# STEADY-STATE BUOYANCY-DRIVEN VISCOUS FLOW WITH MEASURE DATA 

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Dedicated to Prof. J. Nečas on the occasion of his 70th birthday
Abstract. Steady-state system of equations for incompressible, possibly non-Newtonean of the $p$-power type, viscous flow coupled with the heat equation is considered in a smooth bounded domain $\Omega \subset \mathbb{R}^{n}, n=2$ or 3 , with heat sources allowed to have a natural $L^{1}$ structure and even to be measures. The existence of a distributional solution is shown by a fixed-point technique for sufficiently small data if $p>3 / 2$ (for $n=2$ ) or if $p>9 / 5$ (for $n=3$ ).

Keywords: non-Newtonean fluids, heat equation, dissipative heat, adiabatic heat
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## 1. Introduction, Problem formulation

This paper deals with the steady-state buoyancy-driven flow of heat-conductive, possibly non-Newtonean, incompressible fluids. There are various models appearing in literature, cf.e.g. [3], [7], [14], [18] for a genesis of various possibilities. The starting point is always the complete evolutionary compressible fluid system of $n+2$ conservation laws for mass, impulse, and energy; $n$ denotes the spatial dimension. Then, the so-called incompressible limit represents a small perturbation around a stationary homogeneous state, i.e. around constant mass density, constant temperature, and zero velocity; note that small perturbations of velocity $u$ do not necessarily mean small $\nabla u$, which makes it sensible to consider nonlinearity in stress $\tau$ below. This incompressible limit system of $n+1$ equations need not be thermodynamically consistent, however.

We consider $\Omega$ a bounded smooth (namely $C^{3,1}{ }_{-}$) domain in $\mathbb{R}^{n}, n=2$ or 3 ; for $\Omega$ a $C^{0,1}$-domain see Remark 2 below. To cover various possibilities, we consider the
following fairly general system of equations:
(1.1a) $(u \cdot \nabla) u-\operatorname{div} \tau(e(\nabla u))+\nabla \pi=g\left(1-\alpha_{0} \theta\right), \quad e(\nabla u)=\frac{1}{2} \nabla u+\frac{1}{2}(\nabla u)^{\mathrm{T}}$

$$
\begin{align*}
\operatorname{div} u & =0  \tag{1.1b}\\
u \cdot \nabla \theta-\kappa \Delta \theta & =\alpha_{1} \tau(e(\nabla u)): e(\nabla u)+\alpha_{2} \theta g \cdot u+h, \tag{1.1c}
\end{align*}
$$

where $\left[\tau_{i j}\right]:\left[e_{i j}\right]=\sum_{i=1}^{n} \sum_{j=1}^{n} \tau_{i j} e_{i j}, \kappa$ is the heat conductivity, $\alpha_{0}$ is the linearized relative mass density variation with respect to temperature, $\alpha_{1}$ reflects the dissipation effects, $\alpha_{2}$ expresses the adiabatic heat effects, $\tau(e)$ is the viscous stress, $g$ an external (e.g. gravitational or centrifugal) force, and $h=h(x)$ is the external heat source. For simplicity of notation, we normalize the mass density and the heat capacity to 1.

For a rigorous derivation of a system like (1.1) in the evolution case, we refer to Kagei, Růžička and Thäter [7, System (16)] who showed how the coefficient $\alpha_{1}$ depends on Ostrach's dissipation number, while the coefficient $\alpha_{2}$ depends also on the Reynolds and the Prandtl numbers.

The system should be completed by boundary conditions. For simplicity, we will consider a no-slip boundary condition for velocity and the Newton condition with prescribed heat flux $f$ for temperature, i.e.

$$
\begin{equation*}
u=0, \quad \kappa \frac{\partial \theta}{\partial \nu}+b \theta=f \quad \text { on } \Gamma, \tag{1.2}
\end{equation*}
$$

with $\nu$ denoting the unit normal to the boundary $\partial \Omega=: \Gamma$ of $\Omega$ and $b$ denoting the coefficient of the heat transfer through $\Gamma$.

Often a simpler, so-called Oberbeck-Boussinesq model is used for the buoyancydriven flow of heat-conductive incompressible fluids. This model neglects both the dissipative and the adiabatic heat sources, i.e. $\alpha_{1}=\alpha_{2}=0$, and usually considers $\tau(e)=e$ which turns (1.1a,b) into the Navier-Stokes system, cf. e.g. Gebhart et al. [5] or Rajagopal et al. [18], and sometimes it is combined with other phenomena as solidification, see Rodriguez [19]. For a non-Newtonean model coupled with the heat equation we refer to Málek at al. [13] and to Rodriguez and Urbano [20] who allowed the viscosity to depend also on temperature. Temperature dependence of the viscosity tensor $\tau$ was investigated also by Baranger and Mikelić [2] for the special case $\alpha_{1}=1, \alpha_{0}=0$ (i.e. no buoyancy) and $\alpha_{2}=0$, which makes the situation quite different from the buoyancy driven flow. Besides, some buoyancy-driven models include the dissipative heat but not the adiabatic heat sources (i.e. our model (1.1) with $\alpha_{1}>0$ but $\alpha_{2}=0$ ), cf. Landau and Lifshitz [9, Sect. 50] or also, e.g., Kagei [6] or Moseenkov [14].

The measures as heat sources for the buoyancy-driven flow have been investigated for $b=0$ and $f=0$ in [16] in the evolutionary case, which differs from the steadystate case both factually (existence of a non-negative solution holds for arbitrarily large data) and technically ( $L^{1}$-accretivity for the heat equation can be used instead of mere $W^{2,2}$-regularity and interpolation with transposition).

