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# RESEARCH

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# Analysis of Stokes system with solution-dependent subdifferential boundary conditions

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### Abstract

We study the Stokes problem for the incompressible fluid with mixed nonlinear boundary conditions of subdifferential type. The latter involve a unilateral boundary condition, the Navier slip condition, a nonmonotone version of the nonlinear Navier–Fujita slip condition, and the threshold slip and leak condition of frictional type. The weak form of the problem leads to a new class of variational–hemivariational inequalities on convex sets for the velocity field. Solution existence and the weak compactness of the solution set to the inequality problem are established based on the Schauder fixed point theorem.

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## **1** Introduction

In this paper we study the Stokes problem for the incompressible fluid with mixed boundary conditions in a bounded domain of dimension two and three. The problem is formulated as a variational-hemivariational inequality of elliptic type involving both convex and locally Lipschitz, generally nonconvex, superpotentials. The main results concern the existence and compactness of the solution set.

We are motivated by the Stokes system with a frictional type boundary condition with the slip bound threshold value depending on the solution. Such a system has been treated in [10, 11] by the mixed variational formulation with the Lagrange multipliers and then applied to deal with optimum design problems. As mentioned in these papers, it was experimentally observed that the slip bound may depend on the solution itself, e.g., on values of the tangential component of the velocity. This situation may appear in several practical models of flows of polymer melts, blood flow in a vein, fluids on hydrophobic surfaces, problems with multiple interfaces, etc. The simplest threshold slip condition can be described by the Tresca-like condition when the threshold bound is given a priori and it is modeled by a convex potential, see [7, 9]. Note that when the slip bound depends on the

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solution, the model is more complicated, see [11, 14, 15, 17] and the non-stationary case in [8].

The main novelties of the paper are the following. First, we consider a generalization of slip and leak boundary condition of frictional type on diverse parts of the wall of the domain. For the slip condition we study the Clarke subgradient multivalued condition which, in contrast to [17], depends on the tangential velocity. The leak boundary condition can be governed by any convex and lower semicontinuous potential involving the normal velocity. Second, we have supplemented the model with an outflow unilateral boundary condition introduced just recently in [31]. This makes the Stokes problem more involved since the resulting variational-hemivariational inequality has an additional set of unilateral constraints. Third, in comparison with [10, 11], our approach and the proof is different and combines a recent result from the theory of variational-hemivariational inequalities and the Schauder fixed point theorem. In the proof of existence of solution, our main goal is to explore under which conditions concerning coefficients in slip and leak boundary conditions solutions to the Stokes system depend continuously on variations of parameters.

The important feature of the model is the nonlinearity of the form  $k\partial j$ . In such a case we cannot deal with a purely hemivariational inequality since there is not, in general, a potential G with  $G = k\partial j$ . This type of nonlinearity appears in the model twice: in the non-monotone slip boundary condition which is described by the Clarke generalized gradient, and in the generalized leak boundary condition of frictional type governed by the subdifferential of the convex function. The slip bound function k depends on the (norm of the) solution, while the potential  $j: \mathbb{R}^d \to \mathbb{R}$  is a locally Lipschitz function and  $\partial j$  stands for its generalized subgradient. We mention that it is an interesting open problem to generalize the results of this paper to non-Newtonian fluids with various nonlinear constitutive law.

We note that the variational and hemivariational inequalities have been used to solve the fluid flow problems in [18, 19, 22] for the stationary models and in [6] for the evolutionary problems.

The paper is organized as follows. In Sect. 2 we shortly recall our notation and preliminary results. Section 3 contains the classical and variational formulations of the Stokes problem and the statement of the existence theorem. Its proof is delivered in Sect. 4. Finally, in Sect. 5 we demonstrate the weak compactness of the solution set.

#### 2 Preliminary material

In this section we recall the standard notation and definitions from [1, 3, 4, 20, 21].

Throughout the paper, given a Banach space *X*, we denote by  $X^*$  its dual space, by  $\|\cdot\|_X$ a norm in *X*, and by  $\langle \cdot, \cdot \rangle_{X^* \times X}$  the duality brackets between  $X^*$  and *X*. For simplicity of notation, when no confusion arises, we often skip the subscripts. We use the notation  $x_n \to x$  and  $x_n \to x$ , respectively, to denote the strong convergence and weak convergence in various spaces. By  $\mathcal{L}(X_1, X_2)$  we denote the space of linear and bounded operators from a normed space  $X_1$  to a normed space  $X_2$  endowed with the operator norm  $\|\cdot\|_{\mathcal{L}(X_1, X_2)}$ . Given a set  $D \subset X$ , we set  $\|D\|_X = \sup\{\|x\|_X \mid x \in D\}$ .

A single-valued operator  $A: X \to X^*$  is pseudomonotone if it is bounded (it maps bounded subsets of X into bounded subsets of  $X^*$ ), and if  $u_n \to u$  in X and  $\limsup \langle Au_n, u_n - u \rangle \leq 0$  imply  $\langle Au, u - v \rangle \leq \liminf \langle Au_n, u_n - v \rangle$  for all  $v \in X$ . Equivalently, a single-valued operator A is pseudomonotone if and only if it is bounded, and  $u_n \to u$  in X together with  $\limsup \langle Au_n, u_n - u \rangle \leq 0$  yields  $\lim \langle Au_n, u_n - u \rangle = 0$  and  $Au_n \to Au$  in  $X^*$ . Let  $\varphi: X \to \mathbb{R} \cup \{+\infty\}$  be a proper, convex, and lower semicontinuous function. The mapping  $\partial \varphi: X \to 2^{X^*}$  defined by

$$\partial \varphi(x) = \left\{ x^* \in X^* \mid \left\langle x^*, \nu - x \right\rangle \le \varphi(\nu) - \varphi(x) \text{ for all } \nu \in X \right\}$$

is called the convex subdifferential of  $\varphi$ . An element  $x^* \in \partial \varphi(x)$  is called a subgradient of  $\varphi$  in x. Let  $h: X \to \mathbb{R}$  be a locally Lipschitz function on a Banach space X. The generalized (Clarke) directional derivative of h at  $x \in X$  in the direction  $v \in X$  is defined by

$$h^0(x; \nu) = \limsup_{y \to x, \lambda \downarrow 0} \frac{h(y + \lambda \nu) - h(y)}{\lambda}.$$

The generalized (Clarke) gradient of *h* at *x*, denoted by  $\partial h(x)$ , is a subset of the dual space  $X^*$  given by

$$\partial h(x) = \{ \zeta \in X^* \mid h^0(x; \nu) \ge \langle \zeta, \nu \rangle \text{ for all } \nu \in X \}.$$

Finally, we recall an existence and uniqueness result for a class of abstract variational-hemivariational inequalities. Let *X* be a reflexive Banach space. Given an operator  $A: X \to X^*$ , functions  $\varphi: K \times K \to \mathbb{R}$ ,  $j: X \to \mathbb{R}$ , and a set  $K \subset X$ , we consider the following problem.

**Problem 1** Find an element  $u \in K$  such that

$$\langle Au-f, v-u \rangle + \Phi(v) - \Phi(u) + J^0(u; v-u) \ge 0$$
 for all  $v \in K$ .

For this problem, we need the following hypotheses on the data.  $H(A): A: X \to X^*$  is a function such that

- (i) it is pseudomonotone,
- (ii) it is strongly monotone, i.e., there exists  $m_A > 0$  such that

$$\langle Av_1 - Av_2, v_1 - v_2 \rangle \ge m_A ||v_1 - v_2||_X^2$$
 for all  $v_1, v_2 \in X$ .

 $H(\Phi): \Phi: X \to \mathbb{R}$  is convex and lower semicontinuous.

 $H(J): J: X \to \mathbb{R}$  is a function such that

- (i) *J* is locally Lipschitz,
- (ii)  $\|\partial J(\nu)\|_{X^*} \le c_0 + c_1 \|\nu\|_X$  for all  $\nu \in X$  with  $c_0, c_1 \ge 0$ ,
- (iii) there exists  $\alpha_J \ge 0$  such that  $J^0(\nu_1; \nu_2 \nu_1) + J^0(\nu_2; \nu_1 \nu_2) \le \alpha_J ||\nu_1 \nu_2||_X^2$  for all  $\nu_1$ ,  $\nu_2 \in X$ .

H(K): *K* is nonempty, closed and convex subset of *X*. H(f): *f* ∈ *X*<sup>\*</sup>.

Recall the existence and uniqueness result for Problem 1.

**Theorem 2** Assume H(A),  $H(\Phi)$ , H(J), H(K), H(f) and the smallness condition

 $\alpha_J < m_A. \tag{1}$ 

*Then Problem* 1 *has a unique solution*  $u \in K$ *.* 

Theorem 2 represents a particular case of a result proved in [21], where the function  $\Phi$  depends additionally on the solution. Further, in [21, Theorem 18], the authors required additionally that *A* is coercive. However, this assumption is redundant, since if *A* is strongly monotone, then *A* is coercive in the following sense:

$$\langle A\nu, \nu \rangle = \langle A\nu - A0, \nu \rangle + \langle A0, \nu \rangle \ge m_A \|\nu\|_X^2 + \|A0\|_{X^*} \|\nu\|_X$$
(2)

for all  $v \in X$ . Condition H(J)(iii) has been extensively used in the literature for hemivariational inequalities, see [20, 29], and it is equivalent to the relaxed monotone condition

$$\langle z_1 - z_2, \nu_1 - \nu_2 \rangle \ge -\alpha_I \|\nu_1 - \nu_2\|_X^2 \tag{3}$$

for all  $z_i \in \partial J(v_i)$ ,  $z_i \in X^*$ ,  $v_i \in X$ , i = 1, 2. Further, if J is a convex function, then H(J)(iii) or equivalently (3) holds with  $\alpha_J = 0$  and means that the (convex) subdifferential is a monotone map. In this case the smallness condition (1) holds trivially.

