

By (4), we must assume $p-1 \leqslant 4 /(n-2 \mu)$, but this is required for the well-posedness of (1.2) in $\dot{H}^{\mu}$.

In the following, we consider the case $n \geqslant 4$ only since the proofs for the case $n=2,3$ are analogous. We make a classification on $\mu$. The problem is reduced to the existence of the required $1 / q_{i}, i=1,2$.

Case $1.0 \leqslant \mu<(n-3) /(2 n-2)$.
Let $\rho_{i}=0, i=0,1,2$. Let

$$
0 \leqslant 1 / q_{0} \leqslant 1 / 2, \quad 0 \leqslant 1 / q_{i} \leqslant(n-1) \mu /(n+1), \quad i=1,2
$$

Then by Lemma 2.1, we have $\left(1 / q_{0}, 1 / r_{0}\right) \in \Omega_{\varepsilon}$ and $\left(1 / q_{i}, 1 / r_{i}, \rho_{i}\right) \in$ $\Omega_{\varepsilon, \mu}, i=1$, 2, for sufficiently small $\varepsilon>0$. Now $1 / q_{i}, i=0,1,2$, must satisfy (2.9), but the existence of such $1 / q_{i}, i=0,1,2$, is guaranteed if $p$ satisfies

$$
\begin{equation*}
2 /(n-2 \mu) \leqslant p-1 \leqslant(3(n+1)+2(n-1) \mu) /\left(n^{2}+n-4 n \mu\right) \tag{2.11}
\end{equation*}
$$

Case 2. $\mu=(n-3)(2 n-2)$.
In Case 1 , with $1 / q_{0} \leqslant 1 / 2$ replaced with $1 / q_{0}<1 / 2$, we conclude the existence of the required $1 / q_{i}, i=0,1,2$, if $p$ satisfies

$$
\begin{equation*}
2 /(n-2 \mu) \leqslant p-1<4 /(n+1-4 \mu) \tag{2.12}
\end{equation*}
$$

In the following cases, the argument after setting $\rho_{i}, 1 / q_{i}, i=0,1,2$, is similar to that of Case 1 , so that we omit it and write the assumption on $p$ only.

Case 3. $(n-3) /(2 n-2)<\mu \leqslant(n+1) /(2 n-2)$.
Let $\rho_{i}=0, i=0,1,2$. Let

$$
\begin{aligned}
& 0 \leqslant 1 / q_{0} \leqslant(n-1)(1-\mu) /(n+1) \\
& 0 \leqslant 1 / q_{i} \leqslant(n-1) \mu /(n+1), \quad i=1,2
\end{aligned}
$$

for $\mu<(n+1) /(2 n-2)$,

$$
0 \leqslant 1 / q_{i}<1 / 2, \quad i=1,2
$$

for $\mu=(n+1) /(2 n-2)$. The required assumption on $p$ is

$$
2 /(n-2 \mu) \leqslant p-1 \leqslant \begin{cases}4 /(n+1-4 \mu) & \text { if } \mu<1 / 2  \tag{2.13}\\ 4 /(n-2 \mu) & \text { if } \mu \geqslant 1 / 2\end{cases}
$$

Case 4. $(n+1) /(2 n-2)<\mu \leqslant \min (1, \alpha)$.
Let $\rho_{i}=0, i=0,1,2$. Let

$$
0 \leqslant 1 / q_{0} \leqslant(n-1)(1-\mu) /(n+1), \quad 0 \leqslant 1 / q_{i} \leqslant 1 / 2, \quad i=1,2
$$

The assumption on $p$ is

$$
2 /(n-2 \mu) \leqslant p-1 \leqslant 4 /(n-2 \mu)
$$

We refer to the constant $\alpha$ which depends on the spatial dimension. By the condition (2.9), we must assume at least

$$
\ell(1 / 2) \leqslant(n-1)(1-\mu) /(n+1)+1 / 2
$$

which is equivalent to

$$
\begin{equation*}
p-1 \leqslant(2(n-1)(1-\mu) /(n+1)+3) /(n-1-2 \mu) . \tag{2.14}
\end{equation*}
$$