## 2. Distributional solution to (1.1)-(1.2)

We want to treat the system (1.1) in as much general as possible (but still physical) situations. The heat transfer (1.1c) has a natural $L^{1}$-structure, which encourages us to consider the heat sources $h \in L^{1}(\Omega)$ and $f \in L^{1}(\Gamma)$, or even as measures. Then the concept of a weak solution is no longer relevant, and one must speak in terms of distributional solutions, using transposition and $W^{2,2}$-regularity with Hilbertianspace interpolation of the adjoint to the left-hand-side linear operator in (1.1c).

We use the following standard notation for functions spaces: $L^{p}\left(\Omega ; \mathbb{R}^{n}\right)$ denotes the Lebesgue space of measurable functions $\Omega \rightarrow \mathbb{R}^{n}$ whose $p$-power is integrable, $W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$ is the Sobolev space of functions whose gradient is in $L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right)$ and whose trace on $\Gamma$ vanishes, $W_{0, \text { DIV }}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)=\left\{v \in W^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) ; \operatorname{div} v=0\right.$ in the sense of distributions $\}$, and $W^{-1, p^{\prime}}\left(\Omega ; \mathbb{R}^{n}\right) \cong W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)^{*}$ with $p^{\prime}$ denoting the conjugate exponent, i.e. $p^{\prime}=p /(p-1)$. Likewise, $W^{k, p}$ indicates all $k$ th derivatives belonging to the $L^{p}$ space; for $k$ noninteger it refers to a fractional derivative and $W^{k, p}$ then denotes the Sobolev-Slobodetskiĭ space. Let us agree to use the norm $\|u\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}:=\|\nabla u\|_{L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right)}$. Also, "rca" will denote the regular countably additive set functions with respect to a Borel $\sigma$-algebra in question, also called Radon measures.

We will assume the following data qualification:
(2.1a) $\tau$ has a $C^{2}$-potential, $\tau(e): e \geqslant \zeta_{1}|e|^{p},|\tau(e)| \leqslant c\left(|e|^{p-1}+1\right), p>\frac{3 n}{n+2}$,

$$
\begin{align*}
&\left(\tau\left(e_{1}\right)-\tau\left(e_{2}\right)\right):\left(e_{1}-e_{2}\right) \geqslant\left\{\begin{array}{l}
\zeta_{1}\left|e_{1}-e_{2}\right|^{p}+\zeta_{2}\left|e_{1}-e_{2}\right|^{2} \text { if } p \geqslant 2 \\
\zeta_{0}\left(\left|e_{1}\right|+\left|e_{2}\right|\right)^{p-2}\left|e_{1}-e_{2}\right|^{2} \text { if } p<2
\end{array}\right.  \tag{2.1b}\\
& \sum_{i, j, k, l=1}^{n} \frac{\partial \tau_{i j}}{\partial e_{k l}} \xi_{i j} \xi_{k l} \geqslant\left\{\begin{array}{l}
\zeta_{3}\left(1+|e|^{p-2}\right)|\xi|^{2} \text { if } p \geqslant 2 \\
\zeta_{3}|e|^{p-2}|\xi|^{2} \text { if } p<2,
\end{array}\right.  \tag{2.1c}\\
& h \in \operatorname{rca}(\bar{\Omega}), f \in \operatorname{rca}(\Gamma), g \in L^{\infty}\left(\Omega ; \mathbb{R}^{n}\right), b \in C^{0,1}(\Gamma)  \tag{2.1d}\\
& \kappa>0, \alpha_{0}, \alpha_{1}, \alpha_{2} \geqslant 0, b(x) \geqslant b_{0}>0 \tag{2.1e}
\end{align*}
$$

with $\zeta_{i}>0, i=0, \ldots, 3$. An example of $\tau$ satisfying (2.1a-c) is $\tau(e)=\left(1+|e|^{p-2}\right) e$ (if $p \geqslant 2$ ) or $\tau(e)=|e|^{p-2} e$ (if $p \leqslant 2$ ). Let us also recall that (2.1a-c) ensures
(2.2a) $\int_{\Omega} \tau(e(\nabla u)): \nabla u \mathrm{~d} x \geqslant c_{1}\|u\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}^{p}$
$(2.2 \mathrm{~b}) \int_{\Omega}\left(\tau\left(e\left(\nabla u_{1}\right)\right)-\tau\left(e\left(\nabla u_{2}\right)\right)\right): e\left(\nabla u_{1}-\nabla q_{2}\right) \mathrm{d} x$

$$
\geqslant\left\{\begin{array}{l}
\zeta_{1} c_{1, \Omega}\left\|u_{1}-u_{2}\right\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}^{p}+\zeta_{2} c_{2, \Omega}\left\|u_{1}-u_{2}\right\|_{W^{1,2}\left(\Omega ; \mathbb{R}^{n}\right)}^{2} \text { if } p \geqslant 2 \\
\zeta_{0} c_{0, \Omega}\left(\left\|u_{1}\right\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}+\left\|u_{2}\right\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}\right)\left\|u_{1}-u_{2}\right\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}^{2} \text { if } p<2
\end{array}\right.
$$

with some $c_{i, \Omega}>0$ resulting from Korn's inequality, $c_{0, \Omega}(\cdot)$ decreasing; cf. [12, Sect 5.1.2]. Let us also introduce an exponent $q$ by

$$
\frac{2 p}{p-1} \leqslant q\left\{\begin{array}{l}
<\frac{p n}{n-p} \text { if } p<n  \tag{2.3}\\
<+\infty \text { otherwise }
\end{array}\right.
$$

which ensures, in particular, the compact embedding $W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \subset L^{q}\left(\Omega ; \mathbb{R}^{n}\right)$. By using Green's formula once for (1.1a,b) and twice for (1.1c), one gets the following definition:

