# ON APPROXIMATION OF INVERSE PROBLEMS FOR ABSTRACT HYPERBOLIC EQUATIONS 

Dmitry Orlovsky ${ }^{1}$, Sergey Piskarev ${ }^{2, *}$ and Renato Spigler ${ }^{3}$<br>Dedicated to the Memory of Professor Sen-Yen Shaw


#### Abstract

This paper is devoted to the numerical analysis of inverse problems for abstract hyperbolic differential equations. The presentation exploits a general approximation scheme and is based on $C_{0}$-cosine and $C_{0}$-semigroup theory within a functional analysis approach. We consider both discretizations in space as well as in time. The discretization in time is considered under the Krein-Fattorini conditions.


## 1. Introduction

Let $B(E)$ denote the Banach algebra of all linear bounded operators on a complex Banach space $E$. The set of all linear closed densely defined operators in $E$ will be denoted by $\mathcal{C}(E)$.

Let us examine the inverse problem in $E$ consisting of the search for a function $u(\cdot) \in C^{2}([0 ; T] ; E)$ and an element $d \in E$ from the equations

$$
\begin{equation*}
u^{\prime \prime}(t)=A u(t)+\Phi(t) d, \quad 0 \leq t \leq T \tag{1.1a}
\end{equation*}
$$

$$
\begin{array}{r}
u(0)=u^{0}, u^{\prime}(0)=u^{1}, \\
u(T)=u^{T} \tag{1.1c}
\end{array}
$$

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*Corresponding author.
where $A \in \mathcal{C}(E), \Phi(\cdot) \in C^{2}([0 ; T] ; B(E))$ and the elements $u^{0}$, $u^{1}$, $u^{T} \in E$ are given. The cases of parabolic and elliptic equations were considered in [10, 14]. Here we assume that the abstract differential equation in (1.1a) is of the hyperbolic type. This means that the operator $A$ generates a $C_{0}$-cosine operatorfunction $C(\cdot, A)$. Recall that a $C_{0}$-cosine operator-function is used to represent a solution of the abstract Cauchy problem

$$
\left\{\begin{array}{l}
u^{\prime \prime}(t)=A u(t)+f(t), \quad 0 \leq t \leq T  \tag{1.2}\\
u(0)=u^{0}, u^{\prime}(0)=u^{1}
\end{array}\right.
$$

Definition 1.1. A function $u(\cdot)$ is called a classical solution of problem (1.2) if $u(\cdot)$ is twice continuously differentiable, $u(t) \in D(A)$ for all $t \in[0, T]$, and $u(\cdot)$ satisfies the relations in (1.2).

We denote by $\sigma(B)$ the spectrum of the operator $B$, by $\rho(B)$ the resolvent set of $B$.

Proposition 1.1. $[6,18]$. The operator $A$ generates a $C_{0}$-cosine operatorfunction if and only if there are constants $M$ and $\omega$ such that for each $\lambda$ with $R e \lambda>\omega$ the value $\lambda^{2}$ is contained in the resolvent set $\rho(A)$ of the operator $A$ and for the same value $\lambda$ the following estimate holds :

$$
\begin{equation*}
\left\|\frac{d^{n}}{d \lambda^{n}}\left(\lambda R\left(\lambda^{2}, A\right)\right)\right\| \leqslant \frac{M n!}{(\lambda-\omega)^{n+1}}, \quad n=0,1,2, \ldots \tag{1.3}
\end{equation*}
$$

For any strongly continuous $C_{0}$-cosine operator-function $C(\cdot, A)$ the following inequality holds

$$
\begin{equation*}
\|C(t, A)\| \leqslant M \exp (\omega|t|), \quad t \in \mathbb{R} \tag{1.4}
\end{equation*}
$$

In this case we will write $A \in C(M, \omega)$. Furthermore, we introduce the Kisynski space [7]

$$
E^{1}=\left\{x \in E: C(t, A) x \in C^{1}(I R ; E)\right\}
$$

with the norm $\|x\|_{E^{1}}=\|x\|+\sup _{0<t \leq 1}\left\|C^{\prime}(t, A) x\right\|$. This is a Banach space with the norm $\|\cdot\|_{E^{1}}$.

If the operator $A$ generates a $C_{0}$-cosine operator-function $C(\cdot, A)$ and $f(\cdot) \in$ $C([0, T] ; E)$, then for any classical solution of (1.2)

$$
\begin{equation*}
u(t)=C(t, A) u^{0}+S(t, A) u^{1}+\int_{0}^{t} S(t-s, A) f(s) d s, \quad t \in[0, T] \tag{1.5}
\end{equation*}
$$

where $S(t, A):=\int_{0}^{t} C(s, A) d s$ is the corresponding $C_{0}$-sine operator-function. The formula (1.5) is the analog of the variation-of-constants formula for $C_{0}$-semigroups.

As in the case of $C_{0}$-semigroups of operators, the function $u(\cdot)$ given by (1.5) is not a classical solution, in general, since it may be not twice continuously differentiable.

Remark 1.1. According to (1.5), in general, the problem (1.1) is ill-posed. This happens, for instance, if resolvent $(\lambda I-A)^{-1}$ is compact for some $\lambda$. Indeed, in this case the integral operator $\int_{0}^{T} S(T-s, A) \Phi(s) d s$ is compact and thus the equation

$$
\int_{0}^{T} S(T-s, A) \Phi(s) d s d=u(T)-C(T, A) u^{0}-S(T, A) u^{1}
$$

in the space $E$ leads to an ill-posed problem. However, if we consider the operator $\int_{0}^{T} S(T-s, A) \Phi(s) d s$ as the operator from $E$ to $\mathfrak{D}(A)$, where $\mathfrak{D}(A)$ equiped with the norm $\|x\|_{\mathfrak{D}_{(A)}}=\|x\|+\|A x\|$, then the operator $\int_{0}^{T} S(T-$ $s, A) \Phi(s) d s: E \rightarrow \mathfrak{D}(A)$ has a chance to be not compact. Therefore, in case of $u(T), C(T, A) u^{0}, S(T, A) u^{1} \in D(A)$ one can play with formula (4.3) to get a Fredhom equation of the second kind, which is a well-posed problem.

Definition 1.2. The function $u(\cdot) \in C([0, T) ; E)$ given by (1.5) is called a mild solution of problem (1.2).

Proposition 1.2. [6]. Let the operator $A$ be a generator of a $C_{0}$-cosine operatorfunction $C(\cdot, A)$, and let either
(i) $f(\cdot), A f(\cdot) \in C([0, T) ; E)$ and $f(t) \in D(A)$ for $t \in[0, T]$
or
(ii) $f(\cdot) \in C^{1}([0, T] ; E)$.

Then the function $u(\cdot)$ given by (1.5) with $u^{0} \in D(A)$ and $u^{1} \in E^{1}$ is a classical solution of problem (1.2) on $[0, T]$.

If we differentiate both sides of (1.5), we get

$$
u^{\prime}(t)=S(t, A) A u^{0}+C(t, A) u^{1}+\int_{0}^{t} C(t-s, A) f(s) d s
$$

Integrating by parts we obtain an alternative form for the first derivative

$$
\begin{equation*}
u^{\prime}(t)=S(t, A)\left(A u^{0}+f(0)\right)+C(t, A) u^{1}+\int_{0}^{t} S(t-s, A) f^{\prime}(s) d s \tag{1.6}
\end{equation*}
$$

We have to note here that one cannot expect maximal regularity for the problem (1.2), see [4], so in order to get a classical solution the differentiability of $f(\cdot)$ is almost necessary condition. Let us write $v(t)=u^{\prime}(t), v^{0}=u^{1}, v^{1}=A u_{0}+f(0)$, $f_{1}(t)=f^{\prime}(t)$. Then last formula in (1.6) can be written as formula (1.5):

$$
v(t)=C(t, A) v^{0}+S(t, A) v^{1}+\int_{0}^{t} S(t-s, A) f_{1}(s) d s
$$

Proposition 1.2 yields the conditions under which the function $v(\cdot)$ is a classical solution (in particular is twice continuously differentiable) of the problem

$$
\left\{\begin{array}{l}
v^{\prime \prime}(t)=A v(t)+f_{1}(t), \quad 0 \leq t \leq T, \\
v(0)=v^{0}, v^{\prime}(0)=v^{1}
\end{array}\right.
$$

These conditions are that $v^{0} \in D(A), v^{1} \in E^{1}, f_{1}(\cdot) \in C^{1}([0, T] ; E)$, i. e. $u^{0}, u^{1} \in$ $D(A), A u^{0}+f(0) \in E^{1}, f(\cdot) \in C^{2}([0, T] ; E)$. It follows from these conditions that $v(\cdot) \in C^{2}([0, T] ; E)$, i. e. $u(\cdot) \in C^{3}([0, T] ; E)$.

Following the same procedure it is possible to find some sufficient conditions under which the solution of the Cauchy problem becomes as smooth as we like. Set $w(t)=v^{\prime}(t)$. Then, one can write

$$
\begin{equation*}
w(t)=C(t, A) w^{0}+S(t, A) w^{1}+\int_{0}^{t} S(t-s, A) f_{2}(s) d s \tag{1.7}
\end{equation*}
$$

where $w^{0}=v^{1}, w^{1}=A v_{0}+f_{1}(0), f_{2}(t)=f_{1}^{\prime}(t)$.
If $w^{0} \in D(A), w^{1} \in E^{1}$ and $f_{2}(\cdot) \in C^{1}([0, T] ; E)$, then $w(\cdot) \in C^{2}([0, T] ; E)$, i. e. $u(\cdot) \in C^{4}([0, T] ; E)$. This leads us to the next proposition:

Proposition 1.3. Assume that the operator $A \in C(M, \omega)$ and $u^{0}, u^{1} \in D\left(A^{2}\right)$. Suppose also that the following conditions hold
(i) $f(\cdot) \in C^{3}([0, T] ; E)$,
(ii) $A u^{0}+f(0) \in D(A), A u^{1}+f^{\prime}(0) \in E^{1}$.

