# Wave Equations with Time-dependent Spatial Operators of Higher Order 

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We study the initial-boundary value problem for $\partial_{1}^{2} u(t, x)+A(t) u(t, x)+B(t) \partial_{t} u(t, x)=f(t, x)$ on $[0, T]$ $\times \Omega\left(\Omega \subset \mathbb{R}^{n}\right)$ with a homogeneous Dirichlet boundary condition; here $A(t)$ denotes a family of uniformly strongly elliptic operators of order $2 m, B(t)$ denotes a family of spatial differential operators of order less than or equal to $m$, and $u$ is a scalar function. We prove the existence of a unique strong solution $u$. Furthermore, an energy estimate for $u$ is given.

## 1. Introduction

Let $\Omega \subset \mathbb{R}^{n}$ and $m \in \mathbb{N}$ be given. We consider the problem

$$
\left.\begin{array}{ll}
\partial_{t}^{2} u(t, x)+A(t) u(t, x)+B(t) \partial_{t} u(t, x)=f(t, x) & \text { for } t \in[0, T], x \in \Omega,  \tag{1.1}\\
u(t, .) \in \dot{H}^{m}(\Omega) & \text { for } t \in[0, T], \\
u(0, x)=u^{0}(x), \quad \partial_{\imath} u(0, x)=u^{1}(x) & \text { for } x \in \Omega .
\end{array}\right\}
$$

Here $\dot{H}^{m}(\Omega)$ denotes the closure of $C_{0}^{\infty}(\Omega)$ in the $m$ th Sobolev space $H^{m}(\Omega)$, and $A(t), B(t)$ denote families of spatial differential operators of order $2 m$ and less than or equal to $m$, respectively. Problems of this kind appear in the study of fully non-linear wave equations (compare [9]). We make the following assumptions.

Assumption 1.1. (1) The operators $A$ and $B$ are given by

$$
\begin{array}{ll}
A(t) \varphi:=\sum_{|\alpha| \leqslant 2 m}\left[a_{\alpha}(t, .)+\tilde{a}_{\alpha}(t, .)\right] \partial_{x}^{\alpha} \varphi & \text { for } \varphi \in \dot{H}^{m}(\Omega) \cap H^{2 m}(\Omega), \\
B(t) \varphi:=\sum_{|\beta| \leqslant m}\left[b_{\beta}(t, .)+\tilde{b}_{\beta}(t, .)\right] \partial_{x}^{\beta} \varphi \quad \text { for } \varphi \in H^{m}(\Omega), \tag{1.3}
\end{array}
$$

where $a_{a}, b_{\beta} \in C_{b}^{(k-1) m}([0, T] \times \bar{\Omega})$ and

$$
\begin{equation*}
\tilde{a}_{\alpha}, \tilde{b}_{\beta} \in \bigcap_{j=1}^{k-1} C^{j}\left([0, T], H^{(k-1-j) m}(\Omega)\right) \tag{1.4}
\end{equation*}
$$

for $|\alpha| \leqslant 2 m,|\beta| \leqslant m$ and some $k \geqslant[n / 2 m]+3([r]:=\max \{j \in \mathbb{N}: j \leqslant r\})$.
(2) There exist constants $c_{1}, c_{2}>0$ such that

$$
\begin{align*}
& (-1)^{m} \operatorname{Re} \sum_{|\alpha|=2 m}\left[a_{\alpha}(t, x)+\tilde{a}_{\alpha}(t, x)\right] \xi^{\alpha} \geqslant c_{1}|\xi|^{2 m} \\
& \quad \text { for } \xi \in \mathbb{R}^{n}, t \in[0, T], \quad x \in \bar{\Omega},  \tag{1.5}\\
& \left\|\left[A(t)-A^{*}(t)\right] \varphi\right\| \leqslant c_{2}\|\varphi\|_{m} \quad \text { for } \varphi \in \dot{H}^{m}(\Omega) \cap H^{2 m}(\Omega), t \in[0, T], \tag{1.6}
\end{align*}
$$

where $A^{*}(t)$ denotes the formal adjoint to the operator $A(t)$, and $\|\cdot\|,\|\cdot\|_{m}$ denote the norms in $L_{2}(\Omega)$ and $H^{m}(\Omega)$, respectively.
(3) There exists a constant $c_{3} \geqslant 0$ such that

$$
\begin{equation*}
-\operatorname{Re}\langle B(t) \varphi, \varphi\rangle \leqslant c_{3}\|\varphi\|^{2} \quad \text { for } \varphi \in \dot{H}^{m}(\Omega), t \in[0, T] \tag{1.7}
\end{equation*}
$$

(here $\langle.,$.$\rangle denotes the inner product in L_{2}(\Omega)$ ).
Remarks. (1) We admit that the coefficients of $A$ and $B$ are divided into two different parts, one being continuously differentiable with bounded derivatives, and the other lying in some spatial Sobolev space for every $t \in[0, T]$. This is essential for the application to non-linear problems.
(2) Condition (1.6) means that the part of $A(t)$ containing the derivatives $\partial_{x}^{\alpha} \varphi$ with $m+1 \leqslant|\alpha| \leqslant 2 m$ is symmetric. This condition is also used in [2]. An equivalent formulation of (1.6) is used in [3], (3.5), (3.6).
(3) Condition (1.7) is needed for the energy estimate. In the case $m=1$, if $b_{\beta}+\tilde{b}_{\beta}$ is real valued for $|\beta|=1$, (1.7) holds automatically if $k \geqslant[n / 2 m]+4$ in (1.4). This can be shown by integrating by parts (compare [9]). More practical conditions for $B(t)$ guaranteeing (1.7) are given in [9].
(4) By Sobolev's lemma it follows from (1.4) and $k \geqslant[n / 2 m]+3$ that $\tilde{a}_{\alpha}(t) \in C_{b}(\bar{\Omega})$ for $t \in[0, T]$. Hence (1.5) is well defined:

The aim of this paper is to prove the existence of a unique solution

$$
\begin{equation*}
u \in \mathscr{C}_{T}^{k}:=\bigcap_{j=0}^{k} C^{j}\left([0, T], H^{(k-j) m}(\Omega)\right) \tag{1.8}
\end{equation*}
$$

of (1.1). More precisely, we prove the following theorem.
Theorem 1.1. Let Assumption 1.1 be satisfied for some $k \geqslant k_{0}:=[n / 2 m]+4$ and let $2 \leqslant j \leqslant k$. If

$$
\begin{equation*}
f \in \mathscr{C}_{T}^{j-2} \cap C^{j-1}\left([0, T], L_{2}(\Omega)\right) \tag{1.9}
\end{equation*}
$$

and $u^{0} \in H^{j m}(\Omega), u^{1} \in H^{(j-1) m}(\Omega)$ such that $\left(u^{0}, u^{1}, f\right)$ satisfies the compatibility condition (defined in section 2) of order $j$, then (1.1) has a unique solution $u \in \mathscr{C} \mathscr{C}_{\mathbf{T}}^{j}$. Furthermore

$$
\begin{align*}
|u(t)|_{j}:= & \sum_{v=0}^{j}\left\|\partial_{t}^{v} u(t)\right\|_{(j-v) m} \\
\leqslant & \exp \left(C_{1} t\right)\left(D_{1}|u(0)|_{j}+C_{2} \int_{0}^{t}\left[\left\|\partial_{i}^{j-1} f(\tau)\right\|+|f(\tau)|_{j-2}\right] \mathrm{d} \tau\right) \\
& +D_{2}|f(t)|_{j-2} \quad \text { for } t \in[0, T] \tag{1.10}
\end{align*}
$$

where the constants $C_{1}, C_{2}>0$ depend only on $c_{1}, c_{2}, c_{3}$ (of Assumption 1.1), $j$, and

$$
\begin{align*}
& \sup _{|\alpha| \leqslant 2 m,|\beta| \leqslant m} \sup _{|\gamma|+j m \leqslant(k-1) m} \sup _{[0, T] \times \bar{\Omega}}\left(\left|\partial_{i}^{j} \partial_{x}^{\gamma} a_{\alpha}(t, x)\right|+\left|\partial_{t}^{j} \partial_{x}^{\gamma} b_{\beta}(t, x)\right|\right),  \tag{1.11}\\
& \sup _{\sup _{|\alpha| \leqslant 2 m,|\beta| \leqslant m}} \sup _{[0, T]}\left(\left|\tilde{a}_{\alpha}(t) \tilde{\left.\right|_{k-1}}+\left|\tilde{b}_{\beta}(t)\right|_{k-1}\right)\right.  \tag{1.12}\\
& \left(\left|\tilde{a}_{\alpha}(t) \tilde{\left.\right|_{k-1}}:=\left|\partial_{t} \tilde{a}_{\alpha}(t)\right|_{k-2}+\left\|\tilde{a}_{\alpha}(t)\right\|_{(k-2) m}\right) \text {, whereas } D_{1}, D_{2}>0 \text { depend only on } j, c_{1}, c_{2},\right.
\end{align*}
$$ $c_{3},(1.11)$, and

$$
\begin{equation*}
\sup _{|\alpha| \leqslant 2 m .|\beta| \leqslant m} \sup _{[0, T]}\left(\left|\tilde{a}_{\alpha}(t)\right|_{k-2}+\left|\tilde{b}_{\beta}(t)\right|_{k-2}\right) . \tag{1.13}
\end{equation*}
$$

Remarks. (1) The fact that $D_{1}, D_{2}$ depend only on (1.13) and not on (1.12) is essential for the iteration procedure in [9], where this theorem is used. If $D_{1}, D_{2}$ are allowed to depend on (1.12), then the condition $k \geqslant[n / 2 m]+4$ in Theorem 1.1 could be relaxed to $k \geqslant[n / 2 m]+3$. But this would require a more complicated proof.
(2) The term $|u(0)|_{j}$ can be estimated by $\left\|u^{0}\right\|_{j m},\left\|u^{1}\right\|_{(j-1) m}$ and $|f(0)|_{j-2}$ (compare (2.8)).
(3) Condition (1.9) can be relaxed to the assumption that $f \in \mathscr{C}_{T}^{j-2}$ and $\partial_{t}^{j-1} f \in L_{2}$ ( $[0, T], L_{2}(\Omega)$ ). Even (1.4) can be slightly relaxed.

