

ANALYSIS OF THE STOCHASTIC OSKOLKOV SYSTEM WITH
A MULTIPOINT INITIAL-FINAL VALUE CONDITION

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ABSTRACT. Recently, the theory of stochastic equations has been actively developing. Here it is worth noting the classical direction of research by Ito – Stratonovich – Skorokhod. Its main problem is to overcome the difficulties associated with the differentiation of a non-differentiable (in "the usual sense") Wiener process. It is also necessary to note the approach of I.V. Melnikova, within the framework of which stochastic equations are considered in Schwarz spaces using the generalized derivative. Our research will use methods and results of the theory, which is based on the concept of the Nelson – Glicklikh derivative. Most studies consider the Cauchy problem for stochastic equations. In this article, instead of the Cauchy condition, it is proposed to consider a multipoint initial-final value condition. The obtained abstract results are used to analyze the solvability of the stochastic Oskolkov system, which models the dynamics of the velocity and pressure of a viscoelastic incompressible fluid. It is considered with a no-slip boundary condition and a multipoint initial-final value condition. The main result of the article is the proof of the solvability of the posed problem.

Introduction

Let $\Omega \subset \mathbb{R}^n$, $n \in \mathbb{N} \setminus \{1\}$, be a bounded domain with boundary $\partial\Omega$ of the class C^∞ . In the cylinder $\Omega \times \mathbb{R}$ consider the system of equations

$$(1 - \varkappa \nabla^2)v_t = \nu \nabla^2 v - (v \cdot \nabla)v - \nabla p + f, \quad \nabla \cdot v = 0, \quad (0.1)$$

modeling the dynamics of velocity $v = (v_1, v_2, \dots, v_n)$, $v_k = v_k(x, t)$, $k = \overline{1, n}$, and pressure $p = p(x, t)$, $(x, t) \in \Omega \times \mathbb{R}$, viscoelastic incompressible fluid. Here the parameter $\nu \in \mathbb{R}_+$ characterizes the viscous properties of the liquid, and the parameter $\varkappa \in \mathbb{R}$ characterizes the elastic properties of the fluid. The prototype of such a liquid is high-paraffin type of oil, produced, in particular, in the fields of Western Siberia. The system (0.1) was first obtained and studied by A.P. Oskolkov [1]. Therefore, it later received the name "Oskolkov system" [2]. In [3]

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G.A. Sviridyuk began studying a modified system of equations (0.1)

$$(\lambda - \nabla^2)u_t = \nu \nabla^2 u - \alpha(u \cdot \nabla)u - \nabla p, \quad \nabla \cdot u = 0. \quad (0.2)$$

To solve the system of equations (0.2) consider "the condition no-slip" to the boundary of the domain

$$u(x, t) = 0, \quad (x, t) \in \partial\Omega \times \mathbb{R}. \quad (0.3)$$

As is known, Newton rheological relation, which models the dynamics of viscous incompressible fluids, has the form [4], [5], [6]

$$\sigma = 2\nu D - p\mathbb{I}. \quad (0.4)$$

Here σ and D are the stress and strain rate tensors, respectively, $\nu \in \mathbb{R}_+$ is the viscosity coefficient, \mathbb{I} is the unit matrix, p characterizes the pressure. After substituting (0.4) into the equations of motion of a continuous incompressible medium in the Cauchy form

$$v_t = \nabla \cdot \sigma, \quad \nabla \cdot v = 0, \quad (0.5)$$

we obtain the famous Navier – Stokes system of equations

$$v_t = \nu \nabla^2 v - (v \cdot \nabla)v - \nabla p, \quad \nabla \cdot v = 0, \quad (0.6)$$

modeling the evolution of velocity and pressure of a viscous incompressible fluid. Extensive literature is devoted to the study of equations (0.6) in various aspects. Let us note here the fundamental monographs of O.A. Ladyzhenskaya [7] and R. Temam [8].

Various effects (for example, the recoil effect or the fading memory effect) that arise when pumping oil through pipelines and do not fit into the framework of the model (0.6) have prompted many researchers to revise the relationship (0.4). In particular, the rheological relation by V.A. Pavlovskii [9]

$$\sigma = 2\nu D + \varkappa D_t - p\mathbb{I}, \quad (0.7)$$

was offered and experimentally tested [10]. Here the coefficient \varkappa has the physical meaning of relaxation viscosity, and by design it is strictly positive. It immediately follows from (0.7) that the fluid velocity in the absence of stress does not immediately become equal to zero, as in Newton's model (0.4), but tends to zero exponentially, which gives the required relaxation effect.

However, later, in experiments with aqueous solutions of polymers, it turned out that the constant \varkappa can also take negative values [11]. Moreover, for negative values of \varkappa , the model (0.7) demonstrates strong instability. Analysis of the relationship (0.7), carried out in [12] from the standpoint of the rheological theory [4], [5], [6], showed that in the case of $\varkappa \in \mathbb{R}_+$ this model reveals properties characteristic of solids. Consequently, it cannot represent any non-Newtonian fluid [4], and therefore in [12] it is proposed to call such objects *media*.

If the rheological relation (0.7) is substituted into the equation (0.5), we obtain a system of equations (0.1), which, after simple algebraic transformations, will turn into the system (0.2). Here $\alpha = \lambda = \varkappa^{-1}$ are parameters characterizing elastic properties of liquid. In what follows, the system (0.2) will be considered for different values of the parameters α and λ . Initially, various initial-boundary value problems for the equations (0.1) were studied by A.P. Oskolkov [1], [2], [13], [14],

[15], then his students [16], [17] joined the research. The results of their research consisted mainly in proving the unique solvability of various initial-boundary value problems for the system (0.1) and its various generalizations for *positive values* of the parameter \varkappa in bounded and unbounded regions of space \mathbb{R}^n , $n \in \{2, 3, 4\}$.

