

STOCHASTIC PARABOLIC SYSTEMS WITH VARIABLE EXPONENTS OF NONLINEARITY

Some nonlinear parabolic systems of partial differential equations with the white noise is considered. The initial-boundary value problem for a system is investigated and the existence and uniqueness of the weak solution for the problem are proved.

Key words: partial differential equations, stochastic parabolic system, variable exponent of nonlinearity, white noise, weak solution.

Introduction. Let $n, N \in \mathbb{N}$ and $T > 0$ be some fixed numbers, $\Omega \subset \mathbb{R}^n$ be a bounded domain with the smooth boundary $\partial\Omega$, $(\mathbb{S}, \mathcal{F}, \mathbb{P})$ be a complete probability space, where \mathbb{S} is a non-empty set, which is interpreted as a space of states or elementary events, \mathcal{F} is a σ -algebra of the subsets of \mathbb{S} , and $\mathbb{P}: \mathcal{F} \rightarrow \mathbb{R}$ is a probability measure. Let $\mathcal{Q}_{0,T} = \Omega \times (0, T)$, $\Pi_{0,T} = \Omega \times (0, T) \times \mathbb{S}$, $\Theta_{0,T} = (0, T) \times \mathbb{S}$. We intend to find a vector-valued function $u = (u_1, \dots, u_N): \Pi_{0,T} \rightarrow \mathbb{R}^N$ that is dependent on the deterministic variables $x \in \Omega$ and $t \in (0, T)$ and the random variable $\omega \in \mathbb{S}$ such that

$$\begin{aligned} u_{k,t} + \alpha \Delta^2 u_k - \sum_{i=1}^n (a_{ik}(x, t) |u_{x_i}|^{p(x,t)-2} u_{k,x_i})_{x_i} + \\ + (\hat{N}u)_k(x, t, \omega) + \varphi_k((\hat{E}u)_k(x, t, \omega)) = \\ = F_k(x, t, \omega) + b_{k,t}(x, t, \omega), \\ (x, t, \omega) \in \Pi_{0,T}, \quad k = 1, 2, \dots, N, \end{aligned} \quad (1)$$

$$u(x, t, \omega) = \Delta u(x, t, \omega) = 0, \quad x \in \partial\Omega, \quad (t, \omega) \in \Theta_{0,T}, \quad (2)$$

$$u(x, 0, \omega) = u_0(x, \omega), \quad x \in \Omega, \quad \omega \in \mathbb{S}, \quad (3)$$

where $\alpha > 0$ is a number, $\Delta = \partial^2/\partial x_1^2 + \dots + \partial^2/\partial x_n^2$ is the Laplacian, $\Delta^{s+1} = \Delta(\Delta^s)$, $s \in \mathbb{N}$,

$$(\hat{N}u)_k(x, t, \omega) := g_k(x, t) |u(x, t, \omega)|^{q(x,t)-2} u_k(x, t, \omega),$$

$$(\hat{E}u)_k(x, t, \omega) := \int_{\Omega} \mathfrak{J}_k(x, t, y) u_k(y, t, \omega) dy, \quad (x, t, \omega) \in \Pi_{0,T},$$

a_{ik} , g_k , φ_k , \mathfrak{J}_k , F_k , u_0 , b_k , p , q are some functions, $i = 1, 2, \dots, n$, $k = 1, 2, \dots, N$.

We will prove the unique solvability of problem (1)–(3). This work continues the research of [1], where some semilinear stochastic parabolic equations with the variable exponent of nonlinearity are considered. The main conclusion in [1] was based on the methodology shown in [4, 12], that used

[✉] oleh.buhrii@lnu.edu.ua

multivalued functions and their selectors. Similar results for parabolic equations with the constant exponents of the nonlinearity were obtained in [10, 13]. Stochastic problems with the variable exponent of the nonlinearity in a formulation different from ours have been studied in [3, 14]. Nonlinear non-stochastic partial differential equations with the variable exponent of the nonlinearity have been studied, e.g., in [2, 5–9].

In contrast to [1], in this article we make significant use of the uniqueness of the solution to problem (1)–(3) without the random parameter $\omega \in \mathbb{S}$ and the white noise component $b_{k,t}$. Problem (1)–(3) is considered, apparently, for the first time.