There are many papers dealing with problems of the type (1.1) with $B=0$. We only mention Kato [7], Lions and Magenes [10], and Dafermos and Hrusa [3]. In these papers the existence of the solution is proved in two different ways: by the aid of semigroups ([7] and [10]) and by energy methods ([10], §3.8.2, and [3]). In the case $m=1$, if

$$
|\operatorname{Re}\langle B(t) \varphi, \varphi\rangle| \leqslant c\|\varphi\|^{2} \quad \text { for } \varphi \in \dot{H}^{m}(\Omega), t \in[0, T],
$$

Ikawa [5] proved the existence of a solution using semigroups. He considers operators $A(t)$ and $B(t)$ with vanishing $\tilde{a}_{\alpha}$ and $\tilde{b}_{\beta}$. In addition to the Dirichlet boundary condition he studies the Neumann boundary condition. Recently Koch proved in [8] an existence theorem for systems of hyperbolic equations with real-valued coefficients. He assumes that $\Omega$ is bounded and studies the more complicated case of timedependent boundary conditions. Furthermore, he considers a problem similar to (1.1) with $m=1$ (compare (2.12) in [8]). He uses energy methods.

The proof in this paper is also based on energy estimates but differs from the proof in [8]. Section 3 deals with elliptic equations of order $2 m$. In particular, the regularity of a solution of elliptic equations is studied. In section 4 the existence of a unique strong solution $u \in \mathscr{C}_{T}^{2}$ of $(1,1)$ is proved by the method of Faedo-Galerkin, which uses an approximation in finite-dimensional function-spaces. A higher degree of regularity is obtained by induction in section 5 , by solving a system of a simple integral equations and an equation of the type of (1.1) (compare (5.3), (5.4)).

## 2. Notation. The compatibility condition

By $\Omega$ we denote a domain $\Omega \subset \mathbb{R}^{n}$ with $\partial \Omega \in C^{\infty}$ such that $\partial \Omega$ is bounded (or empty) or such that

$$
\begin{equation*}
\boldsymbol{\Omega}=\mathbb{R}^{n^{\prime}} \times \boldsymbol{\Omega}^{\prime} \tag{2.1}
\end{equation*}
$$

with $n^{\prime}<n$ and bounded $\Omega^{\prime} \subset \mathbb{R}^{n-n^{\prime}}$. Let

$$
\partial_{x}^{\alpha} \varphi:=\frac{\partial^{|\alpha|} \varphi}{\partial x_{1}^{\alpha_{1}} \cdots \partial x_{n}^{\alpha_{n}}} \quad \text { for } \alpha \in \mathbb{N}_{0}^{n}
$$

where $|\alpha|=\alpha_{1}+\cdots+\alpha_{n}$. We set
$C_{b}^{k}(\Omega):=\left\{\varphi \in C^{k}(\Omega): \partial_{x}^{\alpha} \varphi\right.$ is bounded in $\Omega$ for $\left.|\alpha| \leqslant k\right\}$.
We use the abbreviation $u^{\prime}=\partial_{t} u, u^{\prime \prime}=\partial_{t}^{2} u$ for time-dependent functions. Furthermore we write $u(t) \in \dot{H}^{m}(\Omega)$ instead of $u(t,.) \in \dot{H}^{m}(\Omega)$. Recalling the definition of $\mathscr{C}_{T}^{k}$ in (1.8) we set

$$
\begin{equation*}
|u(t)|_{k}:=\sum_{j=0}^{k}\left\|\partial_{t}^{j} u(t)\right\|_{(k-j) m} \quad \text { for } t \in[0, T] \tag{2.2}
\end{equation*}
$$

for $u \in \mathscr{C}_{\boldsymbol{T}}^{\mathbf{k}}$. Besides $\mathscr{C}_{\boldsymbol{T}}^{\boldsymbol{k}}$ we use the linear space

$$
\begin{equation*}
\tilde{\mathscr{C}}_{T}^{k}:=\bigcap_{j=1}^{k} C^{j}\left([0, T], H^{(k-j) m}(\Omega)\right) \tag{2.3}
\end{equation*}
$$

Note that $C\left([0, T], H^{k m}(\Omega)\right)$ contains $\mathscr{C}_{T}^{k}$ but not $\tilde{\mathscr{C}}_{T}^{k}$. For $v \in \tilde{\mathscr{C}}_{T}^{k}$ we set

$$
\begin{align*}
|v(t)|_{k} & :=\sum_{j=1}^{k}\left\|\partial_{t}^{j} v(t)\right\|_{(k-j) m}+\|v(t)\|_{(k-1) m} \\
& =\left|v^{\prime}(t)\right|_{k-1}+\|v(t)\|_{(k-1) m} \quad \text { for } t \in[0, T] . \tag{2.4}
\end{align*}
$$

In order to give the compatibility condition we assume that $u \in \mathscr{C}_{T}^{k}$ is a solution of (1.1). From $u(t) \in \dot{H}^{m}(\Omega)$ for $t \in[0, T]$ it follows that $\partial_{i}^{j} u(0) \in \stackrel{H}{H}^{m}(\Omega)$ for $j=0, \ldots$, $k-1$ (compare [9]). Differentiating (1.1) formally ( $j-2$ )-times, we obtain

$$
\begin{align*}
\partial_{t}^{j} u(0)= & \partial_{t}^{j-2} f(0)-\sum_{v=0}^{j-2}\binom{j-2}{v}\left\{\left[\partial_{t}^{v} A(0)\right] \partial_{t}^{j-2-v} u(0)\right. \\
& \left.+\left[\partial_{t}^{v} B(0)\right] \partial_{t}^{j-1-v} u(0)\right\} \tag{2.5}
\end{align*}
$$

where

$$
\begin{equation*}
\partial_{t}^{v} A(0)=\sum_{|\alpha| \leqslant 2 m}\left[\partial_{t}^{v} a_{\alpha}(0, .)+\partial_{t}^{v} \tilde{a}_{z}(0, .)\right] \partial_{x}^{\alpha} \tag{2.6}
\end{equation*}
$$

and $\partial_{t}^{y} B(0)$ is given analogously. We make the following definition.
Definition 2.1. We say that $\left(u^{0}, u^{1}, f\right)$ satisfies the compatibility condition of order $k \in \mathbb{N}$ if $u^{j} \in \dot{H}^{m}(\Omega)$ for $j=0, \ldots, k-1$, where $u^{j}$ is recursively defined by

$$
\begin{equation*}
u^{j}:=\partial_{t}^{j-2} f(0)-\sum_{v=0}^{j-2}\binom{j-2}{v}\left\{\left[\partial_{t}^{v} A(0)\right] u^{j-2-v}+\left[\partial_{t}^{v} B(0)\right] u^{j-1-v}\right\} \tag{2.7}
\end{equation*}
$$

for $j \geqslant 2$.

Remark. Let Assumption 1.1 be satisfied for some $k \geqslant[n / 2 m]+3$ and let $2 \leqslant j \leqslant k$. If

$$
u^{0} \in H^{j m}(\Omega), u^{1} \in H^{(j-1) m}(\Omega), \mathrm{a}_{t}^{\nu-2} f(0) \in H^{(j-v) m}(\Omega) \text { for } v=2, \ldots, j,
$$

then $u^{v}$ is well defined by (2.7) for $v=2, \ldots, j$. Furthermore, it follows from Lemma
8.2 of [9] and induction that

$$
\begin{equation*}
|u(0)|_{j}=\sum_{v=0}^{j}\left\|u^{v}\right\|_{(j-v) m} \leqslant D\left[\left\|u^{0}\right\|_{j m}+\left\|u^{1}\right\|_{(j-1) m}+|f(0)|_{j-2}\right], \tag{2.8}
\end{equation*}
$$

where $D>0$ depends only on $j,(1.11)$ and (1.13).