Then the function $u(\cdot)$ from (1.5) belongs to $C^{4}([0, T] ; E)$. Conversely. Assume that the function $u(\cdot)$ defined by (1.5) belongs to $C^{4}([0, T] ; E)$, i.e. $u(\cdot) \in$ $C^{4}([0, T] ; E)$, and $f(\cdot) \in C^{3}([0, T] ; E)$ with $f(0)=0$. Then $f^{\prime}(0) \in E^{1}$ and so $A S(t, A) f^{\prime}(0) \in C([0, T] ; E)$.

Proof. We prove just second part of Proposition. As it can be seen from (1.7) we have

$$
\begin{gathered}
u^{\prime \prime \prime}(t)=A^{2} S(t, A) u^{0}+C(t, A)\left(A u^{1}+f^{\prime}(0)\right)+\int_{0}^{t} C(t-s, A) f^{\prime \prime}(s) d s \\
=S(t, A) A^{2} u^{0}+C(t, A)\left(A u^{1}+f^{\prime}(0)\right)+S(t, A) f^{\prime \prime}(0)+\int_{0}^{t} S(t-s, A) f^{\prime \prime \prime}(s) d s
\end{gathered}
$$

Now,

$$
\begin{align*}
& u^{\prime \prime \prime \prime}(t)=C(t, A) A^{2} u^{0}+S(t, A) A^{2} u^{1}+A S(t, A) f^{\prime}(0)+C(t, A) f^{\prime \prime}(0) \\
& +\int_{0}^{t} C(t-s, A) f^{\prime \prime \prime}(s) d s \tag{1.8}
\end{align*}
$$

hence the function $A S(t, A) f^{\prime}(0) \in C([0, T] ; E)$.
Let consider the homogenous uniformly well-posed Cauchy problem

$$
\begin{equation*}
u^{\prime \prime}(t)=A u(t), \quad t \in \mathbb{R} ; \quad u(0)=u^{0}, \quad u^{\prime}(0)=u^{1} . \tag{1.9}
\end{equation*}
$$

Define the matrix operator $\mathcal{A}:=\left(\begin{array}{cc}0 & I \\ A & 0\end{array}\right): E^{1} \times E \rightarrow E^{1} \times E$ acting on the element $(x, y) \in E^{1} \times E$ according to the formula $\mathcal{A}(x, y)=(y, A x)$. This operator has the domain $D(\mathcal{A}):=D(A) \times E^{1}$.

Let the uniformly well-posed problem (1.9) have the form

$$
\begin{equation*}
u^{\prime \prime}(t)=\mathfrak{B}^{2} u(t), \quad t \in \mathbb{R} ; \quad u(0)=u^{0}, \quad u^{\prime}(0)=u^{1}, \tag{1.10}
\end{equation*}
$$

where $\mathfrak{B} \in \mathcal{C}(E)$. Then
Definition 1.3. We say that a solution $u(\cdot)$ of problem (1.10) satisfies Condition $(K)$ if

$$
u^{\prime}(\cdot) \in C([0, T] ; \mathfrak{D}(\mathfrak{B}))
$$

Proposition 1.4. [23]. Problem (1.10) has a unique solution satisfying Condition (K) iff the following Cauchy problem:

$$
\binom{u}{v}^{\prime}(t)=\left(\begin{array}{cc}
0 & \mathfrak{B}  \tag{1.11}\\
\mathfrak{B} & 0
\end{array}\right)\binom{u}{v}(t), \quad t \in \mathbb{R}, \quad\binom{u}{v}(0)=\binom{u_{0}}{v_{0}}
$$

is uniformly well posed on the space $E \times E$.
The following Condition ( F ) is analog to Condition (K), which allows to simplify the study of problem (1.9) by using $C_{0}$-semigroups.

Definition 1.4. We say that a $C_{0}$-cosine operator-valued function $C(\cdot, A)$ satisfies Condition $(F)$ if the following conditions hold:
(i) there exists $\mathfrak{B} \in \mathcal{C}(E)$ such that $\mathfrak{B}^{2}=A$, and $\mathfrak{B}$ commutes with any operator from $B(E)$ commuting with $A$;
(ii) the operator $S(t, A)$ maps $E$ into $D(\mathfrak{B})$ for any $t \in \mathbb{R}$;
(iii) the function $\mathfrak{B} S(t, A) x$ is continuous in $t \in \mathbb{R}$ for every fixed $x \in E$.

Proposition 1.5. [6]. Under Condition $(F)$, for each $t \in \mathbb{R}$, we have $\mathfrak{B} S(t, A) \in$ $B(E)$ and $\mathfrak{D}(\mathfrak{B}) \subseteq E^{1}$.

Proposition 1.6. [6]. Pairs of a Banach space $E$ and a $C_{0}$-cosine operatorfunction $C(\cdot, A)$ ( also uniformly bounded) such that Condition $(F)$ does not hold do exist.

We have to note that if $0 \in \rho(A)$, then conditions $(\mathrm{K})$ and $(\mathrm{F})$ are equivalent.
Proposition 1.7. [20]. Let $E$ be a Hilbert space, and let the operator $A$ be self-adjoint and negative-definite. Then $A \in \mathcal{C}(M ; \omega)$, condition $(F)$ is satisfied and the corresponding space $E^{1}$ coincides with $\mathfrak{D}\left((-A)^{1 / 2}\right)$.

Theorem 1.1. [19]. Let $A$ and $\mathfrak{B}$ be operators satisfying condition (i) of Definition 1.4, and let $0 \in \rho(\mathfrak{B})$. The following conditions are equivalent:
(i) the $C_{0}$-cosine operator-function $C(\cdot, A)$ satisfies Condition $(F)$;
(ii) the operator $\mathfrak{B}$ generates a $C_{0-\text { group }} \exp (\cdot \mathfrak{B})$ on $E$;
(iii) the operator $\left(\begin{array}{cc}0 & \mathfrak{B} \\ \mathfrak{B} & 0\end{array}\right)$ with the domain $D(A) \times D(\mathfrak{B})$ generates a $C_{0}$ group on $E \times E$;
(iv) the operator $\mathcal{A}:=\left(\begin{array}{ll}0 & I \\ A & 0\end{array}\right)$ with the domain $D(A) \times D(\mathfrak{B})$ generates a $C_{0}$-group $\exp (\cdot \mathcal{A})$ on $\mathfrak{D}(\mathfrak{B}) \times E$, where $\mathfrak{D}(\mathfrak{B})$ is the Banach space of elements $D(\mathfrak{B})$ endowed with the graph norm;
(v) the embedding $D(\mathfrak{B}) \subseteq E^{1}$ holds;
(vi) $\mathfrak{D}(\mathfrak{B})=E^{1}$.

Proposition 1.8. [19]. Under the conditions of Theorem 1.1, for $t \in \mathbb{R}$, we have
$(i) \exp (t \mathfrak{B})=C(t, A)+\mathfrak{B} S(t, A), C(t, A)=(\exp (t \mathfrak{B})+\exp (-t \mathfrak{B})) / 2$;
(ii) $\exp (t \mathcal{A})=\left(\begin{array}{cc}\mathfrak{B}^{-1} & 0 \\ 0 & I\end{array}\right) \exp \left(t\left(\begin{array}{cc}0 & \mathfrak{B} \\ \mathfrak{B} & 0\end{array}\right)\right)\left(\begin{array}{cc}\mathfrak{B} & 0 \\ 0 & I\end{array}\right)$.

The analog of Proposition 1.2 is given in
Theorem 1.2. [8]. Let the operator $\mathfrak{B}=\sqrt{A}$ in problem (1.2) have a bounded inverse $\mathfrak{B}^{-1} \in B(E)$ and be a generator of a $C_{0}$-group. Assume also that the function $f(\cdot)$ have one of the following properties:
(i) $f(\cdot) \in C^{1}([0, T) ; E)$;
(ii) $\mathfrak{B} f(\cdot) \in C([0, T) ; E)$.

Then for any $u^{0} \in D(A)$ and $u^{1} \in D(\mathfrak{B})$, there exists a unique classical solution of problem (1.2) given by formula (1.5) in the form

$$
\begin{align*}
u(t)= & \frac{1}{2}(\exp (t \mathfrak{B})+\exp (-t \mathfrak{B})) u^{0}+\frac{1}{2}(\exp (t \mathfrak{B})-\exp (-t \mathfrak{B})) \mathfrak{B}^{-1} u^{1} \\
& +\frac{1}{2} \int_{0}^{t}(\exp ((t-s) \mathfrak{B})-\exp (-(t-s) \mathfrak{B})) \mathfrak{B}^{-1} f(s) d s, \quad t \in[0, T] . \tag{1.12}
\end{align*}
$$

## 2. A General Approximation Scheme

A general approximation scheme, due to [21], [22], can be described in the following way. Let $E_{n}$ and $E$ be Banach spaces and $\left\{p_{n}\right\}$ be a sequence of linear bounded operators $p_{n}: E \rightarrow E_{n}, p_{n} \in B\left(E, E_{n}\right), n \in I N=\{1,2, \cdots\}$, with the property:

$$
\left\|p_{n} x\right\|_{E_{n}} \rightarrow\|x\|_{E} \text { as } n \rightarrow \infty \text { for any } x \in E .
$$

Definition 2.1. The sequence of elements $\left\{x_{n}\right\}, x_{n} \in E_{n}, n \in I N$, is said to be $\mathcal{P}$-convergent to $x \in E$ iff $\left\|x_{n}-p_{n} x\right\|_{E_{n}} \rightarrow 0$ as $n \rightarrow \infty$ and we write this $x_{n} \xrightarrow{\mathcal{P}} x$.

