$u(x, t),(x, t) \in \bar{Q}$, that satisfies (in some sense) the homogeneous wave equation
$\frac{\partial^{2} u(x, t)}{\partial t^{2}}-\Delta u(x, t)=0, \quad(x, t) \in Q$,
homogeneous initial conditions
$u(x, 0)=0, \quad \frac{\partial u(x, 0)}{\partial t}=0, \quad x \in \Omega$,
and the Neumann boundary condition
$\partial_{\boldsymbol{v}(x)} u(x, t)=g(x, t), \quad(x, t) \in \Sigma$.
Here $\Delta$ is the Laplacian and $\partial_{v}$ denotes the normal derivative.
To solve the problem (1)-(3) we use the retarded double layer potential
(D) $\lambda$ ) $(x, t):=\left.\frac{1}{4 \pi} \int_{\Gamma} \partial_{\boldsymbol{v}(y)}\left(\frac{\lambda(z, t-|x-y|)}{|x-y|}\right)\right|_{z=y} d \Gamma_{y},(x, t) \in Q$,
where $\lambda: \Gamma \times \mathbb{R} \rightarrow \mathbb{R}$ is a density. It is known (see, e.g., Bamberger and Ha-Duong, 1986a or Polozhyy, 1964) that if an arbitrary function $\lambda(y, \tau),(y, \tau) \in \Gamma \times \mathbb{R}$, is smooth enough and $\lambda(y, \tau)=$ 0 when $y \in \Gamma, \tau \leq 0$, then function
$u(x, t):=(\mathcal{D} \lambda)(x, t), \quad(x, t) \in Q$,
satisfies (in the classical sense) the wave equation and initial conditions (2).

Note that for a sufficiently smooth density $\lambda$ and surface $\Gamma$ we can express the operator $\partial_{\boldsymbol{v}}$ as $\partial_{\boldsymbol{v}(x)} u(x, \cdot)=\boldsymbol{v}(x) \cdot \nabla_{x} u(x, \cdot)$, where $\nabla_{x}$ is the gradient operator. Then there exists the following limit
$(\mathcal{W} \lambda)(x, t):=$
$\left.\frac{1}{4 \pi} \boldsymbol{v}(x) \cdot \lim _{x^{\prime} \rightarrow x} \nabla_{x^{\prime}} \int_{\Gamma} \boldsymbol{v}(y) \cdot \nabla_{y}\left(\frac{\lambda\left(z, t-\left|x^{\prime}-y\right|\right)}{\left|x^{\prime}-y\right|}\right)\right|_{z=y} d \Gamma_{y}$,
where $x^{\prime}:=x-\varepsilon \boldsymbol{v}(x) \in \Omega, \varepsilon>0$, denotes a point close to the point $x \in \Gamma$. We say that $x^{\prime}$ approaches $x, x^{\prime} \rightarrow x$, when $\varepsilon \rightarrow 0$. The function $u$ satisfies the boundary condition (3), if the function $\lambda$ is a solution of the retarded potential integral equation (RPIE)
$(\mathcal{W} \lambda)(x, t)=g(x, t),(x, t) \in \Sigma$.
We shall briefly give the essential notions of the Laguerre transform in the weighted Lebesgue space $L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$ with some parameter $\sigma>0$ and weight $\rho_{\sigma}(t)=e^{-\sigma t}, t \in \mathbb{R}_{+}$. Here $X$ is a Hilbert space with an inner product $(\because \cdot)_{X}$ and an inducted norm $\|\cdot\|_{X}$. Elements $v \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$ are measurable functions $v: \mathbb{R}_{+} \rightarrow X$ such that $\int_{\mathbb{R}_{+}}\|v(t)\|_{X}^{2} e^{-\sigma t} d t<\infty$. This space is equipped with the inner product
$(v, w)_{L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)}:=\int_{\mathbb{R}_{+}}(v(t), w(t))_{X} e^{-\sigma t} d t, v, w \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$,
and the norm
$\|v\|_{L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)}:=\sqrt{(v, v)_{L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)}}, \quad v \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$.
It is well-known (Reed and Simon, 1977) that the space $L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$ is complete. We will assume that the elements of space $L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$ are extended with zero for non-positive arguments.

