

Analysis of a singular convection diffusion system arising in turbulence modelling

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Abstract

We shall study some singular stationary convection diffusion system governing the steady-state of a turbulence model closely related to the $k - \varepsilon$ one. We shall establish existence, positivity and regularity results in a very general framework.

Keywords: degenerate and singular elliptic system, convection-diffusion problem, $k - \varepsilon$ model, Stampacchia estimates.

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1 Introduction

We shall first recall some basic ideas concerning the statistical turbulence modelling for fluids. The reader can consult [23, 26] for a more detailed introduction.

Let \mathbf{u}, p, ρ, T be the velocity, pressure, density and temperature of a newtonien compressible fluid. Let also $\tilde{\Omega} \subset \mathbb{R}^3$ a domain which is assumed to be bounded. Then the motion of the flow in $\tilde{\Omega}$ at a time $t \in \mathbb{R}^+$ can be described by the compressible Navier Stokes equations (see system (C) page 8 in [18]). It is well known that direct simulation based on such a model is harder or even impossible at high reynolds numbers. The reason is that too many points of discretization are necessary and so only very simple configurations can be handled.

Thus engineers and physicists have proposed new sets of equations to describe the average of a turbulent flow. The most famous one is the $k - \varepsilon$ model, introduced by Kolmogorov [13]. We shall briefly present its basic principles in the following. Let v denote a generic physical quantity subject to turbulent

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(i.e. unpredictable at the macroscopic scale), we introduce its mean part (or its esperance) $\langle v \rangle$ by setting:

$$\langle v(x, t) \rangle = \int_{\mathcal{P}} v(\omega, x, t) d\mathcal{P}(\omega),$$

where the integral is taken in a probalistic context which we shall not detail any more here. Note however that the operation $\langle \cdot \rangle$ is more generally called a filter. The probalistic meaning is one but not the only possible filter (see for instance [23] chap.3). We shall then consider the decomposition: $v = \langle v \rangle + v'$, where v' is referred to the non computable or the non relevant part and $\langle v \rangle$ is called the mean part (i.e. the macroscopic part).

The principle of the $k - \varepsilon$ model is to describe the mean flow in terms of the mean quantities $\langle \mathbf{u} \rangle, \langle p \rangle, \langle \rho \rangle, \langle T \rangle$ together with two scalar functions k and ε , which contains relevant informations about the small scales processes (or the turbulent processes). The variable k (SI: $[\frac{m^2}{s^2}]$) is called the turbulent kinetic energy and ε $[\frac{m^2}{s^3}]$ is the rate of dissipation of the kinetic energy. They are defined by:

$$k = \frac{1}{2} |\mathbf{u}'|^2 \quad \varepsilon = \frac{\nu}{2} \langle |\nabla \mathbf{u}' + (\nabla \mathbf{u}')^T|^2 \rangle, \quad (1)$$

where ν is the molecular viscosity of the fluid. The model is then constructed by averaging (i.e. by applying the operator $\langle \cdot \rangle$ on) the Navier-Stokes equations. Under appropriate assumptions (i.e. the Reynolds hypothesis in the incompressible case, and the Favre average in the compressible case) we obtain a closed system of equations for the variables $\langle \mathbf{u} \rangle, \langle p \rangle, \langle \rho \rangle, \langle T \rangle, k$ and ε (see [23] pages 61-62 for the incompressible case, and pages 116-117 in the compressible situation).

Here we shall focus on the equations for k and ε and we consider that the others quantities $\langle \mathbf{u} \rangle, \langle p \rangle, \langle \rho \rangle, \langle T \rangle$ are known. Moreover, in order to simplify the readability we do not use the notation $\langle \cdot \rangle$, i.e. in the sequel we will write \mathbf{u} instead of $\langle \mathbf{u} \rangle$ and ρ instead $\langle \rho \rangle$ to represent the mean velocity and density of the fluid. The equations for k and ε are of convection-diffusion-reaction type:

$$\partial_t k + \mathbf{u} \cdot \nabla k - \frac{c_\nu}{\rho} \operatorname{div} \left(\left(\nu + \rho \frac{k^2}{\varepsilon} \right) \nabla k \right) = c_\nu \frac{k^2}{\varepsilon} F - \frac{2}{3} k D - \varepsilon, \quad (2)$$

$$\partial_t \varepsilon + \mathbf{u} \cdot \nabla \varepsilon - \frac{c_\varepsilon}{\rho} \operatorname{div} \left(\left(\nu + \rho \frac{k^2}{\varepsilon} \right) \nabla \varepsilon \right) = c_1 k F - \frac{2c_1}{3c_\nu} \varepsilon D - c_2 \frac{\varepsilon^2}{k}, \quad (3)$$

where $D(x, t) := \operatorname{div} \mathbf{u}(x, t)$, $F(x, t) := \frac{1}{2} |\nabla \mathbf{u} + (\nabla \mathbf{u})^T|^2 - \frac{2}{3} D(x, t)^2 \geq 0$ (see subsection 6.2 in the Appendix) and $c_\nu, c_\varepsilon, c_1, c_2$ are generally taken as positive constants (see (70) in the Appendix).

Note that equations (2)-(3) are only valid sufficiently far from the walls. In fact, in the vicinity of the walls of the domain $\tilde{\Omega}$, there is a thin domain Σ , called logarithmic layer in which the modulus of the velocity goes from 0 to $\mathcal{O}(1)$. In this layer we can use some wall law or a one equation model (see [23] chap.1 and [24]) instead of (2)-(3). However equations (2)-(3) can be considered even in the logarithmic layer if we allow the coefficients $c_\nu, c_\varepsilon, c_1$ and

c_2 to depend appropriately on some local Reynolds numbers (see [23] pages 59-60 and page 115). In this last situation the system is called Low-Reynolds number k - ε model.