Definition 2.2. The sequence of elements $\left\{x_{n}\right\}, x_{n} \in E_{n}, n \in I N$, is said to be $\mathcal{P}$-compact if for any subset of interges $I N^{\prime} \subseteq I N$ there exist a subset of interges $I N^{\prime \prime} \subseteq I N^{\prime}$ and $x \in E$ such that $x_{n} \xrightarrow{\mathcal{P}} x$, as $n \rightarrow \infty$ in $I^{\prime \prime}$.

Definition 2.3. The sequence of linear bounded operators $B_{n} \in B\left(E_{n}\right), n \in I N$, is said to be $\mathcal{P} \mathcal{P}$-convergent to the bounded linear operator $B \in B(E)$ if for every $x \in E$ and for every sequence $\left\{x_{n}\right\}, x_{n} \in E_{n}, n \in I N$, such that $x_{n} \xrightarrow{\mathcal{P}} x$ one has $B_{n} x_{n} \xrightarrow{\mathcal{P}} B x$. We write this as $B_{n} \xrightarrow{\mathcal{P P}} B$.

For general examples of notions of $\mathcal{P}$-convergence see [21].
Remark 2.1. If we set $E_{n}=E$ and $p_{n}=I$ for every $n \in I N$, where $I$ is the identity operator on $E$, then Definition 2.1 leads to the usual pointwise convergence of bounded linear operators which we denote by $B_{n} \rightarrow B$.

In case of operators which have a compact resolvent it is natural to consider approximating operators which "preserve" the property of compactness. Hence,

Definition 2.4. A sequence of operators $\left\{B_{n}\right\}, B_{n}: E_{n} \rightarrow E_{n}, n \in I N$, converges compactly to an operator $B: E \rightarrow E$ if $B_{n} \xrightarrow{\mathcal{P P}} B$ and the following compactness condition holds:

$$
\left\|x_{n}\right\|_{E_{n}}=O(1) \Longrightarrow\left\{B_{n} x_{n}\right\} \text { is } \mathcal{P} \text {-compact. }
$$

Definition 2.5. The region of stability $\Delta_{s}=\Delta_{s}\left(\left\{A_{n}\right\}\right), A_{n} \in \mathcal{C}\left(B_{n}\right)$, is defined as the set of all $\lambda \in \mathbb{C}$ such that $\lambda \in \rho\left(A_{n}\right)$ for almost all $n$ and such that the sequence $\left\{\left\|\left(\lambda I_{n}-A_{n}\right)^{-1}\right\|\right\}_{n \in N}$ is bounded for almost all $n$. The region of convergence $\Delta_{c}=\Delta_{c}\left(\left\{A_{n}\right\}\right), A_{n} \in \mathcal{C}\left(E_{n}\right)$, is defined as the set of all $\lambda \in \mathbb{C}$ such that $\lambda \in \Delta_{s}\left(\left\{A_{n}\right\}\right)$ and such that the sequence of operators $\left\{\left(\lambda I_{n}-A_{n}\right)^{-1}\right\}_{n \in N}$ is $\mathcal{P} \mathcal{P}$-convergent to some operator $S(\lambda) \in B(E)$.

Definition 2.6. The region of compact convergence of resolvents, $\Delta_{c c}=\Delta_{c c}$ $\left(A_{n}, A\right)$, where $A_{n} \in \mathcal{C}\left(E_{n}\right)$ and $A \in \mathcal{C}(E)$ is defined as the set of all $\lambda \in \Delta_{c} \cap \rho(A)$ such that $\left(\lambda I_{n}-A_{n}\right)^{-1} \xrightarrow{\mathcal{P} \mathcal{P}}(\lambda I-A)^{-1}$ compactly.

In the case of unbounded operators (recall that in general infinitesimal generators are unbounded), we consider the notion of compatibility.

Definition 2.7. The sequence of closed linear operators $\left\{A_{n}\right\}, A_{n} \in \mathcal{C}\left(E_{n}\right), n \in$ $I N$, is said to be compatible with a linear closed operator $A \in \mathcal{C}(E)$ iff for each $x \in D(A)$ there is a sequence $\left\{x_{n}\right\}, x_{n} \in D\left(A_{n}\right) \subseteq E_{n}, n \in I N$, such that $x_{n} \xrightarrow{\mathcal{P}} x$ and $A_{n} x_{n} \xrightarrow{\mathcal{P}} A x$. We write this as $\left(A_{n}, A\right)$ are compatible.

Usually, in practice, the Banach spaces $E_{n}$ are finite-dimensional, although, in general, e.g. in the case of a closed operator $A$, we have $\operatorname{dim} E_{n} \rightarrow \infty$ and $\left\|A_{n}\right\|_{B\left(E_{n}\right)} \rightarrow \infty$ as $n \rightarrow \infty$.

Definition 2.8. A sequence of operators $\left\{B_{n}\right\}, B_{n} \in B\left(E_{n}\right), n \in I N$, is said to be stably convergent to an operator $B \in B(E)$ iff $B_{n} \xrightarrow{\mathcal{P P}} B$ and $\left\|B_{n}^{-1}\right\|_{B\left(E_{n}\right)}=$ $O(1), n \rightarrow \infty$. We will write this as: $B_{n} \xrightarrow{\mathcal{P} \mathcal{P}} B$ stably.

Definition 2.9. A sequence of operators $\left\{B_{n}\right\}, B_{n} \in B\left(E_{n}\right)$, is called regularly convergent to the operator $B \in B(E)$ iff $B_{n} \xrightarrow{\mathcal{P} \mathcal{P}} B$ and the following implication holds:

$$
\left\|x_{n}\right\|_{E_{n}}=O(1) \&\left\{B_{n} x_{n}\right\} \text { is } P \text {-compact } \Longrightarrow\left\{x_{n}\right\} \text { is } P \text {-compact. }
$$

We write this as: $B_{n} \xrightarrow{\mathcal{P} \mathcal{P}} B$ regularly.
Theorem 2.1. [22]. Let $C_{n}, S_{n} \in B\left(E_{n}\right), C, S \in B(E)$ and $\mathcal{R}(S)=E$. Assume also that $C_{n} \xrightarrow{\mathcal{P} \mathcal{P}} C$ compactly and $S_{n} \xrightarrow{\mathcal{P} \mathcal{P}} S$ stably. Then $S_{n}+C_{n} \xrightarrow{\mathcal{P} \mathcal{P}} S+C$ converges regularly.

Theorem 2.2. [22]. For $Q_{n} \in B\left(E_{n}\right)$ and $Q \in B(E)$ the following conditions are equivalent:
(i) $Q_{n} \xrightarrow{\mathcal{P} \mathcal{P}} Q$ regularly, $Q_{n}$ are Fredholm operators of index 0 and $\mathcal{N}(Q)=\{0\}$;
(ii) $Q_{n} \xrightarrow{\mathcal{P} \mathcal{P}} Q$ stably and $\mathcal{R}(Q)=E$;
(iii) $Q_{n} \xrightarrow{\mathcal{P} \mathcal{P}} Q$ stably and regularly;
(iv) if one of conditions (i)-(iii) holds, then there exist $Q_{n}^{-1} \in B\left(E_{n}\right), Q^{-1} \in$ $B(E)$, and $Q_{n}^{-1} \xrightarrow{\mathcal{P P}} Q^{-1}$ regularly and stably.

Theorem 2.3. [5]. Let the operators $A$ and $A_{n}$ generate $C_{0}$-semigroups. The following conditions $(A)$ and $(B)$ are equivalent to condition $(C)$.
(A) Consistency. There exists $\lambda \in \rho(A) \cap \cap_{n} \rho\left(A_{n}\right)$ such that the resolvents converge
$\left(\lambda I_{n}-A_{n}\right)^{-1} \xrightarrow{\mathcal{P P}}(\lambda I-A)^{-1} ;$
(B) Stability. There are some constants $M_{1} \geq 1$ and $\omega_{1} \in \mathbb{R}$ independent of $n$ such that for any $t \geq 0$

$$
\left\|\exp \left(t A_{n}\right)\right\| \leq M_{1} e^{\omega t} \text { for all } n \in I N
$$

(C) Convergence. For any finite $T>0$ we have

$$
\max _{t \in[0, T])}\left\|\exp \left(t A_{n}\right) u_{n}^{0}-p_{n} \exp (t A) u^{0}\right\| \rightarrow 0
$$

as $n \rightarrow \infty$ for any $u^{0} \in E$, whenever $u_{n}^{0} \xrightarrow{\mathcal{P}} u^{0}$.
Usually it is assumed that conditions (A) and (B) for the corresponding $C_{0}{ }^{-}$ semigroup case are satisfied without any loss of generality whatever process of discretization in time is considered. We denote by $T_{n}(\cdot)$ a family of discrete semigroups $T_{n}(t)=T_{n}\left(\tau_{n}\right)^{k_{n}}$, where $k_{n}=\left[\frac{t}{\tau_{n}}\right]$, as $\tau_{n} \rightarrow 0, n \rightarrow \infty$, see [13]. The generator of discrete semigroup is defined by $\breve{A}_{n}=\frac{1}{\tau_{n}}\left(T_{n}\left(\tau_{n}\right)-I_{n}\right) \in B\left(E_{n}\right)$ and hence $T_{n}(t)=\left(I_{n}+\tau_{n} \breve{A}_{n}\right)^{k_{n}}$, where $t=k_{n} \tau_{n}$.

