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# Even Order Differential Equations with Measures as Coefficients 

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#### Abstract

The note deals with differential equations with Borel measures as coefficients. The problem of existence and uniqueness of solutions is discussed. The Ritz-Galerkin method is used for determining approximate solutions.


Key words: Differential equation, Borel measure, Ritz-Galerkin method.

1991 Mathematics Subject Classification: 34A12, 34A45

1. Let us consider the boundary value problem

$$
\begin{gather*}
\sum_{i=1}^{k}(-1)^{i} a_{i} u^{(2 i)}+\mu_{1} u=\mu_{2}  \tag{1}\\
u^{(i)}(a)=u^{(i)}(b)=0 \quad \text { for } i=0, \ldots, k-1, k \geq 1
\end{gather*}
$$

where $\mu_{1}$ and $\mu_{2}$ are real Borel measures, $\mu_{1} \geq 0$ and $a_{i} \geq 0$ for $i=1, \ldots, k$.
It was shown in [1] that for $k=1$ Problem 1 has exactly one solution, which can be approximated by the Ritz-Galerkin method. Here we shall show that Problem 1 can be solved similarly.

Every continuous function $u$ fulfilling the boundary conditions and the differential equation (1) in the distributional sense will be called a solution of Problem 1. This means that

$$
\begin{equation*}
\sum_{i=1}^{k}(-1)^{i} a_{i} \int_{a}^{b} u \varphi^{(2 i)} d x+\int_{a}^{b} u \varphi d \mu_{1}(x)=\int_{a}^{b} \varphi d \mu_{2}(x) \quad \text { for } \varphi \in D(a, b) \tag{2}
\end{equation*}
$$

where $D(a, b)$ denotes the Schwartz space of test functions. Therefore $u^{(2 k-1)} \in$ $B V(a, b) \subset L^{2}(a, b)$. Further it implies that $u \in W^{2 k-1,2}(a, b)\left(W^{n, 2}(a, b), n=\right.$ $1,2, \ldots$ is the Sobolev space of functions $u \in L^{2}(a, b)$ such that the distributional derivative $u^{(n)}$ belongs to $L^{2}(a, b)$; the space $W^{n, 2}(a, b)$ is equipped with the norm $\left.\|u\|:=\left(\sum_{i=0}^{n}\left\|u^{(i)}\right\|_{L^{2}}^{2}\right)^{\frac{1}{2}}\right)$.

Let us define now the space

$$
W_{0}^{n, 2}(a, b):=\left\{u \in W^{n, 2}(a, b): u^{(i)}(a)=u^{(i)}(b)=0 \quad \text { for } i=1, \ldots, n-1\right\} .
$$

It is easy to show that $W_{0}^{n, 2}(a, b)$ is the closure of $D(a, b)$ in $W^{n, 2}(a, b)$ with respect to its norm $\|\cdot\|$.

An easy computation using integration by parts shows that equation (2) can be rewritten in the form

$$
\begin{equation*}
\sum_{i=1}^{k} a_{i} \int_{a}^{b} u^{(i)} \varphi^{(i)} d x+\int_{a}^{b} u \varphi d \mu_{1}(x)=\int_{a}^{b} \varphi d \mu_{2}(x) \quad \text { for } \varphi \in D(a, b) \tag{3}
\end{equation*}
$$

For simplicity of notation we put for $u, \varphi \in W_{0}^{k, 2}(a, b)$

$$
\alpha(u, \varphi):=\sum_{i=1}^{k} a_{i} \int_{a}^{b} u^{(i)} \varphi^{(i)} d x+\int_{a}^{b} u \varphi d \mu_{1}(x), \quad \beta(\varphi):=\int_{a}^{b} \varphi d \mu_{2}(x) .
$$

Thus the original Problem 1 can be written as

$$
\begin{equation*}
\alpha(u, \varphi)=\beta(\varphi) \quad \text { for } \varphi \in W_{0}^{\dot{k}, 2}(a, b) \tag{4}
\end{equation*}
$$

It results from the discussion given above that if equation (4) has a unique solution $u \in W_{0}^{k, 2}(a, b)$, then $u$ is the unique solution of Problem 1 and $u \in$ $W^{2 k-1,2}(a, b)$. Hence we can consider the problem of existence and uniqueness of solutions of equation (4). For this purpose we will define the following two norms in the space $W_{0}^{k, 2}(a, b)$ :

$$
\begin{aligned}
& \|u\|_{\alpha}:=\sqrt{\alpha(u, u)} \\
& \|u\|\|:=\| u^{(k)} \|_{L^{2}}
\end{aligned}
$$