For any $m \in \mathbb{N}$ (set of natural numbers) let us denote the weighted Sobolev space as
$H_{\sigma}^{m}\left(\mathbb{R}_{+} ; X\right):=\left\{v \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right) \mid v^{(k)} \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right), k=\overline{1, m}\right\}$
with the norm $\|v\|_{H_{\sigma}^{m}\left(\mathbb{R}_{+} ; X\right)}:=\left(\sum_{k=0}^{m}\left\|v^{(k)}\right\|_{L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)}^{2}\right)^{1 / 2}$. Here derivatives $v^{k}(k \in \mathbb{N})$ are understood in terms of the space $\mathcal{D}^{\prime}\left(\mathbb{R}_{+} ; X\right)$, elements of which are distributions with values in the space $X$.

We introduce a couple of notations. As the sequence of elements of set $X$ we understand a vector-column $\boldsymbol{v}:=\left(v_{0}, v_{1}, \ldots\right)^{T}$. All possible sequences of elements of the set $X$ are denoted by $X^{\infty}$. We consider the Hilbert space $l^{2}(X):=\left\{\boldsymbol{v} \in X^{\infty} \mid \quad \sum_{j=0}^{\infty}\left\|v_{j}\right\|_{X}^{2}<+\infty\right\}$ with the inner product $(\boldsymbol{v}, \boldsymbol{w})=\sum_{j=0}^{\infty}\left(v_{j}, w_{j}\right)_{X}, \boldsymbol{v}, \boldsymbol{w} \in l^{2}(X)$, and the norm $\|\boldsymbol{v}\|_{l^{2}(X)}:=$ $\left(\sum_{j=0}^{\infty}\left\|v_{j}\right\|_{X}^{2}\right)^{1 / 2}, v \in l^{2}(X)$. Recall that for $X=\mathbb{R}$ we have $l^{2}(\mathbb{R})=l^{2}:=\left\{\left.\boldsymbol{v} \in \mathbb{R}^{\infty}\left|\quad \sum_{j=0}^{\infty}\right| v_{j}\right|^{2}<+\infty\right\}$.

Now let us give the definition of the Laguerre transform and outline some of its properties (Litynskyy and Muzychuk, 2015b). Consider a mapping $\mathcal{L}: L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right) \rightarrow X^{\infty}$ which operates according to the rule
$f_{k}:=\sigma \int_{\mathbb{R}_{+}} f(t) L_{k}(\sigma t) e^{-\sigma t} d t, \quad k \in \mathbb{N}_{0}:=\mathbb{N} \cup\{0\}$,
where $\left\{L_{k}(\sigma \cdot)\right\}_{k \in \mathbb{N}_{0}}$ are Laguerre polynomials, which form orthogonal basis in the space $L_{\sigma}^{2}\left(\mathbb{R}_{+}\right)$(Keilson et al., 1980). Also consider the mapping $\mathcal{L}^{-1}: l^{2}(X) \rightarrow L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$, which maps an arbitrary sequence $\boldsymbol{h}=\left(h_{0}, h_{1}, \ldots, h_{k}, \ldots\right)^{T}$ to a function
$h(t):=\left(\mathcal{L}^{-1} \boldsymbol{h}\right)(t)=\sum_{k=0}^{\infty} h_{k} L_{k}(\sigma t), t \in \mathbb{R}_{+}$.

Proposition 2.1 (theorem 2 from Litynskyy and Muzychuk, 2015b) The mapping $\mathcal{L}: L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right) \rightarrow X^{\infty}$ that maps the arbitrary function $f$ to the sequence $\boldsymbol{f}=\left(f_{0}, f_{1}, \ldots, f_{k}, \ldots\right)^{T}$ according to the formula(10), is injective and its image is the space $l^{2}(X)$, and
$\|f\|_{L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)}^{2}=\frac{1}{\sigma} \sum_{k=0}^{\infty}\left\|f_{k}\right\|_{X}^{2}$.
In addition, for the arbitrary function $f \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$ we have an equality
$\mathcal{L}^{-1} \mathcal{L} f=f$,
where the mapping $\mathcal{L}^{-1}: l^{2}(X) \rightarrow L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$ is the inverse to $\mathcal{L}$ and maps the arbitrary sequence $\boldsymbol{h}=\left(h_{0}, h_{1}, \ldots, h_{k}, \ldots\right)^{T}$ to the function $h$ according to the formula (11).