Definition. We will call $(u, \theta) \in W_{0, \text { DIV }}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \times W^{r, 2}(\Omega)$, with $r \in[0,1]$ satisfying

$$
\begin{equation*}
\frac{2 n-2 p-p n}{2 p}<r<\frac{4-n}{2} \tag{2.4}
\end{equation*}
$$

a distributional solution to (1.1)-(1.2) if

$$
\begin{equation*}
\int_{\Omega}((u \cdot \nabla) u) \cdot v+\tau(e(\nabla u)): e(\nabla v)-g \cdot v\left(1-\alpha_{0} \theta\right) \mathrm{d} x=0 \tag{2.5}
\end{equation*}
$$

for any $v \in W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$, and

$$
\begin{align*}
\int_{\Omega}\left(\left(u \cdot\left(\nabla v+\alpha_{2} g v\right)+\kappa \Delta v\right) \theta\right. & +\alpha_{1} \tau(e(\nabla u)): e(\nabla u) v \mathrm{~d} x  \tag{2.6}\\
& +\int_{\bar{\Omega}} v h(\mathrm{~d} x)+\int_{\Gamma} v f(\mathrm{~d} S)=0
\end{align*}
$$

for any $v$ smooth with $\kappa \frac{\partial}{\partial \nu} v+b v=0$ on $\Gamma$.
Note that (2.1) ensures that $r \in[0,1]$ satisfying (2.4) does exist (recall that $n \leqslant 3$ ); in other words, (2.4) brings no restriction on $p$ if $n \leqslant 3$, as assumed. Let us remark that the inequalities in (2.4) imply respectively $W^{r, 2}(\Omega) \subset L^{q^{\prime}}(\Omega)$ and $W^{2-r, 2}(\Omega) \subset$
$C(\bar{\Omega})$; of course, $q^{\prime}:=q /(q-1)$. Also, (2.1) implies that all integrals in (2.2)-(2.6) have good sense. Also note that (2.1a) indeed enables us to choose $q$ such that $p^{-1}+2 q^{-1} \leqslant 1$, see (2.3), which implies that, e.g., the expression like $|v|^{2} \nabla v$ is integrable for any $v \in W^{1, p}(\Omega)$.

## 3. Existence of the distributional solution

We will prove the existence nonconstructively by using the Schauder fixed point theorem. First, we define the mapping

$$
\begin{equation*}
\mathcal{A}: \vartheta \mapsto u: L^{q^{\prime}}(\Omega) \rightarrow W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \tag{3.1}
\end{equation*}
$$

by $u$ being the weak solution to

$$
\begin{equation*}
(u \cdot \nabla) u-\operatorname{div} \tau(e(\nabla u))+\nabla \pi=g\left(1-\alpha_{0} \vartheta\right), \quad \operatorname{div} u=0,\left.\quad u\right|_{\Gamma}=0 \tag{3.2}
\end{equation*}
$$

For $q<p n /(n-p)$, let us agree to denote by $N_{q}^{1, p}$ the norm of the embedding $W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \subset L^{q}\left(\Omega ; \mathbb{R}^{n}\right)$.

Lemma 1. Assume (2.1). Then there is $R=R\left(p, \Omega, c, \zeta_{0}, \ldots, \zeta_{2}\right)>0$ such that $\mathcal{A}$ is single-valued and (weak,norm)-continuous with respect to the topologies indicated in (3.1) on the set

$$
\begin{equation*}
S_{R}:=\left\{\vartheta \in L^{q^{\prime}}(\Omega) ;\left\|g\left(1-\alpha_{0} \vartheta\right)\right\|_{L^{q^{\prime}}\left(\Omega ; \mathbb{R}^{n}\right)}<R\right\} . \tag{3.3}
\end{equation*}
$$

Proof. Take $\vartheta^{k} \rightharpoonup \vartheta$ in $L^{q^{\prime}}(\Omega)$, which implies $\vartheta^{k} \rightarrow \vartheta$ in $W^{-1, p^{\prime}}(\Omega)$ because $L^{q^{\prime}}(\Omega) \subset W^{-1, p^{\prime}}(\Omega)$ compactly, cf. (2.3). Then denote by $u^{k}$ the weak solution to (3.2) corresponding to $\vartheta^{k}$ in place of $\vartheta$; for the existence of $u^{k}$ we refer to Lions [10, Ch. II, Remark 5.5] after a modification to $\tau$ depending on $e(\nabla u)$ instead of $\nabla u$ or, even for $p \geqslant 2 n /(n+1)$, also Frehse, Málek and Steinhauer [4] or Růžička [22]. By testing with $u^{k}$, we get in a standard way the $a$-priori estimate

$$
\begin{equation*}
\left\|u^{k}\right\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}^{p-1} \leqslant \frac{N_{q}^{1, p}}{\zeta_{1}}\left\|g\left(1-\alpha_{0} \vartheta^{k}\right)\right\|_{L^{q^{\prime}}\left(\Omega ; \mathbb{R}^{n}\right)}<\frac{N_{q}^{1, p} R}{\zeta_{1}}=: R_{0}^{p-1} \tag{3.4}
\end{equation*}
$$

Taking a weakly convergent subsequence in $W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$, it is a standard procedure to show that its limit, denote it by $u$, is a weak solution to (3.2), cf. again [4], [10], [22].