Relevant equations [15] in specially constructed spaces [18] can be reduced to an abstract model

$$L\dot{u} = Mu + N(u),$$

where L , M are linear, and N are nonlinear operators. Let's consider the linear abstract model

$$L\dot{u} = Mu + f,$$

in Banach spaces \mathfrak{U} and \mathfrak{F} , and the operators $L \in \mathcal{L}(\mathfrak{U}; \mathfrak{F})$ (i.e. linear and continuous), $M \in \mathcal{Cl}(\mathfrak{U}; \mathfrak{F})$ (i.e. linear, closed and densely defined). Let the operator M be (L, p) -bounded [19], and its L -spectrum satisfy the condition [20]

$$\left\{ \begin{array}{l} \sigma^L(M) = \bigcup_{j=0}^m \sigma_j^L(M), \quad m \in \mathbb{N}, \quad \text{and } \sigma_j^L(M) \neq \emptyset, \quad \text{there exists} \\ \text{a closed contour } \gamma_j \subset \mathbb{C} \text{ bounding the domain } D_j \supseteq \sigma_j^L(M) \\ \text{such that } \overline{D_j} \cap \sigma_0^L(M) = \emptyset, \quad \overline{D_k} \cap \overline{D_l} = \emptyset \quad \forall j, k, l = \overline{1, m}, k \neq l. \end{array} \right. \quad (0.8)$$

Then there are relative spectral projectors

$$P_j = \frac{1}{2\pi i} \int_{\gamma_j} R_\mu^L(M) d\mu \in \mathcal{L}(\mathfrak{U}); \quad Q_j = \frac{1}{2\pi i} \int_{\gamma_j} L_\mu^L(M) d\mu \in \mathcal{L}(\mathfrak{F}), \quad j = \overline{0, m},$$

where the boundary contours γ_r are defined by (0.8). We will discuss these projectors in more detail in the third paragraph.

The work is devoted to the study of the stochastic linear Sobolev type equation

$$L \overset{\circ}{\eta} = M\eta + N\omega, \quad (0.9)$$

where $\eta = \eta(t)$ is the required one, and $\omega = \omega(t)$ is a given stochastic \mathbf{K} -process (\mathbf{K} -"noise"), with multipoint initial-final value condition

$$\lim_{t \rightarrow 0+} P_0(\eta(t) - \xi_0) = 0, \quad P_j(\eta(\tau_j) - \xi_j) = 0, \quad j = \overline{1, m}. \quad (0.10)$$

A detailed description will be given in the second paragraph.

The article, in addition to the introduction and bibliography, contains three parts. In the first part, spaces of differentiable random processes with values in a separable Hilbert space are constructed. Moreover, by derivative we mean the Nelson – Gliklich derivative [21], [22], [23], [24]. We call random processes that have Nelson – Gliklich derivatives differentiable "noises" [25], [26], [27], [28]. The second part of the article presents results on the solvability of the stochastic problem (0.9), (0.10) under the condition that the operator M , $p \in \{0\} \cup \mathbb{N}$, is (L, p) -bounded, and a condition guaranteeing the existence of relative spectral projectors P_j , $j = \overline{0, n}$, [29]. These results generalize and develop the abstract results of the works [25], [26], [27], [28]. The third part contains applications of the obtained abstract results for the stochastic Oskolkov system. The list of references does not pretend to be complete and reflects only the tastes and preferences of the authors.

1. Spaces of differentiable "noises"

Let $\Omega \equiv (\Omega, \mathcal{A}, \mathbf{P})$ is a complete probability space with a probability measure \mathbf{P} associated with the σ -algebra \mathcal{A} of subsets of the set Ω , and \mathbb{R} is a set of real numbers, endowed with a Borel σ -algebra. A measurable mapping $\xi : \Omega \rightarrow \mathbb{R}$ is called *random variable*. A set of random variables whose mathematical expectation is zero, and the variance is finite, forms the Hilbert space $\mathbf{L}_2 = \{\xi : \mathbf{E}\xi = 0, \mathbf{D}\xi < +\infty\}$ with the scalar product $(\xi_1, \xi_2) = \mathbf{E}\xi_1\xi_2$ and the norm $\|\xi\|_{\mathbf{L}_2}^2 = \mathbf{D}\xi$. Note that in \mathbf{L}_2 the orthogonality of the vectors ξ and η (i.e. $(\xi, \eta) = 0$) is equivalent to correlated random variables ξ and η . Indeed, $0 = \text{cov}(\xi, \eta) = \mathbf{E}\xi\eta = (\xi, \eta) = 0$.

Let us take the set $\mathcal{J} \subset \mathbb{R}$ and consider two mappings: $f : \mathcal{J} \rightarrow \mathbf{L}_2$, which each $t \in \mathcal{J}$ assigns a random variable $\xi \in \mathbf{L}_2$, and $g : \mathbf{L}_2 \times \Omega \rightarrow \mathbb{R}$, which assigns to each pair (ξ, ω) point $\xi(\omega) \in \mathbb{R}$. Display $\eta : \mathcal{J} \times \Omega \rightarrow \mathbb{R}$, which has the form $\eta = \eta(t, \omega) = g(f(t), \omega)$, we call *(one-dimensional) stochastic process*. For every fixed $t \in \mathcal{J}$ value of the stochastic process $\eta = \eta(t, \cdot)$ is a random variable, i.e. $\eta(t, \cdot) \in \mathbf{L}_2$, which we call *cross section* of the stochastic process at point $t \in \mathcal{J}$. For each fixed $\omega \in \Omega$ the function $\eta = \eta(\cdot, \omega)$ is called *(selective) trajectory* of a random process corresponding to the elementary outcome $\omega \in \Omega$. Trajectories are also called *realizations* or *sample functions* of a random process. Usually, when this does not lead to ambiguity, the dependence of $\eta(t, \omega)$ on ω is not indicated and the random process is simply denoted by $\eta(t)$.