Theorem 1 The norms $\|\cdot\|,|||\cdot|||$ and $\|\cdot\|_{\alpha}$ are equivalent on $W_{0}^{k, 2}(a, b)$.
Proof One can show that $\left|u^{(i)}(x)\right| \leq \sqrt{b-a}\left\|u^{(i+1)}\right\|_{L^{2}}$ for $x \in(a, b), u \in$ $W_{0}^{k, 2}(a, b), i=0, \ldots, k-1$. From this we get

$$
\begin{equation*}
\left\|u^{(i)}\right\|_{L^{2}} \leq(b-a)^{k-i} \mid\|u\| \| \tag{5}
\end{equation*}
$$

and

$$
\begin{equation*}
\sup _{x \in(a, b)} u^{2}(x) \leq(b-a)\left\|u^{\prime}\right\|_{L^{2}}^{2} \leq(b-a)^{2 k-1} \mid\|u\| \|^{2} \tag{6}
\end{equation*}
$$

By (5) we obtain

$$
\begin{equation*}
\|u\|\left\|^{2} \leq\right\| u\left\|^{2}=\sum_{i=0}^{k}\right\| u^{(i)}\left\|_{L^{2}}^{2} \leq \sum_{i=0}^{k}(b-a)^{2(k-i)}\right\| u u\| \|^{2}=A \mid\|u\| \|^{2} \tag{7}
\end{equation*}
$$

where $A=\frac{(b-a)^{2 k+2}-1}{(b-a)^{2}-1}$. Therefore the norms $\|\cdot\|$ and $\|\|\cdot\|\|$ are equivalent. Moreover, by virtue of inequalities (6) and (7) there is

$$
\begin{aligned}
a_{k}\| \| u\| \|^{2} & \leq\|u\|_{\alpha}^{2}=\sum_{i=1}^{k} a_{i}\left\|u^{(i)}\right\|_{L^{2}}^{2}+\int_{a}^{b} u^{2} d \mu_{1}(x) \leq \\
& \leq\left(A_{i=1, \ldots, k} \max _{i}+(b-a)^{2 k-1} \mu_{1}([a, b])\right)\|u\| \|^{2}
\end{aligned}
$$

This finishes the proof.
In the space $W_{0}^{k, 2}(a, b)$ besides the usual inner product

$$
(u, \varphi):=\sum_{i=0}^{k}\left(u^{(i)}, \varphi^{(i)}\right)_{L^{2}}
$$

one can consider two other products, namely $\alpha(u, \varphi)$ and $[u, \varphi]:=\left(u^{(k)}, \varphi^{(k)}\right)_{L^{2}}$. It follows from Theorem 1 that $W_{0}^{k, 2}(a, b)$ is a Hilbert space with any of these products.

Theorem 2 Problem 1 has exactly one solution in $W^{2 k-1,2}(a, b)$.
Proof From the previous remarks it follows that it is sufficient to show that equation (4) has exactly one solution in $W_{0}^{k, 2}(a, b) . \beta$ is a continuous linear form on $W_{0}^{k, 2}(a, b),\left(W_{0}^{k, 2}(a, b), \alpha(\cdot, \cdot)\right)$ is a Hilbert space, therefore the Riesz Representation Theorem of a continuous linear form shows that a solution of (4) exists and is unique.
2. In this section we shall show how approximate solutions of Problem 1 can be determined in the space $W_{0}^{k, 2}(a, b)$, if a countable complete system of linear independent functions $\left\{\varphi_{n}\right\}_{n=1}^{\infty}$ is given. It will be noted that $\mathrm{cl}\left(\operatorname{lin}\left\{\varphi_{n} ; n=1,2, \ldots\right\}\right)=W_{0}^{k, 2}(a, b)$. 'Lin' denotes the lineal hull and 'cl' the closure of the set $\left\{\varphi_{n} ; n=1,2, \ldots\right\}$. Let us define the quadratic form

$$
F(u):=\frac{1}{2} \alpha(u, u)-\beta(u) \quad \text { for } u \in W_{0}^{k, 2}(a, b)
$$