To enlarge the right hand side than $4 /(n-2 \mu), \mu$ must satisfy $F(\mu) \geqslant 0$.
But this is guaranteed if $\mu \leqslant \alpha$ since $\alpha$ is the lower root of $F(x)=0$.

$$
\begin{aligned}
& \text { Case 5. } n \geqslant 7, \alpha<\mu<(n-4) / 2 \\
& \text { Let }-\rho_{0}=\rho_{1}=\rho_{2}=\mu-\beta(\mu) \text {. Let } \\
& 0 \leqslant 1 / q_{0} \leqslant(n-1)(1-\beta(\mu)) /(n+1), \quad 0 \leqslant 1 / q_{i} \leqslant 1 / 2, \quad i=1,2
\end{aligned}
$$

The assumption on $p$ is

$$
2 /(n-2 \mu) \leqslant p-1 \leqslant 4 /(n-2 \mu)
$$

Vol. 71, $\mathrm{n}^{\circ}$ 2-1999.

We refer to $\beta(\mu)$ which depends on the spatial dimension and $\mu$. By the condition (2.9), we must assume at least

$$
\begin{equation*}
p-1 \leqslant(2(n-1)(1-\beta(\mu)) /(n+1)+3) /(n-1-2 \mu) \tag{2.15}
\end{equation*}
$$

but the right hand side is equal to $4 /(n-2 \mu)$ by the definition of $\beta(\mu)$.
Case 6. $\max (1,(n-4) / 2) \leqslant \mu<n / 2$.
Let $-\rho_{0}=\rho_{1}=\rho_{2}=\mu-1$. Let

$$
1 / q_{0}=0, \quad 0 \leqslant 1 / q_{i} \leqslant 1 / 2, \quad i=1,2
$$

and

$$
1 / q_{i}<n / 2-\mu, \quad i=1,2
$$

The assumption on $p$ is

$$
2 /(n-2 \mu) \leqslant p-1 \leqslant 4 /(n-2 \mu)
$$

## 3. PROOF OF THEOREM 1.1

We prove Theorem 1.1 in this section.
Proof of Theorem 1.1. - First of all, we recall the following inequalities by Proposition 3. 1 in [10].

$$
\begin{align*}
& \left\|\dot{K}(t) \phi ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant C\left\|\phi ; \dot{H}^{\mu}\right\|  \tag{3.16}\\
& \left\|K(t) \psi ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant C\left\|\psi ; \dot{H}^{\mu-1}\right\|  \tag{3.17}\\
& \left\|\int_{0}^{t} K(t-\tau) h(\tau) d \tau ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant C\left\|h ; L^{q_{0}^{\prime}}\left(I, \dot{B}_{r_{0}^{\prime}}^{-\rho_{0}}\right)\right\| \tag{3.18}
\end{align*}
$$

for any $\mu, \rho, \rho_{0} \in \mathbf{R}$ and

$$
(1 / q, 1 / r),\left(1 / q_{0}, 1 / r_{0}\right) \in \Omega_{\varepsilon}
$$

with

$$
\mu=\rho+n(1 / 2-1 / r)-1 / q=1-\left(\rho_{0}+n\left(1 / 2-1 / r_{0}\right)-1 / q_{0}\right)
$$

where $C$ is a constant independent of $I$.
Let $n, \mu, p, f$ satisfy any of the assumptions in Theorem 1.1. Let $\varepsilon$, $1 / q_{i}, 1 / r_{i}, \rho_{i}, i=0,1,2$, be those in Proposition 2.2. By the above inequalities and Proposition 2.2, we have

$$
\begin{align*}
& \left\|\Phi(u) ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \\
& \quad \leqslant C\|(\phi, \psi)\|_{\mu}+C\left\|f(u) ; L^{q_{0}^{\prime}}\left(I, \dot{B}_{r_{r}^{\prime}}^{-\rho_{0}}\right)\right\| \\
& \quad \leqslant C\|(\phi, \psi)\|_{\mu}+C|I|^{\sigma}\left\|u ; L^{q_{1}}\left(I, \dot{B}_{r_{1}}^{\rho_{1}}\right)\right\|^{p-1}\left\|u ; L^{q_{2}}\left(I, \dot{B}_{r_{2}}^{\rho_{2}}\right)\right\| \tag{3.19}
\end{align*}
$$