Definition 2.2 (Litynskyy and Muzychuk, 2015b) Let $\sigma>0$ and $X$ be a Hilbert space. Mappings $\mathcal{L}: L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right) \rightarrow$ $l^{2}(X)$ and $\mathcal{L}^{-1}: l^{2}(X) \rightarrow L_{\sigma}^{2}\left(\mathbb{R}_{+} ; X\right)$, mentioned in Proposition 2.1, are called, respectively, direct and inverse Laguerre transforms, and the formula (12) is an analogue of the Parseval equality.

Definition 2.3 (Litynskyy et al., 2009) Let $X, Y, Z$ be arbitrary sets and $q: X \times Y \rightarrow Z$ be some mapping. By a $q$-convolution of sequences $\boldsymbol{u} \in X^{\infty}$ and $\boldsymbol{v} \in Y^{\infty}$ we understand the sequence $\boldsymbol{w}:=\left(w_{0}, w_{1}, \ldots, w_{j}, \ldots\right)^{T} \in Z^{\infty}$, whose elements are obtained by the rule
$w_{j}:=\sum_{i=0}^{j} q\left(u_{j-i}, v_{i}\right) \equiv \sum_{i=0}^{j} q\left(u_{i}, v_{j-i}\right), \quad j \in \mathbb{N}_{0} ;$
the $q$-convolution of $\boldsymbol{u}$ and $\boldsymbol{v}$ is shortly written in the form $\boldsymbol{w}=$ $\boldsymbol{u}{ }_{q} \boldsymbol{v}$.

If $X=\mathcal{L}(Y, Z)$ is the space of linear operators acting from the space $Y$ into the space $Z$ and $q(A, v):=A v, A \in \mathcal{L}(Y, Z), v \in Y$, then for components of the q-convolution of arbitrary sequences
$\boldsymbol{A} \in(\mathcal{L}(Y, Z))^{\infty}$ and $\boldsymbol{v} \in Y^{\infty}$ we will have the following formula
$w_{j}=\sum_{i=0}^{j} A_{j-i} v_{i}, \quad j \in \mathbb{N}_{0}$,
and will write $\boldsymbol{w}:=\boldsymbol{A}{ }_{Z} \boldsymbol{v}$.
Now let's consider a sequence of functions
$e_{0}(z):=\frac{e^{-\sigma|z|}}{4 \pi|z|}, \quad e_{k}(z):=\frac{e^{-\sigma|z|}}{4 \pi|z|}\left(L_{k}\left(\sigma|z|-L_{k-1}(\sigma|z|)\right), \quad k \in \mathbb{N}\right.$,
$z \in \mathbb{R}^{3} \backslash\{0\}$.
Based on the above, in the space $H^{1}(\Omega):=\left\{v \in L^{2}(\Omega)| | \nabla v \mid \in\right.$ $\left.L^{2}(\Omega)\right\}$ we can define a function sequence
$\boldsymbol{u}(x):=\left(\boldsymbol{D} \underset{H^{1}(\Omega)}{\circ} \boldsymbol{\lambda}\right)(x), \quad x \in \Omega$,
where $\lambda=\mathcal{L} \lambda$ for any $\lambda \in L_{\sigma}^{2}\left(\mathbb{R}_{+} ; H^{1 / 2}(\Gamma)\right)$ and the sequence $\boldsymbol{D}$ is composed of operators $D_{k}: H^{1 / 2}(\Gamma) \rightarrow H^{1}(\Omega, \Delta), k \in \mathbb{N}_{0}$, given by the rule
$\left(D_{k} \xi\right)(x):=\frac{1}{4 \pi} \int_{\Gamma} \xi(y) \boldsymbol{v}(y) \cdot \nabla_{y} e_{k}(x-y) d \Gamma_{y}$.
Here $H^{1}(\Omega, \Delta):=\left\{v \in H^{1}(\Omega) \mid \Delta v \in L^{2}(\Omega)\right\}, H^{1 / 2}(\Gamma)$ denotes a space of traces of elements of $H^{1}(\Omega)$ on the surface $\Gamma$ and $H^{-1 / 2}(\Gamma):=\left(H^{1 / 2}(\Gamma)\right)^{\prime}$. If $u$ is expressed by the retarded double layer potential (4) with some density $\lambda$ then the sequence (17) represents the transformation $\boldsymbol{u}=\mathcal{L} u$ (Litynskyy and Muzychuk, 2016). Similarly, applying the LT to the equation (7), we obtain a BIE system
$\boldsymbol{W}_{H^{-1 / 2}(\Gamma)}^{\circ} \boldsymbol{\lambda}=\boldsymbol{g}$ in $l^{2}\left(H^{-1 / 2}(\Gamma)\right)$,
where $\boldsymbol{g}=\mathcal{L} g$ and $\boldsymbol{W}: l^{2}\left(H^{1 / 2}(\Gamma)\right) \rightarrow l^{2}\left(H^{-1 / 2}(\Gamma)\right)$ is a sequence of boundary operators
$\left(W_{k} \xi\right)(x):=\frac{1}{4 \pi} \boldsymbol{v}(x) \cdot \lim _{x^{\prime} \rightarrow x} \nabla_{x^{\prime}} \int_{\Gamma} \xi(y) \boldsymbol{v}(y) \cdot \nabla_{y} e_{k}\left(x^{\prime}-y\right) d \Gamma_{y}$,
$k \in \mathbb{N}_{0}$.
After finding the solution $\lambda=\left(\lambda_{0}, \lambda_{1}, \ldots\right)^{T}$ of the BIEs (19), the generalized solution of the problem (1)-(3) can be presented as a sum of the series
$u(x, t)=\sum_{k=0}^{\infty}\left(\sum_{i=0}^{k} D_{k-i} \lambda_{i}(x)\right) L_{k}(\sigma t), \quad(x, t) \in Q$.