Let us now prove uniqueness of $u$ provided $\vartheta \in S_{R}$ from (3.3) with $R$ small enough. Take two weak solutions $u^{1}, u^{2}$ of (3.2), and test the difference of the weak formulation of (3.2) by $u^{12}:=u^{1}-u^{2}$. This gives

$$
\begin{align*}
c\left\|u^{12}\right\|_{W_{0}^{1, \min (2, p)}\left(\Omega ; \mathbb{R}^{n}\right)}^{2} & \leqslant \int_{\Omega}\left(\tau\left(e\left(\nabla u^{1}\right)\right)-\tau\left(e\left(\nabla u^{2}\right)\right): e\left(\nabla u^{12}\right) \mathrm{d} x\right.  \tag{3.5}\\
& =\int_{\Omega}\left(\left(u^{2} \cdot \nabla\right) u^{2}-\left(u^{1} \cdot \nabla\right) u^{1}\right) \cdot u^{12} \mathrm{~d} x \\
& =-\int_{\Omega}\left(\left(u^{12} \cdot \nabla\right) u^{2}\right) \cdot u^{12} \mathrm{~d} x-\int_{\Omega}\left(\left(u^{1} \cdot \nabla\right) u^{12}\right) \cdot u^{12} \mathrm{~d} x \\
& \leqslant\left\|\nabla u^{2}\right\|_{L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right)}\left\|u^{12}\right\|_{L^{2 p^{\prime}}\left(\Omega ; \mathbb{R}^{n}\right)}^{2} \leqslant R_{0}\left\|u^{12}\right\|_{L^{2 p^{\prime}}\left(\Omega ; \mathbb{R}^{n}\right)}^{2}
\end{align*}
$$

with $c=\zeta_{2} c_{2, \Omega}$ (if $p \geqslant 2$ ) or $c=\zeta_{0} c_{0, \Omega}\left(2 R_{0}\right)$ (if $p<2$ ). Then, if $R$ is small enough so that, by (3.4), $R_{0}<c\left(N_{2 p^{\prime}}^{1, \min (2, p)}\right)^{-2}$, we get $u^{12}=0$. This, together with (3.4), gives the bound in (3.3).

Having the uniqueness of $u$, we can conclude that even the whole sequence $\left\{u^{k}\right\}$ converges weakly to $u$. Let us prove the strong convergence: subtracting (3.2) with $u$ and $u^{k}$, testing by $u^{k}-u$, and using Korn's inequality (2.2), we get

$$
\begin{align*}
& \varepsilon\left\|u^{k}-u\right\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}^{\max (2, p)} \leqslant \int_{\Omega}\left(\tau\left(e\left(\nabla u^{k}\right)\right)-\tau(e(\nabla u)): e\left(\nabla u^{k}-\nabla u\right) \mathrm{d} x\right. \\
& \quad=\int_{\Omega}\left(\left(u^{k} \cdot \nabla\right) u^{k}-(u \cdot \nabla) u\right) \cdot\left(u^{k}-u\right)+\alpha_{0}\left(\vartheta^{k}-\vartheta\right) g \cdot\left(u^{k}-u\right) \mathrm{d} x  \tag{3.6}\\
& \quad=: I_{1 k}+I_{2 k}
\end{align*}
$$

with $\varepsilon=\zeta_{1} c_{1, \Omega}$ (if $p \geqslant 2$ ) or $\varepsilon=\zeta_{0} c_{0, \Omega}\left(2 R_{0}\right)$ (if $p<2$ ). By using $\operatorname{div} u^{k}=0=\operatorname{div} u$ and Green's formula, we can calculate

$$
\begin{align*}
I_{1 k} & =\int_{\Omega} \sum_{j=1}^{n}\left(\left(\sum_{i=1}^{n} u_{i}^{k} \frac{\partial}{\partial x_{i}}\right) u_{j}^{k}-\left(\sum_{i=1}^{n} u_{i} \frac{\partial}{\partial x_{i}}\right) u_{j}\right)\left(u_{j}^{k}-u_{j}\right) \mathrm{d} x  \tag{3.7}\\
& =\int_{\Omega} \sum_{i, j=1}^{n} \frac{\partial}{\partial x_{i}}\left(u_{i}^{k} u_{j}^{k}-u_{i} u_{j}\right)\left(u_{j}^{k}-u_{j}\right) \mathrm{d} x \\
& =-\int_{\Omega} \sum_{i, j=1}^{n}\left(u_{i}^{k} u_{j}^{k}-u_{i} u_{j}\right) \frac{\partial}{\partial x_{i}}\left(u_{j}^{k}-u_{j}\right) \mathrm{d} x \\
& =-\int_{\Omega}\left(u^{k} \otimes u^{k}-u \otimes u\right): \nabla\left(u^{k}-u\right) \mathrm{d} x
\end{align*}
$$

Due to the boundedness of $\nabla\left(u^{k}-u\right)$ in $L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right)$, the compact embedding $W^{1, p}(\Omega) \subset L^{q}(\Omega)$ with $p^{-1}+2 q^{-1} \leqslant 1$, and the continuity of the Nemytskiŭ mapping
$u \mapsto u \otimes u: L^{q}\left(\Omega ; \mathbb{R}^{n}\right) \rightarrow L^{q / 2}\left(\Omega ; \mathbb{R}^{n \times n}\right)$, we have $u^{k} \otimes u^{k} \rightarrow u \otimes u$ in $L^{q / 2}\left(\Omega ; \mathbb{R}^{n \times n}\right)$, and eventually we get $I_{1 k} \rightarrow 0$.

Also, the term $I_{2 k}$ converges to zero because $\vartheta^{k} \rightarrow \vartheta$ in $W^{-1, p^{\prime}}(\Omega)$ and $u^{k} \rightharpoonup u$ in $W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$.