Theorem 2.4. (Theorem ABC-discr, [13]). The following conditions $(A)$ and $\left(B^{\prime}\right)$ are equivalent to condition $\left(C^{\prime}\right)$.
(A) Consistency. There exists $\lambda \in \rho(A) \cap \cap_{n} \rho\left(\breve{A}_{n}\right)$ such that the resolvents converge

$$
\left(\lambda I_{n}-\breve{A}_{n}\right)^{-1} \xrightarrow{\mathcal{P} \mathcal{P}}(\lambda I-A)^{-1} ;
$$

( $B^{\prime}$ ) Stability. There are some constants $M \geq 1$ and $\omega_{1} \in \mathbb{R}$ such that

$$
\left\|T_{n}(t)\right\| \leq M \exp \left(\omega_{1} t\right) \text { for } t \in \overline{\mathbb{R}}_{+}=[0, \infty), n \in I N
$$

(C') Convergence. For any finite $T>0$ one has

$$
\max _{t \in[0, T]}\left\|T_{n}(t) u_{n}^{0}-p_{n} \exp (t A) u^{0}\right\| \rightarrow 0
$$

as $n \rightarrow \infty$, whenever $u_{n}^{0} \xrightarrow{\mathcal{P}} u^{0}$ for any $u^{0} \in E, u_{n}^{0} \in E_{n}$.

Theorem 2.5. [13]. Assume that $A \in \mathcal{C}(E), A_{n} \in \mathcal{C}\left(E_{n}\right)$ and let $A, A_{n}$ generate $C_{0}$-semigroups. Assume also that conditions $(A)$ and $(B)$ of Theorem 2.3 hold. Then, the implicit difference scheme

$$
\begin{equation*}
\frac{\bar{U}_{n}\left(t+\tau_{n}\right)-\bar{U}_{n}(t)}{\tau_{n}}=A_{n} \bar{U}_{n}(t+\tau), \bar{U}_{n}(0)=u_{n}^{0} \tag{2.1}
\end{equation*}
$$

is stable, i.e. $\left\|\left(I_{n}-\tau_{n} A_{n}\right)^{-k_{n}}\right\| \leq M_{1} e^{\omega_{1} t}, t=k_{n} \tau_{n} \in \bar{R}_{+}$, and gives an approximation to the $\exp (t A) u_{n}^{0}$, i.e. $\bar{U}_{n}(t) \equiv\left(I_{n}-\tau_{n} A_{n}\right)^{-k_{n}} u_{n}^{0} \xrightarrow{\mathcal{P}} \exp (t A) u_{n}^{0}$ uniformly with respect to $t=k_{n} \tau_{n} \in[0, T]$ as $u_{n}^{0} \xrightarrow{\mathcal{P}} u^{0}, n \rightarrow \infty, k_{n} \rightarrow \infty$, $\tau_{n} \rightarrow 0$.

For $C_{0}$-cosine operator-functions the following ABC Theorem holds:
Theorem 2.6. [13]. Let the operators $A$ and $A_{n}$ be generators of $C_{0}$-cosine operator-functions. Then, the following conditions $(A)$ and $\left(B^{\prime \prime}\right)$ are equivalent to condition $\left(C^{\prime \prime}\right)$ :
(A) Compatability. There exists $\lambda \in \rho(A) \cap \cap_{n} \rho\left(A_{n}\right)$ such that the resolvents converge

$$
\left(\lambda I_{n}-A_{n}\right)^{-1} \xrightarrow{\mathcal{P} \mathcal{P}}(\lambda I-A)^{-1}
$$

( $B$ ") Stability. There are some constants $M_{3} \geq 1$ and $\omega_{3} \geq 0$ such that

$$
\left\|C\left(t, A_{n}\right)\right\| \leq M_{3} e^{\omega_{3} t}, \quad t \geq 0, \quad n \in I N
$$

( $C$ ") Convergence. For any finite $T>0$ one has

$$
\max _{t \in[0, T]}\left\|C\left(t, A_{n}\right) u_{n}^{0}-p_{n} C(t, A) u^{0}\right\| \rightarrow 0
$$

as $n \rightarrow \infty$ for any $u^{0} \in E$, whenever $u_{n}^{0} \xrightarrow{\mathcal{P}} u^{0}$.

## 3. Discretizing in Space and Time

The semidiscrete approximation of (1.2) leads to the following Cauchy problems in the Banach spaces $E_{n}$ :

$$
\begin{align*}
u_{n}^{\prime \prime}(t) & =A_{n} u_{n}(t)+f_{n}(t), t \in[0, T]  \tag{3.1}\\
u_{n}(0) & =u_{n}^{0}, u_{n}^{\prime}(0)=u_{n}^{1}
\end{align*}
$$

with operators $A_{n}$, which generate $C_{0}$-cosine operator-functions, the operators $A_{n}$ and $A$ are compatible, $u_{n}^{0} \xrightarrow{\mathcal{P}} u^{0}, u_{n}^{1} \xrightarrow{\mathcal{P}} u^{1}$ and $f_{n}(\cdot) \xrightarrow{\mathcal{P}} f(\cdot)$ in an appropriate sense. It is natural to assume that conditions $(\mathrm{A})$ and $\left(B^{\prime \prime}\right)$ of Theorem 2.6 for $C_{0}$-cosine operator-functions are satisfied.

The discretization of (3.1) in the time variable has been considered in many papers [ $1,11,17]$. One of the simplest difference scheme is

$$
\begin{align*}
& \frac{U_{n}^{k+1}-2 U_{n}^{k}+U_{n}^{k-1}}{\tau_{n}^{2}}  \tag{3.2}\\
= & A_{n} U_{n}^{k+1}+\varphi_{n}^{k}, k \in\left\{1, \ldots,\left[\frac{T}{\tau_{n}}\right]\right\}, U_{n}^{0}=u_{n}^{0}, U_{n}^{1}=u_{n}^{0}+\tau_{n} u_{n}^{1},
\end{align*}
$$

where, for instance if $f_{n}(\cdot) \in C\left([0, T] ; E_{n}\right)$, one can set $\varphi_{n}^{k}=f_{n}\left(k \tau_{n}\right), k \in$ $\{1, \ldots, K\}, K=\left[\frac{T}{\tau_{n}}\right]$, and in case that $f_{n}(\cdot) \in L^{1}\left([0, T] ; E_{n}\right)$, one can set

$$
\varphi_{n}^{k}=\frac{1}{\tau_{n}} \int_{t_{k-1}}^{t_{k}} f_{n}(s) d s, t_{k}=k \tau_{n}, k \in\{1, \ldots, K\}
$$

The solution to problem (3.2) is given by the formula [16]:

$$
\begin{equation*}
U_{n}^{k}=C_{k}^{(n)} U_{n}^{0}+S_{k}^{(n)} U_{n}^{1}+\tau_{n}^{2} R_{n} \sum_{j=2}^{k} S_{k+1-j}^{(n)} \varphi_{n}^{j-1}, \tag{3.3}
\end{equation*}
$$

where $k \geqslant 2$. Indeed, in order to solve the homogeneous equations associated to (3.2), i.e.

$$
\begin{equation*}
U_{n}^{k+1}-2\left(I_{n}-\tau_{n}^{2} A_{n}\right)^{-1} U_{n}^{k}+\left(I_{n}-\tau_{n}^{2} A_{n}\right)^{-1} U_{n}^{k-1}=0, \tag{3.4}
\end{equation*}
$$

we consider the discrete operator-functions defined by the recurrent relations

$$
\begin{align*}
& C_{k+1}^{(n)}=R_{n}\left(2 C_{k}^{(n)}-C_{k-1}^{(n)}\right), \quad C_{0}^{(n)}=I_{n}, \quad C_{1}^{(n)}=0, \\
& S_{k+1}^{(n)}=R_{n}\left(2 S_{k}^{(n)}-S_{k-1}^{(n)}\right), \quad S_{0}^{(n)}=0, \quad S_{1}^{(n)}=I_{n}, \tag{3.5}
\end{align*}
$$

where $R_{n}=\left(I_{n}-\tau_{n}^{2} A_{n}\right)^{-1}$. Then, the solution of (3.4) is given by

$$
U_{n}^{k}=C_{k}^{(n)} U_{n}^{0}+S_{k}^{(n)} U_{n}^{1}=\left(C_{k}^{(n)}+S_{k}^{(n)}\right) U_{n}^{0}+\tau_{n} S_{k}^{(n)} \frac{U_{n}^{1}-U_{n}^{0}}{\tau_{n}}
$$

To operate with representations of discrete families of operators we give the following

Definition 3.1. [12]. The operators $A_{n}$ of $C_{0}$-cosine operator-valued function $C\left(\cdot, A_{n}\right)$ satisfy the discrete Krein-Fattorini Conditions if the following conditions hold:
(i) there exist $\mathfrak{B}_{n} \in \mathcal{C}\left(E_{n}\right)$ such that $\mathfrak{B}_{n}^{2}=A_{n}$, and $\mathfrak{B}_{n}$ commutes with any operator from $B\left(E_{n}\right)$ commuting with $A_{n}$;
(ii) the operators $\mathfrak{B}_{n}$ generate $C_{0}$-groups such that $\left\|\exp \left( \pm t \mathfrak{B}_{n}\right)\right\| \leq M_{0} e^{\omega_{0}|t|}, t \in$ $\mathbb{R}$;
(iii) the operators $-A_{n}$ are strongly positive, i.e.