In the following we focus on the study in the domain $\Omega := \tilde{\Omega} \setminus \Sigma$ and we assume that its boundary $\partial\Omega$ is Lipschitz¹. We denote by $\mathbf{n}(x)$ the outward normal defined for all points $x \in \partial\Omega$. The boundary conditions for k and ε on $\partial\Omega$ are then on the form:

$$k = k_0 \quad \text{and} \quad \varepsilon = \varepsilon_0 \quad \text{on } \partial\Omega, \quad (4)$$

where k_0 and ε_0 are strictly positive functions which can be calculated by using a wall law (see [23] p.59) or a one equation model (see [24]). In the following we assume that k_0 and ε_0 are given. Moreover we can assume (see again [23] p.59) that:

$$\mathbf{u} \cdot \mathbf{n} = 0 \quad \text{on } \partial\Omega. \quad (5)$$

We shall concentrate in this paper on a modified system obtained after introducing the new variables θ [s] and φ [m^{-2}] given by:

$$\theta = \frac{k}{\varepsilon} \quad \varphi = \frac{\varepsilon^2}{k^3}. \quad (6)$$

These variables have a physical meaning (see [26]): θ represents a characteristic time of turbulence and $L = \varphi^{-1/2}$ is a characteristic turbulent length scale. By using this change of variable in the equations (2)-(3) and after considering some modelisation arguments for the diffusion processes (see the Annexe) we obtain:

$$(P) \quad \begin{cases} \partial_t \theta + \mathbf{u} \cdot \nabla \theta - \frac{1}{\rho} \operatorname{div} \left(\left(\nu + \frac{c_\theta \rho}{\theta \varphi} \right) \nabla \theta \right) = -c_3 \theta^2 F + c_4 \theta D + c_5 \\ \partial_t \varphi + \mathbf{u} \cdot \nabla \varphi - \frac{1}{\rho} \operatorname{div} \left(\left(\nu + \frac{c_\varphi \rho}{\theta \varphi} \right) \nabla \varphi \right) = -\varphi (c_6 \theta F - c_7 D + c_8 \theta^{-1}) \end{cases}$$

where the coefficients c_θ, c_φ and c_i are all positive.

Problem (P) differs from the k - ε one only by the diffusive parts and it is attractive by some stronger mathematical properties. Another model closely related to these systems, and having some popularity, is the k - ω one (with $\omega = \theta^{-1}$, see for instance [26]).

In the papers [16, 20] the authors have established the existence of a weak solution for (P) and a property of positiveness. This last feature takes the model useful in practice: it can be used directly or also as an intermediary stabilization procedure to the k - ε one (see [22]). Another important property attempted for a turbulence model is its capability to predict the possible steady states. In the previous works (except in [21]) only the evolutive version of (P) was studied under very restrictive assumptions. In [21] however the stationary problem is studied but it is simplified by considering a perturbation of the viscosities that artificially cancel the singularity of the system.

¹See [9] p. 127

Hence in this paper we shall study the stationary version of (P) on a bounded domain $\Omega \subset \mathbb{R}^N$, $N=2$ or 3 , on which we impose the boundary conditions $\theta = a, \varphi = b$ on $\partial\Omega$. Remark that by using (4) together with (6) we obtain

$$a = \frac{k_0}{\varepsilon_0} \quad b = \frac{\varepsilon_0^2}{k_0^3}. \quad (7)$$

Hence we can assume that a and b are strictly positive given functions.

We shall establish existence, positivity and regularity results in a very general framework.

2 Main results

2.1 Assumptions and notations

Let (Q) denote the stationary system associated to (P). For simplicity we introduce the new parameters $C_{\text{ind}} := \rho c_{\text{ind}}$ where the subscript 'ind' takes the integer values $3,4,5,6,7,8$ or the letters θ and φ . Then our main model (Q) has the following form:

$$(Q) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \theta - \operatorname{div} \left(\left(\nu + \frac{C_\theta}{\theta \varphi} \right) \nabla \theta \right) = -C_3 F \theta^2 + C_4 \theta D + C_5 & \text{in } \Omega \\ \rho \mathbf{u} \cdot \nabla \varphi - \operatorname{div} \left(\left(\nu + \frac{C_\varphi}{\theta \varphi} \right) \nabla \varphi \right) = -\varphi (C_6 \theta F - C_7 D + C_8 \theta^{-1}) & \text{in } \Omega \\ \theta = a, \varphi = b & \text{on } \partial\Omega \end{cases}$$

For physical reasons we are only interested in positive solutions (θ, φ) for (Q). Note however that even with this restriction, the problem (Q) may be singular (i.e. the viscosities $\nu + \frac{C_\theta}{\theta \varphi}$ and $\nu + \frac{C_\varphi}{\theta \varphi}$ may be unbounded). Moreover, because we allow $\nu \equiv 0$ the equations may degenerate (i.e. the viscosities may vanish). Hence without additional restriction there may be various non equivalent notions of weak solution (see for instance [8]).

In fact a good compromise between respect of the physics, simplification of the mathematical study and obtention of significative results, is to restrict θ and φ to be within the class \mathcal{S} defined by:

$$\mathcal{S} = \left\{ f : \Omega \rightarrow \mathbb{R}^+ \text{ such that } f \in H^1(\Omega) \cap L^\infty(\Omega), f^{-1} \in L^\infty(\Omega) \right\}.$$

In particular, if the parameters appearing in (Q) are sufficiently regular and if we restrict θ and φ to be within the classe \mathcal{S} , then the notion of a weak solution for (Q) is univocally defined: it is a distributional solution (θ, φ) that satisfies the boundary conditions in the sense of the trace.

In this last situation we will tell that (θ, φ) is a **weak solution of (Q) in the class \mathcal{S}** .

In order to can consider such a weak solution for (Q) we shall precise in the following some sufficient conditions of regularity for the data.

Let $N=2$ or 3 denote the dimension of the domain Ω , and r be a fixed number such that:

$$r > \frac{N}{2}. \quad (8)$$

We then have the following continuous injection (see lemma 5):

$$L^r(\Omega) \subset W^{-1,\beta}(\Omega), \quad \text{where } \beta = r^* > N. \quad (9)$$

Recall that $D = \operatorname{div}(\mathbf{u})$. We will consider the following assumptions:

- Assumptions on Ω :

$$\Omega \subset \mathbb{R}^N \quad \text{is open, bounded and it has a Lipschitz boundary } \partial\Omega. \quad (10)$$

- Assumptions on the flow data² \mathbf{u}, F, D, ρ and ν :

$$\nu \geq 0, \quad (11)$$

$$F, \rho : \Omega \rightarrow \mathbb{R}^+, \quad \rho, \rho^{-1} \in L^\infty(\Omega), \quad F \in L^r(\Omega), \quad (12)$$

$$\mathbf{u} \in (L^3(\Omega))^N, D \in L^r(\Omega), \operatorname{div}(\rho\mathbf{u}) = 0, \quad (13)$$

$$\mathbf{u} \cdot \mathbf{n} = 0 \text{ on } \partial\Omega. \quad (14)$$

- Assumptions on the turbulent quantities on the boundary:

$$\begin{aligned} a, b &\in H^{1/2}(\partial\Omega) \cap L^\infty(\partial\Omega), \\ a(x), b(x) &\geq \delta > 0 \text{ a.e. } x \in \partial\Omega, \end{aligned} \quad (15)$$

where $\delta > 0$ is a fixed number.