From the Ritz Theorem ([1], p. 21, see also [2]) one can conclude, that $u$ is a solution of Problem 1 if and only if

$$
\begin{equation*}
F(u)=\inf _{\varphi \in W_{0}^{k, 2}(a, b)} F(\varphi) \tag{8}
\end{equation*}
$$

Assume that $E_{n}$ is a subspace of $W_{0}^{k, 2}(a, b)$ spanned by elements $\varphi_{1}, \ldots, \varphi_{n}$. Let $u_{n}:=\inf _{\varphi \in E_{n}} F(\varphi)$. It is known that $u_{n}$ tends to $u$ with respect to any of the norms considered above. Therefore the problem of determining the approximate solutions of Problem 1 reduces to the possibility of determining functions $u_{n}$. The quadratic form $F$ on the space $E_{n}$ reads as

$$
\begin{gather*}
F(\varphi)=G\left(\lambda_{1}, \ldots, \lambda_{n}\right)= \\
=\frac{1}{2}\left[\sum_{i, j=1}^{n} \sum_{s=1}^{k} \lambda_{i} \lambda_{j} a_{s}\left(\varphi_{i}^{(s)} \varphi_{j}^{(s)}\right)_{L^{2}}+\sum_{i, j=1}^{n} \lambda_{i} \lambda_{j} \int_{a}^{b} \varphi_{i} \varphi_{j} d \mu_{1}(x)\right]- \\
-\sum_{j=1}^{n} \lambda_{j} \int_{a}^{b} \varphi_{j} d \mu_{2}(x) \tag{9}
\end{gather*}
$$

where $\varphi=\lambda_{1} \varphi_{1}+\ldots+\lambda_{n} \varphi_{n}$. The expression (9) may be rewritten in matrix form

$$
G(\Lambda)=\frac{1}{2} \Lambda^{T}\left(\sum_{s=1}^{k} a_{s} \Gamma_{s}+\Delta\right) \Lambda-\Lambda^{T} P
$$

where

$$
\begin{array}{cc}
\Lambda=\left[\lambda_{1} \ldots \lambda_{n}\right]^{T}, & \Gamma_{s}=\left[\left(\varphi_{i}^{(s)}, \varphi_{k}^{(s)}\right)_{L^{2}}\right]_{i, k=1}^{n} \\
\Delta=\left[\int_{a}^{b} \varphi_{i} \varphi_{k} d \mu_{1}(x)\right]_{i, k=1}^{n}, & P=\left[\int_{a}^{b} \varphi_{1} d \mu_{2}(x) \ldots \int_{a}^{b} \varphi_{n} d \mu_{2}(x)\right]^{T} .
\end{array}
$$

It is easy to check that

$$
G\left(\Lambda^{*}\right)=\inf _{\Lambda \in \mathbb{R}^{n}} G(\Lambda)
$$

when

$$
\begin{equation*}
\left(\sum_{i=1}^{k} a_{s} \Gamma_{s}+\Delta\right) \Lambda^{*}=0 \tag{10}
\end{equation*}
$$

Equation (10) may be used for determining approximate solutions of Problem 1.
3. To solve Problem 1 we can apply any complete system of linearly independent functions in the space $W_{0}^{k, 2}(a, b)$, but for calculational reasons systems connected with Haar functions in $L^{2}(a, b)$ are convenient. Let us remind that

$$
\xi_{n}(t)= \begin{cases}2^{\frac{m}{2}} & \text { for } t \in\left[\frac{2 l-2}{2^{m+1}}, \frac{2 l-1}{2^{m+1}}\right] \\ 2^{-\frac{m}{2}} & \text { for } t \in\left(\frac{2-1}{2^{m+1}}, \frac{2 l}{2^{m+1}}\right] \\ 0 & \text { for other } t \in[0,1]\end{cases}
$$

where $n=2^{m}+l, l=1 \ldots 2^{m}, m=0,1, \ldots$ and $\xi_{1}: \equiv 1$ on $[0,1]$ are called Haar functions, they constitute a complete orthonormal system in $L^{2}(0,1)$ ( $[3]$,
p. 132). A complete orthonormal system in $\left(W_{0}^{1,2}(0,1),[.,].\right)$ was constructed in paper [2], which after differentiation gives the Haar functions $\xi_{n}, n=2,3, \ldots$ For various reasons this cannot be repeated for $k>1$. Otherwise it another construction is possible in the space $W_{0}^{k, 2}(0,1)$. Because of calculational difficulties we will discuss this for the case $k=2$ only.

