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## ABOUT THE INFLUENCE OF OSCILLATIONS ON STRICHARTZ-TYPE DECAY ESTIMATES


#### Abstract

Starting with the well-known Strichartz decay estimate for the wave equation we are interested in a similar estimate for wave equations with a time dependent coefficient. The model under consideration is the strictly hyperbolic equation $u_{t t}-a(t) \Delta u=0$. By the aid of an example we illustrate the deteriorating influence of oscillations in $a=a(t)$ on decay estimates. Moreover we prove, that in the case of slow oscillations one gets Strichartz-type decay estimates with a decay rate similar to the classical one.


## 1. Introduction

To prove global existence results (in time) of small data solutions for the Cauchy problem for nonlinear wave equations Strichartz decay estimate [9] plays an important role. That is the following estimate: there exist constants $C$ and $L$ depending on $p$ and $n$ such that

$$
\begin{equation*}
\left\|u_{t}(t, \cdot)\right\|_{L_{q}}+\left\|\nabla_{x} u(t, \cdot)\right\|_{L_{q}} \leq C(1+t)^{-\frac{n-1}{2}\left(\frac{1}{p}-\frac{1}{q}\right)}\|\psi\|_{W_{p}^{L}} \tag{1}
\end{equation*}
$$

where $1 \leq p \leq 2,1 / p+1 / q=1$ and $u=u(t, x)$ is the solution to

$$
u_{t t}-\Delta u=0, \quad u(0, x)=0, \quad u_{t}(0, x)=\psi(x) \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)
$$

If one is interested in such estimates for $u_{t t}-a(t) \Delta u=0$, then one has to explain properties of $a=a(t)$ which take influence on Strichartz-type decay estimates.

In [5] the weakly hyperbolic Cauchy problem

$$
u_{t t}-t^{2 l} \Delta u=0, \quad u(x, 0)=\varphi(x), \quad u_{t}(x, 0)=\psi(x)
$$

was studied and the decay estimates are derived there imply the hypothesis that increasing parts in $a=a(t)$ have an improving influence on decay estimates.
To feel the influence of oscillating parts in $a=a(t)$ let us consider

$$
\begin{equation*}
u_{t t}-(1+\varepsilon \sin t)^{2} \Delta u=0, \quad u(x, 0)=\varphi(x), \quad u_{t}(x, 0)=\psi(x) \tag{2}
\end{equation*}
$$

with a sufficiently small $\varepsilon$. Although the coefficient is near to 1 in usual topologies, decay estimates are in general not valid. Is this a surprise or not? Independent of the point of view we feel, that oscillating parts in $a=a(t)$ have a deteriorating influence on decay estimates.
For (2) this follows from Theorem 1 will be proved in Section 2. In this regard, we have found it necessary from the results of [10] to resort to some specialized representation of the coefficient $a=a(t)$ relating to increasing or oscillating part:

$$
\begin{equation*}
u_{t t}-\lambda(t)^{2} b(t)^{2} \Delta u=0 \tag{3}
\end{equation*}
$$

where

- $\lambda=\lambda(t)$ describes the increasing part,
- $b=b(t)$ describes the oscillating part.

In Section 3 we give the definition of slow oscillations which describes a special interplay between both parts. For the case of slow oscillations Strichartz-type decay estimates will be proved. The main result (Theorem 2) can be applied to (3), where possible $\lambda$ and $b$ are given in the next examples.

Example 1 (Logarithmic Growth). $\lambda(t)=\ln (e+t), b(t)=2+\sin \left((\ln (e+t))^{\beta+1}\right)$, $\beta \in[0,1)$.

Example 2 (potential growth). $\lambda(t)=(1+t)^{\alpha}, \alpha>0, b(t)$ as in Example 1.
EXAMPLE 3 (EXPONENTIAL GROWTH). $\lambda(t)=\exp \left(t^{\alpha}\right), \alpha>1 / 2, b=b(t)$ is periodic, positive, non-constant and smooth.

EXAMPLE 4 (SUPEREXPONENTIAL GROWTH). $\lambda(t)=\exp \left(\exp \left(t^{\alpha}\right)\right), \alpha>0, b=b(t)$ as in Example 3.

## 2. Wave equations with a periodic coefficient

The bad influence of oscillations in the coefficient $a=a(t)$ on decay estimates will be clear by employing Theorem 1. We claim:

Theorem 1. Consider the Cauchy problem

$$
\begin{equation*}
u_{t t}-b(t)^{2} \Delta u=0, \quad u(0, x)=\varphi(x), \quad u_{t}(0, x)=\psi(x) \tag{4}
\end{equation*}
$$

where $b=b(t)$ defined on $\mathbb{R}$, is a 1-periodic, non-constant, smooth, and positive function. Then there are no constants $q, p, C, L$ and a nonnegative function $f$ defined on $\mathbb{N}$ such that for every initial data $\varphi, \psi \in C_{0}^{\infty}(\mathbb{R})$ the estimate

$$
\begin{equation*}
\left\|u_{t}(m, \cdot)\right\|_{L_{q}}+\left\|\nabla_{x} u(m, \cdot)\right\|_{L_{q}} \leq C f(m)\left(\|\varphi\|_{W_{p}^{L}}+\|\psi\|_{W_{p}^{L}}\right) \tag{5}
\end{equation*}
$$

is fulfilled for all $m \in \mathbb{N}$, while $f(m) \rightarrow \infty, \ln f(m)=o(m)$ as $m \rightarrow \infty$.
Remark 1. The conditions for $f=f(m)$ are very near to optimal ones. Indeed it is well-known, that due to Gronwall's inequality one can prove the energy estimate

$$
\left\|u_{t}(t, \cdot)\right\|_{L_{2}}+\left\|\nabla_{x} u(t, \cdot)\right\|_{L_{2}} \leq C \exp \left(C_{0} t\right)\left(\|\varphi\|_{W_{2}^{1}}+\|\psi\|_{L_{2}}\right), \quad t \in[0, \infty)
$$

for the solution of (4). Choosing $t=m, m \in \mathbb{N}, p=q=2, L=1$, we obtain an inequality like (5) with $\ln f(m)=O(m)$ as $m \rightarrow \infty$.