- Assumptions on the model coefficients:

$$C_i : \Omega \rightarrow \mathbb{R}^+, \quad C_5, C_8 \in L^r(\Omega), \text{ for } i \neq 5, 8 : C_i \in L^\infty(\Omega) \quad (16)$$

$$C_{\text{ind}} : \Omega \times (\mathbb{R}^+)^2 \rightarrow \mathbb{R}^+ \text{ are Caratheodory} \quad (17)$$

$$C_{\text{ind}}(x, v, w) \geq \alpha_{\text{ind}} > 0 \quad \forall v, w \in \mathbb{R}^+ \text{ and for a.a. } x \in \Omega \quad (18)$$

$$\forall v, w \in \mathbb{R}^+, \quad x \rightarrow C_{\text{ind}}(x, v, w) \in L^\infty(\Omega), \quad (19)$$

where in (17)-(19) C_{ind} means C_θ or C_φ . The assumption (17) signifies that $\forall v, w \in \mathbb{R}^+, x \rightarrow C_{\text{ind}}(x, u, v)$ is measurable and for a.a. $x \in \Omega : (v, w) \rightarrow C_{\text{ind}}(x, v, w)$ is continuous. This ensures that $C_{\text{ind}}(x, \theta, \varphi)$ is measurable when θ and φ are measurable. The condition (18) means that C_{ind} is uniformly positive, whereas (19) tells that $C_{\text{ind}}(x, \theta, \varphi)$ remains bounded if θ and φ are bounded.

We will study problem (Q) under the **main assumption**:

$$(A_0) : \quad (8) - (19) \text{ are satisfied.}$$

Note that in the main situation (A_0) the assumption (11) made for ν allows the possibility $\nu \equiv 0$. In other words the molecular viscosity ν can be neglected in the model. This is often chosen in practice because the eddy viscosities $\frac{C_\theta}{\theta\varphi}$ and $\frac{C_\varphi}{\theta\varphi}$ are dominant in the physical situations (see [23, 24]). Remark also that the coefficients C_i are allowed to depend on x , and the viscosity parameters C_θ, C_φ may depend on x, θ, φ .

²When $N = 2$ one assumption in (13) can be relaxed: $\mathbf{u} \in (L^p(\Omega))^2$ with $p > 2$ (instead of $p = 3$) is sufficient, but this would not improve any result significantly.

For a given function $f : \Omega \rightarrow \mathbb{R}$, we shall use the notations f^+ and f^- to represent the positive and negative parts of f , that is:

$$f = f^+ + f^-, \quad f^+(x) = \max(f(x), 0) \geq 0, \quad f^-(x) = \min(f(x), 0) \leq 0.$$

We will also consider some assumption of low compressibility of the form:

$$\|D^+\|_{L^r(\Omega)} \leq \tau, \quad (20)$$

for some $\tau > 0$ that will be precised.

This last kind of condition seems to be necessary (see the Appendix) in order to obtain a weak solution for (Q) in the three dimensional case, whereas when $N = 2$ we shall use some particularities of the situation to obtain a weak solution without any assumption of low compressibility. Nevertheless in this case we will assume that in addition to (A_0) the following condition is satisfied:

$$(A_1) : \quad \frac{C_4^2}{C_3} \in L^r, \quad c_6 = 0.014, \quad c_7 = 0.104, \quad c_8 = 0.84.$$

In this last condition the values for c_6, c_7 and c_8 are in fact their classical constant values (see (74) in the Appendix)

In the sequel we denote by *DATA* some quantity depending only on the data under the assumption (A_0) , i.e.:

$$DATA = Const(\Omega, a, b, \|\mathbf{u}\|_{(L^3)^N}, \alpha_\theta, \alpha_\varphi, (\|C_j\|_{L^r})_{j=5,8}, \|F\|_{L^r}, (\|C_i\|_{L^\infty})_{i=3,4,6,7}). \quad (21)$$

Note that *DATA* does not depends on D and ν .

The exact form of the dependency (i.e. the function *Const*) is allowed to change from one part of the text to another.

2.2 Main results

We shall establish two theorems of existence. The first one applies if $N=2$ or 3 and the second one is limited to the case $N=2$.

Theorem 1 *Assume that (A_0) holds. Then there exists $\tau > 0$ such that if $\|D^+\|_{L^r(\Omega)} \leq \tau$ then problem (Q) admits at least one weak solution (θ, φ) in the class \mathcal{S} .*

Theorem 2 *Assume that $N = 2$ and $(A_0), (A_1)$ hold. Then problem (Q) admits at least one weak solution (θ, φ) in the class \mathcal{S} .*

In all the situations we have the following regularity result:

Theorem 3 *Let (θ, φ) be a weak solution of (Q) in the class \mathcal{S} and assume that (A_0) is satisfied. We have:*

i) *If $\mathbf{u} \in (L^\infty(\Omega))^N$ and a, b are Hölder continuous, then $(\theta, \varphi) \in (\mathcal{C}^{0,\alpha}(\overline{\Omega}))^2$, for some $\alpha > 0$*

ii) *Assume that in addition, the following conditions are satisfied:*

$$\partial\Omega, a \text{ and } b \text{ are of class } \mathcal{C}^{2,\alpha}, \quad \rho\mathbf{u} \in (\mathcal{C}^{1,\alpha}(\overline{\Omega}))^N$$

$$F \in \mathcal{C}^{0,\alpha}(\overline{\Omega}), \quad \forall i = 3, \dots, 8 : C_i \in \mathcal{C}^{0,\alpha}(\overline{\Omega})$$

$$C_{ind}(x, v, w) \in \mathcal{C}^{1,\alpha}(\overline{\Omega} \times (\mathbb{R}^+)^2)$$

Then $(\theta, \varphi) \in (\mathcal{C}^{2,\alpha}(\overline{\Omega}))^2$ and it is a classical solution of (Q).

2.3 Discussion on the results

Compared to previous works (see [16, 15, 23, 27]) our basic assumption (A_0) made in theorem 1 is very general. In particular we do not artificially cancel the singularity in the model and we only assume weak regularity on the data. For instance the basic assumption we made on the mean flow is $\mathbf{u} \in (L^3(\Omega))^N$ and $\operatorname{div} \mathbf{u} \in L^r(\Omega)$, with some $r > 3/2$, whereas in the previous works it was assumed $\mathbf{u} \in (L^\infty(\Omega))^N$ and $\operatorname{div} \mathbf{u} \in L^\infty(\Omega)$. Our condition is more interesting from a practical point of view because it is satisfied when \mathbf{u} is a weak solution of the Navier Stokes equations. Hence our work could be used for a future analysis of the full coupled system Navier-Stokes plus (Q).