$$
\left\|\left(\lambda I_{n}-A_{n}\right)^{-1}\right\| \leq \frac{M}{1+|\lambda|}, \quad R e \lambda \geq 0
$$

and $\left\|\mathfrak{B}_{n}^{-1}\right\| \leq C$ as $n \in I N$.
We can obtain explicit representations for the functions $C_{k}^{(n)}, S_{k}^{(n)}$ in the following way. Let us introduce the operators

$$
R_{1, n}=\left(I_{n}-\tau_{n} \mathfrak{B}_{n}\right)^{-1}, \quad R_{2, n}=\left(I_{n}+\tau_{n} \mathfrak{B}_{n}\right)^{-1}
$$

where the operators $\mathfrak{B}_{n}$ are those in the Krein-Fattorini conditions. These operators satisfy the relations

$$
\begin{equation*}
R_{1, n} R_{2, n}=R_{n}, \quad R_{1, n}-R_{2, n}=2 \tau_{n} \mathfrak{B}_{n} R_{n}, \quad R_{1, n}+R_{2, n}=2 R_{n} \tag{3.6}
\end{equation*}
$$

which follow from the well-known Hilbert identity for resolvents. Since under the Krein-Fattorini conditions the operator $\mathfrak{B}_{n}$ generates a $C_{0}$-group one has that $\left\|R_{j, n}^{k}\right\| \leq \operatorname{const}(t), k \tau_{n}=t$ for $j=1,2$.

Simple calculations show that the general solution of (3.4) is given as in [16] by the formula

$$
\begin{equation*}
U_{n}^{k}=R_{1, n}^{k} x+R_{2, n}^{k} y \tag{3.7}
\end{equation*}
$$

where $x$ and $y$ are arbitrary elements of $E_{n}$. Note that the representation (3.7) was established also in [2], [3] without Krein-Fattorini conditions, but in our case we need that $\left\|R_{1, n}^{k_{n}}\right\| \leq M e^{\omega t},\left\|R_{2, n}^{k_{n}}\right\| \leq M e^{\omega t}$ with $k_{n} \tau_{n}=t$. Now if we solve the system

$$
\left\{\begin{array}{l}
x+y=U_{n}^{0} \\
R_{1, n} x+R_{2, n} y=U_{n}^{1}
\end{array}\right.
$$

and insert $x$ and $y$ in (3.7), we obtain by some calculations

$$
C_{k}^{(n)}=-R_{n} \sum_{s=0}^{k-2} R_{1, n}^{s} R_{2, n}^{k-2-s}, \quad S_{k}^{(n)}=\sum_{s=0}^{k-1} R_{1, n}^{s} R_{2, n}^{k-1-s}
$$

From (3.6) we derive

$$
\begin{equation*}
R_{n} \mathfrak{B}_{n} S_{k}^{(n)}=\frac{1}{2 \tau_{n}}\left(R_{1, n}-R_{2, n}\right) \sum_{s=0}^{k-1} R_{1, n}^{s} R_{2, n}^{k-1-s}=\frac{1}{2 \tau_{n}}\left(R_{1, n}^{k}-R_{2, n}^{k}\right) . \tag{3.8}
\end{equation*}
$$

We note also that

$$
\begin{align*}
& R_{1, n}^{k}-R_{1, n}^{k-1}=\tau_{n} \mathfrak{B}_{n} R_{1, n}^{k},  \tag{3.9}\\
& R_{2, n}^{k}-R_{2, n}^{k-1}=-\tau_{n} \mathfrak{B}_{n} R_{2, n}^{k},
\end{align*}
$$

and

$$
\begin{equation*}
R_{1, n}^{k}+R_{2, n}^{k}=2 R_{n}\left(S_{k}^{(n)}-S_{k-1}^{(n)}\right) . \tag{3.10}
\end{equation*}
$$

The equality (3.10) can be proved by induction on $k$. For $k=1$ and $k=2$ it can be checked by direct calculations. For $k>2$,

$$
\begin{aligned}
& R_{1, n}^{k+1}+R_{2, n}^{k+1}=\left(R_{1, n}^{k}+R_{2, n}^{k}\right)\left(R_{1, n}+R_{2, n}\right)-R_{1, n} R_{2, n}\left(R_{1, n}^{k-1}+R_{2, n}^{k-1}\right) \\
= & 2 R_{n}\left(S_{k}^{(n)}-S_{k-1}^{(n)}\right) \cdot 2 R_{n}-R_{n} \cdot\left(S_{k-1}^{(n)}-S_{k-2}^{(n)}\right)=2 R_{n}^{2}\left(2 S_{k}^{(n)}-3 S_{k-1}^{(n)}+S_{k-2}^{(n)}\right) \\
= & 2 R_{n}\left(R_{n}\left(2 S_{k}^{(n)}-S_{k-1}^{(n)}\right)-R_{n}\left(2 S_{k-1}^{(n)}-2 S_{k-2}^{(n)}\right)\right)=2 R_{n}\left(S_{k+1}^{(n)}-S_{k}^{(n)}\right) .
\end{aligned}
$$

From (3.6) and equality $I-R_{1, n}=-\tau_{n} \mathfrak{B}_{n} R_{1, n}$ we have

$$
\begin{aligned}
C_{k}^{(n)}+S_{k}^{(n)} & =-R_{1, n} R_{2, n} \sum_{s=0}^{k-2} R_{1, n}^{s} R_{2, n}^{k-2-s}+\sum_{s=0}^{k-1} R_{1, n}^{s} R_{2, n}^{k-1-s} \\
& =-\sum_{s=0}^{k-2} R_{1, n}^{s+1} R_{2, n}^{k-1-s}+\sum_{s=0}^{k-1} R_{1, n}^{s} R_{2, n}^{k-1-s} \\
& =\sum_{s=0}^{k-2}\left(R_{1, n}^{s}-R_{1, n}^{s+1}\right) R_{2, n}^{k-1-s}+R_{1, n}^{k-1} \\
& =-\tau_{n} \mathfrak{B}_{n} \sum_{s=0}^{k-2} R_{1, n}^{s+1} R_{2, n}^{k-1-s}+R_{1, n}^{k-1} \\
& =-\tau_{n} \mathfrak{B}_{n} R_{1, n} R_{2, n} \sum_{s=0}^{k-2} R_{1, n}^{s} R_{2, n}^{k-2-s}+R_{1, n}^{k-1} \\
& =-\tau_{n} \mathfrak{B}_{n} R_{n} S_{k-1}^{(n)}+R_{1, n}^{k-1} .
\end{aligned}
$$

Using (3.8) we get

$$
\begin{equation*}
\tau_{n} \mathfrak{B}_{n} R_{n} S_{k-1}^{(n)}=\frac{1}{2}\left(R_{1, n}^{k-1}-R_{2, n}^{k-1}\right), \tag{3.11}
\end{equation*}
$$

and consequently

$$
\begin{equation*}
C_{k}^{(n)}+S_{k}^{(n)}=\frac{1}{2}\left(R_{1, n}^{k-1}+R_{2, n}^{k-1}\right) . \tag{3.12}
\end{equation*}
$$

Let consider the inhomogeneous equation (3.2), i.e.
(3.13) $U_{n}^{k+1}-2\left(I_{n}-\tau_{n}^{2} A_{n}\right)^{-1} U_{n}^{k}+\left(I_{n}-\tau_{n}^{2} A_{n}\right)^{-1} U_{n}^{k-1}=\tau_{n}^{2}\left(I_{n}-\tau_{n}^{2} A_{n}\right)^{-1} \varphi_{n}^{k}$.

Using the recurrent relation (3.5) we derive formula (3.3), see [16].

## 4. Existence of Solutions to the Inverse Problem

Consider the inverse problem (1.1) in the following form: for given elements $u^{T}, u^{0}, u^{1} \in D(A)$ find a solution $u(\cdot) \in C^{2}([0, T] ; E)$ and an element $d \in E$ such that

$$
\left\{\begin{array}{l}
u^{\prime \prime}(t)=A u(t)+\Phi(t) d, \quad 0 \leq t \leq T  \tag{4.1}\\
u(0)=u^{0}, u^{\prime}(0)=u^{1} \\
u(T)=u^{T}
\end{array}\right.
$$

Here $A \in C(M ; \omega)$. The problem (4.1) is an inverse problem with overdetermination. Details on such kind of description of problems can be found in [15].