1. Statement of the problem. We incorporate notations from [9]. In addition, we denote $\text{Lip}(\mathbb{R})$ as the set of all Lipschitz functions on \mathbb{R} ; for $m \in \mathbb{N}$ and $\mathcal{O} \subset \mathbb{R}^m$, we introduce the sets: $L^0(\mathcal{O}) := \{v : \mathcal{O} \rightarrow \mathbb{R} \mid v \text{ is a measurable function}\}$, $\mathcal{B}_+(\mathcal{O}) := \{q \in L^\infty(\mathcal{O}) \mid \text{essinf}_{y \in \mathcal{O}} q(y) > 0\}$. For all $q \in \mathcal{B}_+(\mathcal{O})$, we define the following:

$$q_0 := \text{essinf}_{y \in \mathcal{O}} q(y), \quad q^0 := \text{esssup}_{y \in \mathcal{O}} q(y),$$

$$q'(y) := \frac{q(y)}{q(y) - 1}, \quad y \in \mathcal{O}.$$

Also, let $L^r(\mathcal{O})$ be a standard Lebesgue space, where $r \geq 1$ is a number, $L^{q(y)}(\mathcal{O})$ be a generalized Lebesgue space, where $q \in \mathcal{B}_+(\mathcal{O})$ is a function,

$$\mathcal{P}^{\log}(\mathcal{O}) := \{q \in L^\infty(\mathcal{O}) \mid$$

$$q \text{ is a globally log-Hölder continuous, } q_0 \geq 1\}.$$

Let $W = W(t, \omega) : \Theta_{0,T} \rightarrow \mathbb{R}$ denotes the standard Wiener process [1, p. 109] and by $L_r(\mathbb{S}) \equiv L_r(\mathbb{S}, \mathcal{F}, \mathbb{P})$ a random Lebesgue space, consisting of all random variables $\xi : \mathbb{S} \rightarrow \mathbb{R}$ that have a finite r -th absolute momentum. Similarly, let $L_p(\mathbb{S}; B)$ be the random Lebesgue – Bochner spaces of all B -valued random variables $\xi : \mathbb{S} \rightarrow B$, where B is some Banach space.

Let $q \in \mathcal{B}_+(\mathcal{O})$ and $q_0 > 1$. Consider a generalized Lebesgue random space

$$L_{q(y)}(\mathcal{O} \times \mathbb{S}) := \left\{ z : \mathcal{O} \times \mathbb{S} \rightarrow \mathbb{R} \mid$$

$$1) \ z \text{ is a measurable function,}$$

$$2) \ \rho_q(z; \mathcal{O} \times \mathbb{S}) := \int_{\mathcal{O}} \int_{\mathbb{S}} |z(y, \omega)|^{q(y)} \mathbb{P}(d\omega) dy < +\infty \right\}$$

with the Luxembourg norm

$$\|z; L_{q(y)}(\mathcal{O} \times \mathbb{S})\| := \inf \{ \lambda > 0 \mid \rho_q(z/\lambda; \mathcal{O} \times \mathbb{S}) \leq 1 \}.$$

We will substitute \mathcal{O} with $(0, T)$, Ω , $\mathcal{Q}_{0,T}$, etc. below. Assume that the following conditions for the problem coefficients are satisfied:

$$\mathbf{(P)}: \ p, q \in \mathcal{P}^{\log}(\mathcal{Q}_{0,T});$$

$$\mathbf{(Z)}: \ \alpha > 0, \ r_0 = \min\{2, p_0, q_0\}, \ r^0 = \max\{2, p^0, q^0\}, \ s \in \mathbb{N},$$

$$s \geq \frac{1}{2} \max \left\{ 2, 1 + \frac{n(p^0 - 2)}{2p^0}, \frac{n(q^0 - 2)}{2q^0} \right\};$$