for any $(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}$, where $C$ is independent of $I$. Therefore we obtain

$$
\begin{equation*}
\max _{(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}}\left\|\Phi(u) ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant C\|(\phi, \psi)\|_{\mu}+C|I|^{\sigma} R^{p}, \tag{3.20}
\end{equation*}
$$

for any $u \in X_{\varepsilon}(I, R)$. Similarly we have

$$
\begin{align*}
& d(\Phi(u), \Phi(v)) \\
& \quad \leqslant C|I|^{\sigma} R^{p-1} d(u, v)+C|I|^{\sigma} R^{\left[-\rho_{0}\right]+1} d(u, v)^{p-\left[-\rho_{0}\right]-1}, \tag{3.21}
\end{align*}
$$

for any $u, v \in X_{\varepsilon}(I, R)$, where $C$ is independent of $I$ and the second term on the right hand side of (3.21) is disregarded for $p \notin J$. If $p \notin J$, then the unique solution of (1.2) is given by the standard contraction argument on ( $X_{\varepsilon}(I, R), d$ ) with $R$ sufficiently large and $|I|>0$ sufficiently small for the local solution, with $R$ and $\|(\phi, \psi)\|_{\mu}$ sufficiently small for the global solution. If $p \in J$, then we have only to consider the case

$$
(n+1) /(2 n-2)<\mu<n / 2 .
$$

Let $|I|$ and $R$ satisfy

$$
\begin{equation*}
C\|(\phi, \psi)\|_{\mu}+C|I|^{\sigma} R^{p} \leqslant R, \quad C|I|^{\sigma} R^{p-1}<1, \tag{3.22}
\end{equation*}
$$

and let $\|(\phi, \psi)\|_{\mu}$ be sufficiently small for $\sigma=0$. Let $u_{0}=0$ and $u_{i+1} \equiv \Phi\left(u_{i}\right)$ for $i=1,2, \ldots$. Then there is a subsequence $\left\{u_{i_{k}}\right\}_{k} \subset\left\{u_{i}\right\}_{i}$ and $u \in X_{\varepsilon}(I, R)$ such that $u_{i_{k}}$ converges to $u$ in the distribution sense as $k \rightarrow \infty$. On the other hand, let

$$
(n-3) /(2 n-2)<\mu_{0}<(n+1) /(2 n-2)
$$

and let $\lambda>0$ and

$$
\Lambda(\lambda) \equiv\left\{(t, x) \in \mathbf{R}^{n+1}| | x|<\lambda-|t|\},\right.
$$

then we have for sufficiently small $\varepsilon>0$,

$$
\begin{align*}
& \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u_{i+2}-u_{i+1} ; L^{q} L^{r}(\Lambda(\lambda))\right\| \\
& \quad \leqslant C|I|^{\sigma} R^{p-1} \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u_{i+1}-u_{i} ; L^{q} L^{r}(\Lambda(\lambda))\right\| \tag{3.23}
\end{align*}
$$

Indeed, let $w$ and $w_{\lambda}$ satisfy

$$
\left(\partial_{t}^{2}-\Delta\right) w=h, \quad w(0)=\partial_{t} w(0)=0
$$

and

$$
\left(\partial_{t}^{2}-\Delta\right) w_{\lambda}=h \chi_{\Lambda}(\lambda), \quad w_{\lambda}(0)=\partial_{t} w_{\lambda}(0)=0
$$

then $w=w_{\lambda}$ on $\Lambda(\lambda)$, where $\chi_{\Lambda(\lambda)}$ is a characteristic function on $\Lambda(\lambda)$. By this fact and (3.18) and the argument as described in the proof of Proposition 2.2, we obtain the above inequality.