Furthermore, let us consider the Nemytskiǔ-type mapping $\mathcal{N}: W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \times$ $L^{q^{\prime}}(\Omega) \rightarrow \operatorname{rca}(\bar{\Omega})$ defined by

$$
\begin{equation*}
\mathcal{N}:(u, \vartheta) \mapsto h_{1}=\alpha_{1} \tau(e(\nabla u)): e(\nabla u)+\alpha_{2} g \cdot u \vartheta+h, \tag{3.8}
\end{equation*}
$$

and, for $u \in W_{0, \mathrm{DIV}}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$, the linear operator

$$
\begin{equation*}
\mathcal{B}_{u}:\left(h_{1}, f\right) \mapsto \theta: \operatorname{rca}(\bar{\Omega}) \times \operatorname{rca}(\Gamma) \rightarrow W^{r, 2}(\Omega) \tag{3.9}
\end{equation*}
$$

with $\theta$ being the distributional solution to

$$
\begin{equation*}
u \cdot \nabla \theta-\kappa \Delta \theta=h_{1} \quad \text { on } \Omega, \quad \kappa \frac{\partial \theta}{\partial \nu}+b \theta=f \quad \text { on } \Gamma, \tag{3.10}
\end{equation*}
$$

i.e. $\theta \in W^{r, 2}(\Omega)$ satisfies the identity

$$
\begin{equation*}
\int_{\Omega}(u \cdot \nabla v+\kappa \Delta v) \theta \mathrm{d} x+\int_{\bar{\Omega}} v h_{1}(\mathrm{~d} x)+\int_{\Gamma} v f(\mathrm{~d} S)=0 \tag{3.11}
\end{equation*}
$$

for any $v$ smooth with $\kappa \frac{\partial}{\partial \nu} v+b v=0$ on $\Gamma$.
Lemma 2. Let (2.1) be valid. Then the mappings $\mathcal{N}$ and $\mathcal{B}_{u}$ are well defined and both $\mathcal{N}: W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \times L^{q^{\prime}}(\Omega) \rightarrow \operatorname{rca}(\bar{\Omega})$ and $\left(u, h_{1}\right) \mapsto \mathcal{B}_{u}\left(h_{1}, f\right): W_{0, \operatorname{DIV}}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \times$ $\operatorname{rca}(\bar{\Omega}) \rightarrow L^{q^{\prime}}(\Omega)$ are (norm $\times$ weak $^{*}$, weak ${ }^{*}$ )-continuous.

Proof. By the classical result about Nemytskiĭ mappings, $\mathcal{N}_{0}:(\xi, \vartheta) \mapsto$ $\alpha_{1} \tau(e(\xi)): e(\xi)+\alpha_{2} g \cdot u \vartheta: L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right) \times L^{q^{\prime}}(\Omega) \rightarrow L^{1}(\Omega)$ is continuous, so that $\mathcal{N}=\left(\mathcal{N}_{0} \circ \nabla\right)+h$ is continuous, as claimed.

Let us consider the weak solution to the auxiliary linear problem

$$
\begin{equation*}
-u \cdot \nabla v-\kappa \Delta v=\xi \text { on } \Omega, \quad \kappa \frac{\partial v}{\partial \nu}+b v=0 \quad \text { on } \Gamma . \tag{3.12}
\end{equation*}
$$

The existence of $v$ can be proved by the standard energy method by testing (3.12) by $v$; note that

$$
\begin{equation*}
\int_{\Omega}(u \cdot \nabla v) v \mathrm{~d} x=\frac{1}{2} \int_{\Omega} u \cdot \nabla v^{2} \mathrm{~d} x=-\frac{1}{2} \int_{\Omega}(\operatorname{div} u) v^{2} \mathrm{~d} x=0 \tag{3.13}
\end{equation*}
$$

so that we have the estimate $\|v\|_{W^{1,2}(\Omega)} \leqslant K_{1}\|\xi\|_{W^{1,2}(\Omega)^{*}}$ independent of $u$. Moreover, we have also the estimate

$$
\begin{align*}
\int_{\Omega}(u \cdot \nabla v) \Delta v \mathrm{~d} x & \leqslant\|u\|_{L^{q}\left(\Omega ; \mathbb{R}^{n}\right)}\|\nabla v\|_{L^{2 q /(q-2)}\left(\Omega ; \mathbb{R}^{n}\right)}\|\Delta v\|_{L^{2}(\Omega)} \\
(3.14) & \leqslant\|u\|_{L^{q}\left(\Omega ; \mathbb{R}^{n}\right)}\|\nabla v\|_{L^{2}\left(\Omega ; \mathbb{R}^{n}\right)}^{\lambda}\|v v\|_{L^{6}\left(\Omega ; \mathbb{R}^{n}\right)}^{1-\lambda}\|\Delta v\|_{L^{2}(\Omega)}  \tag{3.14}\\
& \leqslant N_{q}^{1, p}\|\nabla u\|_{L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right)} K_{1}^{\lambda}\left(N_{2}^{1,2}\right)^{\lambda}\|\xi\|_{L^{2}(\Omega)}^{\lambda}\left(K_{0} N_{6}^{1,2}\right)^{1-\lambda}\|\Delta v\|_{L^{2}(\Omega)}^{2-\lambda}
\end{align*}
$$

for $\lambda \in(0,1)$ such that $\lambda \frac{1}{2}+(1-\lambda) \frac{1}{6}=\frac{q-2}{2 q}$ which certainly does exist for $p>3 / 2$, and where the constant $K_{0}$ comes from the standard Laplace-operator regularity $\|v\|_{W^{2,2}(\Omega)} \leqslant K_{0}\|\Delta v\|_{L^{2}(\Omega)}$ with the boundary condition $\kappa \frac{\partial v}{\partial \nu}+b v=0$ with $b \in$ $C^{0,1}(\Gamma)$ on the $C^{3,1}$-domain $\Omega$; see Nečas [15]. Then, multiplying (3.12) by $\Delta v$ and integrating over $\Omega$, we get the estimate
$\kappa \int_{\Omega}|\Delta v|^{2} \mathrm{~d} x=-\int_{\Omega}(\xi+u \cdot \nabla v) \Delta v \mathrm{~d} x \leqslant\|\xi\|_{L^{2}(\Omega)}\|\Delta v\|_{L^{2}(\Omega)}$

$$
\begin{equation*}
+N_{q}^{1, p}\|u\|_{W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)} K_{1}^{\lambda}\left(N_{2}^{1,2}\right)^{\lambda}\|\xi\|_{L^{2}(\Omega)}^{\lambda}\left(K_{0} N_{6}^{1,2}\right)^{1-\lambda}\|\Delta v\|_{L^{2}(\Omega)}^{2-\lambda} . \tag{3.15}
\end{equation*}
$$