## 3. Elliptic equations

Consider the elliptic differential equations

$$
\begin{equation*}
A(t) u(t)=f(t) \quad \text { for } t \in[0, T] \tag{3.1}
\end{equation*}
$$

where $A$ satisfies Assumption 1.1 for some $k \geqslant k_{0}=[n / 2 m]+4$. We choose

$$
\left.\begin{array}{l}
c_{\alpha \beta} \in C_{b}^{(k-2) m+|\alpha|}([0, T] \times \bar{\Omega}),  \tag{3.2}\\
\tilde{c}_{\alpha \beta} \in C\left([0, T], H^{(k-3) m+|\alpha|}(\Omega)\right)
\end{array}\right\}
$$

such that

$$
\begin{equation*}
A(t) u(t)=\sum_{|\alpha|,|\beta| \leqslant m}(-1)^{|x|} \partial_{x}^{\alpha}\left\{\left[c_{\alpha \beta}(t)+\tilde{c}_{\alpha \beta}(t)\right] \partial_{x}^{\beta} u(t)\right\} \tag{3.3}
\end{equation*}
$$

(compare, e.g., [1], section 8). It follows from (1.5) that

$$
\begin{equation*}
(-1)^{m} \operatorname{Re} \sum_{|\alpha|=|\beta|=m}\left[c_{\alpha \beta}(t, x)+\tilde{c}_{\alpha \beta}(t, x)\right] \xi^{\alpha+\beta} \geqslant c_{1}|\xi|^{2 m} \tag{3.4}
\end{equation*}
$$

for $\xi \in \mathbb{R}^{n}, t \in[0, T]$ and $x \in \bar{\Omega}$. We set

$$
\begin{align*}
a(t, \varphi, \psi):= & \sum_{|\alpha|=|\beta| \leqslant m}\left\langle\left[c_{\alpha \beta}(t)+\tilde{c}_{\alpha \beta}(t)\right] \partial_{x}^{\beta} \varphi, \partial_{x}^{\alpha} \psi\right\rangle \\
& \text { for } \varphi, \psi \in H^{m}(\Omega), t \in[0, T] \tag{3.5}
\end{align*}
$$

and study instead of (3.1) the equation

$$
\begin{equation*}
a(t, u(t), \varphi)=\langle f(t), \varphi\rangle \quad \text { for every } \varphi \in C_{0}^{\infty}(\Omega), t \in[0, T] \tag{3.6}
\end{equation*}
$$

We prove the following lemma.

Lemma 3.1. Let $k \geqslant[n / 2 m]+4$ and let Assumption 1.1 be satisfied. If $f \in C([0, T]$, $\left.H^{j}(\Omega)\right)$ for some $j \leqslant(k-2) m$ and if $u \in C\left([0, T], \dot{H}^{m}(\Omega)\right)$ such that (3.6) holds, then $u \in C\left([0, T], H^{j+2 m}(\Omega)\right)$. Furthermore

$$
\begin{equation*}
\|u(t)\|_{j+2 m} \leqslant d\left(\|f(t)\|_{j}+\|u(t)\|\right) \quad \text { for } t \in[0, T], \tag{3.7}
\end{equation*}
$$

where $d>0$ depends only on $j$ and on

$$
\begin{equation*}
\sup _{|\alpha| \leqslant 2 m} \sup _{\{0, T]}\left(\sup _{|x| \leqslant(k-2) m} \sup _{x \in \Omega}\left|\partial_{x}^{\gamma} a_{\alpha}(t, x)\right|+\left\|\tilde{a}_{a}(t)\right\|_{(k-2) m}\right) \tag{3.8}
\end{equation*}
$$

(and on $\Omega$ ).
Proof. Let $\Omega$ be of the form (2.1). In the other cases Lemma 3.1 can be shown
analogously. At first we prove the assertion in the case $j=0$. We set

$$
\Omega_{R}:=\{x \in \Omega:|x|<R\}
$$

and choose $R>0$ so large that

$$
\begin{equation*}
\left\{x \in \Omega=\mathbb{R}^{n^{\prime}} \times \Omega^{\prime}:\left|x_{i}\right| \leqslant 1 \quad \text { for } i=1, \ldots, n^{\prime}\right\} \subset \Omega_{R} \tag{3.9}
\end{equation*}
$$

Let

$$
\begin{equation*}
T_{z}\left(\Omega_{R}\right):=\left\{x+z: x \in \Omega_{R}\right\} \quad \text { for } z=\left(z_{1}, \ldots, z_{n^{\prime}}, 0, \ldots, 0\right) \tag{3.10}
\end{equation*}
$$

Note that $(k-3) m \geqslant[n / 2]+1$ and that therefore

$$
c_{\alpha \beta}(t)+\tilde{c}_{\alpha \beta}(t) \in C_{b}^{|\alpha|}(\bar{\Omega}) \quad \text { for }|\alpha|,|\beta| \leqslant m, t \in[0, T]
$$

by the lemma of Sobolev. Hence $a(t, \varphi, \psi)$ is right $m$-smooth for every $t \in[0, T]$ in the sense of Definition 9.1 of [1]. By the proof of Theorem 9.8 in [1] it follows from (3.6) that

$$
\begin{equation*}
\|u(t)\|_{H^{2 m}\left(T_{z}\left(\Omega_{R}\right)\right)} \leqslant d_{1}\left(\|f(t)\|_{L_{2}\left(T_{z}\left(\Omega_{R+i}\right)\right)}+\|u(t)\|_{H^{m}\left(T_{z}\left(\Omega_{R+1}\right)\right)}\right) . \tag{3.11}
\end{equation*}
$$

for every $z \in S:=\mathbb{Z}^{n^{\prime}} \times\{(0, \ldots, 0)\}$, where $d_{1}>0$ can be chosen to be independent of z. Note that

$$
\begin{equation*}
\Omega=\bigcup_{z \in S} T_{z}\left(\Omega_{R}\right)=\bigcup_{z \in S} T_{z}\left(\Omega_{R+1}\right), \tag{3.12}
\end{equation*}
$$

and that every $x \in \Omega$ is contained only in a fixed finite number of sets $T_{z}\left(\Omega_{R+1}\right)$. Hence summation over all $z \in S$ yields $u(t) \in H^{2 m}(\Omega)$ and

$$
\begin{equation*}
\|u(t)\|_{2 m} \leqslant d_{2}\left(\|f(t)\|+\|u(t)\|_{m}\right) \tag{3.13}
\end{equation*}
$$

for $t \in[0, T]$. In view of (3.4) there exist constants $d_{3}, d_{4}>0$ such that

$$
\begin{equation*}
\operatorname{Re} a(t, \varphi, \varphi) \geqslant d_{3}\|\varphi\|_{m}^{2}-d_{4}\|\varphi\|^{2} \quad \text { for } \varphi \in \dot{H}^{m}(\Omega), t \in[0, T] \tag{3.14}
\end{equation*}
$$

Hence it follows from (3.6) (with $\varphi:=u(t)$ ) and (3.13) that (3.7) holds in the case $j=0$. In particular, we have that $\|u(t)\|_{2 m}$ is bounded on $[0, T]$.

If $t, t_{1} \in[0, T]$, we obtain from (3.6) that

$$
\begin{equation*}
a\left(t, u(t)-u\left(t_{1}\right), \varphi\right)=\left\langle g_{1}\left(t, t_{1}^{\prime}\right), \varphi\right\rangle \quad \text { for every } \varphi \in C_{0}^{\infty}(\Omega), t \in[0, T] \tag{3.15}
\end{equation*}
$$

with

$$
\begin{equation*}
g_{1}\left(t, t_{1}\right):=f(t)-f\left(t_{1}\right)+\left[A\left(t_{1}\right)-A(t)\right] u\left(t_{1}\right) \tag{3.16}
\end{equation*}
$$

With Lemma 8.2 in [9] we conclude that $\left\|g_{1}\left(t, t_{1}\right)\right\| \rightarrow 0$ as $t \rightarrow t_{1}$. Hence (3.7) applied to (3.15) yields $\left\|u(t)-u\left(t_{1}\right)\right\|_{2 m} \rightarrow 0$ as $t \rightarrow t_{1}$ and therefore $u \in C\left([0, T], H^{2 m}(\Omega)\right)$. This proves the assertion for $j=0$.

Now let Lemma 3.1 be proved for $j=0, \ldots, J \leqslant(k-2) m-1$. By the induction hypothesis we have $u \in C\left([0, T], H^{j+2 m}(\Omega)\right)$ and (3.7) with $j=J$. Let $\psi \in C_{0}^{\infty}(\Omega)$. From (3.6) and (3.3) we conclude that

$$
\begin{equation*}
a(t, \psi u(t), \varphi)=\left\langle g_{2}(t), \varphi\right\rangle \quad \text { for } \varphi \in C_{0}^{\infty}(\Omega), t \in[0, T] \tag{3.17}
\end{equation*}
$$

with

$$
\begin{equation*}
g_{2}(t):=\psi f(t)+A(t)(\psi u(t))-\psi A(t) u(t) . \tag{3.18}
\end{equation*}
$$

Note that the derivatives of $u$ of order $2 m$ vanish on the right-hand side of (3.18). From $\left.\tilde{a}_{\alpha} \in C(0, T], H^{(k-2) m}(\Omega)\right)$ for $|\alpha| \leqslant 2 m$ and Lemma 8.2 in [9] it follows that $g_{2} \in$
$C\left([0, T], H^{J+1}(\Omega)\right)$. Setting $\varphi=\partial \Phi / \partial x_{i}$ in $(3,17)$ with $\Phi \in C_{0}^{\infty}(\Omega)$, we conclude that

$$
\begin{equation*}
a\left(t, \frac{\partial(\psi u(t))}{\partial x_{i}}, \Phi\right)=\left\langle g_{3}(t), \Phi\right\rangle \quad \text { for } \Phi \in C_{0}^{\infty}(\Omega), t \in[0, T] \tag{3.19}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{3}(t):=\frac{\partial g_{2}(t)}{\partial x_{i}}-\sum_{|\alpha| \leqslant 2 m}\left(\frac{\partial}{\partial x_{i}}\left[a_{a}(t)+\tilde{a}_{\alpha}(t)\right]\right) \partial_{x}^{\alpha}(\psi u(t)) \tag{3.20}
\end{equation*}
$$