Basing on Remark 1.1, we can treat the solution of (4.1) as follows

$$
A \int_{0}^{T} S(T-s, A) \Phi(s) d s d=A u(T)-C(T, A) A u^{0}-A S(T, A) u^{1}
$$

and then use the identities
(4.2) $A \int_{0}^{T} S(T-s, A) \Phi(s) d s=\int_{0}^{T} C(T-s, A) \Phi^{\prime}(s) d s-\Phi(T)+C(T, A) \Phi(0)$
and

$$
\begin{align*}
& A \int_{0}^{T} S(T-s, A) \Phi(s) d s \\
= & \int_{0}^{T} S(T-s, A) \Phi^{\prime \prime}(s) d s-S(T, A) \Phi^{\prime}(0)-\Phi(T)+C(T, A) \Phi(0) \tag{4.3}
\end{align*}
$$

Proposition 4.1. Let $\Phi(\cdot) \in C^{1}([0, T] ; B(E))$, the operator $\Phi(T)$ be invertible, i.e. $\Phi(T)^{-1} \in B(E)$. Then the inverse problem (4.1) is equivalent to that of solving

$$
\begin{equation*}
I d-B_{1} d=g_{1} \tag{4.4}
\end{equation*}
$$

where

$$
B_{1}=\Phi(T)^{-1}\left(\int_{0}^{T}\left(C(T-s, A) \Phi^{\prime}(s)-\lambda S(T-s, A) \Phi(s)\right) d s+C(T, A) \Phi(0)\right)
$$

and

$$
g_{1}:=-\Phi(T)^{-1}(A-\lambda I)\left(u^{T}-C(T, A) u^{0}-S(T, A) u^{1}\right) \text { for } \lambda \in \rho(A)
$$

Proposition 4.2. Let $\Phi(\cdot) \in C^{2}([0, T] ; B(E))$, and assume that the operator

$$
\begin{equation*}
D=\Phi(T)-C(T, A) \Phi(0) \tag{4.5}
\end{equation*}
$$

is invertible, i.e. $D^{-1} \in B(E)$. Then the inverse problem (4.1) is equivalent to that of solving

$$
\begin{equation*}
I d-B_{2} d=g_{2}, \tag{4.6}
\end{equation*}
$$

where

$$
B_{2}:=D^{-1}\left(\int_{0}^{T} S(T-s, A)\left(\Phi^{\prime \prime}(s)-\lambda \Phi(s)\right) d s+S(T, A) \Phi^{\prime}(0)\right)
$$

and

$$
g_{2}:=-D^{-1}(A-\lambda I)\left(u^{T}-C(T, A) u^{0}-S(T, A) u^{1}\right) \text { for } \lambda \in \rho(A) .
$$

Proposition 4.3. [15]. Let the conditions of Proposition 4.1 be satisfied and

$$
\int_{0}^{T}\left(\left\|\Phi^{\prime}(s)\right\|+|\lambda|(T-s)\|\Phi(s)\|\right) e^{\omega(T-s)} d s+\|\Phi(0)\| e^{\omega T}<\frac{1}{M\left\|\Phi(T)^{-1}\right\|}
$$

Then a solution $(u(\cdot), d)$ of the inverse problem (4.1) exists and is unique for any input data $u^{0}$, $u^{T} \in D(A), u^{1} \in E^{1}$.

Proposition 4.4. [15]. Assume that the conditions of Proposition 4.2 and the inequality

$$
\int_{0}^{T}(T-s)\left\|\Phi^{\prime \prime}(s)-\lambda \Phi(s)\right\| e^{\omega(T-s)} d s+T\left\|\Phi^{\prime}(0)\right\| e^{\omega T}<\frac{1}{M\left\|D^{-1}\right\|}
$$

are satisfied. Then a solution $(u(\cdot), d)$ of the inverse problem (4.1) exists and is unique for any input data $u^{0}, u^{T} \in D(A), u^{1} \in E^{1}$.

Proposition 4.5. [15]. Assume that the operator A generates a strongly continuous $C_{0}$-cosine operator-function $C(\cdot, A)$ on the Banach space $E, \Phi(t) \equiv I$ and $0 \in \rho(A)$. Then the inverse problem (4.1) is uniquely solvable for any input data $u^{0}, u^{T} \in D(A), u^{1} \in E^{1}$ if and only if $1 \in \rho(C(T, A))$.

We now assume that $E$ is the Hilbert space and the operator $A$ is selfadjoint and negative. For any real-valued function $\Phi(\cdot)$, the value $\Phi(t)$ will be identified with the operator of multiplication by the number $\Phi(t)$ in the space $E$. The characteristic function $\varphi(\cdot)$ on the negative semi-axis is defined by

$$
\begin{equation*}
\varphi(\lambda)=\frac{1}{\sqrt{-\lambda}} \int_{0}^{T} \Phi(s) \sin (\sqrt{-\lambda}(T-s)) d s \tag{4.7}
\end{equation*}
$$

Note that we might extend the function $\varphi(\cdot)$ from the negative semi-axis to construct an entire function of the complex variable $\lambda$. If, in particular, $\Phi(t) \not \equiv 0$, then the zeroes of the function $\varphi(\cdot)$ are isolated.

In what follows, we denote by $E_{\lambda}$ the spectral decomposition of unity of the operator $A$. With this notation, we can write

$$
A=\int_{0}^{+\infty} \lambda d E(\lambda) .
$$

Theorem 4.1. [15]. If the operator $A$ is self-adjoint and semibounded from above on the Hilbert space $E, \Phi(\cdot) \in C^{1}[0, T]$ and $\Phi(\cdot) \not \equiv 0$, then the following statements hold:
(i) the inverse problem (4.1) with the fixed input data $u^{0}, u^{T} \in D(A), u^{1} \in E^{1}$ is solvable if and only if

$$
\begin{equation*}
\int_{0}^{+\infty}|\varphi(\lambda)|^{-2} d\left(E_{\lambda} g, g\right)<\infty, \tag{4.8}
\end{equation*}
$$

being $g:=u^{T}-C(T, A) u^{0}-S(T, A) u^{1}$;
(ii) if the inverse problem (4.1) is solvable, then its solution is unique if and only if the point spectrum of the operator $A$ contains no zeros of the entire function $\varphi(\cdot)$ defined by (4.7).
Of special interest is the particular case $\Phi(t) \equiv t I$. In this case we have

$$
\varphi(\lambda)= \begin{cases}\frac{\sin (\sqrt{-\lambda} T)-\sqrt{-\lambda} T}{\lambda \sqrt{-\lambda} T} & , \lambda \neq 0 \\ T^{3} / 6 & , \lambda=0\end{cases}
$$

This function has no zeros on the negative semi-axis and $\varphi(\lambda) \sim-\frac{T}{\lambda}$ as $\lambda \rightarrow+\infty$. Hence the convergence of the integral

$$
\int_{0}^{+\infty}|\varphi(\lambda)|^{-2} d\left(E_{\lambda} g, g\right)
$$

is equivalent to that of

$$
\int_{-\infty}^{+\infty}|\lambda|^{2} d\left(E_{\lambda} g, g\right)
$$

This integral converges for every element $g \in D(A)$. We thus proved the following
Proposition 4.6. Assume that the operator $A$ is self-adjoint and negative on the Hilbert space $E, \Phi(t) \equiv t I, t \geq 0$. Then the solution $(u(\cdot), d)$ of the inverse problem (4.1) exists and is unique for any input data $u^{0}, u^{T} \in D(A), u^{1} \in E^{1}$.

## 5. Approximating the Solution of the Inverse Problems

Let us consider the semidiscretization of the inverse problem for the secondorder equation (4.1): for given elements $u_{n}^{T}, u_{n}^{0}, u_{n}^{1} \in D\left(A_{n}\right)$ find a solution $u_{n}(\cdot) \in$ $C^{2}\left([0, T] ; E_{n}\right)$ and an element $d_{n} \in E_{n}$ such that

$$
\left\{\begin{array}{l}
u_{n}^{\prime \prime}(t)=A_{n} u_{n}(t)+\Phi_{n}(t) d_{n}, \quad 0 \leq t \leq T  \tag{5.1}\\
u_{n}(0)=u_{n}^{0}, u_{n}^{\prime}(0)=u_{n}^{1} \\
u_{n}(T)=u_{n}^{T}
\end{array}\right.
$$

where operators $A_{n} \in C(M ; \omega)$, the operators $A, A_{n}$ are consistent, $u_{n}^{T} \xrightarrow{\mathcal{P}} u^{T}$, $u_{n}^{0} \xrightarrow{\mathcal{P}} u^{0}, u_{n}^{1} \xrightarrow{\mathcal{P}} u^{1}$ and $\Phi_{n}(\cdot) \xrightarrow{\mathcal{P} \mathcal{P}} \Phi(\cdot)$ in the sense.

The solution $d_{n}$ of the problem (5.1) must satisfy the equation

$$
\begin{equation*}
I_{n} d_{n}-B_{n, 2} d_{n}=g_{n, 2} \tag{5.2}
\end{equation*}
$$

where

$$
\begin{aligned}
B_{n, 2} & :=D_{n}^{-1}\left(\int_{0}^{T} S\left(T-s, A_{n}\right)\left(\Phi_{n}^{\prime \prime}(s)-\lambda \Phi_{n}(s)\right) d s+S\left(T, A_{n}\right) \Phi_{n}^{\prime}(0)\right), \lambda \in \Delta_{c c} \\
g_{n, 2} & :=-D_{n}^{-1}\left(A_{n}-\lambda I_{n}\right)\left(u_{n}^{T}-C\left(T, A_{n}\right) u_{n}^{0}-S\left(T, A_{n}\right) u_{n}^{1}\right) \\
D_{n} & =\Phi_{n}(T)-C\left(T, A_{n}\right) \Phi_{n}(0)
\end{aligned}
$$

Theorem 5.1. [24]. Let $A, A_{n} \in C(M ; \omega)$. Then $S\left(t, A_{n}\right) \xrightarrow{\mathcal{P} \mathcal{P}} S(t, A)$ compactly for any $t>0$ iff $\Delta_{c c} \neq \emptyset$.