Thus we can see that, if $\xi \in L^{2}(\Omega), \Delta v$ is bounded in $L^{2}(\Omega)$. Then, using again the Laplace-operator regularity, we get $\|v\|_{W^{2,2}(\Omega)} \leqslant K_{u}\|\xi\|_{L^{2}(\Omega)}$ with $K_{u}>0$ depending on $\|u\|_{W_{0, \mathrm{DIV}}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)}$ continuously and increasingly. It is important that this regularity estimate holds uniformly for $u$ ranging over bounded sets in $W_{0, \mathrm{DIV}}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$.

The interpolation between the linear mappings $\xi \mapsto v: W^{1,2}(\Omega)^{*} \rightarrow W^{1,2}(\Omega)$ and $L^{2}(\Omega) \rightarrow W^{2,2}(\Omega)$ gives a mapping $W^{r, 2}(\Omega)^{*} \rightarrow W^{2-r, 2}(\Omega)$ and an estimate $\|v\|_{W^{2-r, 2}(\Omega)} \leqslant K_{u}^{1-r} K_{1}^{r}\|\xi\|_{W^{r, 2}(\Omega)^{*}}$.

Let us rewrite the identity (3.11) into the form $\left\langle B_{u} v, \theta\right\rangle+\langle F, v\rangle=0$ where $B_{u}$ : $W^{2-r, 2}(\Omega) \rightarrow W^{r, 2}(\Omega)^{*}$ and $F \in W^{2-r, 2}(\Omega)^{*}$ are defined by

$$
\begin{equation*}
B_{u} v:=u \cdot \nabla v+\kappa \Delta v, \quad\langle F, v\rangle=\int_{\bar{\Omega}} v h_{1}(\mathrm{~d} x)+\int_{\Gamma} v f(\mathrm{~d} S) \tag{3.16}
\end{equation*}
$$

respectively. Then $\theta=-\left(B_{u}^{*}\right)^{-1} F=-F \circ B_{u}^{-1} \in W^{r, 2}(\Omega)^{* *} \cong W^{r, 2}(\Omega)$ is a solution to $\left\langle B_{u} v, \theta\right\rangle+\langle F, v\rangle=0$. Moreover, because of surjectivity of $B_{u}$, this solution must be unique. Also, we have the estimate $\|\theta\|_{W^{r, 2}(\Omega)} \leqslant K_{u}^{1-r} K_{1}^{r}\|F\|_{W^{2-r, 2}(\Omega)^{*}}$ independent of $u$.

Then we choose $0 \leqslant r \leqslant 1$ so small that $W^{2-r, 2}(\Omega) \subset C(\bar{\Omega})$, i.e. $r<(4-n) / 2$, cf. (2.4). This eventually gives the estimate

$$
\begin{equation*}
\|\theta\|_{W^{r, 2}(\Omega)} \leqslant K_{u}^{1-r} K_{1}^{r}\|F\|_{W^{2-r, 2}(\Omega)^{*}} \leqslant N_{\infty}^{2-r, 2} K_{u}^{1-r} K_{1}^{r}\left\|\left(h_{1}, f\right)\right\|_{\mathrm{rca}(\bar{\Omega}) \times \operatorname{rca}(\Gamma)} \tag{3.17}
\end{equation*}
$$

with $N_{\infty}^{2-r, 2}$ the norm of the embedding $W^{2-r, 2}(\Omega) \subset L^{\infty}(\Omega)$; note that $\left(h_{1}, f\right) \mapsto$ $F: \operatorname{rca}(\bar{\Omega}) \times \operatorname{rca}(\Gamma) \rightarrow W^{2-r, 2}(\Omega)^{*}$ defined by $(3.16)$ is the adjoint mapping to $v \mapsto$ $\left(v,\left.v\right|_{\Gamma}\right): W^{2-r, 2}(\Omega) \rightarrow C(\bar{\Omega}) \times C(\Gamma)$.

To prove continuity of $\left(u, h_{1}\right) \mapsto \mathcal{B}_{u}\left(h_{1}, f\right)$, let us take $h_{1, k} \rightarrow h_{1}$ in rca $(\bar{\Omega})$ weakly* and $u^{k} \rightarrow u$ in $W^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$, and denote by $\theta^{k}$ the distributional solution to (3.10) corresponding to $u^{k}$ and $h_{1, k}$ in place of $u$ and $h_{1}$, respectively. We showed that $\theta^{k}$ does exist and is bounded in $W^{r, 2}(\Omega)$; realize that $\left\{\nabla u^{k}\right\}$ is bounded in $L^{p}\left(\Omega ; \mathbb{R}^{n \times n}\right)$. Then, by Banach-Alaoglu-Bourbaki theorem, we can assume that, possibly up to a subsequence,

$$
\begin{equation*}
\theta^{k} \rightarrow \theta \quad \text { weakly in } W^{r, 2}(\Omega) \tag{3.18}
\end{equation*}
$$

Then we can make the limit passage in the integral identity (3.11), which reads here

$$
\begin{equation*}
\int_{\Omega}\left(u^{k} \cdot \nabla v+\kappa \Delta v\right) \theta^{k} \mathrm{~d} x+\int_{\bar{\Omega}} v h_{1, k}(\mathrm{~d} x)+\int_{\Gamma} v f(\mathrm{~d} S)=0 . \tag{3.19}
\end{equation*}
$$

Note that certainly the term $\theta^{k} u^{k}$ converges to $\theta u$ (even strongly) because, as a consequence of (3.18), $\left\{\theta^{k}\right\}$ converges strongly in $W^{-1, p^{\prime}}(\Omega)$ and $\left\{u^{k}\right\}$ also strongly in $W_{0}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$. Thus $\theta=\mathcal{B}_{u}(h, f)$ and even the whole sequence $\left\{\theta^{k}\right\}$ converges because of the already proved uniqueness of $\theta$.