We apply (3.7) with $j=J$ to (3.19) and obtain

$$
\left.\begin{array}{rl}
\left\|\frac{\partial \psi u(t)}{\partial x_{i}}\right\|_{J+2 m} & \leqslant d\left(\left\|g_{3}(t)\right\|_{J}+\|\psi u(t)\|\right) \\
& \leqslant d_{1}\left(\left\|\frac{\partial(\psi f(t))}{\partial x_{i}}\right\|_{J}+\|u(t)\|_{H}+2 m(\operatorname{supp} \psi)\right. \tag{3.21}
\end{array}\right)
$$

for $t \in[0, T]$ and $i=1, \ldots, n$, where $d_{1}>0$ depends only on (3.8), $J$ and $\psi$.
An analogous procedure can be performed around the boundary points. In fact, let $x_{0} \in \partial \Omega$ and let $U$ be an open neighbourhood of $x_{0}$ such that there exists a $C^{\infty}$ mapping transforming $U \cap \Omega$ in $K_{R}^{+}:=\left\{x \in \mathbb{R}^{n}:|x|<R, x_{n}>0\right\}$ and $\bar{U} \cap \partial \Omega$ in $\Gamma_{R}:=\left\{x \in \mathbb{R}^{n}:|x| \leqslant R, x_{n}=0\right\}$. Since such a mapping preserves the properties of our elliptic equation; we can assume that $U \cap \Omega=K_{R}^{+}, \bar{U} \cap \partial \Omega=\Gamma_{R}$.

For given $R^{\prime} \in(0, R)$ we choose $\psi \in C_{0}^{\infty}\left(K_{R}\right)$ with $\psi=1$ on $K_{R^{\prime}}:=\left\{x \in \mathbb{R}^{n}:|x|<R^{\prime}\right\}$. Then (3.21) follows as above for $i=1, \ldots, n-1$, since $\partial(\psi u(t)) / \partial x_{i} \in \dot{H}^{m}\left(K_{R}^{+}\right)$for $i=1, \ldots, n-1$.

Since $u \in H^{2 m}(\Omega)$, we conclude from (3.17) that

$$
A(t)[\psi u(t)]=g_{2}(t) \quad \text { for } t \in[0, T]
$$

and hence

$$
\begin{equation*}
\left[a_{(0, \ldots, 0,2 m)}(t)+\tilde{a}_{(0, \ldots, 0,2 m)}(t)\right] \partial_{x}^{(0, \ldots, 0,2 m)}(\psi u(t))=g_{4}(t) \tag{3.22}
\end{equation*}
$$

with a suitable $g_{4}$, where according to (3.18) and (3.21)

$$
\left\|g_{4}(t)\right\|_{J+1} \leqslant d_{2}\left[\|\psi f(t)\|_{J+1}+\|u(t)\|_{H^{J+2 m(\text { supp } \psi)}}\right]
$$

Since

$$
\left|a_{(0, \ldots, 0,2 m)}(t, x)+\tilde{a}_{(0 \ldots \ldots 0,2 m)}(t, x)\right| \geqslant c_{1}>0
$$

by (1.5), we conclude from (3.22) and (1.4) that

$$
\begin{equation*}
\mid \partial_{x}^{(0, \ldots, 0,2 m)} \psi u(t) \|_{J+1} \leqslant d_{3}\left(\|\psi f(t)\|_{J+1}+\|u(t)\|_{\left.H^{J+2 m(s u p p} \psi\right)}\right) . \tag{3.23}
\end{equation*}
$$

We choose a finite number of suitable functions $\psi_{1}, \ldots, \psi_{v}$. Then we conclude from (3.21) and (3.23) that

$$
\begin{equation*}
\|u(t)\|_{H^{J+1+2 m\left(T_{z}\left(\Omega_{R}\right)\right)}} \leqslant d_{4}\left(\|f(t)\|_{H^{J+1}\left(T_{x}\left(\Omega_{R+1}\right)\right)}+\|u(t)\|_{\left.H^{J+2 m\left(T_{z}\left(\Omega_{R+1}\right)\right.}\right)}\right) \tag{3.24}
\end{equation*}
$$

for $t \in[0, T]$ and $z=(0, \ldots, 0)$, where $d_{4}>0$ depends only on (3.8) and $J$. Note that we can prove (3.24) for arbitrary $z \in S=\mathbb{Z}^{n} \times\{(0, \ldots, 0)\}$ by the same argument using the functions $\psi_{1}(x-z), \ldots, \psi_{v}(x-z)$. Hence $d_{4}$ can be chosen such that (3.24) holds for every $z \in S$. Using the argument leading to (3.13) we conclude with the induction hypothesis that (3.7) holds for $j=J+1$. Finally, $u \in C\left([0, T], H^{J+1+2 m}(\Omega)\right)$ follows
applying (3.7) with $j=J+1$ to (3.15), since $\left\|g_{1}\left(t, t_{1}\right)\right\|_{J+1} \rightarrow 0$ as $t \rightarrow t_{1}$ by Lemma 8.2 of [9].

## 4. Existence and uniqueness of the solution

In this section we suppose that Assumption 1.1 is satisfied for $k=k_{0}=[n / 2 m]+4$. Let $a(t, \varphi, \psi)$ be defined by (3.2)-(3.5). Note that (1.6) implies that we can choose the coefficients of $a(t, \varphi, \psi)$ such that

$$
\begin{equation*}
c_{\alpha \beta}+\tilde{c}_{\alpha \beta}=\overline{\left(c_{\beta \alpha}+\tilde{c}_{\rho \alpha}\right)} \text { for } m+1 \leqslant|\alpha|+|\beta| \leqslant 2 m \text {. } \tag{4.1}
\end{equation*}
$$

Furthermore it follows from (1.5) that there exist constants $d, d_{1}>0$ such that for

$$
\begin{equation*}
a_{d}(t, \varphi, \psi):=\sum_{\substack{|\alpha|,|\beta| \leqslant m \\|\alpha|+|\beta| \geqslant m+1}}\left\langle\left[c_{\alpha \beta}(t)+\tilde{c}_{\alpha \beta}(t)\right] \partial_{x}^{\beta} \varphi, \partial_{x}^{\alpha} \psi\right\rangle+d\langle\varphi, \psi\rangle \tag{4.2}
\end{equation*}
$$

it holds that

$$
\begin{equation*}
a_{d}(t, \varphi, \varphi)=\operatorname{Re} a_{d}(t, \varphi, \varphi) \geqslant d_{1}\|\varphi\|_{m}^{2} \quad \text { for } \varphi \in \dot{H}^{m}(\Omega), t \in[0, T] \tag{4.3}
\end{equation*}
$$

Here and in the following we denote by $d, d_{1}, d_{2}, \ldots$ positive constants depending only on $c_{1}, c_{2}, c_{3}$ (of Assumption 1.1) and on (1.11), (1.13).

In order to prove the existence of a solution of (1.1) we use the method of Faedo-Galerkin and follow the considerations in [3]. We suppose that

$$
\begin{equation*}
u^{0} \in \dot{H}^{m}(\Omega) \cap H^{2 m}(\Omega), \quad u^{1} \in \dot{H}^{m}(\Omega), \quad f \in H^{1}\left([0, T], L_{2}(\Omega)\right) \tag{4.4}
\end{equation*}
$$

where the last condition means that $f, f^{\prime} \in L_{2}\left([0, T], L_{2}(\Omega)\right)$. Note that this implies that $f \in C\left([0, T], L_{2}(\Omega)\right)$. Let

$$
\left\{\xi_{1}, \xi_{2}, \ldots\right\} \subset \dot{H}^{m}(\Omega) \cap H^{2 m}(\Omega)
$$

such that every finite subset is linearly independent and span $\left\{\xi_{1}, \xi_{2}, \ldots\right\}$ is dense in $L_{2}(\Omega)$. We seek an approximate solution

$$
\begin{equation*}
u_{j}(t)=\sum_{v=1}^{j} u_{j v}(t) \xi_{v} \tag{4.5}
\end{equation*}
$$

of

$$
\left.\begin{array}{l}
\left.\left\langle u_{j}^{\prime \prime}(t)+A(t) u_{j}(t)+B(t) u_{j}^{\prime}(t)-f(t), \xi_{v}\right\rangle=0 \quad \text { for } t \in[0, T], v=1, \ldots, j,\right\}  \tag{4.6}\\
u_{j}(0)=u_{j}^{\prime}, \quad u_{j}^{\prime}(0)=u_{j}^{\mathrm{I}},
\end{array}\right\}
$$

where $u_{j}^{0}, u_{j}^{1} \in \operatorname{span}\left\{\xi_{1}, \ldots, \xi_{j}\right\}$ are chosen so that $\left\|u_{j}^{v}-u^{\nu}\right\|_{(2-v) m} \rightarrow 0$ as $j \rightarrow \infty$ ( $v=0,1$ ). Since (4.6) is a system of ordinary differential equations for the coefficients $u_{j v}(v=1, \ldots, j)$, it follows from standard classical theory that a solution $\left(u_{j 1}, \ldots, u_{j j}\right) \in C^{2}([0, T])$ exists. From (4.6) we conclude that $\left(u_{j 1}, \ldots, u_{j j}\right)$ $\in H^{3}([0, T])$. Hence $u_{j} \in H^{3}\left([0, T], H^{2 m}(\Omega)\right)$.
We set

$$
\begin{equation*}
\left|u_{j}(t)\right|_{E}:=\left[a_{d}\left(t, u_{j}(t), u_{j}(t)\right)+\left\|u_{j}^{\prime}(t)\right\|^{2}\right]^{1 / 2} \quad \text { for } t \in[0, T] . \tag{4.7}
\end{equation*}
$$