- (A):** $a_{ik} \in L^0(\mathbb{Q}_{0,T})$, $0 < a_0 \leq a_{ik}(x, t) \leq a^0 < +\infty$,
for almost every $(x, t) \in \mathbb{Q}_{0,T}$, where $i = 1, \dots, n$, $k = 1, \dots, N$;
- (G):** $g_k \in L^0(\mathbb{Q}_{0,T})$, $0 < g_0 \leq g_k(x, t) \leq g^0 < +\infty$,
for almost every $(x, t) \in \mathbb{Q}_{0,T}$, where $k = 1, \dots, N$;
- (Φ):** $\varphi_k \in \text{Lip}(\mathbb{R})$, $|\varphi_k(\xi)| \leq \varphi^0 |\xi|$,
for all $\xi \in \mathbb{R}$, where $0 \leq \varphi^0 < +\infty$, $k = 1, \dots, N$;
- (E):** $\mathfrak{z}_k \in L^0(\mathbb{Q}_{0,T} \times \Omega)$, $|\mathfrak{z}_k(x, t, y)| \leq \mathfrak{z}^0 < +\infty$,
for almost every $(x, t, y) \in \mathbb{Q}_{0,T} \times \Omega$, where $k = 1, \dots, N$;
- (F):** $F := (F_1, \dots, F_N) \in L_2(\mathbb{S}; L^2(0, T; H))$;
- (U):** $u_0 \in L_2(\mathbb{S}; H)$;
- (W):** W is the standard Wiener process [1, p. 109], $b := (b_1, \dots, b_N)$,
 $b_k(x, t, \omega) = b_{0,k}(x)W(t, \omega)$, $(x, t, \omega) \in \Pi_{0,T}$, $k = 1, \dots, N$,

and $b_0 := (b_{0,1}, \dots, b_{0,N}) \in [C_0^\infty(\Omega)]^N$.

To determine a solution to the problem (1)–(3), let us introduce the following notation:

$$\mathcal{W}_s := \{v \in H^{2s}(\Omega) \mid v|_{\partial\Omega} = \Delta v|_{\partial\Omega} = \dots = \Delta^{s-1}v|_{\partial\Omega} = 0\},$$

$$Z := \mathcal{W}_1, \quad X(t) := W_0^{1,p(x,t)}(\Omega), \quad Y(t) := L^{q(x,t)}(\Omega),$$

$$H := [L^2(\Omega)]^N, \quad V(t) := Z^N \cap [X(t)]^N \cap [Y(t)]^N \cap H,$$

$$X^*(t) := [X(t)]^*, \quad Y^*(t) := [Y(t)]^*, \quad H^* \simeq H, \quad V^*(t) := [V(t)]^*,$$

$$U(\mathbb{Q}_{0,T}) := \left\{ u : (0, T) \rightarrow V(t) \mid D^\alpha u \in [L^2(\mathbb{Q}_{0,T})]^N, |\alpha| = 2, \right.$$

$$u_{x_1}, \dots, u_{x_n} \in [L^{p(x,t)}(\mathbb{Q}_{0,T})]^N,$$

$$\left. u \in [L^{q(x,t)}(\mathbb{Q}_{0,T})]^N \cap [L^2(\mathbb{Q}_{0,T})]^N \right\},$$

$$W(\mathbb{Q}_{0,T}) := \{w \in U(\mathbb{Q}_{0,T}) \mid w_t \in [U(\mathbb{Q}_{0,T})]^*\}.$$

It is known that $W(\mathbb{Q}_{0,T}) \subset C([0, T]; H)$ and that for functions from space $W(\mathbb{Q}_{0,T})$, the formula for integration by parts with respect to the time variable holds (see, for example, Proposition 3.26 [11, p. 95]). Consider the following operators independent of $\omega \in \mathbb{S}$:

$$(Nz)_k(x, t) := g_k(x, t) |z(x, t)|^{q(x,t)-2} z_k(x, t),$$

$$(\mathbf{E}z)_k(x, t) := \int_{\Omega} \mathfrak{Z}_k(x, t, y) z_k(y, t) dy,$$

$$z = (z_1, \dots, z_N), \quad (x, t) \in \mathcal{Q}_{0,T},$$

$$\mathbf{N}z := ((\mathbf{N}z)_1, \dots, (\mathbf{N}z)_N), \quad \mathbf{E}z := ((\mathbf{E}z)_1, \dots, (\mathbf{E}z)_N).$$