REmARK 2. We will construct a special sequence $\left\{\left(\varphi_{M}, \psi_{M}\right)\right\}_{M \in N}$ of data for which (5) is violated at least for large $M$. An interesting question is that for a classification of data: data which allow, don't allow respectively, a decay estimate for the solution of (4).

Proof. The proof is divided into several steps.

1. One lemma for ordinary differential equations with a periodic coefficient

Consider the ordinary differential equation with a periodic coefficient

$$
\begin{equation*}
w_{t t}+\lambda b(t)^{2} w=0 \tag{6}
\end{equation*}
$$

Let the matrix-valued function $X=X\left(t, t_{0}\right)$ depending on $\lambda$ be the solution of the Cauchy problem

$$
d_{t} X=\left(\begin{array}{cc}
0 & -\lambda b(t)^{2}  \tag{7}\\
1 & 0
\end{array}\right) X, \quad X\left(t_{0}, t_{0}\right)=\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right) .
$$

Thus, $X=X\left(t, t_{0}\right)$ gives a fundamental solution to (6). The matrix $X(t+1, t)$ is independent of $t \in \mathbb{N}$. Set

$$
X(1,0)=:\left(\begin{array}{ll}
b_{11} & b_{12} \\
b_{21} & b_{22}
\end{array}\right)
$$

The following lemma follows from the considerations in [1] and provides an important tool for our proof.

Lemma 1 ([10]). If $b=b(t)$ is a non-constant, positive, smooth function on $\mathbb{R}$ which is 1-periodic, then there exists a positive $\lambda_{0}$ such that the corresponding matrix $X\left(t_{0}+1, t_{0}\right)$ of (7) has eigenvalues $\mu_{0}$ and $\mu_{0}^{-1}$ with $\left|\mu_{0}\right|>1$.
2. Lower bound for the energy

We use the periodicity of $b$ and the eigenvalue $\mu_{0}$ of $X(1,0)$ to construct solutions of (6) with prescribed values on a discrete set. Thus we get lower bounds for the energy of the solution on this set.

Lemma 2. Let $w=w(t)$ be the solution of (6) with $\lambda=\lambda_{0}$ and initial data $w(0)=1$, $w_{t}(0)=b_{12} /\left(\mu_{0}-b_{11}\right)$. Then for every positive integer number $M \in \mathbb{N}$ the solution satisfies $w(M)=\mu_{0}^{M}$.

Proof. For the solution $w=w(t)$ we have

$$
\binom{d_{t} w(M)}{w(M)}=(X(1,0))^{M}\binom{d_{t} w(0)}{w(0)} .
$$

The matrix

$$
B:=\left(\begin{array}{cc}
\frac{b_{12}}{\mu_{0}-b_{11}} & 1 \\
1 & \frac{b_{21}}{\mu_{0}^{-1}-b_{22}}
\end{array}\right)
$$

is a diagonalizer for $X(1,0)$, that is,

$$
X(1,0) B=B\left(\begin{array}{cc}
\mu_{0} & 0 \\
0 & \mu_{0}^{-1}
\end{array}\right)
$$

Hence,

$$
\binom{d_{t} w(M)}{w(M)}=\mu_{0}^{M}\binom{\frac{b_{12}}{\mu_{0}-b_{11}}}{1}, \quad w(M)=\mu_{0}^{M}
$$

## 3. Construction of unstable solutions

Let us construct a family of solutions $\left\{u_{M}\right\}$ of (4) with data $\left\{\left(\varphi_{M}, \psi_{M}\right)\right\}$. These functions will violate (5) for sufficiently large $M$. They will be called unstable solutions. With a cut-off function $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right), \chi(x)=1$ when $|x| \leq 1, \chi(x)=0$ when $|x| \geq 2$, let us choose the initial data

$$
\begin{equation*}
\varphi_{M}(x)=e^{i x \cdot y} \chi\left(\frac{x}{M^{2}}\right), \quad \psi_{M}(x)=e^{i x \cdot y} \chi\left(\frac{x}{M^{2}}\right) \frac{b_{12}}{\mu_{0}-b_{11}} . \tag{8}
\end{equation*}
$$

Here $y$ is chosen such that $|y|^{2}=\lambda_{0}$. Using the theory of the Cauchy problem for strictly hyperbolic equations there exists a unique solution $u_{M}=u_{M}(t, x)$ to (4), (8), where $u_{M}(t, \cdot)$ has compact support for every given $t \in[0, \infty)$. Let $B_{1}(0) \subset \mathbb{R}^{n}$ be the ball of radius 1 centered at the origin. Assuming (5) we have

$$
\begin{align*}
\left\|\partial_{t} u_{M}(M, \cdot)\right\|_{L_{q}\left(B_{1}(0)\right)}+\left\|\nabla_{x} u_{M}(M, \cdot)\right\|_{L_{q}\left(B_{1}(0)\right)} & \leq\left\|\partial_{t} u(M, \cdot)\right\|_{L_{q}}+\left\|\nabla_{x} u(M, \cdot)\right\|_{L_{q}} \\
& \leq C f(M)\left\|e^{i x \cdot y} \chi\left(\frac{x}{M^{2}}\right)\right\|_{W_{p}^{L}}  \tag{9}\\
& \leq C_{L} f(M) M^{2 n / p} .
\end{align*}
$$

4. The role of the cone of dependence

If we take into account the cone of dependence for the Cauchy problem (4), then the solution $u_{M}$ is representable in $B_{1}(0)$ at time $t=M$ as

$$
\begin{equation*}
u_{M}(M, x)=e^{i x \cdot y} w(M), \quad \partial_{t} u_{M}(M, x)=e^{i x \cdot y} w_{t}(M) \tag{10}
\end{equation*}
$$

where $w$ is taken from Lemma 2. Indeed, for the set $\{(x, t) ;|x| \leq 1, t=M\}$ let us calculate the lower base at $t=0$ of the truncated cone with the slope $\max _{t \in[0,1]} b(t)$ and height $M$. For sufficiently large $M$ this lower base is contained in the ball

$$
B_{\text {depend }, M}:=\left\{x \in \mathbb{R}^{n}:|x| \leq 2 M \max _{t \in[0,1]} b(t)\right\}
$$

on $t=0$. If $x \in B_{\text {depend, } M}$, then $\left|\frac{x}{M^{2}}\right| \leq 1$ for large $M$. Consequently, $\chi\left(\frac{x}{M^{2}}\right)=1$ on $B_{\text {depend, } M}$ for large $M$.
The function $e^{i x \cdot y} w(t)$ solves (4) and takes at $t=0$ the data (8) if $|x| \leq M^{2}$. The cone of dependence property yields (10).
5. Completion of the proof