It holds that

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t}\left|u_{j}(t)\right|_{E}^{2}=2 \operatorname{Re}\left[a_{d}\left(t, u_{j}(t), u_{j}^{\prime}(t)\right)+\left\langle u_{j}^{\prime \prime}(t), u_{j}^{\prime}(t)\right\rangle\right]+a_{d}^{\prime}\left(t, u_{j}(t), u_{j}(t)\right) \tag{4.8}
\end{equation*}
$$

for $t \in[0, T]$, where

$$
\begin{equation*}
a_{d}^{\prime}(t, \varphi, \psi)=\sum_{m+1 \leqslant|\alpha|+|\beta| \leqslant 2 m}\left\langle\left[c_{\alpha \beta}^{\prime}(t)+\tilde{c}_{\alpha \beta}^{\prime}(t)\right] \partial_{x}^{\beta} \varphi, \partial_{x}^{\alpha} \psi\right\rangle . \tag{4.9}
\end{equation*}
$$

Replacing $\xi_{v}$ in (4.6) by $u_{j}^{\prime}(t) \in \operatorname{span}\left\{\xi_{1}, \ldots, \xi_{j}\right\}$, we conclude that $\frac{\mathrm{d}}{\mathrm{d} t}\left|u_{j}(t)\right|_{E}^{2}$

$$
\begin{aligned}
= & 2 \operatorname{Re}\left\langle f(t)-B(t) u_{j}^{\prime}(t)-\sum_{|\alpha|+|\beta| \leqslant m} \partial_{x}^{\alpha}\left\{\left[c_{\alpha \beta}(t)+\tilde{c}_{\alpha \beta}(t)\right] \partial_{x}^{\beta} u_{j}(t)\right\}+d u_{j}(t), u_{j}^{\prime}(t)\right\rangle \\
& +a_{d}^{\prime}\left(t, u_{j}(t), u_{j}(t)\right) .
\end{aligned}
$$

With (1.7) and (4.3) we obtain

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\left|u_{j}(t)\right|_{E}^{2} \leqslant 2\left[\|f(t)\|+d_{2}\left|u_{j}(t)\right|_{E}\right]\left|u_{j}(t)\right|_{E}
$$

and

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t}\left|u_{j}(t)\right|_{E} \leqslant\|f(t)\|+d_{2}\left|u_{j}(t)\right|_{E} \quad \text { for } t \in[0, T] \tag{4.10}
\end{equation*}
$$

With Gronwall's Lemma it follows that there is a constant $C>0$ such that $\left|u_{j}(t)\right|_{E} \leqslant C$ for $t \in[0, T]$ and $j \in \mathbb{N}$.

In a similar way we prove that $\left|u_{j}^{\prime}(t)\right|_{E}$ is bounded. We differentiate (4.6) with respect to $t$ and obtain

$$
\begin{equation*}
\left\langle u_{j}^{\prime \prime \prime}(t)+A(t) u_{j}^{\prime}(t)+B(t) u_{j}^{\prime \prime}(t)+A^{\prime}(t) u_{j}(t)+B^{\prime}(t) u_{j}^{\prime}(t)-f^{\prime}(t), \xi_{v}\right\rangle=0 \tag{4.11}
\end{equation*}
$$

for $t \in[0, T], v=1, \ldots, j$. We replace $\xi_{v}$ by $u_{j}^{\prime \prime}(t)$ and conclude by the same argument as above that

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t}\left|u_{j}^{\prime}(t)\right|_{E}^{2}=-\left\langle A^{\prime}(t) u_{j}(t), u_{j}^{\prime \prime}(t)\right\rangle+r(t) \tag{4.12}
\end{equation*}
$$

with

$$
\begin{equation*}
|r(t)| \leqslant 2\left\|f^{\prime}(t)\right\|\left|u_{j}^{\prime}(t)\right|_{E}+d_{3}\left|u_{j}^{\prime}(t)\right|_{E}^{2} \leqslant\left\|f^{\prime}(t)\right\|^{2}+\left(d_{3}+1\right)\left|u_{j}^{\prime}(t)\right|_{E}^{2} \tag{4.13}
\end{equation*}
$$

a.e. in $[0, T]$. Let $t_{0} \in[0, T]$. We use

$$
\int_{0}^{t}\left\langle A^{\prime}(\tau) u_{j}(\tau), u_{j}^{\prime \prime}(\tau)\right\rangle \mathrm{d} \tau=\left.\left\langle A^{\prime}(\tau) u_{j}(\tau), u_{j}^{\prime}(\tau)\right\rangle\right|_{\tau=0} ^{t}-\int_{0}^{t}\left\langle\partial_{\tau}\left(A^{\prime}(\tau) u_{j}(\tau)\right), u_{j}^{\prime}(\tau)\right\rangle \mathrm{d} \tau
$$

and

$$
\begin{aligned}
& \left|\left\langle A^{\prime}(t) u_{j}(t), u_{j}^{\prime}(t)\right\rangle\right|=\left\lvert\, a^{\prime}\left(t, u_{j}(t),\left.u_{j}^{\prime}(t)|\leqslant \varepsilon| u_{j}^{\prime}(t)\right|_{E} ^{2}+\frac{d_{4}}{\varepsilon}\left|u_{j}(t)\right|_{E}^{2},\right.\right. \\
& \left\langle\partial_{\varepsilon}\left(A^{\prime}(\tau) u_{j}(\tau)\right), u_{j}^{\prime}(\tau)\right\rangle \leqslant d_{1}^{\prime}\left(\left\|u_{j}(\tau)\right\|_{m}^{2}+\left\|u_{j}^{\prime}(\tau)\right\|_{m}^{2}\right)
\end{aligned}
$$

with $\varepsilon>0$. Here and in the following we denote by $d_{1}^{\prime}, d_{2}^{\prime}, \ldots$ positive constants
depending on $c_{1}, c_{2}$ and $c_{3}$, and on (1.11) and (1.12). Integration of (4.12) yields

$$
\begin{align*}
\left|u_{j}^{\prime}(t)\right|_{E}^{2} \leqslant & d_{5}\left[\left|u_{j}^{\prime}(0)\right|_{E}^{2}+\left|u_{j}(0)\right|_{E}^{2}\right]+\varepsilon\left|u_{j}^{\prime}(t)\right|_{E}^{2}+\frac{d_{4}}{\varepsilon}\left|u_{j}(t)\right|_{E}^{2} \\
& +d_{2}^{\prime} \int_{0}^{t}\left[\left|u_{j}^{\prime}(\tau)\right|_{E}^{2}+\left|u_{j}(\tau)\right|_{E}^{2}+\left\|f^{\prime}(\tau)\right\|^{2}\right] \mathrm{d} \tau . \tag{4.14}
\end{align*}
$$

Note that $\left|u_{j}^{\prime}(0)\right|_{E}^{2}$ is bounded as $j \rightarrow \infty$, since $\left\|\partial_{t}^{v} u_{j}(0)-u^{v}\right\|_{(2-v) m} \rightarrow 0$ as $j \rightarrow \infty$ for $v=0,1,2$, where $u^{2}=f(0)-A(0) u^{0}-B(0) u^{1}$. In fact, (4.6) implies that $u_{j}^{\prime \prime}(0)$ converges weakly to $u^{2}$ in $L_{2}(\Omega)$ and

$$
\varlimsup_{j \rightarrow \infty}\left\|u_{j}^{\prime \prime}(0)\right\|^{2}=\varlimsup_{j \rightarrow \infty}\left\langle f(0)-A(0) u_{j}^{0}-B(0) u_{j}^{1}, u_{j}^{\prime \prime}(0)\right\rangle \leqslant\left\|u^{2}\right\| \varlimsup_{j \rightarrow \infty}\left\|u_{j}^{\prime \prime}(0)\right\| .
$$

Hence $\varlimsup_{j \rightarrow \infty}\left\|u_{j}^{\prime \prime}(0)\right\| \leqslant\left\|u^{2}\right\|$ and it follows that $u_{j}^{\prime \prime}(0) \rightarrow u^{2}$ in $L_{2}(\Omega)$ as $j \rightarrow \infty$.
With Gronwall's lemma we conclude from (4.14) and the boundedness of $\left|u_{j}(t)\right|_{E}^{2}$ that there exists a $C>0$ such that $\left|u_{j}^{\prime}(t)\right|_{E}^{2} \leqslant C$ for $t \in[0, T], j \in \mathbb{N}$. Hence we can extract a subsequence converging weakly to a

$$
u \in H^{2}\left([0, T], L_{2}(\Omega)\right) \cap H^{1}\left([0, T], \dot{H}^{m}(\Omega)\right)
$$

By a standard argument (compare §3.8.2 in [10]) it follows that

$$
\begin{array}{ll}
u^{\prime \prime}(t)+A(t) u(t)+B(t) u^{\prime}(t)=f(t) & \text { a.e. in }[0, T], \\
u(t) \in \stackrel{\circ}{H}^{m}(\Omega) & \text { for } t \in[0, T], \\
u(0)=u^{0}, u^{\prime}(0)=u^{1} . & \tag{4.15c}
\end{array}
$$

From (4.15) we obtain by the argument leading to (4.10) that

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t}|u(t)|_{E} \leqslant\|f(t)\|+d_{2}|u(t)|_{E} \quad \text { a.e. in }[0, T] . \tag{4.16}
\end{equation*}
$$

Using Gronwall's lemma we conclude that

$$
\begin{equation*}
|u(t)|_{E} \leqslant \mathrm{e}^{d^{2} t}\left(|u(0)|_{E}+\int_{0}^{t}\|f(\tau)\| \mathrm{d} \tau\right) \text { for } t \in[0, T] \tag{4.17}
\end{equation*}
$$

and with (4.3) and (4.7) we obtain

$$
\begin{equation*}
|u(t)|_{1} \leqslant d_{5} \mathrm{e}^{d_{2} t}\left(|u(0)|_{1}+\int_{0}^{t}\|f(\tau)\| \mathrm{d} \tau\right) \text { for } t \in[0, T] . \tag{4.18}
\end{equation*}
$$

In the rest of this section we prove that $u \in \mathscr{C}_{T}^{2}$ and derive an estimate for $|u(t)|_{2}$. To this end we introduce a convenient concept of a weak solution.