Proposition 2.4 (theorem 2.4 from Litynskyy and Muzychuk, 2016) Let $g \in H_{\sigma_{0}}^{m+4}\left(\mathbb{R}_{+} ; H^{-1 / 2}(\Gamma)\right)$ for some $\sigma_{0}>0$ and $m \in \mathbb{N}_{0}$. Then there exists a unique generalized solution of the problem (1)-(3), it belongs to the space $H_{\sigma_{0}}^{m+1}\left(\mathbb{R}_{+} ; H^{1}(\Omega)\right)$ and for any $\sigma \geq \sigma_{0}$ such an inequality holds
$\|u\|_{H_{\sigma}^{m+1}\left(\mathbb{R}_{+} ; H^{1}(\Omega)\right)} \leq C\|g\|_{H_{\sigma}^{m+4}\left(\mathbb{R}_{+} ; H^{-1 / 2}(\Gamma)\right)}$,
where $C>0$ is a constant that is not dependent on $g$.
In addition, the generalized solution of the problem (1)-(3) can be obtained by the inverse transform $u=\mathcal{L}^{-1} \boldsymbol{u}$, where $u_{j} \in$ $H^{1}(\Omega, \Delta)\left(j \in \mathbb{N}_{0}\right)$ are the corresponding components of the $q$ convolution (17), and elements of the sequence $\lambda \in l^{2}\left(H^{1 / 2}(\Gamma)\right)$ are solutions of BIE system (19), in which $\boldsymbol{g}=\mathcal{L} g$.

## 3. BEM FOR THE INFINITE BIE SYSTEM

We have now the new representation (21) of the solution of the problem (1)-(3) and the infinite BIE system (19) with unknown functions $\lambda=\left(\lambda_{0}, \lambda_{1}, \ldots\right)^{T}$. It is easy to see that the sys-
tem (19) can be rewritten as a sequence of BIEs