#### **3** Formulation of the Stokes problem

Let  $\Omega \subset \mathbb{R}^d$ , d = 2, 3, be a connected Lipschitz bounded domain occupied by the incompressible fluid. The boundary  $\Gamma = \partial \Omega$  consists of smooth parts  $\Gamma_0$ ,  $\Gamma_1$ ,  $\Gamma_2$ , and  $\Gamma_3$  such that meas( $\Gamma_0$ ) > 0, while the parts  $\Gamma_1$ ,  $\Gamma_2$ , and  $\Gamma_3$  can be empty. The unit outward normal vector exists a.e. on the boundary and is denoted by  $\boldsymbol{\nu}$ . Moreover,  $\mathbb{M}^d$  denotes the class of symmetric  $d \times d$  matrices.

The classical formulation of the Stokes problem with mixed boundary conditions studied in this paper reads as follows.

**Problem 3** Find a flow velocity  $\mathbf{u}: \Omega \to \mathbb{R}^d$ , and a pressure  $p: \Omega \to \mathbb{R}$  such that

$$-\operatorname{Div}(2\widetilde{\mu}\mathbb{D}\boldsymbol{u}) + \nabla p = \boldsymbol{f} \quad in \ \Omega, \tag{4}$$

$$\operatorname{div} \boldsymbol{u} = 0 \quad in \ \Omega, \tag{5}$$

 $\boldsymbol{u} = 0 \quad on \ \Gamma_0, \tag{6}$ 

$$\begin{cases} u_{\nu} = 0, & on \ \Gamma_1, \\ -\boldsymbol{\tau}_{\tau}(\boldsymbol{u}) \in k(\|\boldsymbol{u}_{\tau}\|) \partial j(\boldsymbol{u}_{\tau}), \end{cases}$$
(7)

$$\begin{cases} u_{\nu} + g \ge 0, & \tau_{\nu}(\boldsymbol{u}, p) + h(u_{\nu}) \ge 0, \\ (u_{\nu} + g)(\tau_{\nu}(\boldsymbol{u}, p) + h(u_{\nu})) = 0, & on \ \Gamma_{2}, \\ \boldsymbol{\tau}_{\tau}(\boldsymbol{u}) = 0, \end{cases}$$
(8)

$$\begin{aligned} &-\tau_{\nu}(\boldsymbol{u},p) \in l(\boldsymbol{u}_{\nu})\partial\phi(\boldsymbol{u}_{\nu}), \\ &\boldsymbol{u}_{\tau}=0, \end{aligned} \qquad on \ \Gamma_{3}. \end{aligned}$$

We give a short description of the Stokes system in Problem 3. The extra stress tensor  $\mathbb{S}$  is defined by the linear constitutive law  $\mathbb{S}(\mathbb{D}\boldsymbol{u}) = 2\widetilde{\mu}\mathbb{D}\boldsymbol{u}$  in  $\Omega$ , where  $\widetilde{\mu}$  represents the dynamic viscosity,  $\boldsymbol{f}$  is the volume force, and the deformation-rate tensor is given by  $\mathbb{D}\boldsymbol{u} = \frac{1}{2}(\nabla \boldsymbol{u} + \nabla \boldsymbol{u}^{\top})$ . To simplify presentation, we shall suppose in what follows that  $2\widetilde{\mu} = \mu$ .

The divergence free condition (5) models an incompressible fluid. The divergence operators for tensor and vector-valued functions are defined by  $\text{Div } \mathbb{S} = (\mathbb{S}_{ij,j})$  and  $\text{div } \mathbf{u} = (u_{i,i})$ , where the index that follows a comma represents the partial derivative with respect to the corresponding component of  $\mathbf{x}$ . The homogeneous Dirichlet boundary condition (6) means that the fluid adheres to the wall, see, e.g., [7, 8, 14, 17, 27]. For convenience in notation, in boundary conditions (7)–(9), the traction vector is defined by

$$\boldsymbol{\tau}(\boldsymbol{u},p) = \boldsymbol{\sigma}(\boldsymbol{u},p)\boldsymbol{v} \quad \text{on } \boldsymbol{\Gamma},$$

where the total stress is denoted by

$$\boldsymbol{\sigma}(\boldsymbol{u},\pi) = -p\mathbb{I} + \mathbb{S}(\mathbb{D}\boldsymbol{u}) \quad \text{in } \Omega,$$

and  $\mathbb{I}$  is the identity  $d \times d$  matrix. Hence

$$\tau_{v}(\boldsymbol{u},p) = \boldsymbol{\tau}(\boldsymbol{u},p) \cdot \boldsymbol{v}$$
 and  $\boldsymbol{\tau}_{\tau}(\boldsymbol{u}) = \boldsymbol{\tau}(\boldsymbol{u},p) - \tau_{v}(\boldsymbol{u},p)\boldsymbol{v}$ 

represent normal and tangential components of the traction vector, respectively. Note that  $\boldsymbol{\tau}$  and  $\tau_{\nu}$  do depend on the pressure p, and  $\boldsymbol{\tau}_{\tau}$  is independent of p.

The nonlinear boundary condition (7) describes a generalization, in several directions, of the Navier–Fujita slip condition. The following example of (7) has been studied by Le Roux and Tani in [14, 15]:

$$u_{\nu} = 0 \quad \text{on } \Gamma_1, \tag{10}$$

$$\begin{cases} \|\boldsymbol{\tau}_{\tau}(\boldsymbol{u})\|_{\mathbb{R}^{d}} \leq \alpha + \beta(\|\boldsymbol{u}_{\tau} - \boldsymbol{w}_{\tau}\|_{\mathbb{R}^{d}}), & \text{on } \Gamma_{1}, \\ \boldsymbol{\tau}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{u}_{\tau} - \boldsymbol{w}_{\tau}) = -(\alpha + \beta(\|\boldsymbol{u}_{\tau} - \boldsymbol{w}_{\tau}\|_{\mathbb{R}^{d}}))\|\boldsymbol{u}_{\tau} - \boldsymbol{w}_{\tau}\|_{\mathbb{R}^{d}}, \end{cases}$$
(11)

where  $\alpha \colon \Gamma_1 \to (0,\infty)$  and  $\beta \colon \Gamma_1 \times [0,\infty) \to [0,\infty)$  are prescribed functions such that for a.e.  $\mathbf{x} \in \Gamma_1$ ,  $\beta(\mathbf{x}, r) = 0$  if and only if r = 0, while  $\mathbf{w}_{\tau}$  denotes the tangential velocity of the wall surface at  $\Gamma_1$ . Conditions (10) and (11) in the aforementioned papers are motivated by models of flows of polymer melts during extraction, flows of Newtonian fluids with a moving contact line, and flows of yield-stress fluids, etc. Condition (10) is called the impermeability (no leak) boundary condition, and (11) represents a nonlinear Navier-Fujita slip condition. Physically, condition (11) signifies that for the wall slip to occur the magnitude of the tangential traction has to exceed a prescribed "slip threshold", denoted by  $\alpha$ , independent of the normal stress, and if the slip occurs, the tangential traction equals to a given, nonlinear function of the slip velocity. Condition (11) has been considered as a generalization of three slip boundary conditions: the Navier slip condition in [23] (stating that the tangential velocity  $\mathbf{u}_{\tau}$  is proportional to the shear stress  $\boldsymbol{\sigma}_{\tau}$ ), the nonlinear Naviertype slip condition in [14], and the threshold slip condition of "frictional type" studied in a series of papers by Fujita et al. [7–9, 25, 26]. Note that the nonlinear Navier–Fujita slip condition (11) is a particular case of condition (7) in Problem 3 with functions  $k(\mathbf{x}, \boldsymbol{\xi}) =$  $\alpha(\mathbf{x}) + \beta(\mathbf{x}, \|\mathbf{\xi}\|_{\mathbb{R}^d})$  and  $j(\mathbf{x}, \mathbf{\xi}) = \|\mathbf{\xi}\|_{\mathbb{R}^d}$  for a.e.  $\mathbf{x} \in \Gamma_1$ , all  $\mathbf{\xi} \in \mathbb{R}^d$ . The function *j* satisfies H(j)below with  $c_0(x) = 1$ ,  $c_1 = 0$  since it is convex in the second variable. Further, condition (7) is much more general than (11) since it involves models with nonmonotone slip boundary

condition of frictional type governed by nonconvex locally Lipschitz superpotentials *j*; for concrete examples, see [17, Example (60)] and the references therein.

Condition (8) has been recently suggested and studied in [31] to model a new outflow unilateral boundary condition for blood flow simulations. It has an advantage over the popular do-nothing boundary condition. In (8), *g* is a given positive constant and *h* is a prescribed nonnegative function. It is a counterpart of the Signorini-type condition which models the unilateral contact in the theory of elasticity [13, 28]. Moreover, condition (8) has been studied for the Stokes problem in [27] with  $h(\mathbf{x}; r) = h(\mathbf{x})$ , and recently in [17]. The boundary condition (9) is called the generalized leak boundary condition of frictional type. It describes the lack of slip on the boundary  $\Gamma_3$  and a possible leakage of the fluid through this part of the boundary with a prescribed convex function  $\phi$ . For the choice  $\phi(r) = |r|$  for  $r \in \mathbb{R}$  and further discussion, we refer to [17]. The generalized gradient for *j* and the convex subdifferential of  $\phi$  are always taken with respect to the last variable of a given function.