Furthermore, for $\varrho>0$, we denote the ball of the radius $\varrho$ in $L^{q^{\prime}}(\Omega)$ by

$$
\begin{equation*}
B_{\varrho}:=\left\{\vartheta \in L^{q^{\prime}}(\Omega) ; \quad\|\vartheta\|_{L^{q^{\prime}}(\Omega)} \leqslant \varrho\right\} . \tag{3.20}
\end{equation*}
$$

Proposition 1. Let (2.1) be fulfilled and let $\|g\|_{L^{\infty}\left(\Omega ; \mathbb{R}^{n}\right)}$ be sufficiently small with respect to the other data $\alpha_{0}, \alpha_{1}, \alpha_{2},\|h\|_{\text {rca }(\bar{\Omega})}$ and $\|f\|_{\text {rca }(\Gamma)}$. Then (1.1)-(1.2) has at least one distributional solution $(u, \theta)$.

Proof. We will investigate the mapping $\mathcal{C}: L^{q^{\prime}}(\Omega) \rightarrow L^{q^{\prime}}(\Omega)$ defined by

$$
\begin{equation*}
\mathcal{C}(\vartheta):=\mathcal{B}_{\mathcal{A}(\vartheta)}(\mathcal{N}(\mathcal{A}(\vartheta), \vartheta), f) . \tag{3.21}
\end{equation*}
$$

Note that any fixed point $\theta$ of $\mathcal{C}$ satisfies $\theta=\mathcal{B}_{u}(h, f)$ with $h=\mathcal{N}(u, \theta)$, where $u=\mathcal{A}(\theta)$, which just means that the pair $(u, \theta)$ is the distributional solution to (1.1)-(1.2). We will show that

$$
\begin{equation*}
B_{\varrho} \subset S_{R} \quad \text { and } \quad \mathcal{C}\left(B_{\varrho}\right) \subset B_{\varrho} \tag{3.22}
\end{equation*}
$$

provided $\varrho$ is chosen appropriately and $g$ is small enough. Obviously, $(u, \theta)=$ $(\mathcal{A}(\vartheta), \mathcal{C}(\vartheta))$ solves the decoupled system (3.2) and (3.10) with $u=\mathcal{A}(\vartheta)$ and $h_{1}=$
$h_{u, \vartheta}=\mathcal{N}(u, \vartheta)$. Then, by testing (3.2) by $u$, we get the estimate (3.4) with the subscript $k$ omitted.
Furthermore, using the identity $\int_{\Omega} \tau(e(\nabla u)): e(\nabla u) \mathrm{d} x=\int_{\Omega} g\left(1-\alpha_{0} \vartheta\right) u \mathrm{~d} x$ the source term $h_{u, \vartheta}$ in (3.10) can be estimated as

$$
\begin{align*}
\left\|h_{u, \vartheta}\right\|_{\mathrm{rca}(\bar{\Omega})} \leqslant & \alpha_{1}\|g u\|_{L^{1}(\Omega)}+\mid \alpha_{0} \alpha_{1}-\alpha_{2}\|g \cdot u \vartheta\|_{L^{1}(\Omega)}+\|h\|_{\mathrm{rca}(\bar{\Omega})} \\
& \leqslant  \tag{3.23}\\
& \left(\alpha_{1} N_{1}^{1, p}+\left|\alpha_{0} \alpha_{1}-\alpha_{2}\right| N_{q}^{1, p}\|\vartheta\|_{L^{q^{\prime}}(\Omega)}\right)\|g\|_{L^{\infty}\left(\Omega ; \mathbb{R}^{n}\right)}\|u\|_{W^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)} \\
& +\|h\|_{\mathrm{rca}(\bar{\Omega})} \leqslant \gamma_{1}+\gamma_{2}\left(\|g\|_{L^{\infty}\left(\Omega ; \mathbb{R}^{n}\right)}\right) \varrho^{p^{\prime}},
\end{align*}
$$

where we assume $\vartheta \in B_{\varrho}$ and take into account that $R_{0}=\|g\|_{L^{\infty}\left(\Omega ; \mathbb{R}^{n}\right)}^{1 /(p-1)} \mathcal{O}\left(\|\vartheta\|_{L^{q^{\prime}}(\Omega)}^{1 /(p-1)}\right)$, cf. (3.4); then $\gamma_{1}=\gamma_{1}\left(\alpha_{1}, c, p,\|h\|_{\mathrm{rca}(\bar{\Omega})}\right)$ and $\gamma_{2}(\cdot)$ depends on $\alpha_{0}, \alpha_{2}, p$, and $\zeta_{1}$ and moreover $\lim _{a \rightarrow 0+} \gamma_{2}(a)=0$.
The estimate (3.17) now reads

$$
\left\|\mathcal{B}_{u}\left(h_{u, \vartheta}, f\right)\right\|_{W^{r, 2}(\Omega)} \leqslant N_{\infty}^{2-r, 2} K_{u}^{1-r} K_{1}^{r}\left(\left\|h_{u, \vartheta}\right\|_{\text {rca }(\bar{\Omega})}+\|f\|_{\text {rca }(\Gamma)}\right) .
$$