Considering $\mathcal{J} \subset \mathbb{R}$ to be an interval, we call the stochastic process $\eta = \eta(t)$, $t \in \mathcal{J}$, *continuous*, if a.s. (almost surely) all its trajectories are continuous (i.e. for almost all $\omega \in \mathcal{A}$ trajectories $\eta(\cdot, \omega)$ are continuous functions). A set of continuous stochastic processes forms Banach space, which we denote by the symbol $\mathbf{C}(\mathcal{J}; \mathbf{L}_2)$ with the norm $\|\eta\|_{\mathbf{C}\mathbf{L}_2} = \sup_{t \in \mathcal{J}} (\mathbf{D}\eta(t, \omega))^{1/2}$. Let \mathcal{A}_0 be a σ -subalgebra σ -algebras

\mathcal{A} . Let us construct the subspace $\mathbf{L}_2^0 \subset \mathbf{L}_2$ random variables measurable with respect to \mathcal{A}_0 . Let us denote by $\Pi : \mathbf{L}_2 \rightarrow \mathbf{L}_2^0$ – ortho projector. Let $\xi \in \mathbf{L}_2$, then $\Pi\xi$ is called *conditional mathematical expectation* of the random variable ξ and is denoted by the symbol $\mathbf{E}(\xi|\mathcal{A}_0)$. Let us fix $\eta \in \mathbf{C}(\mathcal{J}; \mathbf{L}_2)$ and $t \in \mathcal{J}$, by \mathcal{N}_t^η we denote the σ -algebra generated by random variable $\eta(t)$, and denote $\mathbf{E}_t^\eta = \mathbf{E}(\cdot|\mathcal{N}_t^\eta)$.

Example 1.1. Wiener process describing Brownian motion in the Einstein – Smoluchowski model (see [22])

$$\beta(t, \omega) = \sum_{k=0}^{\infty} \xi_k(\omega) \sin \frac{\pi}{2}(2k+1)t, \quad t \in \{0\} \cup \mathbb{R}_+,$$

is a continuous stochastic process. Here the coefficients $\{\xi_k = \xi_k(\omega)\} \subset \mathbf{L}_2$ are pairwise uncorrelated Gaussian random variables such that $\mathbf{D}\xi_k^2 = \left[\frac{\pi}{2}(2k+1)\right]^{-2}$, $k \in \{0\} \cup \mathbb{N}$.

Definition 1.2. [21], [22] Let $\eta \in \mathbf{C}(\mathcal{J}; \mathbf{L}_2)$. By the *Nelson – Glicklich derivative* $\overset{\circ}{\eta}$ stochastic process η at point $t \in \mathcal{J}$ a random variable

$$\overset{\circ}{\eta}(t, \cdot) = \frac{1}{2} \lim_{\Delta t \rightarrow 0+} \mathbf{E}_t^\eta \left(\frac{\eta(t + \Delta t, \cdot) - \eta(t, \cdot)}{\Delta t} \right) + \frac{1}{2} \lim_{\Delta t \rightarrow 0+} \mathbf{E}_t^\eta \left(\frac{\eta(t, \cdot) - \eta(t - \Delta t, \cdot)}{\Delta t} \right),$$

is called, if the limit exists in the sense of a uniform metric on \mathbb{R} .

If the Nelson – Gliklich derivatives $\overset{\circ}{\eta}(t, \cdot)$ of the stochastic process $\eta(t, \cdot)$ exist in all (or almost all) points of the interval \mathfrak{J} , then we talk about the existence of the Nelson – Gliklich derivative $\overset{\circ}{\eta}(t, \cdot)$ on \mathfrak{J} (a.s. on \mathfrak{J} .)

Set of continuous stochastic processes having continuous Nelson – Gliklich derivatives $\overset{\circ}{\eta}$ forms a Banach $\mathbf{C}^1(\mathfrak{J}; \mathbf{L}_2)$ space with the norm

$$\|\eta\|_{\mathbf{C}^1 \mathbf{L}_2} = \sup_{t \in \mathfrak{J}} \left(\mathbf{D}\eta(t, \omega) + \mathbf{D} \overset{\circ}{\eta}(t, \omega) \right)^{1/2}.$$

We further define by induction the Banach spaces $\mathbf{C}^l(\mathfrak{J}; \mathbf{L}_2)$, $l \in \mathbb{N}$, stochastic processes whose trajectories a.s. differentiable with respect to Nelson – Gliklich on \mathfrak{J} up to order $l \in \{0\} \cup \mathbb{N}$ inclusive [30]. The norms in them are given by the

formulas $\|\eta\|_{\mathbf{C}^l \mathbf{L}_2} = \sup_{t \in \mathfrak{J}} \left(\sum_{k=0}^l \mathbf{D} \overset{\circ}{\eta}^{(k)}(t, \omega) \right)^{1/2}$. Here we will consider the zero-

order Nelson – Gliklich derivative to be the original random process, i.e. $\overset{\circ}{\eta}^{(0)} \equiv \eta$, and under the Nelson – Gliklich derivative are of order k we will understand the Nelson – Gliklich derivative of the first order from the Nelson – Gliklich derivative of order $k - 1$. For brevity we will call *spaces of differentiable "noises"* (see [25], [26], [27], [28]).

Example 1.3. In [22, 30] it is shown that $\beta \in \mathbf{C}^l(\mathbb{R}_+; \mathbf{L}_2)$, $l \in \{0\} \cup \mathbb{N}$, and $\overset{\circ}{\beta}(t) = \frac{\beta(t)}{2t}$, $t \in \mathbb{R}_+$.

Thus, spaces of random variables \mathbf{L}_2 and spaces of differentiable "noises" $\mathbf{C}^l(\mathfrak{J}; \mathbf{L}_2)$, $l \in \{0\} \cup \mathbb{N}$. Let's move on to constructing a space of random \mathbf{K} -variables. Take \mathfrak{H} is a separable Hilbert space with an orthonormal basis $\{\varphi_k\}$, a monotone sequence $\mathbf{K} = \{\lambda_k\} \subset \mathbb{R}_+$ such that that $\sum_{k=1}^{\infty} \lambda_k^2 < +\infty$, as well as a sequence $\{\xi_k\} = \xi_k(\omega) \subset \mathbf{L}_2$ of random variables such that that $\|\xi_k\|_{\mathbf{L}_2} \leq C$, for all $C \in \mathbb{R}_+$, for all $k \in \mathbb{N}$.