By (3.23), we conclude that $\left\{u_{i}\right\}$ converges to some $v_{\lambda}$ strongly in $L^{q} L^{r}(\Lambda(\lambda))$ for any $(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}$, so that $u=v_{\lambda}$ on $\Lambda(\lambda)$. Therefore we have for any $\lambda>0$

$$
\begin{aligned}
& \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|\Phi(u)-u ; L^{q} L^{r}(\Lambda(\lambda))\right\| \\
& \quad \leqslant \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|\Phi(u)-\Phi\left(u_{i}\right)+u_{i+1}-u ; L^{q} L^{r}(\Lambda(\lambda))\right\| \\
& \leqslant C|I|^{\sigma} R^{p-1} \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u-u_{i} ; L^{q} L^{r}(\Lambda(\lambda))\right\| \\
& \rightarrow 0 \quad \text { as } i \rightarrow \infty,
\end{aligned}
$$

by which we conclude that $u=\Phi(u) a . e(t, x) \in I \times \mathbf{R}^{n}$, namely $u=$ $\Phi(u)$ in $\left(X_{\varepsilon}(I, R), d\right)$.

The uniqueness of the solution also follows from (3.23).
(1) The continuity of the solution $\left(u, \partial_{t} u\right)$ in time with respect to the $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$-norm follows from the Lebesgue convergence theorem. The proof is standard and we omit it.
(2) For the continuous dependence on the initial data of its solution, we consider the case $p \in J$ only since for $p \notin J$ the last term of (3.21) is disregarded, so that we can use the standard argument [3]. By (3.23), we have

$$
\begin{aligned}
& \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u-v ; L^{q} L^{r}(\Lambda(\lambda))\right\| \\
& \leqslant C_{\lambda}\left\|\left(\phi-\phi_{0}, \psi-\psi_{0}\right)\right\|_{\mu}+C|I|^{\sigma} R^{p-1} \\
& \quad \times \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u-v ; L^{q} L^{r}(\Lambda(\lambda))\right\|
\end{aligned}
$$

Annales de l'Institut Henri Poincaré - Physique théorique
where $C_{\lambda}$ is a constant dependent on $\lambda$, but not on $I$. So that we conclude $v \rightarrow u$ in

$$
\bigcap_{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}} L^{q} L^{r}(\Lambda(\lambda))
$$

as $\left(\phi_{0}, \psi_{0}\right)$ tends to $(\phi, \psi)$, by which we conclude that $v$ converges to $u$ in $\mathcal{D}^{\prime}\left(\mathbf{R}^{n+1}\right)$, as required.
(3) Let $\left(\phi_{+}, \psi_{+}\right)$be

$$
\phi_{+} \equiv \phi+\int_{0}^{\infty} K(-\tau) f(u(\tau)) d \tau, \quad \psi_{+} \equiv \psi+\int_{0}^{\infty} \dot{K}(-\tau) f(u(\tau)) d \tau
$$

Then we have

$$
\begin{aligned}
& \left\|u(t)-\dot{K}(t) \phi_{+}-K(t) \psi_{+} ; \dot{H}^{\mu}\right\| \\
& \quad \leqslant\left\|\int_{t}^{\infty} K(t-\tau) f(u(\tau)) d \tau ; \dot{H}^{\mu}\right\| \\
& \quad \leqslant C\left\|u ; L^{q_{1}}\left((t, \infty), \dot{B}_{r_{1}}^{\rho_{1}}\right)\right\|^{p-1}\left\|u ; L^{q_{2}}\left((t, \infty), \dot{B}_{r_{2}}^{\rho_{2}}\right)\right\|,
\end{aligned}
$$

where we have used a similar result to (3.18) and Proposition 2.2, and we can take $1 / q_{i}, i=1,2$, for $1 / q_{i} \neq \infty$ since $p-1=4 /(n-2 \mu)$. Therefore we have

$$
\left\|u(t)-\dot{K}(t) \phi_{+}-K(t) \psi_{+} ; \dot{H}^{\mu}\right\| \rightarrow 0 \quad \text { as } t \rightarrow \infty
$$

(4) For $\left(\phi_{-}, \psi_{-}\right) \in \dot{H}^{\mu} \times \dot{H}^{\mu-1}$, let $\Phi_{-}$be an operator defined by

$$
\begin{equation*}
\Phi_{-}(u)(t) \equiv \dot{K}(t) \phi_{-}+K(t) \psi_{-}+\int_{-\infty}^{t} K(t-\tau) f(u(\tau)) d \tau \tag{3.24}
\end{equation*}
$$