To solve this problem we have to find functions in $W_{0}^{2,2}(0,1)$ which are parts of second order polynomials and are linearly independent. They determine a complete system. After double differentiation they become a linear combination of some Haar functions. Let us begin with the following two functions

$$
y_{3}(t):= \begin{cases}\frac{3 t^{2}}{2} & \text { for } t \in\left[0, \frac{1}{4}\right] \\ \frac{-5}{2}\left(t-\frac{4+\sqrt{6}}{10}\right)\left(t-\frac{4-\sqrt{6}}{10}\right) & \text { for } t \in\left(\frac{1}{4}, \frac{1}{2}\right] \\ \frac{(t-1)^{2}}{2} & \text { for } t \in\left(\frac{1}{2}, 1\right]\end{cases}
$$

and $y_{4}(t):=-y_{3}(1-t), t \in[0,1]$.
One can show that $y_{\alpha}^{\prime \prime}=-\xi_{2}+2 \sqrt{2} \xi_{\alpha}$ for $\alpha=3,4$. Now, we put for $n=2^{m+1}+s, m=0,1, \ldots, s=1, \ldots 2^{m+1}$

$$
y_{2^{m+1}+s}:= \begin{cases}2^{-\frac{3 m}{2}} y_{\alpha}\left[2^{m}\left(t-\frac{4\left\lfloor\frac{s+1}{2}\right\rfloor-4}{2^{m+2}}\right)\right] & \text { for } t \in\left[\frac{4\left\lfloor\frac{s+1}{2}\right\rfloor-4}{2^{m+2}}, \frac{4\left\lfloor\frac{s+1}{2}\right\rfloor}{2^{m+2}}\right]  \tag{11}\\ 0 & \text { for other } t \in[0,1]\end{cases}
$$

where $\alpha=3$ if $s$ is odd and $\alpha=4$ if $s$ is even. By $\lfloor\cdot\rfloor$ the entire part function is denoted. For $t \in[0,1]$ the following formula holds

$$
\begin{gather*}
y_{2^{m+1}+s}^{\prime \prime}(t)=2^{\frac{m}{2}} y_{\alpha}^{\prime \prime}\left[2^{m}\left(t-\frac{4\left\lfloor\frac{s+1}{2}\right\rfloor-4}{2^{m+2}}\right)\right] \\
=2^{\frac{m}{2}}\left\{-\xi_{2}\left[2^{m}\left(t-\frac{4\left\lfloor\frac{s+1}{2}\right\rfloor-4}{2^{m+2}}\right)\right]+2 \sqrt{2} \xi_{\alpha}\left[2^{m}\left(t-\frac{4\left\lfloor\frac{s+1}{2}\right\rfloor-4}{2^{m+2}}\right)\right]\right\} \\
=-\xi_{2^{m}+\left\lfloor\frac{s+1}{2}\right\rfloor}+2 \sqrt{2} \xi_{2^{m+1}+s} \tag{12}
\end{gather*}
$$

Theorem $3 \mathrm{cl}\left(\operatorname{lin}\left\{y_{n}^{\prime \prime} ; n=1,2, \ldots\right\}\right)=L^{2}(0,1)$, where $\hat{y}_{1}^{\prime \prime}:=\xi_{1}$ and $y_{2}^{\prime \prime}:=\xi_{2}$.
Proof We shall show that for every $m, m=0,1, \ldots$ the functions $y_{1}^{\prime \prime}, \ldots y_{2_{m+2}}^{\prime \prime}$ are linearly independent. Let us assume that $\sum_{i=1}^{2^{m+2}} a_{i} y_{i}^{\prime \prime}=0$ for some $a_{i}, i=$ $1 \ldots 2^{m+2}$. By (12) we obtain

$$
\begin{gathered}
0=\sum_{i=1}^{2^{m+1}} a_{i} y_{i}^{\prime \prime}+\sum_{s=1}^{2^{m+1}} a_{2^{m+1}+s} y^{\prime \prime}{ }_{2^{m+1}+s} \\
=\sum_{i=1}^{m} \sum_{j=1}^{2^{i}} a_{2^{i}+j} y^{\prime \prime}{ }_{2^{i}+j}-\sum_{s=1}^{2^{m+1}} a_{2^{m+1}+s} \xi_{2^{m}+\left\lfloor\frac{s+1}{2}\right\rfloor}+2 \sqrt{2} \sum_{s=1}^{2^{m+1}} a_{2^{m+1}+s} \xi_{2^{m+1}+s}
\end{gathered}
$$