From a mathematical point of view the problem we study is a non linear, degenerate and singular elliptic system. Several complications arise for its analysis. In particular the balance between the increase/decrease of the source terms (i.e. the functions C_i, D and F appearing in the second members of (Q)) and the possible explosion/vanishing of the viscosities is difficult to control. The strategy followed here is to first carefully study some elliptic scalar problem (possibly degenerate and singular) of the form:

$$(S) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \zeta - \operatorname{div} \left((\nu + \frac{\kappa}{\zeta}) \nabla \zeta \right) = g(x, \zeta) & \text{in } \Omega \\ \zeta = \zeta_0 & \text{on } \partial\Omega, \end{cases}$$

By using and developing some techniques du to Stampacchia we are able to establish existence, positivity and regularity results for problem (S). These results (see Proposition 9) which also have an independent mathematical interest, are the key ingredients for proving theorem 1. Under the additional assumption $\mathbf{u} \in (L^\infty(\Omega))^N$ we give a Hölder continuity result for k and θ . Moreover, we establish an existence result for a classical solution of (Q) under some smoothness assumptions on the data.

2.4 Organization of the paper

- In section 3 we shall recall some results concerning: the truncature at a fixed level and the Stampacchia's estimates. This last technique takes an important role in our analysis, moreover we shall need a precise control of the estimates. Hence we shall present it with some details and developments.
- In section 4 we introduce a sequence (Q_n) of problems which approximate (Q). For n fixed (Q_n) is a PDE system of two scalar equations of the form (S): one equation for the unknown θ_{n+1} and one for φ_{n+1} . The point is that the unknowns θ_{n+1} and φ_{n+1} are only weakly coupled. The coupling of the two equations is essentially realized through the quantities θ_n and φ_n calculated at the previous step. Hence we shall firstly study carefully the problem (S). The major tool used here are the Stampacchia's estimates. We next prove that (Q_n) is well posed. Hence we obtain an approximate sequence of solutions (θ_n, φ_n) for problem (Q). Moreover, we prove that θ_n and φ_n are uniformly bounded from above and below, which are the key estimates.
- In section 5 we use the uniform bounds established in section 4, in order to extract a converging subsequence from (θ_n, φ_n) . We then prove that the limit (θ, φ) is a weak solution of (Q) in the class \mathcal{S} .

Under the additional assumption $\mathbf{u} \in (L^\infty(\Omega))^3$ we are able to use the De Giorgi-Nash Theorem and we obtain an Hölder continuity result for θ, φ . By assuming in addition some smoothness properties for the data we can iterate the Schauder estimates and prove Theorem 3.

- In section 6 (Appendix) we present the derivation of the $\theta - \varphi$ model from the $k - \varepsilon$ one. Moreover we justify that the choice $\varphi = \frac{\varepsilon^2}{k^3}$ is valid even in the compressible situation. The justification uses in particular a property of positivity of the function F . We also discuss briefly the necessity of the low compressibility assumption when $N = 3$. Finally we recall a generalized version of the chain rule for $G(u)$ where G is a Lipschitz function and u a Sobolev one.

3 Mathematical background

In this section we shall recall some results concerning: the truncature at a fixed level and the Stampacchia's estimates. This last technique takes an important role in our analysis, moreover we shall need a precise control of the estimates. Hence we shall present here the technique with some details and developments. As in the rest on the paper we denote by $\Omega \subset \mathbb{R}^N$ a bounded open Lipschitz domain. These properties for Ω are always implicitly assumed if they are not precised.

3.1 Truncatures and related properties

The technique of Stampacchia is based on the use of special test functions which are constructed by using some truncatures. We shall recall some basic properties of the truncatures used in the paper. An important tool is the generalized chain rule (see Theorem 13 in the Annexe).

Let $l > 0$ we denote by T_l the truncature function $T_l : \mathbb{R} \rightarrow \mathbb{R}$ defined by

$$T_l(s) = \begin{cases} l & \text{if } s > l \\ s & \text{if } -l \leq s \leq l \\ -l & \text{if } s < -l \end{cases} \quad (22)$$

Let $v \in H^1(\Omega)$. By applying Theorem 13 we see that $T_l(v) \in H^1(\Omega) \cap L^\infty(\Omega)$. Moreover, if we denote by $\Omega_{v,l}$ the set $\{x \in \Omega \text{ s.t. } |v(x)| \leq l\}$ then we have:

$$\nabla T_l(v) = \begin{cases} \nabla v & \text{in } \Omega_{v,l} \\ 0 & \text{elsewhere} \end{cases}$$

Note that $T_l(\cdot)$ truncates both the positive and the negative large values. In some cases we need only to truncate in one side. For this reason we introduce the semi-truncatures $T_{l,+}$ and $T_{l,-}$ defined by:

$$T_{l,+}(t) = \begin{cases} l & \text{if } t > l \\ t & \text{elsewhere} \end{cases} \quad T_{l,-}(t) = \begin{cases} -l & \text{if } t < -l \\ t & \text{elsewhere} \end{cases}$$

We then have the decomposition: $T_l = T_{l,+} \circ T_{l,-} = T_{l,-} \circ T_{l,+}$.

For given $v_l \in H^1(\Omega)$ and $s > 0$, we shall also consider

$$\psi_s(v_l) = v_l - T_s(v_l) = \begin{cases} v_l - s & \text{in } A_{s,l}^+ \\ v_l + s & \text{in } A_{s,l}^- \\ 0 & \text{elsewhere} \end{cases} \quad (23)$$

where we have used the notations:

$$A_{s,l}^+ := \{v_l \geq s\}, \quad A_{s,l}^- := \{v_l \leq -s\} \text{ and } A_{s,l} = A_{s,l}^+ \cup A_{s,l}^-. \quad (24)$$

Let also $\psi_{s,\pm}$ be the functions defined above (23) while replacing T_s by $T_{s,+}$ or by $T_{s,-}$. It is easy to verify that $\psi_{s,+}$ (resp. $\psi_{s,-}$) is in fact the positive (resp. the negative) part of ψ_s . In other words, we have :

$$\begin{aligned} \psi_{s,+}(v_l) &= \psi_s^+(v_l) = (v_l - s)\mathbf{1}_{A_{s,l}^+} \geq 0, & \psi_{s,-}(v_l) &= \psi_s^-(v_l) = (v_l + s)\mathbf{1}_{A_{s,l}^-} \leq 0. \\ \psi_s &= \psi_s^+ + \psi_s^-. \end{aligned}$$

The function ψ_s has the following properties:

Lemma 4

- i) $\psi_s \in H^1(\Omega)$ and $\nabla \psi_s = \mathbf{1}_{A_{s,l}} \nabla v_l$
- ii) if $\gamma v_l \in L^\infty(\partial\Omega)$ then for all $s > \|\gamma v_l\|_{L^\infty(\partial\Omega)}$ we have $\psi_s \in H_0^1(\Omega)$,

where $\gamma : H^1(\Omega) \rightarrow H^{1/2}(\partial\Omega)$ denotes the trace function.