We obtained

$$
\begin{aligned}
\left\|\partial_{t} u_{M}(M, \cdot)\right\|_{L_{q}\left(B_{1}(0)\right)} & +\left\|\nabla_{x} u_{M}(M, \cdot)\right\|_{L_{q}\left(B_{1}(0)\right)} \\
& =\left(\left|w_{t}(M)\right|+\lambda_{0}^{1 / 2}|w(M)|\right)\left(\text { meas } B_{1}(0)\right)^{1 / q} .
\end{aligned}
$$

In view of (9) we have

$$
|w(M)|+\left|w_{t}(M)\right| \leq C_{L} f(M) M^{2 n / p}
$$

for all large $M$. But this is a contradiction to the statement of Lemma 2. Thus, Theorem 1 is completely proved.

## 3. Wave equations with slow oscillations in the time-dependent coefficient

### 3.1. Classification of oscillations

Definition 1. Let us suppose that there exists a real $\beta \in[0,1]$ such that the following condition is satisfied:

$$
\begin{equation*}
\left|d_{t} b(t)\right| \leq C \frac{\lambda(t)}{\Lambda(t)}(\ln \Lambda(t))^{\beta} \quad \text { for large } \quad t \tag{11}
\end{equation*}
$$

where $\Lambda(t):=\int_{0}^{t} \lambda(s) d s$. If $\beta \in[0,1), \beta=1$, respectively, we call the oscillations slow oscillations, fast oscillations, respectively. If (11) is not satisfied for $\beta=1$, then we call the oscillations very fast oscillations.

REMARK 3. Very fast oscillations may destroy $L_{p}-L_{q}$ decay estimates. These oscillations give us an exact description of a fairly wide class of equations in which the oscillating part dominates the increasing one. In [6] it is shown that one can prove in this case a statement similar to Theorem 1.

REMARK 4. The case of fast oscillations is studied in [7]. We could derive $L_{p}-L_{q}$ decay estimates only for large dimension $n$. Moreover, the behaviour of $b$ and $\lambda$ and its first two derivatives has an influence on the decay rate.

The goal of the following considerations is to show that for slow oscillations $(\beta \in[0,1)$ in (11)) we have $L_{p}-L_{q}$ decay estimates similar to Strichartz decay estimate (1)

- for any dimension $n \geq 2$,
- with a decay rate which coincides with the classical decay rate,
- with the decay function $1+\Lambda(t)$,
- without an essential influence of the oscillating part.


### 3.2. Main result and philosophy of our approach

Let us study

$$
\begin{equation*}
u_{t t}-\lambda(t)^{2} b(t)^{2} \Delta u=0, \quad u(0, x)=\varphi(x), \quad u_{t}(0, x)=\psi(x), \tag{12}
\end{equation*}
$$

under the following assumptions for the positive coefficient $\lambda(t)^{2} b(t)^{2}$ :

- it holds

$$
\begin{equation*}
\Lambda(t) \rightarrow \infty \quad \text { as } \quad t \rightarrow \infty ; \tag{13}
\end{equation*}
$$

- there exist positive constants $c_{0}, c_{1}$ and $c$ such that

$$
\begin{equation*}
c_{0} \frac{\lambda(t)}{\Lambda(t)} \leq \frac{\lambda^{\prime}(t)}{\lambda(t)} \leq c_{1} \frac{\lambda(t)}{\Lambda(t)} \leq c(\ln \Lambda(t))^{c} \quad \text { for large } \quad t \tag{14}
\end{equation*}
$$

- there exist positive constants $c_{k}$ such that for all $k=2,3, \ldots$ it holds

$$
\begin{equation*}
\left|d_{t}^{k} \lambda(t)\right| \leq c_{k}\left(\frac{\lambda(t)}{\Lambda(t)}\right)^{k} \lambda(t) \quad \text { for large } \quad t \tag{15}
\end{equation*}
$$

- with two positive constants $d_{0}$ and $d_{1}$ we have

$$
\begin{equation*}
d_{0} \leq b^{2}(t) \leq d_{1} \quad \text { for } \quad t \in[0, \infty) \tag{16}
\end{equation*}
$$

- there exist positive constants $d_{k}$ such that for all $k=2,3, \ldots$ it holds

$$
\begin{equation*}
\left|d_{t}^{k} b(t)\right| \leq d_{k}\left(\frac{\lambda(t)}{\Lambda(t)}(\ln \Lambda(t))^{\beta}\right)^{k} \quad \text { for large } \quad t \tag{17}
\end{equation*}
$$

Theorem 2 (Main result). Assume that the conditions (13) to (17) are satisfied with $\beta \in[0,1)$. Then for every $\varepsilon>0$ there exists a constant $C_{\varepsilon}$ such that the decay estimate

$$
\left\|u_{t}(t, \cdot)\right\|_{L_{q}}+\left\|\lambda(t) \nabla_{x} u(t, \cdot)\right\|_{L_{q}} \leq C_{\varepsilon}(1+\Lambda(t))^{1+\varepsilon-\frac{n-1}{2}\left(\frac{1}{p}-\frac{1}{q}\right)}\left(\|\varphi\|_{W_{p}^{L+1}}+\|\psi\|_{W_{p}^{L}}\right)
$$

holds for the solution $u=u(t, x)$ to (12). Here $L=\left[n\left(\frac{1}{p}-\frac{1}{q}\right)\right]+1,1<p<2, \frac{1}{p}+\frac{1}{q}=1$.
Let us explain the philosophy of our approach. By $F, F^{-1}$ we denote the Fourier transform, inverse Fourier transform with respect to $x$, respectively. Applying $F$ to (12) we get with $v=$ $F(u)$ the Cauchy problem

$$
\begin{equation*}
v_{t t}+\lambda(t)^{2} b(t)^{2}|\xi|^{2} v=0, \quad v(0, \xi)=F(\varphi), \quad v_{t}(0, \xi)=F(\psi) \tag{18}
\end{equation*}
$$