Definition 4.1. Let $u^{0} \in \dot{H}^{m}(\Omega), u^{1} \in L_{2}(\Omega)$ and $f \in L_{1}\left([0, T], L_{2}(\Omega)\right)$. We say that $u \in H^{1}\left([0, T], L_{2}(\Omega)\right) \cap L_{2}\left([0, T], \dot{H}^{m}(\Omega)\right)$
is a weak solution of (1.1) if $u(0)=u^{0}$ and

$$
\begin{align*}
& \int_{0}^{T}\left[-a(t, u(t), v(t))+\left\langle B(t) u(t)+u^{\prime}(t), v^{\prime}(t)\right\rangle+\left\langle B^{\prime}(t) u(t)+f(t), v(t)\right\rangle\right] \mathrm{d} t \\
& =\left\langle u^{1}+B u^{0}, v(0)\right\rangle \tag{4.19}
\end{align*}
$$

for every $v \in H^{1}\left([0, T], L_{2}(\Omega)\right) \cap L_{2}\left([0, T], \dot{H}^{m}(\Omega)\right)$ with $v(T)=0($ note that $a(t, u(t)$, $v(t))$ is defined by (3.5)).

Remarks. (1) Since $\partial_{1} v \in L_{2}\left([0, T], L_{2}(\Omega)\right)$ it holds $v \in C\left([0, T], L_{2}(\Omega)\right)$. Hence $v(0)$, $v(T)$ are well defined.
(2) Integration by parts shows that $u \in \mathscr{C}_{T}^{2}$ is a solution of (1.1) if and only if it is a weak solution.

In a first step we prove the uniqueness of a weak solution.
Lemma 4.1. Let Assumption 1.1 be satisfied for $k=k_{0}$. If $u^{0}=u^{1}=0$ and $f(t)=0$ on [ $0, T$ ], and if

$$
u \in H^{1}\left([0, T], L_{2}(\Omega)\right) \cap L_{2}\left([0, T], \dot{H}^{m}(\Omega)\right)
$$

is a weak solution of $(1.1)$, then $u(t)=0$ on $[0, T]$.
Proof. We proceed analogously to the proof of the uniqueness in §3.8.2 of [10]. Let $s \in(0, T)$ and

$$
v(t):= \begin{cases}-\int_{t}^{s} u(\sigma) \mathrm{d} \sigma & \text { for } t<s  \tag{4.20}\\ 0 & \text { for } t \geqslant s\end{cases}
$$

Note that

$$
\begin{aligned}
& \frac{\mathrm{d}}{\mathrm{~d} t}\left[a_{d}(t, v(t), v(t))+2 \operatorname{Re}\left\langle B^{\prime}(t) v(t), v(t)\right\rangle+\|u(t)\|^{2}\right] \\
& \quad=2 \operatorname{Re}\left[a_{d}(t, u(t), v(t))+\left\langle B^{\prime}(t) u(t), v(t)\right\rangle+\left\langle u^{\prime}(t), v^{\prime}(t)\right\rangle\right]+r(t)
\end{aligned}
$$

with

$$
|r(t)| \leqslant d_{3}^{\prime}\left[\|v(t)\|_{m}^{2}+\|u(t)\|^{2}\right]
$$

a.e. in $[0, s)$. We conclude from (4.19) that

$$
\begin{equation*}
a_{d}(0, v(0), v(0))+2 \operatorname{Re}\left\langle B^{\prime}(0) v(0), v(0)\right\rangle+\|u(s)\|^{2} \leqslant d_{4}^{\prime} \int_{0}^{s}\left[\|v(t)\|_{m}^{2}+\|u(t)\|^{2}\right] \mathrm{d} t . \tag{4.21}
\end{equation*}
$$

There exists a $d_{5}^{\prime}>0$ such that

$$
a_{d}(0, v(0), v(0))+2 \operatorname{Re}\left\langle B^{\prime}(0) v(0), v(0)\right\rangle \geqslant \frac{d_{1}}{2}\|v(0)\|_{m}^{2}-d_{5}^{\prime}\|v(0)\|^{2} .
$$

Hence we obtain from (4.21) with the argument proving the uniqueness in §3.8.2 of [10] that $u(t)=0$ for $t \in[0, T]$.

In the next step we prove the existence of a weak solution $u \in \mathscr{C}_{T}^{1}$ of (1.1).
Lemma 4.2. Let Assumption 1.1 be satisfied for $k=k_{0}$. Furthermore let $u^{0} \in \dot{H}^{m}(\Omega)$, $u^{1} \in L_{2}(\Omega)$ and $f \in L_{1}\left([0, T], L_{2}(\Omega)\right)$. Then there exists a weak solution $u \in \mathscr{C}_{T}^{1}$ of (1.1).

Proof. We approximate $u^{0}, u^{1}$ and $f$ by sequences $\left\{u_{j}^{0}\right\}$ in $H^{2 m}(\Omega) \cap \dot{H}^{m}(\Omega),\left\{u_{j}^{1}\right\}$ in $\dot{H}^{m}(\Omega)$ and $\left\{f_{j}\right\}$ in $H^{1}\left([0, T], L_{2}(\Omega)\right)$ such that $\left\|u^{v}-u_{j}^{v}\right\|_{(1-v) m} \rightarrow 0$ as $j \rightarrow \infty$
$(v=0,1)$ and

$$
\begin{equation*}
\int_{0}^{T}\left\|f(\tau)-f_{j}(\tau)\right\| \mathrm{d} \tau \rightarrow 0 \quad \text { as } j \rightarrow \infty . \tag{4.22}
\end{equation*}
$$

Let $u_{j} \in H^{2}\left([0, T], L_{2}(\Omega)\right) \cap H^{1}\left([0, T], \dot{H}^{m}(\Omega)\right) \subset \mathscr{C}_{T}^{1}$ be the solution of (4.15) with data $u_{j}^{0}, u_{j}^{1}$ and $f_{j}$, which we have constructed above. We apply (4.18) to

$$
\begin{align*}
& u_{j}^{\prime \prime}(t)-u_{f}^{\prime \prime}(t)+A(t)\left[u_{j}(t)-u_{f}(t)\right]+B(t)\left[u_{j}^{\prime}(t)-u_{l}^{\prime}(t)\right]=f_{j}(t)-f_{f}(t) \\
& \quad \text { a.e. in }[0, T] \tag{4.23}
\end{align*}
$$

and obtain that $\sup _{\mathbf{N O}_{1}, T_{1}}\left|u_{j}(t)-u_{\ell}(t)\right|_{1} \rightarrow 0$ as $j, \ell \rightarrow \infty$. Let $u \in \mathscr{C} \frac{1}{T}$ be the limit of $\left\{u_{j}\right\}$. Since every $u_{j}$ is a weak solution of (1.1) with data $u_{j}^{0}, u_{j}^{1}$ and $f_{j}$, it follows that $u$ is a weak solution of (1.1).

We note that (4.16) implies

$$
|u(t)|_{E} \leqslant|u(0)|_{E}+\int_{0}^{t}\left[\|f(\tau)\|+d_{2}|u(\tau)|_{E}\right] \mathrm{d} \tau
$$

With (4.3) it follows that

$$
\begin{equation*}
|u(t)|_{1} \leqslant d_{6}\left(|u(0)|_{1}+\int_{0}^{t}\left[\|f(\tau)\|+|u(\tau)|_{1}\right] \mathrm{d} \tau\right) \quad \text { for } t \in[0, T] . \tag{4.24}
\end{equation*}
$$

This estimate holds for every $u_{j}$ used in the above proof and therefore even for the weak solution $u \in \mathscr{C}_{T}^{1}$ of (1.1).