Proof

Point i) is a direct consequence of Theorem 13. Property ii) is proved in [14], lemma 3.3 p.53. \square

3.2 The Stampacchia estimates

The Stampacchia estimates is a general method which allows one to obtain an L^∞ -estimate for the weak solution of a large class of elliptic PDEs of the second order. The L^∞ -estimate presented in the original work [25] depends on various quantities related to the PDE problem studied, but the exact dependency is not established. In our analysis we need a precise control of the L^∞ -estimates with respect to some quantities (in particular with respect to the diffusion coefficient of the PDEs). Hence in the following we take over and detail the technique in order to obtain a more precise L^∞ -estimates.

The Stampacchia estimates are established by using the test functions ψ_s (or ψ_s^\pm) defined previously, where in this case v (resp. v_l) is a weak solution of the problem (resp. the sequence of problems) considered.

For technical reasons we need a classical result concerning some relationship between L^r functions and linear form on Sobolev spaces:

Lemma 5 *Let $1 \leq r < \infty$ and $E \in L^r(\Omega)$. Then*

$$E \in W^{-1,\beta}(\Omega), \text{ with } \beta = r^* = \frac{Nr}{N-r}, \quad (25)$$

and there exists $\tilde{E} \in (L^\beta(\Omega))^3$ such that:

$$\int_{\Omega} E\varphi = \int_{\Omega} \tilde{E} \cdot \nabla \varphi \quad \forall \varphi \in \mathcal{D}(\Omega), \quad \|\tilde{E}\|_{(L^\beta(\Omega))^3} \leq \|E\|_{L^r(\Omega)}. \quad (26)$$

Moreover we have:

$$r > \frac{N}{2} \Rightarrow \beta > N. \quad (27)$$

Proof

Property (25) is easy to prove by using the Sobolev injection Theorem together with the Hölder inequality: $\varphi \rightarrow \int_{\Omega} E\varphi$ is a linear form on $W_0^{1,p}$ if $p^* := \frac{Np}{N-p} \geq \frac{r}{r-1}$. This last condition holds for $p = P := \frac{Nr}{r(N+1)-N}$. Hence $P' = p/(p-1) = \frac{Nr}{N-r} = r^*$.

We next obtain (26) by using a classical result (see [2] Proposition IX.20).

Finally, if we assume that $r > \frac{N}{2}$ then $\beta = r^* > \frac{N^2/2}{N-N/2} = N$. \square

The Stampacchia technique works in two steps: the first one is dependent of the problem (or the sequence of problems) studied and the second one is independent of it. Here the purpose is to present the key ingredients of these two steps. Because the first one is dependent on the problem studied we cannot present it here in its entirety, but we will consider a simple problem which contains the main technical points (in fact this introductory presentation will be useful to treat a more complicated class of problems in Section 4). Let $(v_l) \subset H_0^1(\Omega)$ be a sequence of functions satisfying:

$$\int_{\Omega} \nu_l \nabla v_l \nabla \varphi = \int_{\Omega} g_l \varphi, \quad \forall \varphi \in H_0^1, \quad (28)$$

where (ν_l) is a given sequence of strictly positive bounded functions and $(g_l) \subset L^r(\Omega)$, with $r > \frac{N}{2}$. Let also m_l, M_l denote the bounds from above and below for ν_l , that is:

$$0 < m_l \leq \nu_l \leq M_l < \infty, \quad \text{a.e. in } \Omega.$$

-Step 1

By testing (28) with $\psi_s = \psi_s(v_l)$, we obtain:

$$m_l \int_{\Omega} |\nabla \psi_s|^2 = \int_{\Omega} g_l \psi_s \stackrel{\text{by lemma 5}}{=} \int_{\Omega} E_l \nabla \psi_s, \quad (29)$$

with $E_l \in L^\beta$, $\|E_l\|_{L^\beta} \leq \|g_l\|_{L^r}$ for some $\beta = \beta(r) > N$.

Recall also that $\text{Supp } \psi_s(v_l) \subset A_{l,s}$. Hence by using the Hölder inequality we obtain:

$$\begin{aligned} \int_{\Omega} E_l \nabla \psi_s &\leq \|\nabla \psi_s\|_{L^2} \left(\int_{A_{l,s}} |E_l|^2 \right)^{\frac{1}{2}} \leq \|\nabla \psi_s\|_{L^2} \|E_l\|_{L^\beta} \left(\int_{A_{l,s}} 1 \right)^{\frac{1}{2(\frac{\beta}{2})'}} \\ &\leq \|\nabla \psi_s\|_{L^2} \|E_l\|_{L^\beta} |A_{l,s}|^{\frac{\beta-2}{2\beta}}. \end{aligned}$$

Consequently (29) leads to:

$$\|\psi_s\|_{H_0^1(\Omega)^2}^2 \leq C |A_{l,s}|^\Phi, \quad \text{whith } C > 0, \Phi > \frac{2}{2^*} = \frac{N-2}{N}. \quad (30)$$

This is the key estimate needed to pass at the second step which is independent of the problem studied.

Note that with the particular sequence of problems (28) chosen here, the constants C and Φ are:

$$C = \frac{\|g_l\|_{L^r}^2}{m_l^2}, \quad \Phi = \frac{\beta-2}{\beta}. \quad (31)$$

Hence Φ does not depend on l, s . Moreover if we assume that (g_l) is uniformly bounded in the L^r -norm and that (m_l) is uniformly bounded from above by a strictly positive constant, then C is also independent on l, s .

This is an important point because as we will see it hereafter, an estimate (30) with C and Φ independent on l, s leads to a uniform L^∞ bound for (v_l) .

-*Step 2*

Assume that we have obtained (30). We can obtain an L^∞ -estimate for v_l as follows.