The above expression is a linear combination of the Haar functions and therefore we have $a_{2^{m+1}+s}=0$ for $s=1, \ldots 2^{m+1}$. This implies that $a_{2^{m}+s}=0$ for
$s=1, \ldots 2^{m}$. Continuing the above consideration we can conclude that $a_{i}=0$ for $i=1, \ldots 2^{m+2}$. By induction it can be shown that for $m=0,1, \ldots$, $s=1, \ldots 2^{m+1}$ we have
$\xi_{2^{m+1}+s}=\frac{1}{2 \sqrt{2}} y^{\prime \prime}{ }_{2^{m+1}+s}+\sum_{l=0}^{m-1}\left(\frac{1}{2 \sqrt{2}}\right)^{l+2} y_{2^{\prime \prime}}{ }^{m-l}+\left\lfloor\frac{s+2^{l+1}-1}{2^{l+1}}\right\rfloor+\left(\frac{1}{2 \sqrt{2}}\right)^{m+1} \xi_{2}$.
Therefore every Haar function may be represented as a combination of some functions $y_{n}^{\prime \prime}$. Let $f \in L^{2}(0,1)$. The function $f$ has the Fourier representation

$$
\begin{equation*}
f=\sum_{n=1}^{\infty} c_{n} \xi_{n}=c_{1} \xi_{1}+c_{2} \xi_{2}+\lim _{n \rightarrow \infty} \sum_{m=0}^{2^{n+2}} \sum_{s=1}^{2^{m+1}} c_{2^{m+1}+s} \xi_{2^{m+1}+s} \tag{14}
\end{equation*}
$$

By (13) we obtain

$$
\begin{equation*}
f=c_{1} \xi_{1}+\lim _{n \rightarrow \infty}\left[\sum_{m=3}^{2^{n+2}} \tilde{c}_{m}^{n} y^{\prime \prime}{ }_{m}+A_{n} \xi_{2}\right] \tag{15}
\end{equation*}
$$

where $\tilde{c}_{n}^{m}$ may be determined by formulas (13), (14) and

$$
\begin{equation*}
A_{n}:=c_{2}+\sum_{m=0}^{2^{n+2}} \frac{1}{(2 \sqrt{2})^{m+1}} \sum_{s=1}^{2^{m+1}} c_{2^{m+1}+s}, \quad n=0,1, \ldots \tag{16}
\end{equation*}
$$

This completes the proof of the theorem.
Theorem $4 \mathrm{cl}\left(\operatorname{lin}\left\{y_{n} ; n=3,4 \ldots\right\}\right)=W_{0}^{2,2}(0,1)$.
Proof In virtue of Theorem 1 the space $W_{0}^{2,2}(0,1)$ with the norm ||| $\cdot||\mid$ may be considered. Assume that a function $f \in W_{0}^{2,2}(0,1)$ is given. Then $f^{\prime \prime} \in$ $L^{2}(0,1)$ and the function $f^{\prime \prime}$ may be represented by formulas (14)-(16). Clearly $c_{1}=\int_{0}^{1} f^{\prime \prime} \xi_{1} d x=0$. We shall show that $A:=\lim _{n \rightarrow \infty} A_{n}=0$. It is easy to check that $\int_{0}^{1}\left\{\int_{0}^{t} \xi_{2^{n}+l}(s) d s\right\} d t=\frac{1}{2^{\frac{3 n}{2}+2}}$ for $n=0,1 \ldots, l=1, \ldots 2^{n}$. By (14) integrating two times we get

$$
\begin{aligned}
0 & =\int_{0}^{1}\left(\int_{0}^{t} f^{\prime \prime}(s) d s\right) d t=\sum_{n=2}^{\infty} c_{n} \int_{0}^{1}\left\{\int_{0}^{t} \xi_{n}(s) d s\right\} d t \\
& =\sum_{n=0}^{\infty} \sum_{l=1}^{2^{n}} \frac{c_{2^{n}+l}}{2^{\frac{3 n}{2}+2}}=\frac{c_{2}}{4}+\frac{1}{4} \sum_{m=0}^{\infty} \sum_{l=1}^{2^{m+1}} \frac{c_{2^{m+1}+l}}{2^{\frac{3(m+1)}{2}}}=\frac{A}{4}
\end{aligned}
$$

Now, we can put

$$
\tilde{f}:=\lim _{m \rightarrow \infty}\left[\sum_{i=3}^{2^{m+2}} \tilde{c}_{i}^{m} y_{i}\right], \quad \tilde{f} \in W_{0}^{2,2}(0,1)
$$

This means that

$$
\tilde{f^{\prime \prime}}=\lim _{m \rightarrow \infty}\left[\sum_{i=3}^{2^{m+2}} \tilde{c}_{i}^{m} y^{\prime \prime}{ }_{i}\right]=f^{\prime \prime}
$$

a.e. on $[0,1]$. Hence $f(t)=\tilde{f}(t)$ for each $t \in[0,1]$. This finishes the proof.