Let $\mathbf{S} : U(\mathcal{Q}_{0,T}) \rightarrow [U(\mathcal{Q}_{0,T})]^*$ and $S(t) : V(t) \rightarrow V^*(t)$ be as follows:

$$\begin{aligned} \langle S(t)z, w \rangle_{V(t)} &:= \int_{\Omega} \sum_{k=1}^N \left[\alpha \Delta z_k \Delta w_k + \right. \\ &\quad \left. + \sum_{i=1}^n a_{ik}(x, t) |z_{x_i}|^{p(x,t)-2} z_{k,x_i} w_{k,x_i} + \right. \\ &\quad \left. + (\mathbf{N}z)_k w_k + \varphi_k((\mathbf{E}z)_k) w_k \right] dx, \\ &z, w \in V(t), \quad t \in (0, T); \end{aligned}$$

$$\langle \mathbf{S}u, v \rangle_{U(\mathcal{Q}_{0,T})} := \int_0^T \langle S(t)u(t), v(t) \rangle_{V(t)} dt, \quad u, v \in U(\mathcal{Q}_{0,T}).$$

To introduce the definition of the solution to problem (1)–(3), we formally substitute the unknown function u with \tilde{u} according to the rule

$$u(x, t, \omega) = \tilde{u}(x, t, \omega) + b(x, t, \omega), \quad (4)$$

where function b is taken from condition **(W)**. Since

$$b|_{x \in \partial\Omega} = \Delta b|_{x \in \partial\Omega} = 0, \quad b|_{t=0} = 0,$$

we obtain the following problem for the function $\tilde{u} = (\tilde{u}_1, \dots, \tilde{u}_N) : \Pi_{0,T} \rightarrow \mathbb{R}^N$:

$$\begin{aligned} &\tilde{u}_{k,t} + \alpha \Delta^2 (\tilde{u}_k + b_k) - \\ &\quad - \sum_{i=1}^n \left[a_{ik}(x, t) |\tilde{u}_{x_i} + b_{x_i}(x, t, \omega)|^{p(x,t)-2} (\tilde{u}_k + b_k(x, t, \omega))_{x_i} \right]_{x_i} + \\ &\quad + (\hat{\mathbf{N}}(\tilde{u} + b))_k(x, t, \omega) + \varphi_k((\hat{\mathbf{E}}(\tilde{u} + b))_k(x, t, \omega)) = F_k(x, t, \omega), \\ &\quad (x, t, \omega) \in \Pi_{0,T}, \end{aligned} \quad (5)$$

$$\tilde{u}(x, t, \omega) = \Delta \tilde{u}(x, t, \omega) = 0, \quad x \in \partial\Omega, \quad (t, \omega) \in \Theta_{0,T}, \quad (6)$$

$$\tilde{u}(x, 0, \omega) = u_0(x, \omega), \quad x \in \Omega, \quad \omega \in \mathbb{S}. \quad (7)$$

This problem doesn't contain the white noise component b_t anymore. Now we can determine the solution.

Definition 1. Function $\tilde{u} : \Pi_{0,T} \rightarrow \mathbb{R}^N$ is called a weak solution to problem (5)–(7) if

- 1) $\tilde{u} \in W(\mathcal{Q}_{0,T})$ almost surely;
- 2) function \tilde{u} almost surely satisfies the equality

$$\int_{\mathcal{Q}_{0,T}} \sum_{k=1}^N \left[-\tilde{u}_k v_{k,t} + \alpha \Delta (\tilde{u}_k + b_k) \Delta v_k + \right.$$

$$\begin{aligned}
& + \sum_{i=1}^n a_{ik}(x, t) |\tilde{u}_{x_i} + b_{x_i}|^{p(x,t)-2} (\tilde{u}_{k,x_i} + b_{k,x_i}) v_{k,x_i} + \\
& + (\hat{N}(\tilde{u} + b))_k v_k + \varphi_k((\hat{E}(\tilde{u} + b))_k) v_k \Big] dx dt = \\
& = \int_{Q_{0,T}} \sum_{k=1}^N F_k v_k dx dt
\end{aligned}$$

for all test functions $v \in U(Q_{0,T})$, i.e., in the sense of the spaces $[U(Q_{0,T})]^*$ and $[D^*(Q_{0,T})]^N$ the equality holds almost surely

$$\tilde{u}_t + S(\tilde{u} + b) = F;$$

- 3) \tilde{u} almost surely satisfies condition (7) in the sense of the space $C([0, T]; H)$.

Definition 2. Function $u : \Pi_{0,T} \rightarrow \mathbb{R}^N$ is called a weak solution to problem (1)–(3) if u has form (4) and the function \tilde{u} is a weak solution to problem (5)–(7) in the sense of Definition 1.