Similarly to $\Phi$, we have

$$
\begin{equation*}
\max _{(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}}\left\|\Phi_{-}(u) ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant C\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}+C R^{p} \tag{3.25}
\end{equation*}
$$

$$
\begin{align*}
& d\left(\Phi_{-}(u), \Phi_{-}(v)\right) \\
& \quad \leqslant C R^{p-1} d(u, v)+C|I|^{\sigma} R^{\left[-\rho_{0}\right]+1} d(u, v)^{p-\left[-\rho_{0}\right]-1} \tag{3.26}
\end{align*}
$$

for any $u, v \in X_{\varepsilon}(I, R)$, where the second term on the right hand side of (3.26) is disregarded for $p \notin J$. Therefore for $p \notin J$ we have the Vol. 71, n ${ }^{\circ}$ 2-1999.
unique fixed point of $\Phi_{-}$in $X_{\varepsilon}(I, R)$ by a contraction argument with $\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}$ and $R$ sufficiently small. We show that for $p \in J$ we also have a fixed point of $\Phi_{-}$in $X_{\varepsilon}(I, R)$ with $\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}$ and $R$ sufficiently small in the following. We may assume

$$
(n+1) /(2 n-2)<\mu<n / 2
$$

Let

$$
(n-3) /(2 n-2)<\mu_{0}<(n+1) /(2 n-2)
$$

Let $R_{0}>0$. Let $X_{\varepsilon}\left(I, R, R_{0}\right)$ and $d_{0}$ be

$$
\begin{aligned}
X_{\varepsilon}\left(I, R, R_{0}\right) & \equiv\left\{u \in X_{\varepsilon}(I, R) \mid \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u ; L^{q}\left(I, L^{r}\right)\right\| \leqslant R_{0}\right\} \\
d_{0}(u, v) & \equiv \max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|u-v ; L^{q}\left(I, L^{r}\right)\right\|
\end{aligned}
$$

for any $u, v \in X_{\varepsilon}\left(I, R, R_{0}\right)$. Then similarly to (3.25) and (3.26), we have

$$
\begin{gather*}
\max _{(1 / q, 1 / r, 0) \in \Omega_{\varepsilon, \mu_{0}}}\left\|\Phi_{-}(u) ; L^{q}\left(I, L^{r}\right)\right\| \\
\leqslant C\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu_{0}}+C R^{p-1} R_{0},  \tag{3.27}\\
\max _{(1 / q, 1 / r, \rho) \in \Omega_{\varepsilon, \mu}}\left\|\Phi_{-}(u) ; L^{q}\left(I, \dot{B}_{r}^{\rho}\right)\right\| \leqslant C\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}+C R^{p},  \tag{3.28}\\
d_{0}\left(\Phi_{-}(u), \Phi_{-}(v)\right) \leqslant C R^{p-1} d_{0}(u, v) . \tag{3.29}
\end{gather*}
$$

So that if $\left(\phi_{-}, \psi_{-}\right) \in \dot{H}^{\mu_{0}} \times \dot{H}^{\mu_{0}-1}$ and if $\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}$ and $R$ are sufficiently small and $R_{0}$ sufficiently large, then $\Phi_{-}$becomes a contraction map on $X_{\varepsilon}\left(I, R, R_{0}\right)$ with the metric $d_{0}$. Therefore we obtain the unique fixed point of $\Phi_{-}$. Let $\left\|\left(\phi_{-}, \psi_{-}\right)\right\|_{\mu}$ be sufficiently small. Let $\left\{\left(\phi_{i}, \psi_{i}\right)\right\}_{i=1}^{\infty}$ be a sequence such that

$$
\left(\phi_{i}, \psi_{i}\right) \rightarrow\left(\phi_{-}, \psi_{-}\right)
$$

in $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$ as $i \rightarrow \infty$ and

$$
\left(\phi_{i}, \psi_{i}\right) \in \dot{H}^{\mu_{0}} \times \dot{H}^{\mu_{0}-1}
$$