Remark 1 Let us mention that using $y_{n}, n=3,4$ one can define a complete orthonormal system in $\left(W_{0}^{2,2}(0,1),[.,].\right)$ by putting

$$
\begin{aligned}
\Psi_{2^{m+1}+2 s-1} & :=\frac{1}{2 \sqrt{5}}\left(y_{2^{m+1}+2 s-1}+y_{2^{m+1}+2 s}\right) \\
\Psi_{2^{m+1}+2 s}: & =\frac{1}{4}\left(y_{2^{m+1}+2 s-1}-y_{2^{m+1}+2 s}\right)
\end{aligned}
$$

for $s=1, \ldots 2^{m}, m=0,1, \ldots$
Remark 2 Similarly one can construct a countable complete system of linearly independent functions in the space $W_{0}^{m, 2}(0,1)$ consisting of functions which are parts of polynomials of order $m$ and in some joint points $t_{1}, \ldots, t_{m}$ are at least of class $C^{m-1}$. One can start with function

$$
y(t)= \begin{cases}\frac{a_{0} t^{m}}{m!} & \text { for } t \in\left[0, t_{1}\right]  \tag{17}\\ \frac{a_{i, 1} t^{m}}{m!}+\cdots+a_{i, m+1} & \text { for } t \in\left(t_{i}, t_{i+1}\right], i=1, \ldots m-1 \\ \frac{a_{m}(t-1)^{m}}{m!} & \text { for } t \in\left(t_{m}, 1\right]\end{cases}
$$

which can be always determined (homogeneous system of $m^{2}$ equations with $(m-1)(m+1)+2=m^{2}+1$ unknowns). We assume that $m=2^{n}+l, n=$ $0,1, \ldots, l=0, \ldots\left(2^{n}-1\right)$. Taking $t_{i}, i=1, \ldots m$ from among points $\frac{s}{2^{n+1}}$, $s=1, \ldots\left(2^{n+1}-1\right)$ it can be defined $2^{n+1}-m=2^{n}-l$ linear independent functions according to formula (17). Their $m$-th derivative is a constant in the intervals determined by points $t_{i}, i=1, \ldots m$. Thus, they are a linear combination of Haar functions $\xi_{1}, \ldots \xi_{2^{n+1}}$ (see [3], p. 133). Applying linear transformation appearing in formula (11) one can formulate a complete system of linear independent functions in $W_{0}^{m, 2}(0,1)$ analogously as for $m=2$.
Example 1 For differential equation

$$
\begin{gathered}
y^{(4)}+16 \delta_{\frac{1}{2}} y=1 \\
y(0)=y(1)=y^{\prime}(0)=y^{\prime}(1)=0
\end{gathered}
$$

where $\delta_{\frac{1}{2}}$ is the Dirac measure concentrated at point $t=\frac{1}{2}$ we obtain the exact solution

$$
y(t)= \begin{cases}\frac{1}{24} t^{4}-\frac{49}{600} t^{3}+\frac{97}{2400} t^{2} & \text { for } t \in\left[0, \frac{1}{2}\right] \\ \frac{1}{24} t^{4}-\frac{49}{600} t^{3}+\frac{97}{2400} t^{2}-\frac{1}{300}\left(t-\frac{1}{2}\right)^{3} & \text { for } t \in\left(\frac{1}{2}, 1\right] .\end{cases}
$$

That is $y\left(\frac{1}{2}\right)=\frac{1}{400}=2.5 \cdot 10^{-4}$. Approximate solutions in spaces $E_{2}, E_{6}$ and $E_{14}$ calculated numerically by the method described in this paper have the following values $1.83 \cdot 10^{-4}, 2.23 \cdot 10^{-4}$ and $2.36 \cdot 10^{-4}$ at point $t=\frac{1}{2}$.

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