The main result of the article is the following statement.

Theorem 1. Let conditions **(P)**–**(W)** hold, the constant s is taken from condition **(Z)**, and $\partial\Omega \in C^{2s}$. Then problem (1)–(3) has a unique weak solution u . Furthermore,

$$\begin{aligned}
u & \in L_2(\mathbb{S}; C([0, T]; H) \cap L^2(0, T; Z)), \\
u & \in L_{q(x,t)}(\Pi_{0,T}), \quad u_{x_1}, \dots, u_{x_N} \in L_{p(x,t)}(\Pi_{0,T}). \tag{8}
\end{aligned}$$

Similar to (1)–(3), the problems without random parameter were considered in [9]. White noise-perturbed integro-differential systems (1) are considered, apparently, for the first time.

P r o o f. First, we introduce the random argument function

$$\mathbb{k} : \mathbb{S} \rightarrow H \times L^2(0, T; H) \times U(Q_{0,T})$$

according to the rule

$$\mathbb{k}(\omega) := \{u_0(\cdot, \omega), F(\cdot, \cdot, \omega), b(\cdot, \cdot, \omega)\}, \quad \omega \in \mathbb{S}. \tag{9}$$

Let

$$\mathfrak{R} : H \times L^2(0, T; H) \times U(Q_{0,T}) \rightarrow C([0, T]; H) \cap U(Q_{0,T})$$

be a function such that

$$\mathfrak{R}\{u_0, F, b\} = \tilde{u}, \tag{10}$$

where \tilde{u} is the solution of problem (5)–(7) with a fixed random parameter $\omega \in \mathbb{S}$. Using the Faedo – Galerkin method, we prove the existence of the solution to this problem. Since the nonlinear components corresponding to the exponents p and q are monotonous, and the integral component belongs to system (1) with the Lipschitz nonlinearity, we can use the standard methods to prove that the solution of problem (5)–(7) with a fixed random parameter $\omega \in \mathbb{S}$ continuously depends on the input data u_0, F, b . Therefore, the function \mathfrak{R} is correctly defined and continuous.

Using (9) and (10) we define the function $\mathfrak{P} : \mathbb{S} \rightarrow C([0, T]; H) \cap U(Q_{0,T})$ according to the following rule:

$$\mathfrak{P}(\omega) := \mathfrak{R} \circ \mathbb{k}(\omega) = \mathfrak{R}(\mathbb{k}(\omega)), \quad \omega \in \mathbb{S}. \quad (11)$$

Therefore, for each $\omega \in \mathbb{S}$ the value of $\mathfrak{P}(\omega)$ is equal to $\tilde{u}(\cdot, \cdot, \omega)$, where \tilde{u} is the solution to problem (5)–(7) without the random parameter ω . Since the function \mathbb{k} from (9) is measurable, and \mathfrak{R} is continuous, the function \mathfrak{P} from (11) is measurable. Thus, \tilde{u} will be a $C([0, T]; H) \cap U(Q_{0, T})$ -valued random variable. In addition, \tilde{u} satisfies the following estimates:

$$\int_{\Omega} |\tilde{u}(x, \tau, \omega)|^2 dx \leq C_1 \mathbf{F}(\tau, \omega), \quad (12)$$

$$\begin{aligned} \int_{Q_{0, \tau}} \left[|\Delta \tilde{u}(x, t, \omega)|^2 + \sum_{i=1}^n |\tilde{u}_{x_i}(x, t, \omega)|^{p(x, t)} + \right. \\ \left. + |\tilde{u}(x, t, \omega)|^{q(x, t)} \right] dx dt \leq C_1 \mathbf{F}(\tau, \omega), \end{aligned} \quad (13)$$

where the constant $C_1 > 0$ does not depend on \tilde{u} , ω , u_0 , F , b ,

$$\begin{aligned} \mathbf{F}(\tau, \omega) = \int_{\Omega} |u_0(x, \omega)|^2 dx + \int_{Q_{0, \tau}} \left[|F(x, t, \omega)|^2 + |\Delta b(x, t, \omega)|^2 + \right. \\ \left. + \sum_{i=1}^n |b_{x_i}(x, t, \omega)|^{p(x, t)} + |b(x, t, \omega)|^{q(x, t)} + \right. \\ \left. + |b(x, t, \omega)|^2 \right] dx dt, \end{aligned}$$

$\tau \in (0, T]$, $\omega \in \mathbb{S}$. Let's take the maximum for τ in (12) and $\tau = T$ in (13). Since the function $\mathbf{F}(T, \cdot)$ belongs to $L_1(\mathbb{S})$, we can obtain (by integrating the obtained equalities with respect to the variable $\omega \in \mathbb{S}$) the existence of integrals, that guarantee that (8) holds.