In a particular case, Problem 3 has been studied in [10, 11] under the hypotheses  $\Gamma_2 = \Gamma_3 = \emptyset$  and the convex potential  $j(\xi) = \|\xi\|_{\mathbb{R}^d}$  for  $\xi \in \mathbb{R}^d$ . There, a weak form of Problem 3 was treated as a mixed variational formulation which couples a nonlinear variational inequality and an equation for the multiplier. This formulation is quite different to the one we study in the present paper. Due to the presence of both convex and nonconvex potentials, the weak formulation of Problem 3 leads to a variational –hemivariational inequality of elliptic type for the velocity field.

Next, we introduce the following spaces needed for the weak formulation.

$$\widetilde{V} = \left\{ \boldsymbol{v} \in \mathcal{C}^{\infty}(\overline{\Omega}; \mathbb{R}^{d}) \mid \operatorname{div} \boldsymbol{v} = 0 \text{ in } \Omega, \boldsymbol{v} = 0 \text{ on } \Gamma_{0}, \nu_{\nu} = 0 \text{ on } \Gamma_{1}, \boldsymbol{v}_{\tau} = 0 \text{ on } \Gamma_{3} \right\},$$

$$V = \operatorname{closure of} \widetilde{V} \text{ in } H^{1}(\Omega; \mathbb{R}^{d}).$$
(12)

The space *V* is endowed with the standard norm  $\|\mathbf{v}\| = \|\mathbf{v}\|_{H^1(\Omega;\mathbb{R}^d)}$ , and we also consider the norm given by  $\|\mathbf{v}\|_V = \|\mathbb{D}\mathbf{v}\|_{L^2(\Omega:\mathbb{M}^d)}$  for  $\mathbf{v} \in V$ . Using the Korn inequality

$$c_{K} \|\boldsymbol{\nu}\|_{H^{1}(\Omega;\mathbb{R}^{d})} \leq \|\mathbb{D}\boldsymbol{\nu}\|_{L^{2}(\Omega;\mathbb{M}^{d})} \quad \text{for } \boldsymbol{\nu} \in V \text{ with } c_{K} > 0,$$
(13)

see [5, Theorem 8], it is known that  $\|\cdot\|_{H^1(\Omega;\mathbb{R}^d)}$  and  $\|\cdot\|_V$  are the equivalent norms on V. The set of unilateral constraints for Problem 3 is given by

$$K = \{ \mathbf{v} \in V \mid v_{\nu} + g \ge 0 \text{ on } \Gamma_2 \} \quad \text{with } g > 0.$$

$$(14)$$

We recall that, for a bounded Lipschitz domain with boundary  $\Gamma$ , there exists a unique trace operator  $\gamma_0: H^1(\Omega; \mathbb{R}^d) \to H^{1/2}(\Gamma; \mathbb{R}^d)$  which is linear, bounded, surjective, and such that  $\gamma_0(\mathbf{v}) = \mathbf{v}|_{\Gamma}$  for all  $\mathbf{v} \in C^{\infty}(\overline{\Omega}; \mathbb{R}^d)$ , see, e.g., [24, Chap. 2, Theorems 5.5 and 5.7]. There are various definitions of  $H^{1/2}(\Gamma; \mathbb{R}^d)$  which are equivalent for Lipschitz boundaries, for details, see, e.g., [30, Chap. II.1.2, p. 47], [2, Chap. 3.2]. It follows from [16, p. 93] that the injection map  $i: H^{1/2}(\Gamma; \mathbb{R}^d) \to L^2(\Gamma; \mathbb{R}^d)$  is compact. We use also the trace operator  $\gamma = i \circ \gamma_0: V \subset H^1(\Omega; \mathbb{R}^d) \to L^2(\Gamma; \mathbb{R}^d)$  which is linear, bounded, compact, and its norm in the space  $\mathcal{L}(V, L^2(\Gamma; \mathbb{R}^d))$  is denoted by  $\|\gamma\|$ . It is clear that the (tangential and normal) trace operator

$$H^{1}(\Omega; \mathbb{R}^{d}) \ni \boldsymbol{w} \mapsto (\boldsymbol{w}_{\tau}, \boldsymbol{w}_{\nu}) \in H^{1/2}(\Gamma; \mathbb{R}^{d}) \times H^{1/2}(\Gamma)$$
(15)

is linear and bounded, where for a vector-valued function  $\boldsymbol{w}$  on  $\Gamma$ , we denote by  $w_{\nu}$  and  $\boldsymbol{w}_{\tau}$  its normal and tangential components defined by  $w_{\nu} = \boldsymbol{w} \cdot \boldsymbol{\nu}$  and  $\boldsymbol{w}_{\tau} = \boldsymbol{w} - w_{\nu} \boldsymbol{\nu}$ , respectively. Further, instead of  $\gamma \boldsymbol{v}$ , for simplicity, we retain the notation  $\boldsymbol{v}$ .

We now introduce hypotheses on the dynamic viscosity  $\mu$ , the nonconvex superpotential *j*, the convex potential  $\phi$ , the functions *h*, *k*, and *l*, and the external body force density *f* in Problem 3.

 $H(\mu)$ :  $\mu \in L^{\infty}(\Omega)$  and  $0 < \mu_0 \le \mu(\mathbf{x})$  for a.e.  $\mathbf{x} \in \Omega$ .

 $H(h): h: \Gamma_2 \times \mathbb{R} \to \mathbb{R}$  is a function such that

(i)  $h(\cdot, r)$  is measurable for all  $r \in \mathbb{R}$ ,

- (ii)  $h(\mathbf{x}, \cdot)$  is continuous and nondecreasing for a.e.  $\mathbf{x} \in \Gamma_2$ ,
- (iii)  $0 \le h(\mathbf{x}, r) \le h_0$  for all  $r \in \mathbb{R}$  and a.e.  $\mathbf{x} \in \Gamma_2$ ,
- (iv)  $h(\boldsymbol{x}, 0) = 0$  for a.e.  $\boldsymbol{x} \in \Gamma_2$ .
- $H(j): j: \Gamma_1 \times \mathbb{R}^d \to \mathbb{R}$  is a function such that
- (i)  $j(\cdot, \boldsymbol{\xi})$  is measurable on  $\Gamma_1$  for all  $\boldsymbol{\xi} \in \mathbb{R}^d$  and there is  $e \in L^2(\Gamma_1; \mathbb{R}^d)$  such that  $j(\cdot, e(\cdot)) \in L^1(\Gamma_1)$ ,
- (ii)  $j(\mathbf{x}, \cdot)$  is locally Lipschitz for a.e.  $\mathbf{x} \in \Gamma_1$ ,
- (iii)  $\|\partial j(\mathbf{x}, \mathbf{\xi})\|_{\mathbb{R}^d} \le c_0(x) + c_1 \|\mathbf{\xi}\|_{\mathbb{R}^d}$  for all  $\mathbf{\xi} \in \mathbb{R}^d$ , a.e.  $\mathbf{x} \in \Gamma_1$  with  $c_0 \in L^2(\Gamma_1), c_0, c_1 \ge 0$ ,
- (iv) there exists  $\alpha_j \ge 0$  such that  $j^0(\boldsymbol{\xi}_1; \boldsymbol{\xi}_2 \boldsymbol{\xi}_1) + j^0(\boldsymbol{\xi}_2; \boldsymbol{\xi}_1 \boldsymbol{\xi}_2) \le \alpha_j \|\boldsymbol{\xi}_1 \boldsymbol{\xi}_2\|_{\mathbb{R}^d}^2$  for all  $\boldsymbol{\xi}_1, \boldsymbol{\xi}_2 \in \mathbb{R}^d$ .

 $H(k): k: \Gamma_1 \times \mathbb{R} \to \mathbb{R}$  is a function such that

- (i)  $k(\cdot, r)$  is measurable on  $\Gamma_1$  for all  $r \in \mathbb{R}$ ,
- (ii)  $k(\mathbf{x}, \cdot)$  is continuous on  $\mathbb{R}$  for a.e.  $\mathbf{x} \in \Gamma_1$ ,
- (iii)  $0 < k_0 \le k(\mathbf{x}, r) \le k_1$  for all  $r \in \mathbb{R}$ , a.e.  $\mathbf{x} \in \Gamma_1$ .
- $H(l): l: \Gamma_3 \times \mathbb{R} \to \mathbb{R}$  is a function such that
- (i)  $l(\cdot, r)$  is measurable on  $\Gamma_3$  for all  $r \in \mathbb{R}$ ,
- (ii)  $l(\mathbf{x}, \cdot)$  is continuous on  $\mathbb{R}$  for a.e.  $\mathbf{x} \in \Gamma_3$ ,
- (iii)  $0 < l_0 \le l(\mathbf{x}, r) \le l_1$  for all  $r \in \mathbb{R}$ , a.e.  $\mathbf{x} \in \Gamma_1$ .
- $H(\phi): \phi: \mathbb{R} \to \mathbb{R}$  is convex and lower semicontinuous.

$$H(f): \mathbf{f} \in V^*.$$

 $(H_0): \alpha_j k_1 \|\gamma\|^2 < \mu_0.$ 

The condition ( $H_0$ ) means that a "small" decrease (with respect to the lower bound  $\mu_0$  for the viscosity) of the graph of the subdifferential of *j* is permissible. As mentioned before, when  $j(\mathbf{x}, \cdot)$  is a convex potential, then ( $H_0$ ) holds trivially with  $\alpha_j = 0$ .