Let us construct a \mathfrak{H} -valued random \mathbf{K} -variable $\xi(\omega) = \sum_{k=1}^{\infty} \lambda_k \xi_k(\omega) \varphi_k$. Completion of the linear hull of the set $\{\lambda_k \xi_k \varphi_k\}$ by the norm

$$\|\eta\|_{\mathbf{H}_{\mathbf{K}} \mathbf{L}_2}^2 = \left(\sum_{k=1}^{\infty} \lambda_k^2 \mathbf{D} \xi_k \right)^{1/2}$$

is called the *space of (\mathfrak{H} -valued) random \mathbf{K} -variables* and is denoted by the symbol $\mathbf{H}_{\mathbf{K}} \mathbf{L}_2$. How easy it is to see the space $\mathbf{H}_{\mathbf{K}} \mathbf{L}_2$ is Hilbertian, and the random \mathbf{K} -variable constructed above $\xi = \xi(\omega) \in \mathbf{H}_{\mathbf{K}} \mathbf{L}_2$. Likewise, Banach space (\mathfrak{H} -valued) \mathbf{K} - "noises" $\mathbf{C}^l(\mathfrak{J}; \mathbf{H}_{\mathbf{K}} \mathbf{L}_2)$, $l \in \{0\} \cup \mathbb{N}$, we define as the completion of the linear

hull of the set $\{\lambda_k \eta_k \varphi_k\}$ by the norm

$$\|\eta\|_{\mathbf{C}^l \mathbf{H}_K \mathbf{L}_2}^2 = \sup_{t \in \mathcal{J}} \left(\sum_{k=1}^{\infty} \lambda_k^2 \sum_{m=1}^l \mathbf{D} \overset{\circ}{\eta}_k^m \right)^{1/2},$$

where the sequence of "noises" $\{\eta_k\} \subset \mathbf{C}^l(\mathcal{J}; \mathbf{L}_2)$, $l \in \{0\} \cup \mathbb{N}$. As is easy to see, the vector $\eta(t, \omega) = \sum_{k=1}^{\infty} \lambda_k \eta_k(t, \omega) \varphi_k$ lies in the space $\mathbf{C}^l(\mathcal{J}; \mathbf{H}_K \mathbf{L}_2)$, if a sequence of vectors $\{\eta_k\} \subset \mathbf{C}^l(\mathcal{J}; \mathbf{L}_2)$ and all their Nelson – Glicklich derivatives up to order $l \in \{0\} \cup \mathbb{N}$ inclusive are uniformly bounded by the norm $\|\cdot\|_{\mathbf{C}^l \mathbf{L}_2}$.

Example 1.4. Vector lying in all spaces $\mathbf{C}^l(\mathbb{R}_+; \mathbf{H}_K \mathbf{L}_2)$, $l \in \{0\} \cup \mathbb{N}$,

$$W_{\mathbf{K}}(t, \omega) = \sum_{k=1}^{\infty} \lambda_k \beta_k(t, \omega) \varphi_k,$$

where $\{\beta_k\} \subset \mathbf{C}^l(\mathcal{J}; \mathbf{L}_2)$ is sequence of Brownian motions, called (*st-valued*) *Wiener K-process*.

2. The multipoint initial-final value condition

Now let $\mathfrak{U}(\mathfrak{F})$ be a real separable Hilbert space with an orthonormal basis $\{\varphi_k\}$ ($\{\psi_k\}$). Let us introduce into consideration a monotone sequence $\mathbf{K} = \{\lambda_k\} \subset \{0\} \cup \mathbb{R}$ such that $\sum_{k=1}^{\infty} \lambda_k^2 < +\infty$. The symbol $\mathbf{U}_K \mathbf{L}_2$ ($\mathbf{F}_K \mathbf{L}_2$) denotes the Hilbert space, which is the completion of the linear hull of *random K-variables*

$$\xi = \sum_{k=1}^{\infty} \lambda_k \xi_k \varphi_k, \quad \xi_k \in \mathbf{L}_2, \quad \left(\zeta = \sum_{k=1}^{\infty} \mu_k \zeta_k \psi_k, \quad \zeta_k \in \mathbf{L}_2 \right),$$

according to the norm $\|\eta\|_{\mathfrak{U}}^2 = \sum_{k=1}^{\infty} \lambda_k^2 \mathbf{D} \xi_k$ ($\|\omega\|_{\mathfrak{F}}^2 = \sum_{k=1}^{\infty} \mu_k^2 \mathbf{D} \zeta_k$). Note that in different spaces ($\mathbf{U}_K \mathbf{L}_2$ and $\mathbf{F}_K \mathbf{L}_2$) the sequence \mathbf{K} can be different ($\mathbf{K} = \{\lambda_k\}$ in $\mathbf{U}_K \mathbf{L}_2$ and $\mathbf{K} = \{\mu_k\}$ in $\mathbf{F}_K \mathbf{L}_2$), however, all sequences marked with \mathbf{K} must be monotonic and summable with square. All results, generally speaking, will be true for different sequences $\{\lambda_k\}$ and $\{\mu_k\}$, but for the sake of simplicity we will limit ourselves to the case $\lambda_k = \mu_k$.