We prove the uniqueness of the solution to the problem by contradiction. Suppose that \tilde{u}^1 and \tilde{u}^2 are two solutions to problem (5)–(7), $\tilde{u} = \tilde{u}^1 - \tilde{u}^2$. Then we have

$$\begin{aligned} \frac{1}{2} \int_{\Omega} |\tilde{u}(x, \tau, \omega)|^2 dx + \int_{Q_{0, \tau}} \sum_{k=1}^N \left[\alpha |\Delta \tilde{u}_k|^2 + \right. \\ \left. + \sum_{i=1}^n a_{ik}(x, t) (|\tilde{u}_{x_i}^1 + b_{x_i}|^{p(x, t)-2} (\tilde{u}_{k, x_i}^1 + b_{k, x_i}) - \right. \\ \left. - |\tilde{u}_{x_i}^2 + b_{x_i}|^{p(x, t)-2} (\tilde{u}_{k, x_i}^2 + b_{k, x_i})) \times \right. \\ \left. \times ((\tilde{u}_{k, x_i}^1 + b_{k, x_i}) - (\tilde{u}_{k, x_i}^2 + b_{k, x_i})) + \right. \\ \left. + g_k(x, t) (|\tilde{u}^1 + b|^{q(x, t)-2} (\tilde{u}_k^1 + b_k) - \right. \\ \left. - |\tilde{u}^2 + b|^{q(x, t)-2} (\tilde{u}_k^2 + b_k)) \times \right. \\ \left. \times ((\tilde{u}_k^1 + b_k) - (\tilde{u}_k^2 + b_k)) + \left. \left\{ \varphi_k((\mathbf{E}(\tilde{u}^1 + b))_k) - \right. \right. \right. \end{aligned}$$

$$- \varphi_k((\mathbf{E}(\tilde{u}^2 + b))_k) \Big\} \tilde{u}_k \Big] dx dt = 0,$$

$$\tau \in (0, T], \quad \omega \in \mathbb{S}.$$

Discarding the monotonous components, we obtain the following estimate

$$\frac{1}{2} \int_{\Omega} |\tilde{u}(x, \tau, \omega)|^2 dx + \int_{Q_{0,\tau}} \alpha |\Delta \tilde{u}|^2 dx dt \leq \int_{Q_{0,\tau}} \sum_{k=1}^N |\varphi_k((\mathbf{E}(\tilde{u}^1 + b))_k) - \varphi_k((\mathbf{E}(\tilde{u}^2 + b))_k)| |\tilde{u}_k| dx dt. \quad (14)$$

From (14), conditions **(Φ)**–**(E)**, and the Gronwall – Belman lemma, we get, in particular, that

$$\int_{\Omega} |\tilde{u}(x, \tau, \omega)|^2 dx \leq 0, \quad \tau \in (0, T), \quad \omega \in \mathbb{S}.$$

Integrating by $\tau \in (0, T)$ and $\omega \in \mathbb{S}$, we obtain that $\tilde{u} = 0$, i.e. $\tilde{u}^1 = \tilde{u}^2$, for example, in the sense of the space $L_2(\Pi_{0,T})$. Theorem 1 is proved. \blacklozenge