Setting $V=\left(V_{1}, V_{2}\right)^{T}:=\left(\lambda(t)|\xi| v, D_{t} v\right), D_{t}:=d / d t$, the differential equation can be transformed to the system of first order

$$
D_{t} V-\left(\begin{array}{cc}
0 & \lambda(t)|\xi|  \tag{19}\\
\lambda(t) b(t)^{2}|\xi| & 0
\end{array}\right) V-\frac{D_{t} \lambda}{\lambda}\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right) V=0 .
$$

Our main object is the fundamental solution $X=X(t, \tau, \xi) \in C^{\infty}\left([\tau, \infty) \times \mathbb{R}^{n}\right)$ of (19), that is the solution of

$$
D_{t} X-\left(\begin{array}{cc}
0 & \lambda(t)|\xi|  \tag{20}\\
\lambda(t) b(t)^{2}|\xi| & 0
\end{array}\right) X-\frac{D_{t} \lambda}{\lambda}\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right) X=0,
$$

$$
\begin{equation*}
X(\tau, \tau, \xi)=I \tag{21}
\end{equation*}
$$

with $\tau \geq 0$. We prove that $X=X(t, 0, \xi)$ can be represented in the form

$$
\begin{align*}
X(t, 0, \xi)= & X^{+}(t, 0, \xi) \exp \left(i \int_{0}^{t} \lambda(s) b(s) d s|\xi|\right) \\
& +X^{-}(t, 0, \xi) \exp \left(-i \int_{0}^{t} \lambda(s) b(s) d s|\xi|\right) \tag{22}
\end{align*}
$$

where $X^{+}$and $X^{-}$have connections to symbol classes. Using this representation we obtain the solution of (12) in the form

$$
\begin{align*}
u(t, x) & =F^{-1}\left(\frac{\lambda(0)}{\lambda(t)} X_{11}(t, 0, \xi) F(\varphi)+\frac{1}{\lambda(t)|\xi|} X_{12}(t, 0, \xi) F(\psi)\right),  \tag{23}\\
D_{t} u(t, x) & =F^{-1}\left(\lambda(0)|\xi| X_{21}(t, 0, \xi) F(\varphi)+X_{22}(t, 0, \xi) F(\psi)\right), \tag{24}
\end{align*}
$$

where $X_{j k}$ are the elements of $X$. For these Fourier multipliers $L_{p}-L_{q}$ decay estimates are derived in Section 3.5.
We intend to obtain the representation (22) in $\left\{(t, \xi) \in[0, \infty) \times \mathbb{R}^{n} \backslash\{0\}\right\}$ by splitting this set into two zones.

Definition 2. We define the pseudodifferential zone $Z_{p d}(\beta, N)$ by

$$
Z_{p d}(\beta, N):=\left\{(t, \xi) \in[0, \infty) \times\left(\mathbb{R}^{n} \backslash\{0\}\right):\left(e^{4}+\Lambda(t)\right)|\xi| \leq N\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}\right\}
$$

the hyperbolic zone $Z_{\text {hyp }}(\beta, N)$ by

$$
Z_{\mathrm{hyp}}(\beta, N):=\left\{(t, \xi) \in[0, \infty) \times\left(\mathbb{R}^{n} \backslash\{0\}\right):\left(e^{4}+\Lambda(t)\right)|\xi| \geq N\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}\right\}
$$

The positive constant $N$ will be chosen later.
For $|\xi| \in\left(0, p_{0}\right], p_{0}=4^{\beta} N / e^{4}$, we define the function $t_{\xi}=t(|\xi|)$ as the solution of $\left(e^{4}+\Lambda\left(t_{\xi}\right)\right)|\xi|=N\left(\ln \left(e^{4}+\Lambda\left(t_{\xi}\right)\right)\right)^{\beta}$.

Lemma 3. The derivatives $\partial_{|\xi|}^{k} t_{\xi}$ can be estimated in the following way:

$$
\left|\partial_{|\xi|}^{k} t_{|\xi|}\right| \leq C_{k}|\xi|^{-k} \frac{\left(e^{4}+\Lambda\left(t_{\xi}\right)\right)}{\lambda\left(t_{\xi}\right)} \quad \text { for all } \quad \xi \in \mathbb{R}^{n},|\xi| \in\left(0, p_{0}\right]
$$

### 3.3. The fundamental solution in $Z_{p d}(\beta, N)$

Denoting

$$
A(t,|\xi|):=\left(\begin{array}{cc}
0 & \lambda(t)|\xi| \\
\lambda(t) b(t)^{2}|\xi| & 0
\end{array}\right)+\frac{D_{t} \lambda}{\lambda}\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right)
$$

the fundamental solution $X(t, 0, \xi)$ can be written explicitly in the form

$$
\begin{equation*}
X(t, 0,|\xi|)=I+\sum_{k=1}^{\infty} \int_{0}^{t} A\left(t_{1},|\xi|\right) \cdots \int_{0}^{t_{k-1}} A\left(t_{k},|\xi|\right) d t_{k} \cdots d t_{1} \tag{25}
\end{equation*}
$$

for $|\xi| \in\left(0, p_{0}\right]$. For a given positive number $T$ let us distinguish two cases.
a) $t_{\xi} \leq T$ : in this case we have

$$
\int_{0}^{t}\|A(s,|\xi|)\| d s \leq C(T) \quad \text { for all } \quad t \leq t_{\xi}
$$

b) $T \leq t_{\xi}$ : in this case we have

$$
\begin{aligned}
\int_{0}^{t}\|A(s,|\xi|)\| d s & \leq C(T)+C_{b} \int_{T}^{t} \lambda(s)|\xi| d s+\int_{T}^{t} \frac{\lambda^{\prime}(s)}{\lambda(s)} d s \\
& \leq C(T)+C_{b} \Lambda(t)|\xi|+\ln \frac{\lambda(t)}{\lambda(T)} \\
& \leq C(T)+C_{b} N\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}+\ln \frac{\lambda(t)}{\lambda(T)} \\
& \leq C(T)+C_{b} N\left(\ln \left(e^{4}+\Lambda(T)\right)\right)^{\beta-1} \ln \left(e^{4}+\Lambda(t)\right)+\ln \frac{\lambda(t)}{\Lambda(T)}
\end{aligned}
$$

for all $t \leq t_{\xi}$. Consequently,

$$
\exp \left(\int_{0}^{t}\|A(s,|\xi|)\| d s\right) \leq C(T) \lambda(t)\left(e^{4}+\Lambda(t)\right)^{\frac{C_{b} N}{\left(\ln \left(e^{4}+\Lambda(T)\right)\right)^{1-\beta}}}
$$