Lemma 2.1. *Operator $A \in \mathcal{L}(\mathfrak{U}; \mathfrak{F})$ exactly when $A \in \mathcal{L}(\mathbf{U}_K \mathbf{L}_2; \mathbf{F}_K \mathbf{L}_2)$.*

How easy it is to see

$$\|A\xi\|_{\mathfrak{F}} \leq \sum_{k=1}^{\infty} \lambda_k^2 \mathbf{D} \xi_k \|A\varphi_k\|_{\mathfrak{F}}^2 \leq \text{const} \sum_{k=1}^{\infty} \lambda_k^2 \mathbf{D} \xi_k = \text{const} \|\xi\|_{\mathfrak{U}}.$$

Lemma 2.2. *Operator $M \in \mathcal{L}(\mathfrak{U}; \mathfrak{F})$ is σ -bounded with respect to operator $L \in \mathcal{L}(\mathfrak{U}; \mathfrak{F})$ exactly when $M \in \mathcal{L}(\mathbf{U}_K \mathbf{L}_2; \mathbf{F}_K \mathbf{L}_2)$ is σ -bounded with respect to the operator $L \in \mathcal{L}(\mathbf{U}_K \mathbf{L}_2; \mathbf{F}_K \mathbf{L}_2)$. Moreover, the L -spectrum of the operator M coincide in both cases.*

The proof of Lemma 2.2 is similar to the proof of Lemma 2.1 and is therefore omitted. Ideas and methods of theory regarding σ -bounded operators can be found, for example, in [19]. According to this theory, in the case of a (L, σ) -bounded operator M there is a pair of relative spectral projectors

$$P = \frac{1}{2\pi i} \int_{\gamma} R_{\mu}^L(M) d\mu \in \mathcal{L}(\mathfrak{U}), \quad Q = \frac{1}{2\pi i} \int_{\gamma} L_{\mu}^L(M) d\mu \in \mathcal{L}(\mathfrak{F}), \quad (2.1)$$

where $\gamma \subset \mathbb{C}$ is the contour bounding the region containing the L -spectrum $\sigma^L(M)$ of the operator M , and $R_{\mu}^L(M) = (\mu L - M)^{-1} L$ ($L_{\mu}^L(M) = L(\mu L - M)^{-1}$ – right (left) L -resolvent operator M). By Lemma 2.1, we transfer the projectors P and Q to the spaces $\mathbf{U}_{\mathbf{K}}\mathbf{L}_2$ and $\mathbf{F}_{\mathbf{K}}\mathbf{L}_2$ respectively and introduce into consideration the subspaces $\mathbf{U}_{\mathbf{K}}^0\mathbf{L}_2 = \ker P$, $\mathbf{F}_{\mathbf{K}}^0\mathbf{L}_2 = \ker Q$ and $\mathbf{U}_{\mathbf{K}}^1\mathbf{L}_2 = \text{im } P$, $\mathbf{F}_{\mathbf{K}}^1\mathbf{L}_2 = \text{im } Q$. We denote by $L_k(M_k)$ the restriction of the operator $L(M)$ to $\mathbf{U}_{\mathbf{K}}^k\mathbf{L}_2$ ($\mathbf{F}_{\mathbf{K}}^k\mathbf{L}_2$), $k = 0, 1$. Fair

Theorem 2.3. (*splitting theorem*). [28] *Let $L, M \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}\mathbf{L}_2)$, operator M (L, σ) -bounded. Then*

- (i) *operators $L_0, M_0 \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}^0\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}^0\mathbf{L}_2)$; $L_1, M_1 \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}^1\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}^1\mathbf{L}_2)$;*
- (ii) *operators $M_0^{-1} \in \mathcal{L}(\mathbf{F}_{\mathbf{K}}^0\mathbf{L}_2; \mathbf{U}_{\mathbf{K}}^0\mathbf{L}_2)$, $L_1^{-1} \in \mathcal{L}(\mathbf{F}_{\mathbf{K}}^1\mathbf{L}_2; \mathbf{U}_{\mathbf{K}}^1\mathbf{L}_2)$.*

Let us construct the operator $H = L_0^{-1}M_0 \in \mathcal{L}(\mathfrak{U}^0)$ ($G = M_0L_0^{-1} \in \mathcal{L}(\mathfrak{F}^0)$). We call a (L, σ) -bounded operator M (L, p) -bounded if there exists a number $p \in \{0\} \cup \mathbb{N}$ such, that $H^p \neq \mathbb{O}$, and $H^{p+1} = \mathbb{O}$. Note that the number $p \in \{0\} \cup \mathbb{N}$ does not change if the operator H is replaced by the operator G . Now let us take the (L, p) -bounded operator M , $p \in \{0\} \cup \mathbb{N}$, and construct families of operators

$$U^t = \frac{1}{2\pi i} \int_{\gamma} R_{\mu}^L(M) e^{\mu t} d\mu, \quad F^t = \frac{1}{2\pi i} \int_{\gamma} L_{\mu}^L(M) e^{\mu t} d\mu, \quad (2.2)$$

where the contour $\gamma \subset \mathbb{C}$ is the same as in (2.1). On \mathbb{R} the families (2.2) are commutative groups whose units are the projectors (2.1). Moreover, the families (2.2) extend to the entire complex plane while preserving the group property.

Let $\tau_0 = 0$, $\tau_j \in \mathbb{R}_+$, $(\tau_{j-1} < \tau_j)$, $j = \overline{1, m}$. Let condition (0.8) be satisfied.

Theorem 2.4. [31] *Let the operator M be (L, p) -bounded, $p \in \{0\} \cup \mathbb{N}$, and the condition (0.8) holds. Then there exist families of operators*

$$U_j^t = \frac{1}{2\pi i} \int_{\gamma_j} R_{\mu}^L(M) e^{\mu t} d\mu, \quad F_j^t = \frac{1}{2\pi i} \int_{\gamma_j} L_{\mu}^L(M) e^{\mu t} d\mu, \quad j = \overline{1, m},$$

$$U_0^t = U^t - \sum_{k=1}^m U_k^t, \quad F_0^t = F^t - \sum_{k=1}^m F_k^t, \quad t \in \mathbb{R}.$$

Moreover

- (i) $U_j^t U_j^s = U_j^s U_j^t = U_j^{s+t}$, $F_j^t F_j^s = F_j^s F_j^t = F_j^{s+t}$ for all $s, t \in \mathbb{R}$, $j = \overline{0, m}$;
- (ii) $U_k^t U_l^s = U_l^s U_k^t = \mathbb{O}$, $F_k^t F_l^s = F_l^s F_k^t = \mathbb{O}$ for all $s, t \in \mathbb{R}$, $k, l = \overline{0, m}$, $k \neq l$.