1. Бугрій Н., Бугрій О., Доманська О. Напівлінійне стохастичне параболічне рівняння зі змінним показником нелінійності // Вісн. Львів. ун-ту. Сер. мех.-мат. – 2022. – Вип. 93. – С. 108–121. – <http://doi.org/10.30970/vmm.2022.93.108-121>.
2. Бугрій О. М. Про існування слабких розв'язків мішаних задач для напівлінійних параболічних за Петровським систем зі змінними показниками нелінійності // Укр. мат. журн. – 2014. – **66**, № 4. – С. 435–444.
Engl. translation: Buhrii O. M. On the existence of mild solutions of the initial-boundary-value problems for the Petrovskii-type semilinear parabolic systems with variable exponents of nonlinearity // Ukr. Math. J. – 2014. – **66**, No. 4. – P. 487–498. – <https://doi.org/10.1007/s11253-014-0947-2>.
3. Bauzet C., Vallet G., Wittbold P., Zimmermann A. On a $p(t, x)$ -Laplace evolution equation with a stochastic force // Stoch. PDE: Anal. Comp. – 2013. – **1**. – P. 552–570. – <https://doi.org/10.1007/s40072-013-0017-z>.
4. Bensoussan A., Temam R. Equations stochastiques du type Navier – Stokes // J. Funct. Anal. – 1973. – **13**, No. 2. – P. 195–222. – [https://doi.org/10.1016/0022-1236\(73\)90045-1](https://doi.org/10.1016/0022-1236(73)90045-1).
5. Bokalo M. M., Domanska O. V. Initial-boundary value problem for higher-orders nonlinear elliptic-parabolic equations with variable exponents of the nonlinearity in unbounded domains without conditions at infinity // Mat. studii. – 2023. – **59**, № 1. – С. 86–105. – <https://doi.org/10.30970/ms.59.1.86-105>.
6. Bokalo M., Buhrii O., Hryadii N. Initial-boundary value problems for nonlinear elliptic-parabolic equations with variable exponents of nonlinearity in unbounded domains without conditions at infinity // Nonlinear Anal. – 2020. – **192**. – Article No. 111700. – <https://doi.org/10.1016/j.na.2019.111700>.
7. Buhrii O. M. Visco-plastic, Newtonian, and dilatant fluids: Stokes equations with variable exponent of nonlinearity // Mat. studii. – 2018. – **49**, № 2. – С. 165–180. – <https://doi.org/10.15330/ms.49.2.165-180>.
8. Buhrii O. M., Kholyavka O. T., Bokalo T. M. Nonlocal hyperbolic Stokes system with variable exponent of nonlinearity // Mat. studii. – 2023. – **60**, № 2. – С. 173–179. – <https://doi.org/10.30970/ms.60.2.173-179>.
9. Buhrii O., Buhrii N. Integro-differential systems with variable exponents of nonlinearity // Open Math. – 2017. – **15**, No. 1. – P. 859–883. – <https://doi.org/10.1515/math-2017-0069>.
10. Coayla-Teran E. A., Ferreira J., Magalhães P. M. D. Weak solutions for random nonlinear parabolic equations of nonlocal type // Random Oper. Stoch. Equat. – 2008. – **16**, No. 3. – P. 213–223. – <https://doi.org/10.1515/ROSE.2008.011>.
11. Kaltenbach A. Pseudo-monotone operator theory for unsteady problems with variable exponents. – Ser. Lect. Notes Math. – Vol. 2329. – Cham: Springer, 2023. – xiii+358 p.

12. Kuratowski K., Ryll-Nardzewski C. A general theorem on selectors // Bull. Acad. Polon. Sci. Sér. Math. Sci. Astronom. Phys. – 1965. – **13**. – P. 397–403.
13. Pardoux É. Stochastic Partial differential equations: An introduction. – Ser. Springer Briefs in Mathematics. – Cham: Springer, 2021. – VIII+74 p.
– <https://doi.org/10.1007/978-3-030-89003-2>.
14. Ren J., Röckner M., Wang F. Y. Stochastic generalized porous media and fast diffusion equations // J. Diff. Equat. – 2007. – **238**, No. 1. – P. 118–152.
– <https://doi.org/10.1016/j.jde.2007.03.027>.

СТОХАСТИЧНІ ПАРАБОЛІЧНІ СИСТЕМИ ЗІ ЗМІННИМИ ПОКАЗНИКАМИ НЕЛІНІЙНОСТІ

Розглянуто деякі нелінійні параболічні системи диференціальних рівнянь з частинними похідними з білим шумом. Досліджено початково-крайову задачу для системи та доведено існування та єдиність слабкого розв'язку задачі.

Ключові слова: диференціальні рівняння з частинними похідними, стохастична параболічна система, змінний показник нелінійності, білий шум, слабкий розв'язок.

Ivan Franko National University of Lviv, Lviv

Received
07.10.24