Theorem 5.2. Assume that $\Phi(\cdot) \in C^{3}([0, T] ; B(E)), \Phi_{n}(\cdot) \in C^{3}([0, T] ;$ $\left.B\left(E_{n}\right)\right), D_{n}^{-1} \xrightarrow{\mathcal{P P}} D^{-1}$, the resolvents $\left(\lambda I_{n}-A_{n}\right)^{-1},(\lambda I-A)^{-1}$ are compact, $(A),\left(B^{\prime \prime}\right)$ and (1.4) are satisfied, $\Phi_{n}^{(j)}(t) \xrightarrow{\mathcal{P} \mathcal{P}} \Phi^{(j)}(t)$ uniformly in $t \in[0, T]$ for $j \in \overline{1,3}$, and $\Delta_{c c} \neq \emptyset$. Assume also that the problem (4.1) has a unique solution for any $u^{T} \in D(A)$. Then there are solutions to problems (5.1) for almost all $n$ and they converge to solution of problem (4.1), i.e. $u_{n}(t) \xrightarrow{\mathcal{P}} u(t)$ uniformly in $t \in[0, T]$ and $d_{n} \xrightarrow{\mathcal{P}} d$ as $n \in I N$, whenever $A_{n} u_{n}^{0} \xrightarrow{\mathcal{P}} A u^{0}, A_{n} u_{n}^{1} \xrightarrow{\mathcal{P}} A u^{1}, A_{n} u_{n}^{T} \xrightarrow{\mathcal{P}} A u^{T}$.

Proof. We, first, show that the solutions of equations (5.2) converge to the solution of equation (4.6). Since $D_{n}^{-1} \xrightarrow{\mathcal{P} \mathcal{P}} D^{-1}$, it is clear that $g_{n, 2} \xrightarrow{\mathcal{P}} g_{2}$. If $B_{n, 2} \xrightarrow{\mathcal{P} \mathcal{P}} B_{2}$ compactly, then by Theorems 2.1 and 2.2 it follows that $d_{n} \xrightarrow{\mathcal{P}} d$ and Theorem 5.2 is proved.

Using Theorem 5.1 one can show that the operators $B_{n, 2} \xrightarrow{\mathcal{P} \mathcal{P}} B_{2}$ compactly. To see this recall that operators $B_{n, 2}, B_{2}$ can be split into two parts. The first term

$$
D_{n}^{-1} S\left(T, A_{n}\right) \Phi_{n}^{\prime}(0) \xrightarrow{\mathcal{P} \mathcal{P}} D^{-1} S(T, A) \Phi^{\prime}(0)
$$

converges compactly and the second term

$$
\begin{aligned}
& \quad D_{n}^{-1}\left(\lambda I_{n}-A_{n}\right)^{-1}\left(\lambda I_{n}-A_{n}\right) \\
& \int_{0}^{T} S\left(T-s, A_{n}\right) \Phi_{n}^{\prime \prime}(s) d s \xrightarrow{\mathcal{P} \mathcal{P}} D^{-1}(\lambda I-A)^{-1}(\lambda I-A) \\
& \int_{0}^{T} S(T-s, A) \Phi^{\prime \prime}(s) d s
\end{aligned}
$$

also converges compactly, since $\Delta_{c c} \neq \emptyset$ and

$$
\left(\lambda I_{n}-A_{n}\right) \int_{0}^{T} S\left(T-s, A_{n}\right) \Phi_{n}^{\prime \prime}(s) d s \xrightarrow{\mathcal{P} \mathcal{P}}(\lambda I-A) \int_{0}^{T} S(T-s, A) \Phi^{\prime \prime}(s) d s
$$

The last statement can be derived from the representation like (4.2). Therefore from Theorems 2.1 and 2.2 it follows that $d_{n} \xrightarrow{\mathcal{P}} d$. The convergence of solutions $u_{n}(t) \xrightarrow{\mathcal{P}} u(t)$ uniformly in $t \in[0, T]$ then follows from representation formulae like (1.5).

Consider the discretization of (5.1) in time

$$
\begin{align*}
& \frac{U_{n}^{k+1}-2 U_{n}^{k}+U_{n}^{k-1}}{\tau_{n}^{2}}=A_{n} U_{n}^{k+1}+\Phi_{n}\left(k \tau_{n}\right) \tilde{d}_{n}  \tag{5.3}\\
& k \in\left\{1, \ldots,\left[\frac{T}{\tau_{n}}\right]\right\}, U_{n}^{0}=u_{n}^{0}, U_{n}^{1}=u_{n}^{0}+\tau_{n} u_{n}^{1}
\end{align*}
$$

According to (3.3) one can write its solution as

$$
\begin{equation*}
U_{n}^{k}=C_{k}^{(n)} U_{n}^{0}+S_{k}^{(n)} U_{n}^{1}+\tau_{n}^{2} R_{n} \sum_{j=2}^{k} S_{k+1-j}^{(n)} \Phi_{n}^{j-1} \tilde{d}_{n} \tag{5.4}
\end{equation*}
$$

where we wrote $\Phi_{n}^{j}=\Phi_{n}\left(j \tau_{n}\right)$. Using (3.8) and (3.9) and summing by parts we have

$$
\begin{align*}
& R_{n} A_{n} \tau_{n}^{2} \sum_{j=2}^{k} S_{k+1-j}^{(n)} \Phi_{n}^{j-1}=R_{n} \mathfrak{B}_{n}^{2} \tau_{n}^{2} \sum_{j=2}^{k} S_{k+1-j}^{(n)} \Phi_{n}^{j-1} \\
= & \frac{\tau_{n} \mathfrak{B}_{n}}{2} \sum_{j=2}^{k}\left(R_{1, n}^{k+1-j}-R_{2, n}^{k+1-j}\right) \Phi_{n}^{j-1}=\frac{\tau_{n} \mathfrak{B}_{n}}{2} \sum_{j=2}^{k} R_{1, n}^{k+1-j} \Phi_{n}^{j-1} \\
& -\frac{\tau_{n} \mathfrak{B}_{n}}{2} \sum_{j=2}^{k} R_{2, n}^{k+1-j} \Phi_{n}^{j-1}  \tag{5.5}\\
& =\frac{1}{2} \sum_{j=2}^{k}\left(R_{1, n}^{k+1-j}-R_{1, n}^{k-j}\right) \Phi_{n}^{j-1}+\frac{1}{2} \sum_{j=2}^{k}\left(R_{2, n}^{k+1-j}-R_{2, n}^{k-j}\right) \Phi_{n}^{j-1}
\end{align*}
$$

$$
\begin{aligned}
= & \frac{1}{2}\left(\sum_{j=2}^{k} R_{1, n}^{k+1-j}\left(\Phi_{n}^{j-1}-\Phi_{n}^{j-2}\right)-\Phi_{n}^{k-1}+R_{1, n}^{k-1} \Phi_{n}^{0}\right) \\
& +\frac{1}{2}\left(\sum_{j=2}^{k} R_{2, n}^{k+1-j}\left(\Phi_{n}^{j-1}-\Phi_{n}^{j-2}\right)-\Phi_{n}^{k-1}+R_{2, n}^{k-1} \Phi_{n}^{0}\right) \\
= & \frac{1}{2} \sum_{j=2}^{k}\left(R_{1, n}^{k+1-j}+R_{2, n}^{k+1-j}\right)\left(\Phi_{n}^{j-1}-\Phi_{n}^{j-2}\right)-\Phi_{n}^{k-1}+\frac{1}{2}\left(R_{1, n}^{k-1}+R_{2, n}^{k-1}\right) \Phi_{n}^{0} .
\end{aligned}
$$

Using (3.10) and summing by parts we obtain again

$$
\begin{align*}
& \frac{1}{2} \sum_{j=2}^{k}\left(R_{1, n}^{k+1-j}+R_{2, n}^{k+1-j}\right)\left(\Phi_{n}^{j-1}-\Phi_{n}^{j-2}\right) \\
& =R_{n} \sum_{j=2}^{k}\left(S_{k+1-j}^{(n)}-S_{k-j}^{(n)}\right)\left(\Phi_{n}^{j-1}-\Phi_{n}^{j-2}\right)  \tag{5.6}\\
& =R_{n}\left(\sum_{j=2}^{k} S_{k+1-j}^{(n)}\left(\Phi_{n}^{j-1}-2 \Phi_{n}^{j-2}+\Phi_{n}^{j-3}\right)+S_{k-1}^{(n)}\left(\Phi_{n}^{0}-\Phi_{n}^{-1}\right)\right) .
\end{align*}
$$

From (5.5) and (5.6) we get the next identity valid for any solution of (5.3)

$$
\begin{align*}
& A_{n} U_{n}^{k}=A_{n} C_{k}^{(n)} U_{n}^{0}+A_{n} S_{k}^{(n)} U_{n}^{1} \\
+ & R_{n} \sum_{j=2}^{k} S_{k+1-j}^{(n)}\left(\Phi_{n}^{j-1}-2 \Phi_{n}^{j-2}+\Phi_{n}^{j-3}\right) \tilde{d}_{n}  \tag{5.7}\\
+ & \left(R_{n} \tau_{n} S_{k-1}^{(n)} \frac{\Phi_{n}^{0}-\Phi_{n}^{-1}}{\tau_{n}}-\Phi_{n}^{k-1}+\frac{1}{2}\left(R_{1, n}^{k-1}+R_{2, n}^{k-1}\right) \Phi_{n}^{0}\right) \tilde{d}_{n}, k \geq 2 .
\end{align*}
$$

As in (4.5) define the operator

$$
D_{n, k_{n}}=\Phi_{n}^{k_{n}-1}-\left(C_{k_{n}}^{(n)}+S_{k_{n}}^{(n)}\right) \Phi_{n}^{0} .
$$

Then we have the following
Theorem 5.3. Assume that $\Phi(\cdot) \in C^{4}([0, T] ; B(E)), \Phi_{n}(\cdot) \in C^{4}\left([0, T] ; B\left(E_{n}\right)\right)$, $\mathfrak{B}_{n}^{-1} \xrightarrow{\mathcal{P} \mathcal{P}} \mathfrak{B}^{-1}$ compactly, $\left(D_{n, k_{n}}\right)^{-1} \xrightarrow{\mathcal{P} \mathcal{P}} D^{-1}, k_{n} \tau_{n}=T$, the resolvents $\left(\lambda I_{n}-\right.$ $\left.A_{n}\right)^{-1},(\lambda I-A)^{-1}$ are compact, $\left(B^{\prime \prime}\right)$ and $(1.4)$ are satisfied and $\Phi_{n}^{(l)}(t) \xrightarrow{\mathcal{P P}} \Phi^{(l)}(t)$ uniformly in $t \in[0, T]$ for $l=\overline{1,4}$. Assume also that the problem (4.1) has a unique solution for any $u^{T} \in D(A)$ and the Krein-Fattorini conditions are satisfied. Then there are solutions of the problem (5.3) for almost all $n$ and they converge to the solution of problem (4.1), i.e.

$$
U_{n}(t) \xrightarrow{\mathcal{P}} u(t) \text { uniformly in } t \in[0, T]
$$

and $\tilde{d}_{n} \xrightarrow{\mathcal{P}} d$ as $n \in I N$, whenever $A_{n} u_{n}^{T} \xrightarrow{\mathcal{P}} A u^{T}, A_{n} u_{n}^{0} \xrightarrow{\mathcal{P}} A u^{0}, \mathfrak{B}_{n} u_{n}^{1} \xrightarrow{\mathcal{P}} \mathfrak{B} u^{1}$.