Remark 2.5. Consider the units $P_j = U_j^0$, $Q_j = F_j^0$, $j = \overline{0, m}$. It can be shown, based on the results of Theorem 2.4, that

- (i) $PP_j = P_jP = P_j$, $QQ_j = Q_jQ = Q_j$, $j = \overline{0, m}$;
- (ii) $P_kP_l = P_lP_k = \mathbb{O}$, $Q_kQ_l = Q_lQ_k = \mathbb{O}$ $k, l = \overline{0, m}$, $k \neq l$.

By Lemma 2.1, we transfer the projectors $P_j, Q_j, j = \overline{0, m}$, to the spaces $\mathbf{U}_{\mathbf{K}}\mathbf{L}_2$ and $\mathbf{F}_{\mathbf{K}}\mathbf{L}_2$ respectively and introduce subspaces into consideration $\mathbf{U}_{\mathbf{K}}^{1j}\mathbf{L}_2 = \text{im } P_j, \mathbf{F}_{\mathbf{K}}^{1j}\mathbf{L}_2 = \text{im } Q_j, j = \overline{0, m}$. By construction $\mathbf{U}_{\mathbf{K}}\mathbf{L}_2 = \bigoplus_{j=0}^m \mathbf{U}_{\mathbf{K}}^{1j}\mathbf{L}_2,$

$\mathbf{F}_{\mathbf{K}}\mathbf{L}_2 = \bigoplus_{j=0}^m \mathbf{F}_{\mathbf{K}}^{1j}\mathbf{L}_2.$ We denote by $L_{1j} (M_{1j})$ the restriction of the operator $L (M)$ to $\mathbf{U}_{\mathbf{K}}^{1j}\mathbf{L}_2 (\mathbf{F}_{\mathbf{K}}^{1j}\mathbf{L}_2), j = \overline{0, m}.$

Theorem 2.6. (Generalized splitting theorem). [28] *Let $L, M \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}\mathbf{L}_2),$ operator M is (L, p) -bounded, $p \in \{0\} \cup \mathbb{N},$*

(i) operators $L_{1j} \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}^{1j}\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}^{1j}\mathbf{L}_2), M_{1j} \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}^{1j}\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}^{1j}\mathbf{L}_2), j = \overline{0, m};$

(ii) there are operators $L_{1j}^{-1} \in \mathcal{L}(\mathbf{F}_{\mathbf{K}}^{1j}\mathbf{L}_2; \mathbf{U}_{\mathbf{K}}^{1j}\mathbf{L}_2), j = \overline{0, m}.$

Let the operators $L, M, N \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}\mathbf{L}_2),$ Let us consider a linear stochastic equation of Sobolev type (0.9). Let us supply the equation (0.9) multipoint initial-final value condition (0.10)

Let us call the stochastic \mathbf{K} -process $\eta \in \mathbf{C}^1(\mathbb{R}_+; \mathbf{L}_2)$ (*classical*) *solution of the equation (0.9),* if a.s. all of it trajectories satisfy the equation (0.9) with some \mathbf{K} -”noise” $\omega \in \mathbf{C}(\mathbb{R}_+; \mathbf{L}_2)$ and all $t \in \mathbb{R}_+.$ The solution $\eta = \eta(t)$ of the equation (0.9) will be called a *solution of the problem (0.9), (0.10)* if the condition (0.10) for some random \mathbf{K} -variables $\xi_k \in \mathbf{U}_{\mathbf{K}}\mathbf{L}_2, k = \overline{0, l}.$

Theorem 2.7. [28] *Let the operator M is (L, p) -bounded, $p \in \{0\} \cup \mathbb{N},$ and condition (0.8) is satisfied. Then for any $\tau_j \in \mathbb{R}_+, j = \overline{1, m},$ operator $N \in \mathcal{L}(\mathfrak{L}; \mathfrak{F}),$*

monotonic sequence $K \subset \{\lambda_k\}$ such that $\sum_{k=1}^{\infty} \lambda_k^2 < +\infty,$ \mathbf{K} - ”noise” $\omega = \omega(t)$ such

that $(\mathbb{I} - Q)N\omega \in \mathbf{C}^{p+1}(\mathbb{R}_+; \mathbf{U}_{\mathbf{K}}\mathbf{L}_2)$ and $QN\omega \in \mathbf{C}(\mathbb{R}_+; \mathbf{U}_{\mathbf{K}}\mathbf{L}_2),$ and random \mathbf{K} -variables $\xi_j \in \mathbf{U}_{\mathbf{K}}\mathbf{L}_2, j = \overline{0, m},$ independent of $\omega,$ there is a unique solution $\eta \in \mathbf{C}^1(\mathbb{R}_+; \mathbf{U}_{\mathbf{K}}\mathbf{L}_2),$ problem (0.9), (0.10), having the form

$$\begin{aligned} \eta(t) = & - \sum_{q=0}^p H^q M_0^{-1} (\mathbb{I} - Q) \overset{\circ}{\omega}^{(q)}(t) + \\ & + \sum_{j=0}^m \left[U_j^{t-\tau_j} \xi_j + \int_{\tau_j}^t U_j^{s-\tau_j} L_{1j}^{-1} Q_j N \omega(s) ds \right], t \in \mathfrak{I}. \end{aligned}$$

Corollary 2.8. *Let all the conditions of Theorem 2.7 be satisfied and $\omega(t) = \overset{\circ}{W}_{\mathbf{K}}(t).$ Then for any random \mathbf{K} -variables $\xi_j \in \mathbf{U}_{\mathbf{K}}\mathbf{L}_2, j = \overline{0, m},$ there is a unique solution to the problem (0.9), (0.10), having the form*

$$\begin{aligned} \eta(t) = & - \sum_{q=0}^p H^q M_0^{-1} (\mathbb{I} - Q) \overset{\circ}{W}_{\mathbf{K}}^{(q+1)}(t) + \\ & + \sum_{j=0}^m \left[U_j^{t-\tau_j} \xi_j + L_{1j}^{-1} Q_j N W_{\mathbf{K}}(t) - S_j P_j \int_{\tau_j}^t U_j^{s-\tau_j} L_{1j}^{-1} Q_j N W_{\mathbf{K}}(s) ds \right], t \in \overline{\mathbb{R}}_+. \end{aligned}$$