We now derive the variational formulation of Problem 3. We assume now that  $\boldsymbol{u}$  and p are sufficiently smooth functions which satisfy (4)–(9). Let  $\boldsymbol{v} \in K$ . We multiply equation (4) by  $\boldsymbol{v} - \boldsymbol{u}$  and integrate over  $\Omega$  to find that

$$\int_{\Omega} \left( -\operatorname{Div} \mathbb{S}(\mathbb{D}\boldsymbol{u}) \right) \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx + \int_{\Omega} \nabla p \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx = \int_{\Omega} \boldsymbol{f} \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx.$$
(16)

By using a second Green-type formula, see [20, Theorem 2.25], the fact that functions in V are divergence free,  $\mathbf{v} = \mathbf{u} = 0$  on  $\Gamma_0$  and  $v_v = u_v = 0$  on  $\Gamma_1$ , we obtain

$$\begin{split} \int_{\Omega} \nabla p \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx &= -\int_{\Omega} \left( \operatorname{div}(\boldsymbol{v} - \boldsymbol{u}) \right) p \, dx + \int_{\Gamma_0} (v_v - u_v) p \, d\Gamma \\ &+ \int_{\Gamma_1} (v_v - u_v) p \, d\Gamma + \int_{\Gamma_2} (v_v - u_v) p \, d\Gamma + \int_{\Gamma_3} (v_v - u_v) p \, d\Gamma \end{split}$$

$$= \int_{\Gamma_2} (v_\nu - u_\nu) p \, d\Gamma + \int_{\Gamma_3} (v_\nu - u_\nu) p \, d\Gamma. \tag{17}$$

By the Green formula in [20, Theorems 2.24], we get

$$\int_{\Omega} \left( -\operatorname{Div} \mathbb{S}(\mathbb{D}\boldsymbol{u}) \right) \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx$$
$$= \int_{\Omega} \mathbb{S}(\mathbb{D}\boldsymbol{u}) : \mathbb{D}(\boldsymbol{v} - \boldsymbol{u}) \, dx - \int_{\partial \Omega} \left( \mathbb{S}(\mathbb{D}\boldsymbol{u})\boldsymbol{v} \right) \cdot (\boldsymbol{v} - \boldsymbol{u}) \, d\Gamma.$$

Further, we employ the decomposition formula, see [20, relation (6.33)], condition  $\boldsymbol{u} = \boldsymbol{v} = 0$  on  $\Gamma_0$ , and

$$\mathbb{S}_{\nu}(\boldsymbol{u}) = \tau_{\nu}(\boldsymbol{u}, p) + p, \qquad \mathbb{S}_{\tau}(\boldsymbol{u}) = \boldsymbol{\tau}_{\tau}(\boldsymbol{u}) \quad \text{on } \Gamma,$$

to obtain

$$\begin{split} &\int_{\partial\Omega} (\mathbb{S}\boldsymbol{v}) \cdot (\boldsymbol{v} - \boldsymbol{u}) d\Gamma \\ &= \int_{\Gamma_1} \mathbb{S}_{\nu}(\boldsymbol{u})(\nu_{\nu} - u_{\nu}) + \mathbb{S}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma + \int_{\Gamma_2} \mathbb{S}_{\nu}(\boldsymbol{u})(\nu_{\nu} - u_{\nu}) \\ &+ \mathbb{S}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma + \int_{\Gamma_3} \mathbb{S}_{\nu}(\boldsymbol{u})(\nu_{\nu} - u_{\nu}) + \mathbb{S}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma \\ &= \int_{\Gamma_1} \boldsymbol{\tau}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma + \int_{\Gamma_2} (\boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p) + p)(\nu_{\nu} - u_{\nu}) d\Gamma \\ &+ \int_{\Gamma_3} (\boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p) + p)(\nu_{\nu} - u_{\nu}) d\Gamma = \int_{\Gamma_1} \boldsymbol{\tau}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma \\ &+ \int_{\Gamma_2} (\boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p) + h(u_{\nu}))(\nu_{\nu} + g) - (\boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p) + h(u_{\nu}))(g + u_{\nu}) \\ &- h(u_{\nu})(\nu_{\nu} - u_{\nu}) + p(\nu_{\nu} - u_{\nu}) d\Gamma + \int_{\Gamma_3} (\boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p) + p)(\nu_{\nu} - u_{\nu}) d\Gamma. \end{split}$$

Now, taking into account the relations

$$\big(\tau_{\nu}(\boldsymbol{u},p)+h(u_{\nu})\big)(\nu_{\nu}+g)\geq 0, \big(\tau_{\nu}(\boldsymbol{u},p)+h(u_{\nu})\big)(u_{\nu}+g)=0 \quad \text{on } \Gamma_{2},$$

we deduce

$$\int_{\partial\Omega} (\mathbb{S}\boldsymbol{\nu}) \cdot (\boldsymbol{\nu} - \boldsymbol{u}) d\Gamma$$

$$\geq \int_{\Gamma_1} \boldsymbol{\tau}_{\tau} (\boldsymbol{u}) \cdot (\boldsymbol{\nu}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma$$

$$+ \int_{\Gamma_2} p(\nu_{\nu} - u_{\nu}) - h(u_{\nu})(\nu_{\nu} - u_{\nu}) d\Gamma + \int_{\Gamma_3} (\boldsymbol{\tau}_{\nu} (\boldsymbol{u}, p) + p)(\nu_{\nu} - u_{\nu}) d\Gamma.$$
(18)

Therefore, combining (16), (17), and (18), we have

$$\int_{\Omega} \mathbb{S}(\mathbb{D}\boldsymbol{u}) : \mathbb{D}(\boldsymbol{v} - \boldsymbol{u}) \, dx + \int_{\Gamma_2} (v_v - u_v) p \, d\Gamma + \int_{\Gamma_3} (v_v - u_v) p \, d\Gamma$$

$$\geq \int_{\Omega} \boldsymbol{f} \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx + \int_{\Gamma_1} \boldsymbol{\tau}_{\tau} (\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) \, d\Gamma$$
$$+ \int_{\Gamma_2} p(v_{\nu} - u_{\nu}) - h(u_{\nu})(v_{\nu} - u_{\nu}) \, d\Gamma + \int_{\Gamma_3} (\boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p) + p)(v_{\nu} - u_{\nu}) \, d\Gamma,$$

and finally

$$\int_{\Omega} \mathbb{S}(\mathbb{D}\boldsymbol{u}) : \mathbb{D}(\boldsymbol{v} - \boldsymbol{u}) \, dx$$
  

$$\geq \int_{\Omega} \boldsymbol{f} \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx + \int_{\Gamma_1} \boldsymbol{\tau}_{\tau}(\boldsymbol{u}) \cdot (\boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) \, d\Gamma$$
  

$$- \int_{\Gamma_2} h(u_{\nu})(v_{\nu} - u_{\nu}) \, d\Gamma + \int_{\Gamma_3} \boldsymbol{\tau}_{\nu}(\boldsymbol{u}, p)(v_{\nu} - u_{\nu}) \, d\Gamma.$$

We utilize the subgradient boundary conditions (7) and (9) to obtain the following variational formulation of Problem 3.

**Problem 4** Find a velocity  $\mathbf{u} \in K$  such that

$$\begin{split} \int_{\Omega} \mu \mathbb{D} \boldsymbol{u} : \mathbb{D} (\boldsymbol{v} - \boldsymbol{u}) \, dx + \int_{\Gamma_2} h(u_{\nu})(v_{\nu} - u_{\nu}) \, d\Gamma + \int_{\Gamma_1} k \big( \|\boldsymbol{u}_{\tau}\| \big) j^0(\boldsymbol{u}_{\tau}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) \, d\Gamma \\ &+ \int_{\Gamma_3} l(u_{\nu}) \big( \phi(v_{\nu}) - \phi(u_{\nu}) \big) \, d\Gamma \\ &\geq \int_{\Omega} \boldsymbol{f} \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx \quad \text{for all } \boldsymbol{v} \in K. \end{split}$$

Problem 4 is called a variational–hemivariational inequality on a convex set. The main existence result on Problem 4 is stated below, and its proof will be given in Sect. 4.

**Theorem 5** Under hypotheses  $H(\mu)$ , H(h), H(j), H(k), H(l),  $H(\phi)$ , H(f), and  $(H_0)$ , Problem 4 has a solution.

#### 4 Proof of Theorem 5

The proof is based on an application of Theorem 2 and a fixed point argument. We divide the proof into four steps.