Let $2^* = \frac{2N}{N-2}$ be the Sobolev exponent associated to 2 in dimension N . By using the Poincaré-Sobolev inequality we have:

$$\left(\int_{A_{l,s}} |\psi_s|^{2^*} \right)^{2/2^*} \leq C_1 \|\psi_s\|_{H_0^1(\Omega)}^2, \quad C_1 = C_1(|\Omega|) \quad (32)$$

Let now $t > s$. It is clear that $A_{l,t} \subset A_{l,s}$ and consequently

$$\left(\int_{A_{l,s}} |\psi_s|^{2^*} \right)^{2/2^*} \geq \left(\int_{A_{l,t}} |\psi_s|^{2^*} \right)^{2/2^*} \geq \left(\int_{A_{l,t}} |t-s|^{2^*} \right)^{2/2^*} \geq |t-s|^2 |A_{l,t}|^{2/2^*} \quad (33)$$

We set

$$\chi(s) := |A_{l,s}|, \quad \forall s \geq 0. \quad (34)$$

For fixed l , χ is a decreasing function, and from the estimates (30)-(33), we obtain

$$\chi(t) \leq C_2 |\chi(s)|^{2^*\Phi/2} (t-s)^{-2^*} \quad \forall t > s \geq 0, \quad C_2 = (C_1 C)^{2^*/2}.$$

Recall that we have assumed in (30) that $\Phi > \frac{2}{2^*}$. Hence $2^*\Phi/2 > 1$ and by using Lemma 4.1 in [25] we obtain:

$$\chi(d) = 0, \quad d = 2^{2^*\Phi/2} C_2^{1/2^*} |\Omega|^{\frac{\Phi}{2} - \frac{1}{2^*}} < \infty. \quad (35)$$

This property tells exactly that:

$$\|v_l\|_{L^\infty(\Omega)} \leq d, \quad d = C_3(\Phi, N, |\Omega|) \sqrt{C}. \quad (36)$$

In particular d does not depend on l if the constants C and Φ appearing in (30) are independent of l . For instance with the particular sequence of problems (28) the constants C and Φ are given by (31), and if we assume that $\|g_l\|_{L^r} \leq K$, $m_l \geq m > 0$ we obtain:

$$d = \tilde{C} \frac{K}{m}, \quad \tilde{C} = \tilde{C}(\Omega, N, r).$$

Remark 6

- If you are only interested in obtaining an uniform majoration or minoration for v_l then instead of (30) it suffices to have

$$\|\psi_s^\pm\|_{H_0^1(\Omega)^2}^2 \leq C^\pm |A_{l,s}^\pm|^{\Phi^\pm}, \quad \text{with } C^\pm > 0, \Phi^\pm > \frac{N-2}{N}. \quad (37)$$

In fact in this case we consider $\chi_\pm(s) := |A_{l,s}^\pm|$ instead of (34). This function is decreasing and we obtain again (35). But now this property tells:

$$\pm v_l(x) \leq d^\pm \text{ a.e. in } \Omega, \quad d^\pm = \tilde{C}^\pm(\Phi^\pm, N, |\Omega|) \sqrt{C^\pm}.$$

- Let again (v_l) be a sequence of functions satisfying (28) and assume that $m_l \geq m > 0$. Then we have:

$$g_l \leq h_l \text{ and } H := \sup_l \|h_l\|_{L^r} < \infty \Rightarrow v_l \leq d_1 \text{ a.e. in } \Omega, \quad (38)$$

$$g_l \geq k_l \text{ and } K := \sup_l \|k_l\|_{L^r} < \infty \Rightarrow v_l \geq -d_2 \text{ a.e. in } \Omega, \quad (39)$$

$$\text{with } d_1, d_2 > 0, \quad d_1 = \tilde{C} \frac{H}{m}, \quad d_2 = \tilde{C} \frac{K}{m}, \quad \tilde{C} = \tilde{C}(\Phi, N, |\Omega|).$$

The proof of (38) and (39) are obtained by taking over the first step of the technique of Stampacchia, but by using ψ_s^\pm as test function instead of ψ_s .

In fact, the function ψ_s^+ is positive. Hence instead of (29) we have:

$$m \int_{\Omega} |\nabla \psi_s^+|^2 = \int_{\Omega} g_l \psi_s^+ \leq \int_{\Omega} h_l \psi_s^+ = \int_{\Omega} E_{l,+} \nabla \psi_s^+,$$

with $E_{l,+} \in L^\beta$, $\|E_{l,+}\|_{L^\beta} \leq \|h_l\|_{L^r} \leq H$. Consequently, in this case we obtain (37) for the function ψ_s^+ , with $C^+ = \frac{H^2}{m^2}$, and the uniform majoration (38) follows.

The relation (39) can be proven by using ψ_s^- as test function in (28). In fact we remark now that ψ_s^- is negative. Hence instead of (29) we obtain:

$$m \int_{\Omega} |\nabla \psi_s^-|^2 = \int_{\Omega} g_l \psi_s^- \leq \int_{\Omega} k_l \psi_s^- = \int_{\Omega} E_{l,-} \nabla \psi_s^-,$$

with $E_{l,-} \in L^\beta$, $\|E_{l,-}\|_{L^\beta} \leq \|k_l\|_{L^r} \leq K$. Consequently we now obtain (37) for the function ψ_s^- , with $C^- = \frac{K^2}{m^2}$. This implies: $-v_l \leq d_2$ and consequently $v_l \geq -d_2$.

4 Approximate sequence and estimates

Let $g \in H^{1/2}(\partial\Omega)$, we denote by $\mathcal{R}g$ its harmonic lifting, that is:

$$\mathcal{R}g \in H^1(\Omega), \mathcal{R}g = g \text{ on } \partial\Omega \text{ and } \Delta \mathcal{R}g = 0 \text{ on } \Omega.$$

We define the functions θ_0 and φ_0 by the formula:

$$\theta_0 := \mathcal{R}a \quad \varphi_0 := \mathcal{R}b. \quad (40)$$

Hence, by the using the maximum principle (see [2] p.189 and [11]) together with the condition (15) we obtain:

$$0 < \delta \leq \theta_0 \leq \|a\|_{L^\infty(\partial\Omega)}, \quad \delta \leq \varphi_0 \leq \|b\|_{L^\infty(\partial\Omega)}. \quad (41)$$

Let now $n \geq 0$, (θ_n, φ_n) be given and

$$C_\theta^{(n)}(\cdot) := C_\theta(\cdot, \theta_n(\cdot), \varphi_n(\cdot)), \quad C_\varphi^{(n)}(\cdot) := C_\varphi(\cdot, \theta_n(\cdot), \varphi_n(\cdot)).$$

In order to construct an approximate solution $(\theta_{n+1}, \varphi_{n+1})$ for problem (Q), we introduce the following system:

$$(Q_n) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \theta_{n+1} - \operatorname{div} \left(\left(\nu + \frac{C_\theta^{(n)}}{\theta_{n+1} \varphi_n} \right) \nabla \theta_{n+1} \right) = g_\theta^{(n)} & \text{in } \Omega, \\ \rho \mathbf{u} \cdot \nabla \varphi_{n+1} - \operatorname{div} \left(\left(\nu + \frac{C_\varphi^{(n)}}{\theta_{n+1} \varphi_{n+1}} \right) \nabla \varphi_{n+1} \right) = g_\varphi^{(n)} & \text{in } \Omega, \\ \theta_{n+1} = a \text{ and } \varphi_{n+1} = b & \text{on } \partial\Omega, \end{cases}$$

where we used the notations:

$$g_\theta^{(n)} := C_5 - C_3 F \theta_{n+1}^2 + C_4 \theta_{n+1} D, \quad g_\varphi^{(n)} := -\varphi_n (C_8 \theta_{n+1}^{-1} + C_6 \theta_{n+1} F - C_7 D).$$

For $n \in \mathbb{N}$ we denote by (H_n) the following condition:

$$(H_n) \quad \begin{cases} \theta_n, \varphi_n, \theta_n^{-1}, \varphi_n^{-1} \in L^\infty(\Omega), \\ \theta_n, \varphi_n \geq 0. \end{cases}$$

Let φ_{\max} be a fixed real number such that:

$$\varphi_{\max} > \|b\|_{L^\infty(\partial\Omega)}. \quad (42)$$

We shall also consider the condition: $(K_n) := (H_n) + (\varphi_n \leq \varphi_{\max})$.