3. Linear stochastic Oskolkov system

Let $\Omega \subset \mathbb{R}^n$, $n \in \mathbb{N} \setminus \{1\}$, be a bounded domain with boundary $\partial\Omega$ of the class C^∞ . In the cylinder $\Omega \times \mathbb{R}$ consider the linear Oskolkov system of equations

$$(1 - \varkappa \nabla^2)v_t = \nu \nabla^2 v - \nabla p + f, \quad \nabla \cdot v = 0. \quad (3.1)$$

Based on the results of points 2 and 3, we will reduce the system (3.1) and the condition (0.3) to the equation (0.9). Following [18], [32], [33], we denote by $\mathbb{H}^2 = (W_2^2)^n$, $\mathring{\mathbb{H}}^1 = (\mathring{W}_2^1)^n$, $\mathbb{L}^2 = (L_2)^n$ space vector-functions $v = (v_1, v_2, \dots, v_n)$ defined on Ω . Consider the lineal $\mathfrak{L} = \{v \in (C_n^\infty)^n : \nabla \cdot v = 0\}$ of vector-functions, solenoidal and finite in the domain Ω . We denote the closure of \mathfrak{L} with respect to the norm of the space \mathbb{L}^2 by \mathbb{H}_σ . The space \mathbb{H}_σ is Hilbert with the scalar product $\langle \cdot, \cdot \rangle$ inherited from \mathbb{L}^2 ; in addition, there is a splitting $\mathbb{L}^2 = \mathbb{H}_\sigma \oplus \mathbb{H}_\pi$, where \mathbb{H}_π is the orthogonal complement of \mathbb{H}_σ . We denote by $\Sigma : \mathbb{L}^2 \rightarrow \mathbb{H}_\sigma$ the corresponding orthoprojector. The restriction of the projector Σ to $\mathbb{H}^2 \cap \mathring{\mathbb{H}}^1$ is a continuous operator, $\Sigma : \mathbb{H}^2 \cap \mathring{\mathbb{H}}^1 \rightarrow \mathbb{H}^2 \cap \mathring{\mathbb{H}}^1$. Let us therefore represent the space $\mathbb{H}^2 \cap \mathring{\mathbb{H}}^1$ as a direct sum $\mathbb{H}^2 \cap \mathring{\mathbb{H}}^1 = \mathbb{H}_\sigma^2 \oplus \mathbb{H}_\pi^2$, where $\mathbb{H}_\sigma^2 = \text{im } \Sigma$, $\mathbb{H}_\pi^2 = \text{ker } \Sigma$. There are continuous and dense embeddings $\mathbb{H}_\sigma^2 \hookrightarrow \mathbb{H}_\sigma$ and $\mathbb{H}_\pi^2 \hookrightarrow \mathbb{H}_\pi$. The space \mathbb{H}_π^2 consists of vector functions that are equal to zero on $\partial\Omega$ and are gradients of functions from $W_2^3(\Omega)$.

Lemma 3.1. [19]

(i) By the formula $A = (-\nabla^2)^n : \mathbb{H}^2 \cap \mathring{\mathbb{H}}^1 \rightarrow \mathbb{L}^2$ defines a linear continuous operator with a positive discrete finite multiple spectrum $\sigma(A)$ condensing only to the point $+\infty$, and the mapping $A : \mathbb{H}_{\sigma(\pi)}^2 \rightarrow \mathbb{H}_{\sigma(\pi)}$ is bijective.

(ii) Formula $B : v \rightarrow -\nabla(\nabla \cdot v)$ a linear continuous surjective operator is given $B : \mathbb{H}^2 \cap \mathring{\mathbb{H}}^1 \rightarrow \mathbb{H}_\pi$, and $\text{ker } B = \mathbb{H}_\sigma^2$.

Let $\mathfrak{U} = \mathbb{H}_\sigma^2 \times \mathbb{H}_\pi^2 \times \mathbb{H}_p$, $\mathfrak{F} = \mathbb{H}_\sigma \times \mathbb{H}_\pi \times \mathbb{H}_p$ are real separable Hilbert spaces with an orthonormal basis $\{\varphi_k\}$ and $\{\psi_k\}$, respectively. Moreover $\mathbb{H}_p = \mathbb{H}_\pi$, $A_\lambda = \lambda \mathbb{I} + A$. Let's construct operators

$$L = \begin{pmatrix} \Sigma A_\lambda & \mathbb{O} & \mathbb{O} \\ \mathbb{O} & \Pi A_\lambda & \mathbb{O} \\ \mathbb{O} & \mathbb{O} & \mathbb{O} \end{pmatrix}, \quad M = \begin{pmatrix} -\nu \Sigma A & \mathbb{O} & \mathbb{O} \\ \mathbb{O} & -\nu \Pi A & -Pi \\ \mathbb{O} & \Pi B & \mathbb{O} \end{pmatrix}.$$

Obviously, $L, M \in \mathcal{L}(\mathfrak{U}; \mathfrak{F})$, and $\text{im } L = \mathbb{H}_\sigma \times \mathbb{H}_\pi \times \{0\}$, $\text{ker } L = \{0\} \times \{0\} \times \mathbb{H}_p$.