This gives the next statement:

Lemma 4. To each small positive $\varepsilon$ there exists a constant $C_{\varepsilon}(N)$ such that for all $(t, \xi) \in$ $Z_{p d}(\beta, N)$ it holds

$$
\|X(t, 0, \xi)\| \leq C_{\varepsilon}(N) \lambda(t)\left(e^{4}+\Lambda(t)\right)^{\varepsilon}, \quad\|X(t, 0, \xi)\| \leq C_{\varepsilon}(N) \lambda(t)|\xi|^{-\varepsilon}
$$

respectively.
To continue the solution from $Z_{p d}(\beta, N)$ to $Z_{\text {hyp }}(\beta, N)$ for $|\xi| \in\left(0, p_{0}\right]$ and to study its properties in $Z_{\text {hyp }}(\beta, N)$ we need the behaviour of $\partial_{t}^{k} \partial_{\xi}^{\alpha} X(t, 0, \xi)$, too. It is obtained among other things from (25) and (14).

Theorem 3. To each small positive $\varepsilon$ and each $k$ and $\alpha$ there exists a constant $C_{\varepsilon, k, \alpha}(N)$ such that

$$
\left\|\partial_{t}^{k} \partial_{\xi}^{\alpha} X(t, 0, \xi)\right\| \leq C_{\varepsilon, k, \alpha}(N) \lambda(t)\left(\lambda(t)|\xi|+\frac{\lambda(t)}{e^{4}+\Lambda(t)}\right)^{k}|\xi|^{-|\alpha|-\varepsilon}
$$

for all $(t, \xi) \in Z_{p d}(\beta, N)$.

### 3.4. The fundamental solution in $Z_{\text {hyp }}(\beta, N)$

The hyperbolic zone $Z_{\text {hyp }}(\beta, N)$ can be represented as the union of the two sets $\{(t, \xi):|\xi| \in$ $\left.\left(0, p_{0}\right]:\left(e^{4}+\Lambda(t)\right)|\xi| \geq N(\ln (e+\Lambda(t)))^{\beta}\right\}$ and $\left\{(t, \xi) \in[0, \infty) \times\left\{|\xi| \geq p_{0}\right\}\right\}$. We restrict ourselves to the first set and sketch at the end of this section the approach in the second set.
In $Z_{\text {hyp }}(\beta, N)$ we apply a diagonalization procedure to the first order system (19). To carry out this procedure we need the following classes of symbols.

Definition 3. For given real numbers $m_{1}, m_{2}, m_{3}, \beta \in[0,1)$ and for positive $N$ we denote by $S_{\beta, N}\left\{m_{1}, m_{2}, m_{3}\right\}$ the set of all symbols $a=a(t, \xi) \in C^{\infty}\left(Z_{\mathrm{hyp}}(\beta, N):|\xi| \in\right.$ $\left.\left(0, p_{0}\right]\right)$ satisfying there

$$
\left|\partial_{t}^{k} \partial_{\xi}^{\alpha} a(t, \xi)\right| \leq C_{k, \alpha}|\xi|^{m_{1}-|\alpha|} \lambda(t)^{m_{2}}\left(\frac{\lambda(t)}{e^{4}+\Lambda(t)}\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}\right)^{m_{3}+k}
$$

These classes of symbols are related to the Definitions 1 and 2. To understand that the diagonalization procedure improves properties of the remainder (as usually) one takes into consideration the following rules of the symbolic calculus:

- $\quad S_{\beta, N}\left\{m_{1}, m_{2}, m_{3}\right\} \subset S_{\beta, N}\left\{m_{1}+k, m_{2}+k, m_{3}-k\right\}, \quad k \geq 0 ;$
- $\quad a \in S_{\beta, N}\left\{m_{1}, m_{2}, m_{3}\right\}, \quad b \in S_{\beta, N}\left\{n_{1}, n_{2}, n_{3}\right\}$, then $a b \in S_{\beta, N}\left\{m_{1}+n_{1}, m_{2}+n_{2}, m_{3}+n_{3}\right\} ;$
- $\quad a \in S_{\beta, N}\left\{m_{1}, m_{2}, m_{3}\right\}$, then $\partial_{t} a \in S_{\beta, N}\left\{m_{1}, m_{2}, m_{3}+1\right\}$;
- $\quad a \in S_{\beta, N}\left\{m_{1}, m_{2}, m_{3}\right\}$, then $\partial_{\xi}^{\alpha} a \in S_{\beta, N}\left\{m_{1}-|\alpha|, m_{2}, m_{3}\right\}$.

Let us define the matrices

$$
M^{-1}(t):=\frac{1}{\sqrt{\lambda(t) b(t)}}\left(\begin{array}{cc}
1 & 1 \\
-b(t) & b(t)
\end{array}\right), \quad M(t):=\frac{1}{2} \sqrt{\frac{\lambda(t)}{b(t)}}\left(\begin{array}{cc}
b(t) & -1 \\
b(t) & 1
\end{array}\right) .
$$

Substituting $X=M^{-1} Y$ some calculations transform (20) into the first order system

$$
\begin{equation*}
D_{t} Y-D Y+B Y=0 \tag{26}
\end{equation*}
$$

where

$$
\begin{array}{ll}
D(t, \xi):=\left(\begin{array}{cc}
\tau_{1}(t, \xi) & 0 \\
0 & \tau_{2}(t, \xi)
\end{array}\right), & B(t, \xi):=-\frac{D_{t}(\lambda(t) b(t))}{2 \lambda(t) b(t)}\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right) \\
\tau_{1}(t, \xi):=-\lambda(t) b(t)|\xi|+\frac{D_{t} \lambda(t)}{\lambda(t)}, & \tau_{2}(t, \xi):=\lambda(t) b(t)|\xi|+\frac{D_{t} \lambda(t)}{\lambda(t)}
\end{array}
$$

Without difficulties one can prove $D \in S_{\beta, N}\{1,1,0\}, B \in S_{\beta, N}\{0,0,1\}$. To prove $L_{p}-L_{q}$ decay estimates for the solution of (12) we need further steps of the diagonalization of (26). This is carried out in the next lemma.