with recursive right-hand sides
$\tilde{g}_{k}:=g_{k}-\sum_{i=0}^{k-1} W_{k-i} \lambda_{i}, \quad k \in \mathbb{N}$.
For every $k \in \mathbb{N}_{0}$ the corresponding $k$-th equation of (23) is a hypersingular equation that has the form
$W_{0} \eta=f$ in $H^{-1 / 2}(\Gamma)$.
It has a unique solution $\eta \in H^{1 / 2}(\Gamma)$ for an arbitrary function $f \in H^{-1 / 2}(\Gamma)$ (Hsiao and Wendland, 2008). We can choose (by some criteria) value $N$ and find from (23) the first components $\lambda_{0}, \lambda_{1}, \ldots, \lambda_{N}$. Then the approximate solution of the problem (1)-
(3) is the partial sum
$u^{N}(x, t)=\sum_{j=0}^{N}\left(\sum_{i=0}^{k} D_{k-i} \lambda_{i}(x)\right) L_{k}(\sigma t), \quad(x, t) \in Q$.
We are now in position to apply the BEM for finding unknown functions. Consider in $H^{1 / 2}(\Gamma)$ a sequence of finite-dimensional subspaces $X^{M} \subset H^{1 / 2}(\Gamma), M \in N$, assuming that $\left\{\phi_{i}\right\}_{i=1}^{M}$ is a basis of $X^{M}$. The numerical solution of the equation (25) can be presented as a linear combination
$\eta^{M}:=\sum_{i=1}^{M} \eta_{i} \phi_{i} \in X^{M}$
that is a solution of such variational equation
$\left\langle W_{0} \eta^{M}, \eta\right\rangle_{\Gamma}=\langle f, \eta\rangle_{\Gamma} \quad \forall \eta \in X^{M}$.
Applying the Galerkin method, that is taking the elements of the basis $\phi_{j}$ as test functions in order to find the vector of unknown coefficients $\boldsymbol{\eta}^{[M]}:=\left\{\eta_{i}\right\}_{i=1}^{M} \in \mathbb{R}^{M}$, we will get a system of linear algebraic equations (SLAE)
$W_{0}{ }^{[M]} \boldsymbol{\eta}^{[M]}=\boldsymbol{f}^{[M]}$.
where $W_{0}^{[M]}[j, i]:=\left\langle W_{0} \phi_{i}, \phi_{j}\right\rangle_{\Gamma}, \quad f_{j}^{[M]}:=\left\langle f, \phi_{j}\right\rangle_{\Gamma}, \quad i, j=\overline{1, M}$. As a result of the $H^{1 / 2}(\Gamma)$-ellipticity of the operator $W_{0}$, the matrix $W_{0}^{[M]}$ is positive definite (Costabel M., 1988; Hsiao and Wedland, 1977; Schtainbah, 2008). Therefore, the system (29) has a unique solution for an arbitrary right-hand side $\forall M \in N$.

Let $\Gamma_{\widetilde{M}}=\cup_{l=1}^{\widetilde{M}} \bar{\tau}_{l}$ be some approximation of the boundary $\Gamma$ composed of triangular boundary elements $\left\{\tau_{l}\right\}_{l=1}^{\tilde{M}}$ with vertices $\left\{x^{\left[l_{1}\right]}, x^{\left[l_{2}\right]}, x^{\left[l_{3}\right]}\right\}$. We assume that vertices of all triangles have a global numbering $\left\{x_{k}\right\}_{k=1}^{M^{*}}$ and for each point $x_{k}$ there exists an associated set $\mathcal{J}(k)$ of numbers of those triangles that have this point as a vertex. We treat the value $h:=\max _{l=1, \tilde{M}}\left(\int_{\tau_{l}} d s\right)^{1 / 2}$ as the parameter of the spatial approximation.

Note that each triangle can be projected on a "standard" triangle $\tau:=\left\{\xi:=\left(\xi_{1}, \xi_{2}\right) \in R^{2}: 0<\xi_{1}<1,0<\xi_{2}<1-\xi_{1}\right\}$. Following Dautray and Lions (1992) and Schtainbah (2008), functions $\phi_{1}^{1}(\xi)=1-\xi_{1}-\xi_{2}, \phi_{2}^{1}(\xi)=\xi_{1}$ and $\phi_{3}^{1}(\xi)=\xi_{2}$, defined locally on the triangle $\tau$, form a set $\left\{\varphi_{i}^{1}\right\}_{i=1}^{M}, M=M^{*}$, that contains linearly-independent on $\Gamma_{\widetilde{M}}$ functions. Moreover, $\operatorname{supp} \varphi_{i}^{1}=$ $\mathrm{U}_{l \in \mathcal{J}(i)} \bar{\tau}_{l}=: \tau_{i}^{*}$. Then the numerical solution $\lambda^{h}:=\left(\lambda_{0}^{h}, \lambda_{1}^{h}, \ldots\right)^{T}$ of the BIE system (23) is sought in the form
$\lambda_{k}^{h}=\sum_{l=1}^{M} \lambda_{k, l}^{h} \varphi_{l}^{1} \quad \in S_{h}^{1}(\Gamma):=\operatorname{span}\left\{\varphi_{i}^{1}\right\}_{i=1}^{M}, k \in N_{0}$,