*Step 1.* Let  $Y_1 = H^{1/2}(\Gamma_1; \mathbb{R}^d)$  and  $Y_2 = H^{1/2}(\Gamma_3)$ . Let  $(\eta, \xi) \in Y_1 \times Y_2$  be fixed. Consider the auxiliary problem

$$P(\eta,\xi) \begin{cases} \text{find } \boldsymbol{u} \in K \text{ such that} \\ \int_{\Omega} \mu \mathbb{D} \boldsymbol{u} : \mathbb{D}(\boldsymbol{v} - \boldsymbol{u}) \, dx + \int_{\Gamma_2} h(u_{\nu})(v_{\nu} - u_{\nu}) \, d\Gamma \\ + \int_{\Gamma_1} k(\|\boldsymbol{\eta}\|_{\mathbb{R}^d}) j^0(\boldsymbol{u}_{\tau}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) \, d\Gamma + \int_{\Gamma_3} l(\xi)(\phi(v_{\nu}) - \phi(u_{\nu})) \, d\Gamma \\ \geq \int_{\Omega} \boldsymbol{f} \cdot (\boldsymbol{v} - \boldsymbol{u}) \, dx \quad \text{for all } \boldsymbol{v} \in K. \end{cases}$$

We shall prove that problem  $P(\eta, \xi)$  has a unique solution. With this problem, we associate the following inequality: find  $\mathbf{u} \in K$  such that

$$\langle A\boldsymbol{u} - \boldsymbol{f}, \boldsymbol{v} - \boldsymbol{u} \rangle + \Phi(\boldsymbol{v}) - \Phi(\boldsymbol{u}) + J^0(\boldsymbol{u}; \boldsymbol{v} - \boldsymbol{u}) \ge 0 \quad \text{for all } \boldsymbol{v} \in K,$$
(19)

where  $A: V \to V^*$  and  $\Phi, J: V \to \mathbb{R}$  are defined by

$$\langle A\boldsymbol{u},\boldsymbol{v}\rangle = \int_{\Omega} \mu \mathbb{D}\boldsymbol{u} : \mathbb{D}\boldsymbol{v} \, dx + \int_{\Gamma_2} h(u_v) v_v \, d\Gamma, \quad \boldsymbol{u}, \boldsymbol{v} \in V,$$
(20)

$$J(\mathbf{v}) = \int_{\Gamma_1} k \big( \|\boldsymbol{\eta}\|_{\mathbb{R}^d} \big) j(\mathbf{v}_{\tau}) \, d\Gamma, \quad \mathbf{v} \in V,$$
<sup>(21)</sup>

$$\Phi(\mathbf{v}) = \int_{\Gamma_3} l(\xi) \phi(\nu_{\nu}) \, d\Gamma, \quad \mathbf{v} \in V.$$
(22)

We shall verify conditions H(A), H(J),  $H(\Phi)$ , H(K), H(f), and  $(H_0)$  of Theorem 2. The operator  $A_1: V \to V^*$  defined by

$$\langle A_1 \boldsymbol{u}, \boldsymbol{v} \rangle = \int_{\Omega} \mu \mathbb{D} \boldsymbol{u} : \mathbb{D} \boldsymbol{v} \, dx, \quad \boldsymbol{u}, \boldsymbol{v} \in V,$$

is linear, bounded, strongly monotone with constant  $m_A = \mu_0 > 0$ , and so also coercive. In particular, the function  $V \ni \mathbf{v} \mapsto \langle A_1 \mathbf{v}, \mathbf{v} \rangle \in \mathbb{R}$  for  $\mathbf{v} \in V$  is strictly convex and lower semicontinuous, see, e.g., [28, Proposition 1.30]. Hence, it is weakly lower semicontinuous on V which implies

$$\limsup \langle A_1 \boldsymbol{u}^n, \boldsymbol{v} - \boldsymbol{u}^n \rangle \leq \langle A_1 \boldsymbol{u}, \boldsymbol{v} - \boldsymbol{u} \rangle \quad \text{for all } \boldsymbol{u}^n \rightharpoonup \boldsymbol{u} \text{ in } V, \boldsymbol{v} \in V.$$
(23)

Next, we note that the nonlinear operator

$$\langle A_2 \boldsymbol{u}, \boldsymbol{v} \rangle = \int_{\Gamma_2} h(u_v) v_v \, d\Gamma, \quad \boldsymbol{u}, \boldsymbol{v} \in V,$$

is bounded, continuous (we use the compactness of the normal trace operator  $V \mapsto L^2(\Gamma_2)$ ), and monotone (by hypothesis H(h)(ii)). Further, by H(h)(ii) and the Hölder inequality, we get

$$\lim \langle A_2 \boldsymbol{u}^n, \boldsymbol{v} - \boldsymbol{u}^n \rangle = \langle A_2 \boldsymbol{u}, \boldsymbol{v} - \boldsymbol{u} \rangle \quad \text{for all } \boldsymbol{u}^n \rightharpoonup \boldsymbol{u} \text{ in } V, \boldsymbol{v} \in V.$$
(24)

In conclusion, the operator *A* given by (20) is bounded, continuous, and strongly monotone with constant  $m_A = \mu_0$ , so by [20, Theorem 3.69], it is pseudomonotone. Condition H(A) is verified.

Hypothesis H(j)(i)-(ii) guarantees that the function *J*, given by (21), is well defined, and conditions H(J)(i) and (ii) follow from H(j)(i)-(iii), respectively, see [20, Theorem 3.47]. Condition H(J)(ii) is a consequence of the inequality

$$J^{0}(\boldsymbol{u};\boldsymbol{v}) \leq \int_{\Gamma_{1}} k(\|\boldsymbol{\eta}\|_{\mathbb{R}^{d}}) j^{0}(\boldsymbol{u}_{\tau};\boldsymbol{v}_{\tau}) d\Gamma \quad \text{for } \boldsymbol{u},\boldsymbol{v} \in V$$
(25)

combined with hypothesis H(j)(iv). We obtain that H(J)(iii) holds with  $\alpha_I = \alpha_i k_1 ||\gamma||^2$ .

Furthermore, taking into account hypothesis  $H(\phi)$ , we deduce  $H(\Phi)$ . The convexity of  $\Phi$  is obvious. Let  $\{\mathbf{v}^n\} \subset V$  be such that  $\mathbf{v}^n \to \mathbf{v}$  in V. We use the continuity of the normal trace operator  $V \ni \mathbf{v} \mapsto \nu_{\nu} \in L^2(\Gamma_3)$  to get  $\nu_{\nu}^n \to \nu_{\nu}$  in  $L^2(\Gamma_3)$ , and next, at least for a

subsequence, we have  $\nu_{\nu}^{n}(\mathbf{x}) \rightarrow \nu_{\nu}(\mathbf{x})$  for a.e.  $\mathbf{x} \in \Gamma_{3}$ . From H(l)(iii) and  $H(\phi)$ , we obtain

$$l(\boldsymbol{\xi}(\boldsymbol{x}))\boldsymbol{\phi}(\boldsymbol{v}_{\boldsymbol{\nu}}(\boldsymbol{x})) \leq \liminf l(\boldsymbol{\xi}(\boldsymbol{x}))\boldsymbol{\phi}(\boldsymbol{v}_{\boldsymbol{\nu}}^{n}(\boldsymbol{x})) \quad \text{a.e. } \boldsymbol{x} \in \Gamma_{3}.$$

$$(26)$$

It is clear from [3, Proposition 5.2.25] that there are  $a, b \in \mathbb{R}$  such that  $\phi(r) \ge ar + b$ for all  $r \in \mathbb{R}$ . Hence  $l(\xi(\mathbf{x}))\phi(\nu_{\nu}(\mathbf{x})) \ge \varrho_1(\mathbf{x})$  for a.e.  $\mathbf{x} \in \Gamma_3$  and all  $\mathbf{v} \in V$  with  $\varrho_1(\mathbf{x}) :=$  $-l_1(|a||\nu_{\nu}(\mathbf{x})| + |b|)$  a.e.  $\mathbf{x} \in \Gamma_3$ . By a direct calculation, we get  $\int_{\Gamma_3} \varrho_1 d\Gamma \ge -m_1 \|\mathbf{v}\|_V - m_2$ for  $\mathbf{v} \in V$  with  $m_1, m_2 \ge 0$ . Therefore  $\varrho_1 \in L^1(\Gamma_3)$ . Using the latter and integrating both sides of (26), by Fatou's lemma, it yields

$$\int_{\Gamma_3} l(\xi)\phi(v_{\nu})\,d\Gamma \leq \int_{\Gamma_3}\liminf l(\xi)\phi(v_{\nu}^n)\,d\Gamma \leq \liminf \int_{\Gamma_3} l(\xi)\phi(v_{\nu}^n)\,d\Gamma,$$

which shows that  $\Phi$  is a lower semicontinuous function on *V*. Thus,  $H(\Phi)$  is proved.

Subsequently, it is easy to see that the set of unilateral constraints (14) is closed and convex in *V* with  $\mathbf{0} \in K$  and, thus, condition H(K) is satisfied. Finally, the hypothesis  $(H_0)$  assures that the smallness condition (1) holds. Having verified the hypotheses of Theorem 2, we deduce that there exists a unique solution  $\mathbf{u} = \mathbf{u}_{\xi\eta} \in K$  to problem (19). Using inequality (25), we know that  $\mathbf{u} = \mathbf{u}_{\xi\eta} \in K$  is also the solution to problem  $P(\eta, \xi)$ .

By the direct computation, we can show that the solution to problem  $P(\eta, \xi)$  is unique. Indeed, let  $\boldsymbol{u}_1, \boldsymbol{u}_2 \in V$  be solutions to problem  $P(\eta, \xi)$ . Choosing the test function  $\boldsymbol{v} = \boldsymbol{u}_2 \in K$  in the inequality satisfied by  $\boldsymbol{u}_1$ , and  $\boldsymbol{v} = \boldsymbol{u}_1 \in K$  in the inequality satisfied by  $\boldsymbol{u}_2$ , and then adding the two resulting inequalities, we obtain

$$\langle A\boldsymbol{u}_1 - A\boldsymbol{u}_2, \boldsymbol{u}_1 - \boldsymbol{u}_2 \rangle \leq \int_{\Gamma_1} k \big( \|\boldsymbol{\eta}\|_{\mathbb{R}^d} \big) \big( j^0(\boldsymbol{u}_{1\tau}; \boldsymbol{u}_{2\tau} - \boldsymbol{u}_{1\tau}) + j^0(\boldsymbol{u}_{2\tau}; \boldsymbol{u}_{1\tau} - \boldsymbol{u}_{2\tau}) \big) d\Gamma$$

Taking into account the strong monotonicity of the operator *A* and hypotheses H(j)(iv), H(k)(iii) as well as the continuity of the trace operator  $\gamma$ , we have

$$\mu_0 \| \boldsymbol{u}_1 - \boldsymbol{u}_2 \|_V^2 \le k_1 \alpha_j \| \boldsymbol{u}_{1\tau} - \boldsymbol{u}_{2\tau} \|_{L^2(\Gamma_1; \mathbb{R}^d)}^2 \le \alpha_j k_1 \| \boldsymbol{\gamma} \|^2 \| \boldsymbol{u}_1 - \boldsymbol{u}_2 \|_V^2.$$

Finally, due to hypothesis ( $H_0$ ), we deduce  $\boldsymbol{u}_1 = \boldsymbol{u}_2$ . In conclusion, problem  $P(\eta, \xi)$  has a unique solution.