Note that (41) shows that the condition (K_n) is satisfied for $n = 0$. We will prove in the sequel that under condition (K_n) we can obtain a weak solution $(\theta_{n+1}, \varphi_{n+1})$ for problem (Q_n) , with moreover $(\theta_{n+1}, \varphi_{n+1})$ satisfying the condition (K_{n+1}) . This last property ensures the right definition of an approximate sequence. More precisely, we have:

Proposition 7 *Let $n \in \mathbb{N}$ be given and assume that (A_0) is satisfied. Let also (θ_n, φ_n) be given and satisfying condition (K_n) . There exists $\tau > 0$ depending only on DATA such that if $\|D^+\|_{L^r(\Omega)} \leq \tau$ then problem (Q_n) admits at least one weak solution $(\theta_{n+1}, \varphi_{n+1}) \subset (H^1(\Omega) \times L^\infty(\Omega))^2$.*

Moreover $(\theta_{n+1}, \varphi_{n+1})$ satisfies condition (K_{n+1}) and the estimates:

$$0 < \theta_{\min} \leq \theta_{n+1} \leq \theta_{\max}, \quad (43)$$

$$0 < \varphi_{\min} \leq \varphi_{n+1} \leq \varphi_{\max}, \quad (44)$$

where φ_{\max} was fixed in (42) and $\varphi_{\min}, \theta_{\min}, \theta_{\max}$ are positive numbers depending on DATA, but not on n .

Remark 8 *Proposition 7 is the key result that will be used later on to prove Theorem 1, whereas for Theorem 2 we shall establish and use a more simple version of this proposition (see Subsection 5.2).*

In order to prove the proposition we establish intermediate results.

4.1 Auxiliary results

Let $n \in \mathbb{N}$ and (θ_n, φ_n) be given and satisfying (K_n) . We want to obtain $(\theta_{n+1}, \varphi_{n+1})$ by solving (Q_n) and in order to iterate the algorithm we also want that $(\theta_{n+1}, \varphi_{n+1})$ satisfies (K_{n+1}) .

Remark that the system (Q_n) is composed of two coupled scalar elliptic equation in divergence form, with a possible singular and degenerate structure. Hence the goal of this subsection is to study this last kind of scalar problem.

In order to do this, we first introduce a weight $\kappa : \mathbb{R}^+ \rightarrow \mathbb{R}^+$ which is assumed to be measurable and satisfying:

$$0 < \kappa_0 \leq \kappa \leq \kappa_1 \quad \text{a.e. in } \Omega, \quad (45)$$

where κ_0 and κ_1 are two given reals.

Let also $g : \Omega \times \mathbb{R}^+ \rightarrow \mathbb{R}$ be a Caratheodory function (i.e. $\forall u \in \mathbb{R}^+ : x \rightarrow g(x, u)$ is measurable, and for a.a. $x \in \Omega : u \rightarrow g(x, u)$ is continuous).

Let us consider the scalar problem:

$$(S) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \zeta - \operatorname{div} \left(\left(\nu + \frac{\kappa}{\zeta} \right) \nabla \zeta \right) = g(x, \zeta) & \text{in } \Omega \\ \zeta = \zeta_0 & \text{on } \partial\Omega, \end{cases}$$

where $\zeta_0 \in H^{1/2}(\partial\Omega) \cap L^\infty(\partial\Omega)$, $\zeta_0 \geq \delta > 0$ a.e. in $\partial\Omega$, is given. We always assume that $\rho, \mathbf{u}, \nu, \Omega$ which appear in (S) satisfy their corresponding conditions in (A_0) .

Recall that we allow $\nu \equiv 0$ in (A_0) . Hence problem (S) may degenerate (i.e. the viscosity $\nu + \frac{\kappa}{\zeta}$ may vanish) when $\zeta \rightarrow \infty$. Moreover (S) is singular (i.e. the viscosity tends to infinity) when $\zeta \rightarrow 0$.

We want now to find sufficient additional conditions for g that guarantee the existence of a bounded positive weak solution for problem (S). Hence we shall consider:

$$\gamma(x) := \sup_{u \in [0,1]} |g^-(x, u)| \in L^r(\Omega), \quad (46)$$

$$g^+(x, u) \leq \gamma_1(x) + \gamma_2(x)h(u), \quad (47)$$

where $\gamma_1, \gamma_2 \in L^r(\Omega)$ and $h : \mathbb{R}^+ \rightarrow \mathbb{R}^+$ is continuous. In fact, more than establishing only the existence of a bounded positive solution for (S), we are interested in obtaining some uniform (with respect to κ_1) bounds from above and below and some regularity results. We have:

Proposition 9

i) Let κ satisfying (45) and $g : \Omega \times \mathbb{R}^+ \rightarrow \mathbb{R}$ be a Caratheodory function satisfying (46), (47). There exists a real $\tau > 0$ depending on κ_0, h such that if $\|\gamma_2\|_{L^r(\Omega)} \leq \tau$ then there exists a weak solution $\zeta \in H^1(\Omega) \cap L^\infty(\Omega)$ for problem (S). Moreover we have:

$$0 < \zeta_{min} \leq \zeta \leq \zeta_{max}, \quad \zeta_{min} = e^{-\frac{C}{\kappa_0}}, \zeta_{max} = e^{\frac{C}{\kappa_0}}, \quad (48)$$

where C depends only on $\gamma, \gamma_1, \Omega, r, N, \zeta_0$. In particular ζ_{min} and ζ_{max} are independent of κ_1 .