LEMMA 5. For a given nonnegative integer $M$ there exist matrix-valued functions $N_{M}=$ $N_{M}(t, \xi) \in S_{\beta, N}\{0,0,0\}, F_{M}=F_{M}(t, \xi) \in S_{\beta, N}\{-1,-1,2\}$ and $R_{M}=R_{M}(t, \xi) \in$ $S_{\beta, N}\{-M,-M, M+1\}$ such that the following operator-valued identity holds:

$$
\left(D_{t}-D+B\right) N_{M}=N_{M}\left(D_{t}-D+F_{M}-R_{M}\right)
$$

where $F_{M}$ is diagonal while $N_{M}$ is invertible and its inverse $N_{M}^{-1}$ belongs as $N_{M}$ to $S_{\beta, N}$ $\{0,0,0\}$.

REMARK 5. The invertibility of the diagonalizer $N_{M}=N_{M}(t, \xi)\left(\bmod S_{\beta, N}\{-M,-M\right.$, $M+1\}$ ) is essential. This property follows by a special choice of the positive constant $N$ in Definition 2. We need only a finite number of steps of diagonalization (cf. proof of Theorem 2), thus $N$ can be fixed after carrying out these steps.

Now let us devote to the system

$$
\begin{equation*}
\left(D_{t}-D+F_{M}-R_{M}\right) Z=0, \quad Z=Z(t, r, \xi) \tag{27}
\end{equation*}
$$

where $t_{\xi} \leq r \leq t$. Let $E_{2}=E_{2}(t, r, \xi) ; t, r \geq t_{\xi}$, is defined by

$$
E_{2}(t, r, \xi):=\frac{\lambda(t)}{\lambda(r)}\left(\begin{array}{cc}
\exp \left(-i \int_{r}^{t} \lambda(s) b(s) d s|\xi|-i \int_{r}^{t} F_{M}^{(1,1)}(s, \xi) d s\right) & 0 \\
0 & \exp \left(i \int_{r}^{t} \lambda(s) b(s) d s|\xi|-i \int_{r}^{t} F_{M}^{(2,2)}(s, \xi) d s\right)
\end{array}\right)
$$

be the solution of the Cauchy problem $\left(D_{t}-D+F_{M}\right) Z=0, Z(r, r, \xi)=I$. Let us denote

$$
P_{M}(t, r, \xi):=E_{2}(r, t, \xi) R_{M}(t, \xi) E_{2}(t, r, \xi)
$$

By the aid of $P_{M}$ we define the matrix-valued function

$$
\begin{align*}
Q_{M}(t, r, \xi):= & \sum_{k=1}^{\infty} i^{k} \int_{r}^{t} P_{M}\left(t_{1}, r, \xi\right) \int_{r}^{t_{1}} P_{M}\left(t_{2}, r, \xi\right) \cdots  \tag{28}\\
& \cdots \int_{r}^{t_{k-1}} P_{M}\left(t_{k}, r, \xi\right) d t_{k} \ldots d t_{1}
\end{align*}
$$

The function $Q_{M}=Q_{M}(t, r, \xi)$ solves the Cauchy problem

$$
D_{t} Q-P_{M} Q-P_{M}=0, \quad Q(r, r, \xi)=0 \quad \text { for } \quad t, r \geq t \xi
$$

Using these auxiliary functions it is easy to prove the next result.

Lemma 6. The matrix-valued function $Z(t, r, \xi)=E_{2}(t, r, \xi)\left(I+Q_{M}(t, r, \xi)\right)$ solves the Cauchy problem (27) for $t, r \geq t_{\xi}$.

Now we can go back to (20), (21) and obtain as its solution

$$
\begin{align*}
X(t, 0, \xi)= & M^{-1}(t) N_{M}(t, \xi) E_{2}\left(t, t_{\xi}, \xi\right)\left(I+Q_{M}\left(t, t_{\xi}, \xi\right)\right) \cdot \\
& \cdot N_{M}^{-1}\left(t_{\xi}, \xi\right) M\left(t_{\xi}\right) X\left(t_{\xi}, 0, \xi\right) . \tag{29}
\end{align*}
$$

We write $\exp \left(-i \int_{t_{\xi}}^{t} \lambda(s) b(s) d s\right)=\exp \left(-i \int_{0}^{t} \lambda(s) b(s) d s-i \int_{t_{\xi}}^{0} \lambda(s) b(s) d s\right)$ in correspondence with our goal (22) and include the second factor in the amplitudes. The matrices $M$ and $M^{-1}$ are given in an explicit form. The properties of $N_{M}$ and $N_{M}^{-1}$ are described by Lemma 5 using Definition 3. To estimate $X\left(t_{\xi}, 0, \xi\right)$ we use Theorem 3. Consequently, it remains to estimate $E_{2}\left(0, t_{\xi}, \xi\right)$ and $Q_{M}\left(t, t_{\xi}, \xi\right)$.

Lemma 7. For every positive small $\varepsilon$ and every $\alpha$ the following estimate in $Z_{p d}(\beta, N)$ holds:

$$
\left|\partial_{\xi}^{\alpha} \exp \left(i \int_{t_{\xi}}^{t} \lambda(s) b(s) d s\right)\right| \leq C_{\varepsilon, \alpha|\xi|^{-|\alpha|-\varepsilon},}
$$

where $C_{\varepsilon, \alpha}=C_{\varepsilon, \alpha}(\beta, N)$.
Lemma 8. For every positive small $\varepsilon$ and every $\alpha$ the following estimate in $Z_{\text {hyp }}(\beta, N)$ holds, $|\xi| \in\left(0, p_{0}\right]$ :

$$
\left|\partial_{\xi}^{\alpha} \exp \left(-i \int_{t_{\xi}}^{t} F_{M}^{(k, k)}(s, \xi) d s\right)\right| \leq C_{\varepsilon, \alpha}|\xi|^{-|\alpha|-\varepsilon}, \quad k=1,2
$$