*Step 2.* We shall show the estimate on the unique solution  $\mathbf{u} = \mathbf{u}_{\xi\eta} \in K$  to  $P(\eta, \xi)$ . We choose  $\mathbf{v} = \mathbf{0} \in K$  in the inequality  $P(\eta, \xi)$  to get

$$\langle A\boldsymbol{u},\boldsymbol{u}\rangle \leq \Phi(\boldsymbol{0}) - \Phi(\boldsymbol{u}) + \left| \int_{\Gamma_1} k \big( \|\boldsymbol{\eta}\|_{\mathbb{R}^d} \big) j^0(\boldsymbol{u}_{\tau};-\boldsymbol{u}_{\tau}) \, d\Gamma \right| + \left| \langle \boldsymbol{f},\boldsymbol{u} \rangle \right|.$$
(27)

From (2), H(l), and  $H(\phi)$ , it is clear that

$$\langle A\boldsymbol{u}, \boldsymbol{u} \rangle \ge \mu_0 \|\boldsymbol{u}\|_V^2 + \|A\boldsymbol{0}\|_{V^*} \|\boldsymbol{u}\|_V \ge \mu_0 \|\boldsymbol{u}\|_V^2 \quad \text{and} \quad \Phi(\boldsymbol{0}) \le l_1 |\phi(0)|.$$
 (28)

By [3, Proposition 5.2.25], it follows that  $\Phi$  is bounded from below by an affine function, i.e., there are  $q \in V^*$  and  $\beta \in \mathbb{R}$  such that  $\Phi(\mathbf{v}) \ge \langle q, \mathbf{v} \rangle + \beta$  for all  $\mathbf{v} \in V$ . Hence

$$-\Phi(\mathbf{u}) \le \|q\|_{V^*} \|\mathbf{u}\|_{V} + |\beta|.$$
<sup>(29)</sup>

Subsequently, from H(j)(iii), (iv) and the inequality  $\|\boldsymbol{u}_{\tau}\|_{\mathbb{R}^d} \leq \|\boldsymbol{u}\|_{\mathbb{R}^d}$ , we deduce

$$\left|j^{0}(\boldsymbol{u}_{\tau};-\boldsymbol{u}_{\tau})\right| \leq \alpha_{j} \|\boldsymbol{u}_{\tau}\|_{\mathbb{R}^{d}}^{2} + \left|j^{0}(0;\boldsymbol{u}_{\tau})\right| \leq \alpha_{j} \|\boldsymbol{u}\|_{\mathbb{R}^{d}}^{2} + c_{0}(\boldsymbol{x}) \|\boldsymbol{u}\|_{\mathbb{R}^{d}}.$$

Due to H(k), the continuity of the trace operator  $V \mapsto L^2(\Gamma_1; \mathbb{R}^d)$ , and the Hölder inequality, we have

$$\begin{aligned} \left| \int_{\Gamma_{1}} k(\|\boldsymbol{\eta}\|_{\mathbb{R}^{d}}) j^{0}(\boldsymbol{u}_{\tau}; -\boldsymbol{u}_{\eta}) d\Gamma \right| \qquad (30) \\ &\leq \int_{\Gamma_{1}} \left| k(\|\boldsymbol{\eta}\|_{\mathbb{R}^{d}}) \right| \left| j^{0}(\boldsymbol{u}_{\tau}; -\boldsymbol{u}_{\tau}) \right| d\Gamma \\ &\leq \int_{\Gamma_{1}} k_{1}(\alpha_{j} \|\boldsymbol{u}\|^{2} + c_{0}(\boldsymbol{x}) \|\boldsymbol{u}\|) d\Gamma \leq \alpha_{j} k_{1} \|\boldsymbol{u}\|_{L^{2}(\Gamma_{1};\mathbb{R}^{d})}^{2} + k_{1} \|c_{0}\|_{L^{2}(\Gamma_{1})} |\Gamma_{1}|^{1/2} \|\boldsymbol{u}\|_{L^{2}(\Gamma_{1};\mathbb{R}^{d})} \\ &\leq \alpha_{j} k_{1} \|\boldsymbol{\gamma}\|^{2} \|\boldsymbol{u}\|_{V}^{2} + k_{1} \|c_{0}\|_{L^{2}(\Gamma_{1})} |\Gamma_{1}|^{1/2} \|\boldsymbol{\gamma}\| \|\boldsymbol{u}\|_{V}. \end{aligned}$$

Combining (28)–(30) with (27), we obtain

$$(\mu_{0} - \alpha_{j}k_{1} \|\gamma\|^{2}) \|\boldsymbol{u}\|_{V}^{2} \leq (\|q\|_{V^{*}} + \|\boldsymbol{f}\|_{V^{*}} + k_{1}\|c_{0}\|_{L^{2}(\Gamma_{1})}|\Gamma_{1}|^{1/2}\|\gamma\|) \|\boldsymbol{u}\|_{V}$$

$$+ l_{1} |\phi(0)| + |\beta|.$$

$$(31)$$

We use the relation

$$x^2 \le ax + b \implies x^2 \le a^2 + 2b$$
, for all  $a, b, x \ge 0$ 

to conclude that

$$\|\boldsymbol{u}\|_{V} \le M,\tag{32}$$

where the constant M > 0 is such that

$$M^{2} = \left(\frac{\|q\|_{V^{*}} + \|\mathbf{f}\|_{V^{*}} + k_{1}\|c_{0}\|_{L^{2}(\Gamma_{1})}|\Gamma_{1}|^{1/2}\|\gamma\|}{\mu_{0} - \alpha_{j}k_{1}\|\gamma\|^{2}}\right)^{2} + 2\frac{l_{1}|\phi(0)| + |\beta|}{\mu_{0} - \alpha_{j}k_{1}\|\gamma\|^{2}}$$

*Step 3.* Let  $\Lambda : Y_1 \times Y_2 \to Y_1 \times Y_2$  be defined by

$$\Lambda(\boldsymbol{\eta}, \boldsymbol{\xi}) = (\boldsymbol{u}_{\eta \boldsymbol{\xi} \tau}, \boldsymbol{u}_{\eta \boldsymbol{\xi} \nu}) \quad \text{for } (\boldsymbol{\eta}, \boldsymbol{\xi}) \in Y_1 \times Y_2,$$

where  $\boldsymbol{u} = \boldsymbol{u}_{\eta\xi} \in K$  is the unique solution to problem  $P(\eta, \xi)$ . Note that  $\Lambda$  is well defined, and by (32), the continuity of the trace operator (15), and the elementary inequalities  $\|\boldsymbol{u}_{\tau}\|_{\mathbb{R}^d} \leq \|\boldsymbol{u}\|_{\mathbb{R}^d}$  and  $u_{\nu} \leq |u_{\nu}| \leq \|\boldsymbol{u}\|_{\mathbb{R}^d}$ , we have

$$\|\Lambda(\eta,\xi)\|_{Y_1\times Y_2} = \|\boldsymbol{u}_{\eta\xi\tau}\|_{Y_1} + \|u_{\eta\xi\nu}\|_{Y_2} \le 2\|\boldsymbol{\gamma}\|\|\boldsymbol{u}_{\eta\xi}\|_V \le 2M\|\boldsymbol{\gamma}\|.$$
(33)

In other words, due to the uniform bounds for k and l, the image of the whole space  $Y_1 \times Y_2$  remains in a closed ball in  $Y_1 \times Y_2$ , that is,

$$\Lambda(Y_1 \times Y_2) \subset \overline{B}(0, 2M \|\gamma\|). \tag{34}$$

Next, we shall show that  $\Lambda$  is a weakly-weakly continuous map. Let  $(\eta_n, \xi_n), (\eta, \xi) \in Y_1 \times Y_2$  such that

$$\eta_n \rightharpoonup \eta \quad \text{in } Y_1 \quad \text{and} \quad \xi_n \rightharpoonup \xi \quad \text{in } Y_2.$$
 (35)

Using the compactness of the embedding  $H^{1/2}(\Gamma; \mathbb{R}^d)$  into  $L^2(\Gamma; \mathbb{R}^d)$ , (35) implies

$$\eta_n \to \eta \quad \text{in } L^2(\Gamma_1; \mathbb{R}^d) \quad \text{and} \quad \xi_n \to \xi \quad \text{in } L^2(\Gamma_3),$$
(36)

and therefore we can assume that

$$\eta_n \to \eta$$
 a.e. on  $\Gamma_1$  and  $\xi_n \to \xi$  a.e. on  $\Gamma_3$ . (37)