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where $\left\{\lambda_{k, l}^{h}\right\}_{l=1}^{M}$ are the unknown coefficients. In order to find the vector $\lambda_{k}^{h}:=\left\{\lambda_{k, l}^{h}\right\}_{l=1}^{M} \in \mathbb{R}^{M}$ we obtain the following SLAE from (29)
$\boldsymbol{W}_{0}^{h} \lambda_{k}^{h}=\boldsymbol{g}_{k}^{h}-\sum_{i=0}^{k-1} \boldsymbol{W}_{k-i}^{h} \lambda_{i}^{h}, \quad k \in N$,
where
$W_{j}^{h}[i, l]=\int_{\tau_{i}^{*}} \varphi_{i}^{1}(x) \partial_{v(x)} \int_{\tau_{l}^{*}} \varphi_{l}^{1}(y) \partial_{v(y)} E_{j}(x-y) d s_{y} d s_{x}, j=$
$\overline{0, k}$,
$g_{k}^{h}[i]=\int_{\tau_{i}^{*}} \varphi_{i}^{1}(x) \tilde{g}_{k}(x) d s_{x}, \quad i, l=\overline{1, M}$.
After finding the consequent vector $\lambda_{\mathrm{k}}^{\mathrm{h}}$ we can compute the corresponding element of the sequence $\boldsymbol{u}^{h}:=\left(u_{0}^{h}, u_{1}^{h}, \ldots, u_{k}^{h}, \ldots, u_{N}^{h}, 0, \ldots\right)^{T}$ :
$u_{k}^{h}(x)=\sum_{j=0}^{k}\left(D_{j} \lambda_{k-j}^{h}\right)(x), x \in \Omega$.
Having $\boldsymbol{u}^{h}$, the numerical solution of the Neumann problem is found by the formula
$u^{N, h}(x, t)=\sum_{j=0}^{N} u_{k}^{h}(x) L_{k}(\sigma t), \quad(x, t) \in Q$.
Let us obtain, following Hsiao and Wedland (1977), an a priory error estimate of the numerical solution after the introduction of Sobolev spaces of piecewise-smooth on the boundary $\Gamma$ functions. Let $\Gamma$ be a union $\Gamma=\mathrm{U}_{i=1}^{\widetilde{N}} \bar{\Gamma}_{i}$ of surfaces $\Gamma_{i}\left(\Gamma_{i} \cap \Gamma_{j}=\varnothing\right.$ when $i \neq j$ ), each of which has a sufficiently smooth parameterization $\Gamma_{i}:=\left\{x \in R^{3}: x=\tilde{\chi}_{i}(\xi), \xi \in \tilde{\tau}_{i} \subset R^{2}\right\}$. Then, using the set of the non-negative functions $\phi_{i} \in C_{0}^{\infty}\left(R^{3}\right)$ such that $\sum_{i=1}^{\widetilde{N}} \phi_{i}(x)=1 \quad \forall x \in \Gamma, \quad \phi_{i}(x)=0 \quad \forall x \in \Gamma \backslash \Gamma_{i}$, an arbitrary piecewise-smooth function can be given in a form
$v(x)=\sum_{i=1}^{\widetilde{N}} \phi_{i}(x) v(x)=\sum_{i=1}^{\widetilde{N}} v_{i}(x) \quad \forall x \in \Gamma$,
where $v_{i}(x):=\phi_{i}(x) v(x) \quad \forall x \in \Gamma_{i}$. Taking into account the parameterization of the fragments $\Gamma_{i}$, we consider Sobolev spaces $H^{s}\left(\tilde{\tau}_{i}\right)$, elements of which are functions $\tilde{v}_{i}(\xi):=v_{i}\left(\tilde{\chi}_{i}(\xi)\right)$ when $\xi \in \tilde{\tau}_{i}$, with corresponding norms and semi-norms
$\left\|\tilde{v}_{i}\right\|_{H^{s}\left(\tilde{\tau}_{i}\right)}:=\left(\sum_{|\alpha| \leq m}\left\|\partial^{\alpha} \tilde{v}_{i}\right\|_{L^{2}\left(\tilde{\tau}_{i}\right)}^{2}\right)^{\frac{1}{2}}$,