Altogether,

$$
\begin{align*}
\|\mathcal{C}(\vartheta)\|_{L^{\prime}(\Omega)} & \leqslant N_{q^{\prime}, 2}^{r,}\|\mathcal{C}(\vartheta)\|_{W^{r, 2}(\Omega)}  \tag{3.24}\\
& \leqslant N_{q^{q^{2}}}^{r,} N_{\infty}^{2-r, 2} K_{u}^{1-r} K_{1}^{r}\left(\gamma_{1}+\gamma_{2}\left(\|g\|_{L^{\infty}\left(\Omega ; \mathbb{R}^{n}\right)}\right) \varrho^{p^{\prime}}+\|f\|_{\mathrm{rca}(\Gamma)}\right) .
\end{align*}
$$

If $g$ is small, one can find $\varrho>N_{q^{\prime}}^{r, 2} N_{\infty}^{2-r, 2} K_{u}^{1-r} K_{1}^{r}\left(\gamma_{1}+\|f\|_{\text {rca( }(\Gamma)}\right)$ small enough so that (3.24) implies $\|\mathcal{C}(\vartheta)\|_{L^{q^{\prime}}(\Omega)} \leqslant \varrho$. In other words, we have proved $\mathcal{C}\left(B_{\varrho}\right) \subset B_{\varrho}$ for such $\varrho$. Moreover, if $g$ is small enough, we have also $B_{\varrho} \subset S_{R}$.

We endow $B_{\varrho}$ with the weak (or, if $q^{\prime}=+\infty$, weak*) topology of $L^{q^{\prime}}(\Omega)$, which makes $B_{\varrho}$ compact (note that, due to (2.3), always $q^{\prime}>1$ ). By Lemmas 1 and 2 and by (3.22), $\mathcal{C}$ maps $B_{\varrho}$ (weak,weak)-continuously into itself. Then, by Schauder's theorem, it has a fixed point $\theta$ on $B_{\varrho}$.

Remark 1. The interpolation/transposition method in Hilbert-space setting was thoroughly presented by Lions and Magenes [11]. Here, however, we did not assume infinitely smooth $\Gamma$ or the coefficients $u$ and $b$ in (3.12) and, moreover, it was important to derive the estimate (3.17) uniformly for $u$ from bounded sets in $W_{0, \mathrm{DIV}}^{1, p}\left(\Omega ; \mathbb{R}^{n}\right)$.

Remark 2. Under a quite restrictive assumption $p>2 n$, we can alternatively use a continuous imbedding of $\mathcal{W}:=\left\{v \in W^{1,2}(\Omega) ; \Delta v \in L^{n / 2+\varepsilon}(\Omega), \frac{\partial}{\partial \nu} v \in\right.$ $\left.L^{n-1+\varepsilon}(\Gamma)\right\}$ with $\varepsilon>0$ into $C^{0}(\bar{\Omega})$, proved by Alibert and Raymond [1] even for Lipschitz domains. Indeed, for $u \in W^{1, p}\left(\Omega ; \mathbb{R}^{n}\right) \subset L^{q}\left(\Omega ; \mathbb{R}^{n}\right)$ with $q$ satisfying (2.3) and for $v \in W^{1,2}(\Omega)$, we have $u \cdot \nabla v \in L^{n / 2+\varepsilon}(\Omega)$, which enables us to
get the auxiliary mapping $\xi \mapsto v: L^{n / 2+\varepsilon}(\Omega) \rightarrow \mathcal{W}$ in the proof of Lemma 1. Then $B_{u}: \mathcal{W} \rightarrow L^{n / 2+\varepsilon}(\Omega)$ and all above considerations work equally for $\theta$ in $L^{n /(n-2)-\varepsilon}(\Omega)$ instead of $W^{r, 2}(\Omega)$. Beside the $C^{0,1}$-domain $\Omega$, this modification enables us also to consider $b$ from $L^{4 / 3+\varepsilon}(\Gamma)$ (if $n=2$ ) or from $L^{6+\varepsilon}(\Gamma)$ (if $n=3$ ) because then $b v \in L^{n-1+\varepsilon}(\Gamma)$ for any $v \in W^{1,2}(\Omega)$.

Remark 3. Contrary to the evolution case (cf. [16]), if $\alpha_{2}>0$, it does not seem possible to prove $\theta \geqslant 0$ for some solution obtained in Proposition 1 even if one assumes $h \geqslant 0$ and $f \geqslant 0$. Yet, negative temperature need not be interpreted as non-physical solution because $\theta$ is a "small" deviation from some constant reference temperature rather than the absolute temperature. Nevertheless, this holds true if the adiabatic effect can be neglected, i.e. $\alpha_{2}=0$. Then, assuming $h \geqslant 0$ and $f \geqslant 0$ and regularizing (1.1c) by a term $\varepsilon \theta$ on the left-hand side, we can prove existence of the "regularized" solution $\left(u_{\varepsilon}, \theta_{\varepsilon}\right)$ again by Proposition 1 with all estimates independent of $\varepsilon>0$ and then nonnegativity $\theta_{\varepsilon} \geqslant 0$ by testing $\varepsilon \theta_{\varepsilon}+u_{\varepsilon} \cdot \nabla \theta_{\varepsilon}-\kappa \Delta \theta_{\varepsilon}=h_{1} \geqslant 0$ by $\operatorname{signum}\left(\theta_{\varepsilon}\right)-1$ or, more rigorously, by a regularization of this test function. Then, passing with $\varepsilon \rightarrow 0$, one gets $\theta \geqslant 0$.

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