Let us introduce into consideration a monotone sequence $\mathbf{K} = \{\lambda_k\} \subset \{0\} \cup \mathbb{R}$ such that $\sum_{k=1}^{\infty} \lambda_k^2 < +\infty$. Let us consider Hilbert spaces, which is the completion of the linear hull of *random K-variables*

$$\mathbf{U}_\mathbf{K} \mathbf{L}_2 = \left\{ \xi = \sum_{k=1}^{\infty} \lambda_k \xi_k \varphi_k, \quad \xi_k \in \mathbf{L}_2, \quad \varphi_k \in \mathbb{H}_\sigma^2 \times \mathbb{H}_\pi^2 \times \mathbb{H}_p \right\},$$

$$\mathbf{F}_\mathbf{K} \mathbf{L}_2 = \left\{ \zeta = \sum_{k=1}^{\infty} \mu_k \zeta_k \psi_k, \quad \zeta_k \in \mathbf{L}_2, \quad \psi_k \in \mathbb{H}_\sigma \times \mathbb{H}_\pi \times \mathbb{H}_p \right\}.$$

Lemma 3.2. [19] For any $\lambda \in \mathbb{R} \setminus \sigma(A)$, $\nu \in \mathbb{R}_+$ the operator M is $(L, 1)$ -bounded.

If we set $N\omega = \text{col}(\Sigma f, \Pi f, 0)$, and $f(t) = \overset{\circ}{W}_{\mathbf{K}}(t)$, then the reduction of the problem (0.3), (3.1) to the equation (0.9) finished.

Let the operators $L, M, N \in \mathcal{L}(\mathbf{U}_{\mathbf{K}}\mathbf{L}_2; \mathbf{F}_{\mathbf{K}}\mathbf{L}_2)$. Consider a linear stochastic Sobolev type equation (0.9) with multipoint initial-final value condition (0.10).

Let us call the stochastic \mathbf{K} -process $\eta \in \mathbf{C}^1(\mathbb{R}_+; \mathbf{L}_2)$ (classical) solution of the equation (0.9), if a.s. all of it trajectories satisfy the equation (0.9) for some \mathbf{K} - "noise" $\omega \in \mathbf{C}(\mathbb{R}_+; \mathbf{L}_2)$ and all $t \in \mathbb{R}_+$. Solution $\eta = \eta(t)$ to the equation (0.9) let's call solution to the problem (0.9), (0.10), if the condition (0.10) is met for some random \mathbf{K} -variables $\xi_k \in \mathbf{U}_{\mathbf{K}}\mathbf{L}_2$, $k = \overline{0, l}$.

Let us find the L -spectrum $\sigma^L(M)$ of the operator M . As is easy to see, the operator $\mu L - M$ is invertible exactly when the operator

$$\Sigma(\mu\lambda\mathbb{I} - (\mu - \nu)A) : \mathbb{H}_{\sigma}^2 \rightarrow \mathbb{H}_{\sigma}$$

is invertible. Let \tilde{A} denote the restriction of the operator A to \mathbb{H}_{σ}^2 . The spectrum of the operator $\tilde{A} \in \mathcal{L}(\mathbb{H}_{\sigma}^2; \mathbb{H}_{\sigma})$ is positive, discrete, finitely multiple and condenses only to $+\infty$ (Solonnikov - Vorovich - Yudovich theorem [31]). Let $\{\lambda_k\}$ denote the set of eigenvalues of the operator \tilde{A} , numbered in non-decreasing order taking into account multiplicity. Then

$$\sigma^L(M) = \left\{ \mu_k = \frac{\nu\lambda_k}{\lambda_k - \lambda} : \lambda_k \in \sigma(\tilde{A}) \setminus \{\lambda\} \right\}.$$

It is clear that for such a set one can select contours $\gamma_j \subset \mathbb{C}$. Let's construct

$$U_j^t = \begin{pmatrix} \sum_{\lambda_k \in \sigma_j^L(M)} e^{\lambda_k t} \langle \cdot, \varphi_k \rangle_{\sigma} \varphi_k & \mathbb{O} & \mathbb{O} \\ \mathbb{O} & \mathbb{O} & \mathbb{O} \\ \mathbb{O} & \mathbb{O} & \mathbb{O} \end{pmatrix}, j = \overline{0, m}.$$

It follows from Lemma 3.2 that under the conditions of this lemma the condition (0.8) is satisfied.

Theorem 3.3. Let the operators L and M be defined as in Lemma 3.2. Then for any $\tau_j \in \mathbb{R}_+$, $j = \overline{1, m}$, operator $N \in \mathcal{L}(\mathfrak{U}; \mathfrak{F})$, monotonic sequence $\mathbf{K} \subset \{\lambda_k\}$ such that $\sum_{k=1}^{\infty} \lambda_k^2 < +\infty$, \mathbf{K} - "noise" $\omega = \omega(t)$ such that $(\mathbb{I} - Q)N\omega \in \mathbf{C}^{p+1}(\mathbb{R}_+; \mathbf{U}_{\mathbf{K}}\mathbf{L}_2)$

and $QN\omega \in \mathbf{C}(\mathbb{R}_+; \mathbf{U}_{\mathbf{K}}\mathbf{L}_2)$, $\omega(t) = \overset{\circ}{W}_{\mathbf{K}}(t)$ and random \mathbf{K} -variables $\xi_j \in \mathbf{U}_{\mathbf{K}}\mathbf{L}_2$, $j = \overline{0, m}$, independent of ω , there is a unique solution $\eta \in \mathbf{C}^1(\mathbb{R}_+; \mathbf{U}_{\mathbf{K}}\mathbf{L}_2)$, problem (0.9), (0.10), having the form

$$\begin{aligned} \eta(t) = & \sum_{j=0}^m \left[U_j^{t-\tau_j} \xi_j + L_{1j}^{-1} Q_j N W_{\mathbf{K}}(t) - S_j P_j \int_{\tau_j}^t U_j^{s-\tau_j} L_{1j}^{-1} Q_j N W_{\mathbf{K}}(s) ds \right] - \\ & - \sum_{q=0}^p H^q M_0^{-1} (\mathbb{I} - Q) \overset{\circ}{W}_{\mathbf{K}}^{(q+1)}(t), t \in \mathbb{R}_+. \end{aligned}$$

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