$\left|\tilde{v}_{i}\right|_{H^{s}\left(\tilde{\tau}_{i}\right)}:=\left(\sum_{|\alpha|=m}\left|\partial^{\alpha} \tilde{v}_{i}\right|_{L^{2}\left(\tilde{\tau}_{i}\right)}^{2}\right)^{\frac{1}{2}}, \quad s=m \in N ;$
$\left|\tilde{v}_{i}\right|_{H^{s}\left(\tilde{\tau}_{i}\right)}:=\left(\sum_{|\alpha|=m} \int_{\tilde{\tau}_{i}} \int_{\tilde{\tau}_{i}} \frac{\left|\partial^{\alpha} \tilde{v}_{i}(\xi)-\partial^{\alpha} \tilde{v}_{i}(\eta)\right|^{2}}{|\xi-\eta|^{2+2 \sigma}} d s_{\xi} d s_{\eta}\right)^{\frac{1}{2}}$,
$s=m+\sigma, \sigma \in(0,1)$.
Here $\partial^{\alpha}$ is a notation of the partial derivative with a multi-index $\alpha=\left(\alpha_{1}, \alpha_{2}\right)$. Moreover, for functions that are defined on the whole boundary $\Gamma$, we use the semi-norm
$|v|_{H^{s}(\Gamma)}:=\left(\sum_{i=1}^{\widetilde{N}}\left|\tilde{v}_{i}\right|_{H^{s}\left(\tilde{\tau}_{i}\right)}^{2}\right)^{1 / 2}$.
Lemma 3.1 Let $\lambda \in\left(H^{s}(\Gamma)\right)^{\infty}$ for some $\boldsymbol{s} \in\left[\frac{1}{2}, 2\right]$ be the exact solution of the BIE system (23), which satisfies an inequality
$\sum_{j=0}^{\infty}\left|\lambda_{j}\right|_{H^{s}(\Gamma)}<+\infty$.
Then for any values of parameters $N \in \mathbb{N}_{0}$ and $T \in \mathbb{R}_{+}$the following error estimates hold
$\left\|\lambda^{N}(\cdot, t)-\lambda^{N, h}(\cdot, t)\right\|_{H^{\frac{1}{2}}(\Gamma)} \leq \tilde{C}_{N, T} h^{s-\frac{1}{2}} \sum_{k=0}^{N}\left|\lambda_{k}\right|_{H^{s}(\Gamma)}$,
$t \in(0, T)$,
$\left|u^{N}(x, t)-u^{N, h}(x, t)\right| \leq \tilde{C}_{N, T}^{*} h^{s-\frac{1}{2}} \sum_{k=0}^{N}\left|\lambda_{k}\right|_{H^{s}(\Gamma)}$,
$x \in \Omega, \quad t \in(0, T)$,
where $\tilde{C}_{N, T}$ and $\tilde{C}_{N, T}^{*}$ are constants independent of $h$.
Proof. Let us fix two arbitrary values $N \in \mathbb{N}_{0}$ and $T \in \mathbb{R}_{+}$and consider $\delta_{N, T}:=\left\|\lambda^{N}(\cdot, t)-\lambda^{N, h}(\cdot, t)\right\|_{H^{1 / 2}(\Gamma)}=\| \sum_{k=0}^{N}\left(\lambda_{k}(\cdot)-\right.$ $\left.\lambda_{k}^{h}(\cdot)\right) L_{k}(\sigma t) \|_{H^{1 / 2}(\Gamma)} \quad$ for any $t \in(0, T)$. Setting $\quad C_{N, T}:=$ $\max ^{(0, T], k=0, N}\left|L_{k}(\sigma t)\right|$ we can write
$\delta_{N, T} \leq C_{N, T} \sum_{k=0}^{N}\left\|\lambda_{k}-\lambda_{k}^{h}\right\|_{H^{1 / 2}(\Gamma)}$.
Note that for any function $\lambda_{k}$, which satisfies an equation like (25), the inequality (39) yields the following estimate (Schtainbah, 2008)
$\left\|\lambda_{k}-\lambda_{k}^{h}\right\|_{H^{1 / 2}(\Gamma)} \leq \tilde{C}_{k} h^{s-1 / 2}\left|\lambda_{k}\right|_{H^{s}(\Gamma)}, \quad k \in N_{0}$,
where $\tilde{C}_{k}$ are constants independent of $h$. Using this inequality and setting $\tilde{C}_{N, T}:=C_{N, T} \max _{k=0, N}\left\{\tilde{C}_{k}\right\}$ we obtain (40) from (42).