Then by the above argument, there exists $u_{i} \in X_{\varepsilon}\left(I, R, R_{0}\right)$ which satisfies

$$
\begin{equation*}
u_{i}=\dot{K}(t) \phi_{i}+K(t) \psi_{i}+\int_{-\infty}^{t} K(t-\tau) f\left(u_{i}(\tau)\right) d \tau \tag{3.30}
\end{equation*}
$$

for $i$ sufficiently large. We can take a subsequence of $\left\{u_{i}\right\}$ which converges to some $u$ in the distribution sense. This $u$ is the required fixed point of $\Phi_{-}$in $X_{\varepsilon}(I, R)$. For details, we refer to the discussion before Lemma 7.1 and itself in [15]. The result

$$
\left\|u(t)-\dot{K}(t) \phi_{-}-K(t) \psi_{-} ; \dot{H}^{\mu}\right\| \rightarrow 0 \quad \text { as } t \rightarrow-\infty
$$

now follows similarly to the proof of (3).
Next we show that the scattering map

$$
\left(\phi_{-}, \psi_{-}\right) \rightarrow\left(\phi_{+}, \psi_{+}\right)
$$

is continuous in the neighborhood at the origin in $\dot{H}^{\mu} \times \dot{H}^{\mu-1}$ for $p \notin J$. By the proof of (3) and (4), we have the following relation between $\left(\phi_{-}, \psi_{-}\right)$and $\left(\phi_{+}, \psi_{+}\right)$as

$$
\begin{align*}
& \phi_{+}=\phi_{-}+\int_{-\infty}^{\infty} K(-\tau) f(u(\tau)) d \tau \\
& \psi_{+}=\psi_{-}+\int_{-\infty}^{\infty} \dot{K}(-\tau) f(u(\tau)) d \tau \tag{3.31}
\end{align*}
$$

where $u$ is the solution of $u=\Phi_{-}(u)$. Let

$$
\left(\left(\tilde{\phi}_{-}, \tilde{\psi}_{-}\right), \tilde{u},\left(\tilde{\phi}_{+}, \tilde{\psi}_{+}\right)\right)
$$

be another triplet. It suffices to show that

$$
\begin{equation*}
\left\|\left(\phi_{+}-\widetilde{\phi}_{+}, \psi_{+}-\widetilde{\psi}_{+}\right)\right\|_{\mu} \rightarrow 0 \quad \text { as }\left\|\left(\phi_{-}-\widetilde{\phi}_{-}, \psi_{-}-\widetilde{\psi}_{-}\right)\right\|_{\mu} \rightarrow 0 \tag{3.32}
\end{equation*}
$$

Similarly to the proof of (3.21), we have

$$
\begin{equation*}
\left\|\phi_{+}-\widetilde{\phi}_{+} ; \dot{H}^{\mu}\right\| \leqslant\left\|\phi_{-}-\widetilde{\phi}_{-} ; \dot{H}^{\mu}\right\|+C R^{p-1} d(u, \widetilde{u}) \tag{3.33}
\end{equation*}
$$

and

$$
\begin{equation*}
d(u, \tilde{u}) \leqslant C\left\|\left(\phi_{-}-\widetilde{\phi}_{-}, \psi_{-}-\widetilde{\psi}_{-}\right)\right\|_{\mu}+C R^{p-1} d(u, \tilde{u}) \tag{3.34}
\end{equation*}
$$

Since $C R^{p-1}<1$, we conclude that

$$
\left\|\phi_{+}-\widetilde{\phi}_{+} ; \dot{H}^{\mu}\right\| \rightarrow 0
$$

Vol. 71, $\mathrm{n}^{\circ}$ 2-1999.
as $\left\|\left(\phi_{-}-\widetilde{\phi}_{-}, \psi_{-}-\widetilde{\psi}_{-}\right)\right\|_{\mu}$ tends to zero. For $\left\|\psi_{+}-\widetilde{\psi}_{+} ; \dot{H}^{\mu-1}\right\|$, the proof is analogous.

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