Proof. First we apply the operator $\left(A_{n}-\lambda I_{n}\right)$ to (5.4) for $\lambda \in \Delta_{c c}$. Using (5.7) we get equation

$$
\begin{align*}
& \left(\Phi_{n}^{k-1}-\left(C_{k}^{(n)}+S_{k}^{(n)}\right) \Phi_{n}^{0}\right) \tilde{d}_{n}-\left[R_{n} \sum_{j=2}^{k} S_{k+1-j}^{(n)}\left(\Phi_{n}^{j-1}-2 \Phi_{n}^{j-2}+\Phi_{n}^{j-3}\right)\right. \\
- & \left.\lambda R_{n} \tau_{n}^{2} \sum_{j=2}^{k} S_{k+1-j}^{(n)} \Phi_{n}^{j-1}+R_{n} S_{k-1}^{(n)}\left(\Phi_{n}^{0}-\Phi_{n}^{-1}\right)\right] \tilde{d}_{n}  \tag{5.8}\\
= & \left(A_{n}-\lambda I_{n}\right)\left[C_{k}^{(n)} U_{n}^{0}+S_{k}^{(n)} U_{n}^{1}-U_{n}^{k}\right]
\end{align*}
$$

Since

$$
\left(D_{n, k_{n}}\right)^{-1}=\left(\Phi_{n}^{k_{n}-1}-\left(C_{k_{n}}^{(n)}+S_{k_{n}}^{(n)}\right) \Phi_{n}^{0}\right)^{-1} \xrightarrow{\mathcal{P} \mathcal{P}} D^{-1}
$$

we can rewrite (5.8) in the form

$$
\begin{equation*}
I_{n} \tilde{d}_{n}-B_{n, 3} \tilde{d}_{n}=g_{n, 3} \tag{5.9}
\end{equation*}
$$

where

$$
\begin{aligned}
B_{n, 3}:= & \left(D_{n, k_{n}}\right)^{-1}\left[R_{n} \sum_{j=2}^{k_{n}} S_{k_{n}+1-j}^{(n)}\left(\Phi_{n}^{j-1}-2 \Phi_{n}^{j-2}+\Phi_{n}^{j-3}\right)\right. \\
& \left.-\lambda R_{n} \tau_{n}^{2} \sum_{j=2}^{k_{n}} S_{k_{n}+1-j}^{(n)} \Phi_{n}^{j-1}+R_{n} S_{k_{n}-1}^{(n)}\left(\Phi_{n}^{0}-\Phi_{n}^{-1}\right)\right]
\end{aligned}
$$

and

$$
g_{n, 3}:=\left(D_{n, k_{n}}\right)^{-1}\left(A_{n}-\lambda I_{n}\right)\left(C_{k_{n}}^{(n)} U_{n}^{0}+S_{k_{n}}^{(n)} U_{n}^{1}-u_{n}^{T}\right), \quad k_{n} \tau_{n}=T
$$

To show that $B_{n, 3} \xrightarrow{\mathcal{P} \mathcal{P}} B_{2}$ compactly we split the operators $B_{n, 3}, B_{2}$ into two parts. Compact convergence

$$
\frac{1}{2}\left(R_{1, n}^{k-1}-R_{2, n}^{k-1}\right) \mathfrak{B}_{n}^{-1} \xrightarrow{\mathcal{P} \mathcal{P}} \frac{1}{2}(\exp (t \mathfrak{B})-\exp (-t \mathfrak{B})) \mathfrak{B}^{-1}
$$

because of (3.11) and Theorem 2.5, implies that

$$
\left(D_{n, k_{n}}\right)^{-1} R_{n} \tau_{n} S_{k_{n}-1}^{(n)} \frac{\Phi_{n}^{0}-\Phi_{n}^{-1}}{\tau_{n}} \xrightarrow{\mathcal{P} \mathcal{P}} D^{-1} S(T, A) \Phi^{\prime}(0)
$$

compactly. The other parts of the operators $B_{n, 3}$ also converge compactly to the corresponding parts of $B_{2}$. One can see by the same reasons as in (3.11) and (1.12) that

$$
\mathfrak{B}_{n}\left(B_{n, 3}-\left(D_{n, k_{n}}\right)^{-1} R_{n} \tau_{n} S_{k_{n}-1}^{(n)} \frac{\Phi_{n}^{0}-\Phi_{n}^{-1}}{\tau_{n}}\right) \xrightarrow{\mathcal{P} \mathcal{P}} \mathfrak{B}\left(B_{2}-D^{-1} S(T, A) \Phi^{\prime}(0)\right)
$$

and this implies that $B_{n, 3} \xrightarrow{\mathcal{P} \mathcal{P}} B_{2}$ compactly.
The convergence of the finite differences to derivatives follows, e.g., from [9], p. 409. Therefore, from Theorems 2.1 and 2.2 it follows that $\tilde{d}_{n} \xrightarrow{\mathcal{P}} d$. The convergence of solutions $U_{n}(t) \xrightarrow{\mathcal{P}} u(t)$ uniformly in $t \in[0, T]$ follows from the representation formulas (5.4) and (1.5).

Remark 5.1. In case of Hilbert space and negative self-adjoint operators $A$ in Theorem 5.3, one can omit the condition that $\mathfrak{B}_{n, 3} \xrightarrow{\mathcal{P P}} \mathfrak{B}_{2}$ compactly and just claim the condition $\Delta_{c c} \neq \emptyset$. Indeed, then one can get the compact convergence of square roots of operators as in [10] and then get the compact convergence $B_{n, 3} \xrightarrow{\mathcal{P P}} B_{2}$ as before.

Remark 5.2. In case of a Banach space in Theorem 5.3 one can also omit the condition of compact convergence $\mathfrak{B}_{n}^{-1} \xrightarrow{\mathcal{P P}} \mathfrak{B}^{-1}$ and just use the condition $\Delta_{c c} \neq \emptyset$. In this case one should assume that problem (4.1) possesses some extra smoothness condition. More precisely, assume that $u(\cdot) \in C^{4}([0, T] ; B(E)), \Phi(\cdot) \in$ $C^{3}([0, T] ; B(E)), \Phi(0)=0$ and $u^{0}, u^{1} \in D\left(A^{2}\right)$. Then from Proposition 1.3 follows that $A S(\cdot, A) \Phi^{\prime}(0) \in C([0, T] ; E)$. Moreover, if we assume that for the problems (5.1) and (4.1) $U_{n}^{(4)}(t) \xrightarrow{\mathcal{P}} u^{(4)}(t)$ uniformly in $t \in[0, T]$ for any $\tilde{d}_{n} \xrightarrow{\mathcal{P}} d$, then from the discrete analog of (1.8) follows that

$$
\begin{equation*}
A_{n} R_{n} \tau_{n} S_{k_{n}-1}^{(n)} \frac{\Phi_{n}^{0}-\Phi_{n}^{-1}}{\tau_{n}} \tilde{d}_{n} \xrightarrow{\mathcal{P}} A S(T, A) \Phi^{\prime}(0) d . \tag{5.10}
\end{equation*}
$$

This means that without loss of generality one can assume that

$$
A_{n} R_{n} \tau_{n} S_{k_{n}-1}^{(n)} \frac{\Phi_{n}^{0}-\Phi_{n}^{-1}}{\tau_{n}} \xrightarrow{\mathcal{P} \mathcal{P}} A S(T, A) \Phi^{\prime}(0),
$$

and then the compact convergence $B_{n, 3} \xrightarrow{\mathcal{P} \mathcal{P}} B_{2}$ can be established directly from the convergence $A_{n} B_{n, 3} \xrightarrow{\mathcal{P P}} A B_{2}$.

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Dmitry Orlovsky<br>Department of Mathematics,<br>"MEPhI" National Nuclear Research University,<br>Kashyrskoye shosse 31,<br>Moscow 115409,<br>Russia<br>E-mail: odg@bk.ru<br>Sergey Piskarev<br>Scientific Research Computer Center,<br>Lomonosov Moscow State University,<br>Leninskie Gory,<br>Moscow 119991,<br>Russia<br>E-mail: piskarev@gmail.com<br>Renato Spigler<br>Department of Mathematics,<br>University Roma Tre,<br>Roma, Italy<br>E-mail: spigler@mat.uniroma3.it