Since in the case of Lipschitz boundary all functions $E_{j}(x-\cdot)$ are bounded and infinitely continuously differentiable on $\Gamma$ for any fixed point $x \in \Omega$, we get inequality $\left\|\partial_{\nu(\cdot)} E_{j}(x-\cdot)\right\|_{H^{-1 / 2}(\Gamma)} \leq$ $c_{j}^{*}=$ const. Taking it and (43) into account by the Generalized Cauchy-Schwarz inequality we obtain $\left|u_{k}(x)-u_{k}^{h}(x)\right|=$ $\left|\sum_{i=0}^{k}\left\langle\partial_{\nu(\cdot)} E_{k-i}(x-\cdot),\left(\lambda_{i}-\lambda_{i}^{h}\right)\right\rangle_{\Gamma}\right| \leq \tilde{c}_{k} h^{s-1 / 2} \sum_{i=0}^{k}\left|\lambda_{i}\right|_{H^{s}(\Gamma)}$, where $\tilde{c}_{k}$ are constants independent of $h$. Using this estimate, the rest of the proof for (41) can be carried out analogously to the proof for (40).

## 4. RESULTS OF THE COMPUTATIONAL EXPERIMENT

Let us demonstrate the suggested method to solve some model problem and assess the accuracy of numerical solutions. Let the domain $\Omega:=R^{3} \backslash \Omega^{-}$be outside of a cube $\Omega^{-}:=[-1,1] \times$ $[-1,1] \times[-1,1]$ and the function $g(x, t):=-\partial_{v(x)} v(x, t)$, $(x, t) \in \Sigma$, in the boundary condition (3) is defined by means of a spherical impulse $v(x, t):=f_{3}(t-|x|+1)|x|^{-1},(x, t) \in Q$, where $f_{3}$ is a cubic $B$-spline.

Tab. 1. Convergence behavior of $u_{0}^{h}(x)$ at points $x=\left(x_{1}, 0,0\right)$

|  | $u_{0}^{h}(x)$ |  |  | $\boldsymbol{u}_{\mathbf{0}}(\boldsymbol{x})$ |
| :---: | :--- | :--- | :--- | :--- |
| $\boldsymbol{x}_{\mathbf{1}}$ | $\overline{\boldsymbol{M}}=\mathbf{5 8 8}$ | $\overline{\boldsymbol{M}}=\mathbf{1 2 0 0}$ | $\overline{\boldsymbol{M}}=\mathbf{1 7 2 8}$ |  |
| 1.2 | 4.97567 <br> $\times 10^{-1}$ | 5.16241 <br> $\times 10^{-1}$ | 5.23443 <br> $\times 10^{-1}$ | 5.58600 <br> $\times 10^{-1}$ |
| 1.5 | 2.19226 <br> $\times 10^{-1}$ | 2.26607 <br> $\times 10^{-1}$ | 2.29583 <br> $\times 10^{-1}$ | 2.45252 <br> $\times 10^{-1}$ |
| 2.0 | 6.07058 <br>  <br> $\times 10^{-2}$ | 6.25240 <br> $\times 10^{-2}$ | 6.32929 <br> $\times 10^{-2}$ | 6.76676 <br> $\times 10^{-2}$ |
| 3.0 | 5.50191 | 5.64833 <br> $\times 10^{-3}$ | 5.71311 <br> $\times 10^{-3}$ | 6.10521 <br> $\times 10^{-3}$ |
| 4.0 | 5.59834 |  |  |  |
|  | $\times 10^{-4}$ | 5.73967 <br> $\times 10^{-4}$ | 5.80360 <br> $\times 10^{-4}$ | 6.19688 <br> $\times 10^{-4}$ |

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