Let  $\boldsymbol{u}_{\eta_n\xi_n}^n \in K$ ,  $n \in \mathbb{N}$ , be the unique solution to problem  $P(\eta_n, \xi_n)$  corresponding to  $(\boldsymbol{\eta}_n, \boldsymbol{\xi}_n)$ . For simplicity, we write  $\boldsymbol{u}^n = \boldsymbol{u}_{\eta_n\xi_n}^n$  and we have

$$\langle A\boldsymbol{u}^{n} - \boldsymbol{f}, \boldsymbol{v} - \boldsymbol{u}^{n} \rangle + \int_{\Gamma_{1}} k \left( \|\boldsymbol{\eta}_{n}\|_{\mathbb{R}^{d}} \right) j^{0} \left( \boldsymbol{u}_{\tau}^{n}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^{n} \right) d\Gamma$$
  
+ 
$$\int_{\Gamma_{3}} l(\xi_{n}) \left( \phi(v_{\nu}) - \phi\left(\boldsymbol{u}_{\nu}^{n}\right) \right) d\Gamma \geq 0 \quad \text{for all } \boldsymbol{v} \in K,$$
 (38)

where *A* is the operator defined by (20). By estimate (32) in Step 2, we know that  $\{\boldsymbol{u}^n\}$  remains uniformly bounded in *V*. Thus, by the reflexivity of *V*, there is a subsequence of  $\{\boldsymbol{u}^n\}$ , denoted in the same way, and  $\boldsymbol{u} \in V$  such that

$$\boldsymbol{u}^n \rightharpoonup \boldsymbol{u} \quad \text{in } V. \tag{39}$$

Since  $\{\boldsymbol{u}^n\} \subset K$  and K is weakly closed in V (being closed and convex), we have  $\boldsymbol{u} \in K$ . From the compactness of the trace operator  $\gamma$  and the converse Lebesgue dominated convergence theorem [20, Theorem 2.39], we may suppose that  $\boldsymbol{u}^n \to \boldsymbol{u}$  in  $L^2(\Gamma; \mathbb{R}^d)$  and

$$\boldsymbol{u}^{n}(\boldsymbol{x}) \rightarrow \boldsymbol{u}(\boldsymbol{x}) \quad \text{a.e. } \boldsymbol{x} \in \Gamma, \qquad \left\| \boldsymbol{u}^{n}(\boldsymbol{x}) \right\|_{\mathbb{R}^{d}} \le \rho(\boldsymbol{x}) \quad \text{a.e. } \boldsymbol{x} \in \Gamma$$
 (40)

with  $\rho \in L^2(\Gamma)$ . Next, let  $\mathbf{v} \in K$ . We use H(j)(iii), H(k)(iii), and [20, Proposition 3.23 (iii)] to obtain

$$\begin{aligned} &k \big( \|\boldsymbol{\eta}_n\|_{\mathbb{R}^d} \big) j^0 \big( \boldsymbol{u}_{\tau}^n; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^n \big) \\ &\leq k_1 \|\partial j \big( \boldsymbol{u}_{\tau}^n \big) \|_{\mathbb{R}^d} \| \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^n \|_{\mathbb{R}^d} \\ &\leq k_1 \big( c_0(\boldsymbol{x}) + c_1 \| \boldsymbol{u}_{\tau}^n \|_{\mathbb{R}^d} \big) \big( \| \boldsymbol{v}(x) \|_{\mathbb{R}^d} + \| \boldsymbol{u}_{\tau}^n \|_{\mathbb{R}^d} \big) \\ &\leq k_1 \big( c_0(\boldsymbol{x}) + c_1 \rho(\boldsymbol{x}) \big) \big( \| \boldsymbol{v}(x) \|_{\mathbb{R}^d} + \rho(\boldsymbol{x}) \big) =: \rho(\boldsymbol{x}), \end{aligned}$$

where  $\rho \in L^1(\Gamma_1)$ . In view of the latter, by Fatou's lemma, we have

$$\limsup \int_{\Gamma_1} k \big( \|\boldsymbol{\eta}_n\|_{\mathbb{R}^d} \big) j^0 \big( \boldsymbol{u}_{\tau}^n; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^n \big) \, d\Gamma$$

$$\leq \int_{\Gamma_{1}} \limsup k(\|\boldsymbol{\eta}_{n}\|_{\mathbb{R}^{d}}) j^{0}(\boldsymbol{u}_{\tau}^{n}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^{n}) d\Gamma$$
  
$$\leq \int_{\Gamma_{1}} \limsup (k(\|\boldsymbol{\eta}_{n}\|_{\mathbb{R}^{d}}) - k(\|\boldsymbol{\eta}\|_{\mathbb{R}^{d}})) \varrho d\Gamma$$
  
$$+ \int_{\Gamma_{1}} k(\|\boldsymbol{\eta}\|_{\mathbb{R}^{d}}) \limsup j^{0}(\boldsymbol{u}_{\tau}^{n}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^{n}) d\Gamma.$$
(41)

On the other hand, we employ (40), the upper semicontinuity of the function  $j^0(\cdot; \cdot)$ , see [20, Proposition 3.23(ii)], and get

$$\int_{\Gamma_1} k\big(\|\boldsymbol{\eta}\|_{\mathbb{R}^d}\big) \limsup j^0\big(\boldsymbol{u}_{\tau}^n; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^n\big) d\Gamma \leq \int_{\Gamma_1} k\big(\|\boldsymbol{\eta}\|_{\mathbb{R}^d}\big) j^0(\boldsymbol{u}_{\tau}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) d\Gamma.$$
(42)

Combining (35), (41), and (42), it follows

$$\limsup \int_{\Gamma_1} k \big( \|\boldsymbol{\eta}_n\|_{\mathbb{R}^d} \big) j^0 \big( \boldsymbol{u}_{\tau}^n; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^n \big) \, d\Gamma \leq \int_{\Gamma_1} k \big( \|\boldsymbol{\eta}\|_{\mathbb{R}^d} \big) j^0 (\boldsymbol{u}_{\tau}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}) \, d\Gamma \tag{43}$$

for all  $\mathbf{v} \in V$ . In a similar way to (43), we can use (35), H(l),  $H(\phi)$ , and Fatou's lemma to obtain

$$\limsup \int_{\Gamma_3} l(\xi_n) \big( \phi(\nu_\nu) - \phi(u_\nu^n) \big) \, d\Gamma \le \int_{\Gamma_3} l(\xi) \big( \phi(\nu_\nu) - \phi(u_\nu) \big) \, d\Gamma \tag{44}$$

for all  $\mathbf{v} \in V$ .

We are now in a position to pass to the limit in (38), as  $n \to \infty$ . To this end, we use convergences (23), (24), (39), (43), and (44) to deduce

$$\langle A\boldsymbol{u}-\boldsymbol{f},\boldsymbol{v}-\boldsymbol{u}\rangle+\int_{\Gamma_1}k\big(\|\boldsymbol{\eta}\|_{\mathbb{R}^d}\big)j^0(\boldsymbol{u}_{\tau};\boldsymbol{v}_{\tau}-\boldsymbol{u}_{\tau})\,d\Gamma+\int_{\Gamma_3}l(\xi)\big(\phi(\boldsymbol{v}_{\nu})-\phi(\boldsymbol{u}_{\nu})\big)\,d\Gamma\geq 0$$

for all  $\mathbf{v} \in K$ . This means that  $\mathbf{u} \in K$  is the solution to problem  $P(\eta, \xi)$  corresponding to  $(\eta, \xi)$ . Since the problem admits a unique solution, the whole sequence  $\{\mathbf{u}^n\}$  converges weakly in *V* to  $\mathbf{u}$ . Since the (tangential and normal) trace operator (15) is linear and continuous, we have

$$\boldsymbol{u}_{\tau}^{n} \rightarrow \boldsymbol{u}_{\tau}$$
 in  $Y_{1}$  and  $\boldsymbol{u}_{\nu}^{n} \rightarrow \boldsymbol{u}_{\nu}$  in  $Y_{2}$ 

and  $\Lambda(\eta_n, \xi_n) \rightarrow \Lambda(\eta, \xi)$  in  $Y_1 \times Y_2$ . This proves the weak-to-weak continuity of the map  $\Lambda$ .

*Step 4*. We apply the Schauder fixed point theorem in its weak variant, see, e.g., [12, p. 185].

**Theorem 6** Let X be a separable reflexive Banach space and  $C \subset X$  be nonempty, convex, and weakly compact in X. Let S:  $C \rightarrow C$  be weak-to-weak continuous. Then S has a fixed point.

From Step 2, we know that the map  $\Lambda$  is continuous in the weak topology and it satisfies (34). By the reflexivity of  $Y_1 \times Y_2$ , it is clear that the closed ball is weakly compact and obviously convex. Thus  $\Lambda$  has a fixed point  $(\eta^*, \xi^*) \in Y_1 \times Y_2$ , that is,

$$\boldsymbol{\eta}^* = \boldsymbol{u}_{\eta^* \xi^* \tau}, \quad \xi^* = u_{\eta^* \xi^* \nu},$$

where  $\boldsymbol{u}_{\eta^*\xi^*} \in K$  solves problem  $P(\eta^*, \xi^*)$ . We conclude that  $\boldsymbol{u}_{\eta^*\xi^*} \in K$  is a solution to Problem 4. The proof is complete.

#### 5 Compactness and uniqueness results

We conclude the paper with the following compactness and uniqueness results for Problem 4.