where $C_{\varepsilon, \alpha}=C_{\varepsilon, \alpha}(\beta, N)$.
Proof. The statement for $|\alpha|=0$ follows from

$$
\begin{aligned}
\left|\int_{t_{\xi}}^{t} F_{M}^{(k, k)}(s, \xi) d s\right| \leq & \int_{t_{\xi}}^{t}\left|F_{M}^{(k, k)}(s, \xi)\right| d s \\
\leq & C_{k} \int_{t_{\xi}}^{t} \frac{\lambda(s)\left(\ln \left(e^{4}+\Lambda(s)\right)\right)^{2 \beta}}{|\xi|\left(e^{4}+\Lambda(s)\right)^{2}} d s \\
\int_{t_{\xi}}^{t} \frac{\lambda(s)\left(\ln \left(e^{4}+\Lambda(s)\right)\right)^{2 \beta}}{\left(e^{4}+\Lambda(s)\right)^{2}} d s \leq & \frac{\left(\ln \left(e^{4}+\Lambda\left(t_{\xi}\right)\right)\right)^{2 \beta}}{e^{4}+\Lambda\left(t_{\xi}\right)} \\
& +\int_{t_{\xi}}^{t} \frac{2 \beta}{\ln \left(e^{4}+\Lambda(s)\right)} \frac{\lambda(s)\left(\ln \left(e^{4}+\Lambda(s)\right)\right)^{2 \beta}}{\left(e^{4}+\Lambda(s)\right)^{2}} d s
\end{aligned}
$$

Definition 2 and

$$
\frac{\left(\ln \left(e^{4}+\Lambda\left(t_{\xi}\right)\right)\right)^{\beta}}{N}=\left(\ln \left(e^{4}+\Lambda\left(t_{\xi}\right)\right)\right)^{\frac{1}{N\left(\ln \left(e^{4}+\Lambda\left(t_{\xi}\right)\right)\right)^{1-\beta}}} .
$$

By induction we prove the statement for $|\alpha|>0$ by using $F_{M} \in S_{\beta, N}\{-1,-1,2\}$ and Lemma 3.

More problems appear if we derive an estimate for $Q_{M}$. Here we refer the reader to [8].
LEMMA 9. The matrix-valued function $P_{M}=P_{M}\left(t, t_{\xi}, \xi\right)$ satisfies for every $l$ and $\alpha$ in $Z_{\text {hyp }}(\beta, N) \cap\left\{|\xi| \in\left(0, p_{0}\right]\right\}$ the estimates

$$
\begin{aligned}
\left\|\partial_{t}^{l} \partial_{\xi}^{\alpha} P_{M}\left(t, t_{\xi}, \xi\right)\right\| \leq & C_{M, l, \alpha}(\lambda(t)|\xi|)^{l}\left(e^{4}+\Lambda(t)\right)^{|\alpha|} \frac{\lambda(t)}{e^{4}+\Lambda(t)} \\
& \cdot\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}\left(\frac{\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}}{\left(e^{4}+\Lambda(t)\right)|\xi|}\right)^{M}
\end{aligned}
$$

Lemma 10. For every positive small $\varepsilon$ and every $\alpha,|\alpha| \leq M-1$, it holds the following estimate in $Z_{\text {hyp }}(\beta, N),|\xi| \in\left(0, p_{0}\right]$ :

$$
\left\|\partial_{\xi}^{\alpha} Q_{M}(t, t \xi, \xi)\right\| \leq C_{\varepsilon, \alpha}(1+\Lambda(t))^{\varepsilon}|\xi|^{-|\alpha|-\varepsilon}
$$

where $C_{\varepsilon, \alpha}=C_{\varepsilon, \alpha}(\beta, N)$.
Proof. We use the representation (28) with $r=t_{\xi}$ and form the derivatives $\partial_{\xi}^{\alpha} Q_{M}\left(t, t_{\xi}, \xi\right)$. For $|\alpha|=0$ the statement from Lemma 9 and similar calculations as in the proof of Lemma 8 imply the estimate for $\left\|Q_{M}\left(t, t_{\xi}, \xi\right)\right\|$. If we differentiate for $|\alpha|=1$ inside of the integrals, then the estimate follows immediately. If we differentiate the lower integral bound in $\int_{t_{\xi}}^{t_{k-1}} P_{M}\left(t_{k}, t_{\xi}, \xi\right) d t_{k}$, then there appears a term of the form

$$
P_{M}\left(t_{\xi}, t_{\xi}, \xi\right) \frac{\partial t_{\xi}}{\partial \xi_{l}} \int_{t_{\xi}}^{t} P_{M}\left(t_{1}, t_{\xi}, \xi\right) \cdots \int_{t_{\xi}}^{t_{k-2}} P_{M}\left(t_{k-1}, t_{\xi}, \xi\right) d t_{k-1} \ldots d t_{1}
$$

Using Lemma 9 and Lemma 3 gives the desired estimate in this case, too. But we can only get estimates for $|\alpha| \leq M-1$. In this case we can have an integrand of the form

$$
\frac{\lambda(t)}{\left(e^{4}+\Lambda(t)\right)^{2}|\xi|}\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta(M+1)}
$$

The term $\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta M}$ can be estimated by $C(1+\Lambda(t))^{\varepsilon}$, the other factor is integrable and can be estimated by $C$. An induction procedure yields the statement for $|\alpha| \leq M-1$ (see [8]).