Consider the solution $u \in H^{2}\left([0, T], L_{2}(\Omega)\right) \cap H^{1}\left([0, T], \stackrel{\circ}{H}^{m}(\Omega)\right)$ of $(4.15)$ constructed above. It holds that

$$
\begin{equation*}
A(t) u(t)=f(t)-u^{\prime \prime}(t)-B(t) u^{\prime}(t) \tag{4.25}
\end{equation*}
$$

a.e. in $[0, T]$. We consider a fixed $t \in[0, T]$ and conclude by Lemma 3.1 (applied to functions being constant in $t$ ) that

$$
\begin{equation*}
\|u(t)\|_{2 m} \leqslant d_{7}\left[\|f(t)\|+\left\|u^{\prime \prime}(t)\right\|+\left\|u^{\prime}(t)\right\|_{m}+\|u(t)\|\right] . \tag{4.26}
\end{equation*}
$$

Since this inequality holds a.e. in $[0, T]$, we obtain that $u \in L_{2}\left([0, T], H^{2 m}(\Omega)\right)$.
From (4.11), $u_{j}^{\prime}(0) \rightarrow u^{1}, u_{j}^{\prime \prime}(0) \rightarrow u^{2}$ and the construction of the solution $u$ of (4.15) it follows by a standard argument (compare §3.8.2 of [10]) that $v:=u^{\prime}$ is a weak solution of

$$
\begin{align*}
& v^{\prime \prime}(t)+A(t) v(t)+B(t) v^{\prime}(t)=f^{\prime}(t)-A^{\prime}(t) u(t)-B^{\prime}(t) u^{\prime}(t) \text { for } t \in[0, T],  \tag{4.27a}\\
& v(t) \in \dot{H}^{m}(\Omega) \text { for } t \in[0, T],  \tag{4.27b}\\
& v(0)=u^{1}, \quad v^{\prime}(0)=u^{2} . \tag{4.27c}
\end{align*}
$$

Since $f^{\prime}-A^{\prime} u-B^{\prime} u^{\prime} \in L_{2}\left([0, T], L_{2}(\Omega)\right)$, Lemma 4.2 yields the existence of a weak solution $v \in \mathscr{C}{ }_{T}^{1}$ of (4.27). From the uniqueness of $v$ we obtain $u^{\prime}=v$ and hence $u^{\prime} \in \mathscr{C} \mathscr{C}_{T}^{1}$. With $u \in \mathscr{C}{ }_{T}^{1}$, (4.25) and Lemma 3.1 we conclude that $u \in \mathscr{C}_{T}^{2}$. We apply (4.24) to (4.27) and obtain

$$
\begin{equation*}
\left|u^{\prime}(t)\right|_{1} \leqslant d_{6}\left(\left|u^{\prime}(0)\right|_{1}+\int_{0}^{t}\left\|f^{\prime}(t)-A^{\prime}(t) u(t)-B^{\prime}(t) u^{\prime}(t)\right\| \mathrm{d} t\right) . \tag{4.28}
\end{equation*}
$$

With (4.24) and (4.26) it follows that

$$
\begin{equation*}
|u(t)|_{2} \leqslant d_{8}\left(|u(0)|_{2}+\int_{0}^{t}\left[\left\|f^{\prime}(\tau)\right\|+\|f(\tau)\|+|u(\tau)|_{2}\right] \mathrm{d} \tau+\|f(t)\|\right) \tag{4.29}
\end{equation*}
$$

for $t \in[0, T]$. Thus we have proved the following lemma.
Lemma 4.3. Let Assumption 1.1 be satisfied for $k=k_{0}$. If $u^{0} \in H^{2 m}(\Omega) \cap \dot{H}^{m}(\Omega)$, $u^{1} \in \dot{H}^{m}(\Omega)$ and $f \in H^{1}\left([0, T], L_{2}(\Omega)\right)$, then (1.1) has a unique solution $u \in \mathscr{C}_{T}^{2}$. Furthermore (4.29) holds, where $d_{8}>0$ depends only on $c_{1}, c_{2}, c_{3}$ of Assumption 1.1 and on (1.11), (1.13).

## 5. Higher regularity

We prove Theorem 1.1 by induction with respect to $j$. Instead of (1.10) we prove

$$
\begin{align*}
|u(t)|_{j} \leqslant & D_{3}\left(|u(0)|_{j}+\int_{0}^{t}\left[\left\|\partial_{i}^{j-1} f(\tau)\right\|+|f(\tau)|_{j-2}+C_{3}|u(\tau)|_{j}\right] \mathrm{d} \tau\right) \\
& +D_{4}|f(t)|_{j-2} \text { for } t \in[0, T] \tag{5.1}
\end{align*}
$$

where $C_{3}>0$ depends only on $j, c_{1}, c_{2}$ and $c_{3}$ (of Assumption 1.1) and on (1.11), (1.12), while $D_{3}, D_{4}>0$ depend only on $j, c_{1}, c_{2}, c_{3}$ and on (1.11), (1.13). Note that it follows from (5.1) by Gronwall's lemma that

$$
\begin{align*}
|u(t)|_{j} \leqslant & \mathrm{e}^{D_{3} c_{3}}\left(2 D_{3}|u(0)|_{j}+D_{3}\left(2+C_{3} D_{4}\right) \int_{0}^{t}\left[\left\|\partial_{t}^{j-1} f(\tau)\right\|+|f(\tau)|_{j-2}\right] \mathrm{d} \tau\right) \\
& +D_{4}|f(t)|_{j-2} \text { for } t \in[0, T] . \tag{5.2}
\end{align*}
$$

Hence (5.1) implies (1.10).
For $j=2$, Theorem 1.1 is proved by Lemma 4.3. In order to prove Theorem 1.1 for $\mathrm{j} \geqslant 3$, we differentiate (1.1) with respect to $t$ (formally) and set $v:=u^{\prime}$. This yields

$$
\begin{align*}
& v^{\prime \prime}(t)+A(t) v(t)+B(t) v^{\prime}(t)=f^{\prime}(t)-A^{\prime}(t) u(t)-B^{\prime}(t) u^{\prime}(t) \text { for } t \in[0, T]  \tag{5.3a}\\
& v(t) \in \dot{H}^{m}(\Omega) \text { for } t \in[0, T]  \tag{5.3b}\\
& v(0)=u^{1}, \quad v^{\prime}(0)=u^{2}  \tag{5.3c}\\
& u(t)=u^{0}+\int_{0}^{t} v(\tau) \mathrm{d} \tau \quad \text { for } t \in[0, T] \tag{5.4}
\end{align*}
$$

with $u^{2}=f(0)-A(0) u^{0}-B(0) u^{1}$. On the other hand, if $(u, v) \in \tilde{\mathscr{C}}_{T}^{j+1} \times \mathscr{C}_{T}^{j}$ is a solution of (5.3), (5.4) (with $j \geqslant 2$ ), then $u \in \mathscr{C}_{T}^{j+1}$ and $u$ solves (1.1). In fact, it follows from $u^{0} \in \dot{H}^{m}(\Omega), v(t) \in \dot{H}^{m}(\Omega)$ for $t \in[0, T]$ and $(5.4)$ that $u(t) \in \dot{H}^{m}(\Omega)$ for $t \in[0, T]$. Furthermore (5.4) implies $u^{\prime}=v$ and $u(0)=u^{0}, u^{\prime}(0)=u^{1}$. Integrating (5.3) with respect to $t$ we obtain that (1.1) holds. Finally, it follows from

$$
\begin{equation*}
A(t) u(t)=f(t)-u^{\prime \prime}(t)-B(t) u^{\prime}(t) \text { for } t \in[0, T] \tag{5.5}
\end{equation*}
$$

by Lemma 3.1 that $u \in C\left([0, T], H^{(j+1) m}(\Omega)\right)$. This and $u \in \tilde{\mathscr{C}}_{T}^{j+1}$ imply $u \in \mathscr{C}_{T}^{j+1}$.
Let Theorem 1.1 be proved for $j=2, \ldots, J \leqslant k-1$. We prove the existence of a solution $(u, v) \in \tilde{\mathscr{C}}_{\boldsymbol{T}}^{J+1} \times \mathscr{C}_{T}^{J}$ of (5.3), (5.4) by the method of successive approximations
(as in [5]). Let $u_{0} \in \tilde{\mathscr{C}}_{T}^{J+1}$ such that $\partial_{t}^{v} u_{0}(0)=u^{v}$ for $v=0, \ldots, J$ (compare Lemma 8.8 of [9]); here $u^{v}$ is defined by (2.7) for $v \geqslant 2$. We define $u_{\mu}$ for $\mu \geqslant 1$ by

$$
\begin{align*}
& v_{\mu}^{\prime \prime}(t)+A(t) v_{\mu}(t)+B(t) v_{\mu}^{\prime}(t)=f^{\prime}(t)-A^{\prime}(t) u_{\mu-1}(t)-B^{\prime}(t) u_{\mu-1}^{\prime}(t) \text { for } t \in[0, T],  \tag{5.6a}\\
&  \tag{5.6b}\\
& \quad v_{\mu}(t) \in \dot{H}^{m}(\Omega) \quad \text { for } t \in[0, T],  \tag{5.6c}\\
& v_{\mu}(0)=u^{1}, \quad v_{\mu}^{\prime}(0)=u^{2},  \tag{5.7}\\
& \\
& u_{\mu}(t):=u^{0}+\int_{0}^{t} v_{\mu}(\tau) \mathrm{d} \tau \quad \text { for } t \in[0, T] .
\end{align*}
$$

Then $u_{\mu} \in \tilde{\mathscr{C}}_{r}^{J+1}$ for $\mu=1,2, \ldots$, as can be seen in the following way: assume that $u_{\mu-1} \in \tilde{\mathscr{C}}_{T}^{J+1}$ and $\partial_{t}^{v} u_{\mu-1}(0)=u^{v}$ for $v=0, \ldots, J$. Then $\left(u^{1}, u^{2}, f^{\prime}-A^{\prime} u_{\mu-1}-B^{\prime} u_{\mu-1}^{\prime}\right)$ satisfies the compatibility condition for (5.6) of order J. Corollary 8.3 of [9] yields

$$
\begin{equation*}
f^{\prime}-A^{\prime} u_{\mu-1}-B^{\prime} u_{\mu-1}^{\prime} \in \tilde{\mathscr{C}}_{T}^{J-1} \subset \mathscr{C}_{T}^{J-2} \cap C^{J-1}\left([0, T], L_{2}(\Omega)\right) . \tag{5.8}
\end{equation*}
$$

By the induction hypothesis we obtain the existence of a unique solution $v_{\mu} \in \mathscr{C}_{T}^{J}$ of (5.6), and (5.7) yields $u_{\mu} \in \tilde{\mathscr{C}}_{T}^{J+1}$. Finally it follows from (5.6) and (5.7) that $\partial_{T}^{v} u_{\mu}(0)=u^{v}$ for $v=0, \ldots, J$.