**Theorem 7** Under the hypotheses of Theorem 5, the set of solutions to Problem 4 is weakly compact in V.

*Proof* Let *S* be the solution set to Problem 4. It is a nonempty subset of *V* by Theorem 5. Let  $\{u^n\} \subset S$ , so  $u^n \in K$  and

$$\int_{\Omega} \mu \mathbb{D} \boldsymbol{u}^{n} : \mathbb{D} (\boldsymbol{v} - \boldsymbol{u}^{n}) \, dx + \int_{\Gamma_{2}} h(\boldsymbol{u}_{\nu}^{n}) (\boldsymbol{v}_{\nu} - \boldsymbol{u}_{\nu}^{n}) \, d\Gamma + \int_{\Gamma_{1}} k(\|\boldsymbol{u}_{\tau}^{n}\|) j^{0}(\boldsymbol{u}_{\tau}^{n}; \boldsymbol{v}_{\tau} - \boldsymbol{u}_{\tau}^{n}) \, d\Gamma \\ + \int_{\Gamma_{3}} l(\boldsymbol{u}_{\nu}^{n}) (\phi(\boldsymbol{v}_{\nu}) - \phi(\boldsymbol{u}_{\nu}^{n})) \, d\Gamma \ge \int_{\Omega} \boldsymbol{f} \cdot (\boldsymbol{v} - \boldsymbol{u}^{n}) \, dx \quad \text{for all } \boldsymbol{v} \in K.$$

Next, we follow the arguments used in Steps 2 and 3 of Theorem 5 to demonstrate that  $\{u^n\}$  is uniformly bounded in *V* and, at least for a subsequence,  $u^n \rightarrow u^0$  weakly in *V*, where  $u^0 \in K$  solves Problem 4. Hence  $u^0 \in S$ . This completes the proof.

The uniqueness of a solution to Problem 4 under the general hypotheses of Theorem 5 is an open problem. However, it can be obtained, for instance, under the additional assumptions stated below.

**Theorem 8** Assume the hypotheses of Theorem 5 and the following conditions:

$$\begin{cases} the functions j(\mathbf{x}, \cdot), k(\mathbf{x}, \cdot), l(\mathbf{x}, \cdot) \text{ and } \phi \text{ are Lipschitz continuous} \\ for a.e. \mathbf{x} \text{ with constants } L_j, L_k, L_l, L_{\phi} > 0, respectively, and \\ (\alpha_j k_1 + L_j L_k + L_l L_{\phi}) \|\gamma\|^2 < \mu_0. \end{cases}$$
(45)

Then Problem 4 has a unique solution.

*Proof* Let  $u_1, u_2 \in K$  be solutions to Problem 4. We take the test function  $v = u_2$  in the inequality satisfied by  $u_1$ , and  $v = u_1$  in the inequality satisfied by  $u_2$ . Next, we add the resulting inequalities to get

$$\langle A \boldsymbol{u}_{1} - A \boldsymbol{u}_{2}, \boldsymbol{u}_{2} - \boldsymbol{u}_{1} \rangle$$

$$+ \int_{\Gamma_{1}} \left( k \big( \| \boldsymbol{u}_{1\tau} \|_{\mathbb{R}^{d}} \big) j^{0} (\boldsymbol{u}_{1\tau}; \boldsymbol{u}_{2\tau} - \boldsymbol{u}_{1\tau}) + k \big( \| \boldsymbol{u}_{2\tau} \|_{\mathbb{R}^{d}} \big) j^{0} (\boldsymbol{u}_{2\tau}; \boldsymbol{u}_{1\tau} - \boldsymbol{u}_{2\tau}) \big) d\Gamma$$

$$+ \int_{\Gamma_{3}} \big( l(u_{1\nu}) \big( \phi(u_{2\nu}) - \phi(u_{1\nu}) \big) + l(u_{2\nu}) \big( \phi(u_{1\nu}) - \phi(u_{2\nu}) \big) \big) d\Gamma \ge 0,$$

where the operator A is defined by (20). Hence, we have

$$\langle A\boldsymbol{u}_1 - A\boldsymbol{u}_2, \boldsymbol{u}_1 - \boldsymbol{u}_2 \rangle \le I_1 + I_2 \tag{46}$$

with

$$I_{1} := \int_{\Gamma_{1}} \left( k \left( \| \boldsymbol{u}_{1\tau} \|_{\mathbb{R}^{d}} \right) - k \left( \| \boldsymbol{u}_{2\tau} \|_{\mathbb{R}^{d}} \right) \right) j^{0}(\boldsymbol{u}_{1\tau}; \boldsymbol{u}_{2\tau} - \boldsymbol{u}_{1\tau}) \\ + k \left( \| \boldsymbol{u}_{2\tau} \|_{\mathbb{R}^{d}} \right) \left( j^{0}(\boldsymbol{u}_{1\tau}; \boldsymbol{u}_{2\tau} - \boldsymbol{u}_{1\tau}) + j^{0}(\boldsymbol{u}_{2\tau}; \boldsymbol{u}_{1\tau} - \boldsymbol{u}_{2\tau}) \right) d\Gamma, \\ I_{2} := \int_{\Gamma_{3}} \left( l(\boldsymbol{u}_{1\nu}) - l(\boldsymbol{u}_{2\nu}) \right) \left( \phi(\boldsymbol{u}_{2\nu}) - \phi(\boldsymbol{u}_{1\nu}) \right) d\Gamma.$$

Note that  $j(\mathbf{x}, \cdot)$  is Lipschitz with constant  $L_j > 0$  for a.e.  $\mathbf{x} \in \Gamma_1$  if and only if  $\|\partial j(\boldsymbol{\xi})\|_{\mathbb{R}^d} \le L_j$  for all  $\boldsymbol{\xi} \in \mathbb{R}^d$ . Under the hypotheses, we use the continuity of the trace operator  $\gamma$ , and by a direct calculation, we deduce

$$\begin{aligned} |I_1| &\leq \alpha_j k_1 \|\gamma\|^2 \|\boldsymbol{u}_1 - \boldsymbol{u}_2\|_V^2 + L_j L_k \|\gamma\|^2 \|\boldsymbol{u}_1 - \boldsymbol{u}_2\|_V^2, \\ |I_2| &\leq L_l L_\phi \|\gamma\|^2 \|\boldsymbol{u}_1 - \boldsymbol{u}_2\|_V^2. \end{aligned}$$

Combining these inequalities with (46) and the strong monotonicity of the operator A imply

$$(\mu_0 - (\alpha_j k_1 + L_j L_k + L_l L_{\phi}) \|\gamma\|^2) \|\boldsymbol{u}_1 - \boldsymbol{u}_2\|_V^2 \leq 0.$$

Finally, from (45), we deduce  $\boldsymbol{u}_1 = \boldsymbol{u}_2$ . The proof is complete.

*Example* 9 We provide examples of functions *j* and  $\phi$ . Let *j*:  $\mathbb{R}^d \to \mathbb{R}$  be defined by

$$j(\boldsymbol{\xi}) = (a-1)e^{-\|\boldsymbol{\xi}\|_{\mathbb{R}^d}} + a\|\boldsymbol{\xi}\|_{\mathbb{R}^d} \quad \text{for } \boldsymbol{\xi} \in \mathbb{R}^d$$

with a constant  $a \in [0, 1)$ . Then the generalized gradient has the form

$$\partial j(\boldsymbol{\xi}) = \begin{cases} \overline{B}(\mathbf{0},1) & \text{if } \boldsymbol{\xi} = \mathbf{0}, \\ ((1-a)e^{-\|\boldsymbol{\xi}\|_{\mathbb{R}^d}} + a)\frac{\boldsymbol{\xi}}{\|\boldsymbol{\xi}\|_{\mathbb{R}^d}} & \text{if } \boldsymbol{\xi} \neq \mathbf{0} \end{cases} \quad \text{for } \boldsymbol{\xi} \in \mathbb{R}^d.$$

The function *j* is nonconvex and  $\|\partial j(\boldsymbol{\xi})\|_{\mathbb{R}^d} \leq 1$  for all  $\boldsymbol{\xi} \in \mathbb{R}^d$ . Further, it satisfies H(j) with  $c_0 = 1, c_1 = 0$  and  $\alpha_j = 1$ , and *j* is Lipschitz continuous with  $L_j = 1$ . Let  $\phi \colon \mathbb{R} \to \mathbb{R}$  be defined by

$$\phi(r) = \begin{cases} |r| & \text{if } |r| \le 2, \\ 2|r|-2 & \text{if } |r| > 2 \end{cases} \quad \text{for } r \in \mathbb{R}.$$

Then the function  $\phi$  is convex, continuous, and

$$\partial \phi(r) = \begin{cases} -2 & \text{if } r \le -2, \\ [-2, -1] & \text{if } r = -2, \\ -1 & \text{if } r \in (-2, 0), \\ [-1, 1] & \text{if } r = 0, \\ 1 & \text{if } r \in (0, 2), \\ [1, 2] & \text{if } r = 2, \\ 2 & \text{if } r \ge 2 \end{cases}$$

So it satisfies  $H(\phi)$ ,  $|\partial \phi(r)| \le 2$  for all  $r \in \mathbb{R}$ , and it is Lipschitz with  $L_{\phi} = 2$ .

Finally, we observe that any smooth solution to Problem 4 is a solution to Problem 3. In fact, we are able to recover the pressure from the weak formulation. For the proof, see [17, Proposition 17].

**Proposition 10** If  $\mathbf{u} \in K$  is a solution to Problem 4 given by Theorem 5 which is smooth  $\mathbf{u} \in H^2(\Omega; \mathbb{R}^d)$ , then  $\mathbf{u}$  satisfies the conditions of Problem 3. In particular, there exists a unique  $p \in L^2(\Omega)$  such that  $\int_{\Omega} p \, dx = 0$  which satisfies (4)–(9).

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#### Availability of data and materials

Not applicable.

#### Declarations

#### **Competing interests**

The authors declare that they have no competing interests.

#### Authors' contributions

SM developed the mathematical methods. SD discussed the results and contributed to the final manuscript. JZ contributed to typing and corrections. All authors read and approved the final manuscript.

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