Now we have all tools to get an estimate for (29).
THEOREM 4. The fundamental solution $X=X(t, 0, \xi)$ can be represented in $Z_{\mathrm{hyp}}(\beta, N)$ $\cap\left\{|\xi| \in\left(0, p_{0}\right]\right\}$ as follows:
$X(t, 0, \xi)=X^{+}(t, 0, \xi) \exp \left(i \int_{0}^{t} \lambda(s) b(s) d s|\xi|\right)+X^{-}(t, 0, \xi) \exp \left(-i \int_{0}^{t} \lambda(s) b(s) d s|\xi|\right)$,
where the matrix-valued amplitudes $X^{-}, X^{+}$satisfy for all $|\alpha| \leq M-1$ and all positive small $\varepsilon$ the estimates

$$
\begin{equation*}
\left\|\partial_{\xi}^{\alpha} X^{ \pm}(t, 0, \xi)\right\| \leq C_{M, \varepsilon} \sqrt{\lambda(t) \lambda\left(t_{\xi}\right)}|\xi|^{-|\alpha|-\varepsilon} \tag{30}
\end{equation*}
$$

REMARK 6. There are no new difficulties to derive a corresponding estimate to (30) in $\left\{(t, \xi) \in[0, \infty) \times\left\{|\xi| \geq p_{0}\right\}\right\}$ which belongs to $Z_{\text {hyp }}(\beta, N)$ completely. We obtain for all $|\alpha| \leq M-1$ the estimates

$$
\begin{equation*}
\left\|\partial_{\xi}^{\alpha} X^{ \pm}(t, 0, \xi)\right\| \leq C_{M} \sqrt{\lambda(t)}|\xi|^{-|\alpha|} \tag{31}
\end{equation*}
$$

### 3.5. Fourier multipliers

The statements of Theorems 3, 4 and Remark 6 imply together with (23), (24) representations of the solution of (12) and its derivatives by the aid of Fourier multipliers. To estimate $\lambda=\lambda\left(t_{\xi}\right)$ in (30) we use assumption (14), especially $\frac{\lambda(t)}{\Lambda(t)} \leq c(\ln \Lambda(t))^{c}$ for large $t$. This corresponds to the Examples 3 and 4. To study Examples 1 and 2 we can follow our strategies in the same way. To get $L_{p}-L_{q}$ decay estimates we divide our consideration into two steps in accordance with two completely different ideas:
Hardy-Littlewood inequality [2] and Littman's lemma [3].
Instead of $Z_{p d}(\beta, N)$ from Definition 2 we use now

$$
Z_{p d}(\beta, N)=\left\{(t, \xi) \in[0, \infty) \times\left(\mathbb{R}^{n} \backslash\{0\}\right):\left(e^{4}+\Lambda(t)\right)|\xi| \leq 4 N\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}\right\}
$$

It is clear that the statement of Theorem 3 remains unchanged.
Let us choose a function $\psi \in C^{\infty}(\mathbb{R})$ satisfying $\psi(s) \equiv 0$ for $|s| \leq 1 / 2, \psi(s) \equiv 1$ for $|s| \geq 1$ and $0 \leq \psi(\xi) \leq 1$. Moreover, we define

$$
K(t):=\frac{2 N\left(\ln \left(e^{4}+\Lambda(t)\right)\right)^{\beta}}{e^{4}+\Lambda(t)}
$$

Following the approach of [8] which generalizes that one of [4] one can prove the next two results.

## a) Fourier multipliers with amplitudes supported in the pseudodifferential zone

THEOREM 5 (APPLICATION OF HARDY - LITTLEWOOD INEQUALITY). Let us consider Fourier multipliers depending on the parameter $t \in[0, \infty)$ which are defined by
$I_{1}:=F^{-1}\left(e^{i \int_{0}^{t} \lambda(s) b(s) d s|\xi|}\left(1-\psi\left(\frac{\xi}{K(t)}\right)\right)|\xi|^{-2 r} a(t, \xi) F\left(u_{0}\right)(\xi)\right), \quad u_{0} \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$.
Suppose that $|a(t, \xi)| \leq C_{\varepsilon} \lambda(t)|\xi|^{-\varepsilon}$ in $Z_{p d}(\beta, N)$. Then we have

$$
\begin{gathered}
\qquad\left\|I_{1}\right\|_{L_{q}} \leq C_{r, N, \varepsilon} \lambda(t)\left(e^{4}+\Lambda(t)\right)^{2 r+\varepsilon-n\left(\frac{1}{p}-\frac{1}{q}\right)}\left\|u_{0}\right\|_{L_{p}} \\
\text { provided that } 1<p<2,1 / p+1 / q=1,2 r+\varepsilon \leq n\left(\frac{1}{p}-\frac{1}{q}\right) .
\end{gathered}
$$

b) Fourier multipliers with amplitudes supported in the hyperbolic zone

Theorem 6 (Application of Littman's Lemma). Let us consider

$$
I_{2}:=F^{-1}\left(e^{i \int_{0}^{t} \lambda(s) b(s) d s|\xi|} \psi\left(\frac{\xi}{K(t)}\right)|\xi|^{-2 r} a(t, \xi) F\left(u_{0}\right)(\xi)\right), \quad u_{0} \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)
$$

Suppose that $\left|\partial_{\xi}^{\alpha} a(t, \xi)\right| \leq C \sqrt{\lambda(t)}|\xi|^{-\frac{1}{2}-|\alpha|-\varepsilon}$ for $|\alpha| \leq n$ in $Z_{\mathrm{hyp}}(\beta, N)$ from Definition 2. Then we have

$$
\left\|I_{2}\right\|_{L_{q}} \leq C_{r, N, n, \varepsilon} \sqrt{\lambda(t)}\left(e^{4}+\Lambda(t)\right)^{2 r+\frac{1}{2}+\varepsilon-n\left(\frac{1}{p}-\frac{1}{q}\right)}\left\|u_{0}\right\|_{L_{p}}
$$

provided that $1<p<2,1 / p+1 / q=1, \frac{n+1}{2}\left(\frac{1}{p}-\frac{1}{q}\right) \leq 2 r$.
Proof of Theorem 2. Theorem 6 tells us that we have to carry out $M=n+1$ steps of the perfect diagonalization. Here we use $2 r=\frac{n+1}{2}\left(\frac{1}{p}-\frac{1}{q}\right)$. In Theorem 5 we use $2 r=\left[n\left(\frac{1}{p}-\frac{1}{q}\right)\right]+1$. But from (24) we know that this is the supposed regularity for $\nabla \varphi, \psi$, respectively. Consequently, the regularity for $\varphi$ and $\psi$ from Theorem 2 can be understood. The decay function $1+\Lambda(t)$ with the rate $1+\varepsilon-\frac{n-1}{2}\left(\frac{1}{p}-\frac{1}{q}\right)$ follows immediately from the statements of Theorems 5 and 6. Due to (23) the estimate for $\lambda(t) \nabla u$ coincides with that for $D_{t} u$.

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