Consider $w_{\mu}:=u_{\mu+1}-u_{\mu}, \tilde{w}_{\mu}:=v_{\mu+1}-v_{\mu}$. It holds

$$
\begin{align*}
& \tilde{w}_{\mu}^{\prime \prime}(t)+A(t) \tilde{w}_{\mu}(t)+B(t) \tilde{w}_{\mu}^{\prime}(t)=-A^{\prime}(t) w_{\mu-1}(t)-B^{\prime}(t) w_{\mu-1}^{\prime}(t) \\
& \quad \text { for } t \in[0, T],  \tag{5.9a}\\
& \tilde{w}_{\mu}(t) \in \dot{H}^{m}(\Omega) \text { for } t \in[0, T],  \tag{5.9b}\\
& \partial_{t}^{\prime} \tilde{w}_{\mu}(0)=0 \quad \text { for } v=0,1, \ldots, J,  \tag{5.9c}\\
& w_{\mu}(t)=\int_{0}^{t} \tilde{w}_{\mu}(\tau) \mathrm{d} \tau \quad \text { for } t \in[0, T] \tag{5.10}
\end{align*}
$$

for $\mu \geqslant 1$. By Corollary 8.3 in [9] we obtain

$$
\begin{align*}
& \left|A^{\prime}(t) w_{\mu-1}(t)+B^{\prime}(t) w_{\mu-1}^{\prime}(t)\right|_{J-1} \leqslant d_{1}^{\prime}\left|w_{\mu-1}(t)\right|_{J+1},  \tag{5.11}\\
& \left|A^{\prime}(t) w_{\mu-1}(t)+B^{\prime}(t) w_{\mu-1}^{\prime}(t)\right|_{J-2} \leqslant d_{2}^{\prime}\left|w_{\mu-1}(t)\right|_{J} \tag{5.12}
\end{align*}
$$

for $t \in[0, T]$. Here and in the following we denote by $d_{1}^{\prime}, d_{2}^{\prime}, \ldots$, positive constants that may depend on $c_{1}, c_{2}, c_{3}$ (of Assumption 1.1), $J$ and on (1.11), (1.12). We apply the induction hypothesis and (5.2) to (5.9). This yields

$$
\begin{align*}
\left|\tilde{w}_{\mu}(t)\right|_{J} & \leqslant d_{3}^{\prime} \mathrm{e}^{d_{4}^{\prime} t} \int_{0}^{t}\left|w_{\mu-1}(\tau)\right|_{J+1} \mathrm{~d} \tau+d_{5}^{\prime}\left|w_{\mu-1}(t)\right|_{J} \\
& \leqslant\left(d_{3}^{\prime}+d_{5}^{\prime}\right) \mathrm{e}^{d_{j} \tau} \int_{0}^{t} \mid w_{\mu-1}(\tau) \tilde{J}_{J+1} \mathrm{~d} \tau \tag{5.13}
\end{align*}
$$

since $\partial_{t}^{v} w_{\mu}(0)=0$ for $v=0, \ldots, J$ and therefore

$$
\begin{equation*}
\left|w_{\mu-1}(t)\right|_{J} \leqslant \int_{0}^{t}\left|w_{\mu-1}^{\prime}(\tau)\right|_{J} \mathrm{~d} \tau \leqslant \int_{0}^{t}\left|w_{\mu-1}(\tau)\right|_{J+1} \mathrm{~d} \tau \tag{5.14}
\end{equation*}
$$

for $t \in[0, T]$.

From (5.10) and (5.13) we conclude that

$$
\begin{align*}
\left|w_{\mu}(t)\right|_{J+1} & =\left\|w_{\mu}(t)\right\|_{J_{m}}+\left|w_{\mu}^{\prime}(t)\right|_{J} \\
& \leqslant \int_{0}^{t}\left\|\tilde{w}_{\mu}(\tau)\right\|_{J m} \mathrm{~d} \tau+\left|\tilde{w}_{\mu}(t)\right|_{J} \\
& \leqslant\left(d_{3}^{\prime}+d_{5}^{\prime}\right)(1+T) \mathrm{e}^{d_{4} T} \int_{0}^{t}\left|w_{\mu-1}(\tau)\right|_{J+1} \mathrm{~d} \tau . \tag{5.15}
\end{align*}
$$

Let $K(T):=\left(d_{3}^{\prime}+d_{5}^{\prime}\right)(1+T) \mathrm{e}^{d_{i} T}$. By induction it follows that

$$
\begin{equation*}
\left|w_{\mu}(t) \tilde{J}_{J+1} \leqslant K(T)^{\mu} \frac{t^{\mu-1}}{(\mu-1)!} \int_{0}^{T}\right| w_{0}(\tau) \tilde{J}+1 \mathrm{~d} \tau \quad \text { for } t \in[0, T] . \tag{5.16}
\end{equation*}
$$

Note that $w_{\mu}=u_{\mu+1}-u_{\mu}, w_{\mu}^{\prime}=v_{\mu+1}-v_{\mu}$. Hence we obtain from (5.16) that

$$
\begin{align*}
\left|v_{\mu}(t)-v_{v}(t)\right|_{J} & \leqslant \mid u_{\mu}(t)-u_{v}(t) \tilde{J}_{J+1} \\
& \leqslant C \sum_{l=v}^{\mu-1} K(T)^{l} \frac{T^{l-1}}{(l-1)!} \quad \text { for } \mu>v, t \in[0, T] \tag{5.17}
\end{align*}
$$

with some $C>0$. Hence $\left\{\left(u_{\mu}, v_{\mu}\right)\right\}$ converges in $\tilde{\mathscr{C}}_{T}^{J+1} \times \mathscr{C}_{T}^{J}$. From (5.6), (5.7) it follows that the limit $(u, v) \in \tilde{\mathscr{C}}_{T}^{J+1} \times \mathscr{C}_{T}^{J}$ is a solution of (5.3), (5.4). Thus $u \in \mathscr{C}_{T}^{J+1}$ and $u$ is a solution of (1.1) by the considerations following (5.4).

It remains to prove (5.1) for $j=J+1$. In order to apply (5.1) with $j=J$ to (5.3) we note that

$$
\begin{align*}
&\left\|\partial_{t}^{J-1}\left[A^{\prime}(t) u(t)+B^{\prime}(t) u(t)\right]\right\|+\left|A^{\prime}(t) u(t)+B^{\prime}(t) u(t)\right|_{J-2} \\
& \leqslant\left|A^{\prime}(t) u(t)+B^{\prime}(t) u(t)\right|_{J-1} \leqslant d_{6}^{\prime}|u(t)|_{J+1},  \tag{5.18}\\
&\left|i^{\prime}(t) u(t)+B^{\prime}(t) u(t)\right|_{J-2} \leqslant d_{1}|u(t)|_{J} \\
& \leqslant d_{1}\left(|u(0)|_{J}+\int_{0}^{t} \mid u(\tau) \tilde{\left.\right|_{J+1}} \mathrm{~d} \tau\right) \tag{5.19}
\end{align*}
$$

by Corollary 8.3 of [9]; here $d_{1}$ depends on $J, c_{1}, c_{2}, c_{3}$ and on (1.11), (1.13). Note that $v=u^{\prime}$ in (5.3), (5.4). Applying (5.1) with $j=J$ to (5.3) and using (5.18), (5.19), we conclude that

$$
\begin{align*}
\left|u^{\prime}(t)\right|_{J} \leqslant & \left(D_{3}+D_{4} d_{1}\right)|u(0)|_{J}+D_{3} \int_{0}^{t}\left[\left\|\partial_{1}^{J} f(\tau)\right\|+\left|f^{\prime}(\tau)\right|_{J-2}\right] \mathrm{d} \tau \\
& +\left[D_{3}\left(C_{3}+d_{6}^{\prime}\right)+D_{4} d_{1}\right] \int_{0}^{t}|u(\tau)|_{J+1} \mathrm{~d} \tau+D_{4}\left|f^{\prime}(t)\right|_{J-2} \\
& \text { for } t \in[0, T] . \tag{5.20}
\end{align*}
$$

From (5.5) we obtain by Lemma 8.2 of [9] and by Lemma 3.1 that

$$
\begin{equation*}
\|u(t)\|_{(J+1) m} \leqslant d_{2}\left[\|f(t)\|_{(J-1) m}+\left|u^{\prime}(t)\right|_{J-1}+\|u(t)\|\right] \tag{5.21}
\end{equation*}
$$

for $t \in[0, T]$, where $d_{2}>0$ depends on $J, c_{1}, c_{2}, c_{3}$ and (1.11), (1.13). Combining (5.1) (with $j=J$ ), (5.20) and (5.21) we conclude that (5.1) holds with $j=J+1$.

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