In addition, the following extended ³ maximum principle holds:

$$\|\zeta\|_{L^\infty(\Omega)} \leq \|\zeta_0\|_{L^\infty(\partial\Omega)} + \tilde{C} \frac{\|g^+\|_{L^r}}{\kappa_0} \|\zeta\|_{L^\infty(\Omega)}, \quad \tilde{C} = \tilde{C}(\Omega, N, r, \gamma_1). \quad (49)$$

³When $g^+ \equiv 0$ it is a maximum principle

- ii) Assume that in addition $\mathbf{u} \in (L^\infty(\Omega))^N$ and ζ_0 is Hölder continuous. Then $\zeta \in \mathcal{C}^{0,\alpha}(\bar{\Omega})$ for some $0 < \alpha < 1$. Moreover if $\partial\Omega$ is of class $\mathcal{C}^{2,\alpha}$, $g \in \mathcal{C}^{0,\alpha}(\bar{\Omega} \times \mathbb{R}^+)$, $\rho\mathbf{u} \in (\mathcal{C}^{1,\alpha}(\bar{\Omega}))^N$, $\kappa \in \mathcal{C}^{1,\alpha}(\bar{\Omega})$ and $\zeta_0 \in \mathcal{C}^{2,\alpha}(\partial\Omega)$, then $\zeta \in \mathcal{C}^{2,\alpha}(\bar{\Omega})$ and it is a classical solution of (S).

Before proving Proposition 9 we establish an intermediate result. In a first step we consider the change of variable $v = \ln \zeta$ in (S), and for $l \in \mathbb{N}$ we introduce a truncated version (S_l) of the system obtained:

$$(S_l) \quad \begin{cases} \rho\mathbf{u}\nabla e^{T_l(v)} - \operatorname{div}((\nu e^{T_l(v)} + \kappa)\nabla v) = g(x, e^{T_l(v)}) & \text{in } \Omega \\ v = \ln \zeta_0 & \text{on } \partial\Omega \end{cases}$$

We then establish:

Lemma 10

- i) Let κ satisfying (45) and $g : \Omega \times \mathbb{R}^+ \rightarrow \mathbb{R}$ be a Caratheodory function satisfying (46), (47). Then, for any $l \in \mathbb{N}$ there exists a weak solution $v = v_l \in H^1(\Omega) \cap L^\infty$ for the problem (S_l) .
- ii) Let (v_l) be the sequence given in i) and $l \geq 1$ be a fixed integer. Then there exists $\tau_l > 0$ such that if the function γ_2 in (46) satisfies $\|\gamma_2\|_{L^r} \leq \tau_l$ then we have:

$$\|v_l\|_{L^\infty(\Omega)} \leq C \frac{\|\gamma + \gamma_1\|_{L^r}}{\kappa_0}, \quad C = C(\text{DATA}, \zeta_0). \quad (50)$$

In particular C is independent of κ, ν and l .

Proof

- i) By using the divergence formula we obtain, for all $w \in H^1(\Omega)$:

$$\begin{aligned} \int_{\Omega} \rho\mathbf{u}\nabla e^{T_l(v)} \cdot w &= - \int_{\Omega} e^{T_l(v)} \operatorname{div}(\rho\mathbf{u}w) + \underbrace{\int_{\partial\Omega} e^{T_l(v)} \rho w \mathbf{u} \cdot \mathbf{n} d\sigma}_{=0} \\ &= - \int_{\Omega} e^{T_l(v)} \rho\mathbf{u} \cdot \nabla w - \int_{\Omega} e^{T_l(v)} \underbrace{\operatorname{div}(\rho\mathbf{u})}_{=0} w. \end{aligned}$$

Let $v_0 := \ln(\mathcal{R}\zeta_0)$ and consider the change of variable $\tilde{v} := v - v_0$. Then problem (S_l) is equivalent to find $\tilde{v} \in H_0^1(\Omega)$ such that:

$$- \operatorname{div} \sigma(x, \tilde{v}, \nabla \tilde{v}) = f(x, \tilde{v}) \quad \text{in } \mathcal{D}'(\Omega), \quad (51)$$

where $\sigma : \Omega \times \mathbb{R} \times \mathbb{R}^N \rightarrow \mathbb{R}^N$ and $f : \Omega \times \mathbb{R} \rightarrow \mathbb{R}$ are defined by:

$$\begin{aligned} \sigma(x, w, \mathbf{G}) &= (\nu e^{T_l(w+v_0)} + \kappa(x))\mathbf{G} - \rho(x)\mathbf{u}(x)e^{T_l(w+v_0(x))}, \\ f(x, w) &= g(x, e^{T_l(w+v_0(x))}). \end{aligned}$$

We now remark that (51) is a quasilinear equation in divergence form. Moreover, it is easy to see that f and σ satisfy the classical growth

assumptions and σ satisfies also the classical coercivity condition. Note that:

$$\begin{aligned} \langle \sigma(x, w, \mathbf{G}) - \sigma(x, w, \mathbf{G}'), \mathbf{G} - \mathbf{G}' \rangle &= (\nu e^{T_l(w+v_0)} + \kappa) |\mathbf{G} - \mathbf{G}'|^2 \\ &\geq \kappa_0 |\mathbf{G} - \mathbf{G}'|^2. \end{aligned}$$

Hence σ is strictly monotone in the third variable. We then conclude (see for instance [1] Theorem 1.5, or [14] Theorem 8.8 page 311) that there exists a weak solution $\tilde{v} \in H_0^1(\Omega)$ for (51).

Consequently $v_l := \tilde{v} + v_0$ is a weak solution for (S_l) , that is $\forall w \in H_0^1(\Omega)$:

$$\int_{\Omega} (\nu e^{T_l(v_l)} + \kappa) \nabla v_l \nabla w + \int_{\Omega} \rho \mathbf{u} \nabla e^{T_l(v_l)} w = \int_{\Omega} g(x, e^{T_l(v_l)}) w. \quad (52)$$

By applying Theorem 4.2 page 108 in [25] we obtain: $v_l \in L^\infty(\Omega)$.

ii)

$$\text{Let } \tau_l := \frac{1}{h(e^l)} > 0 \text{ and assume now that } \|\gamma_2\|_{L^r} \leq \tau_l. \quad (53)$$

With this additional assumption, we are able to obtain a useful estimation for $\|v_l\|_{L^\infty}$. Technically we will detail a method due to Stampacchia (see Subsection 3.2 for the notations and for an introduction of the method. Here only the first step of the technique will be developed further). Let $s > |\ln \|\zeta_0\|_{L^\infty(\partial\Omega)}|$, we consider the function $\psi_s = v_l - T_s(v_l)$. We have (see Lemma 4) $\psi_s \in H_0^1(\Omega) \cap L^\infty(\Omega)$, and by testing (52) with ψ_s we obtain:

$$\int_{\Omega} (\nu e^{T_l(v_l)} + \kappa) |\nabla \psi_s|^2 + \underbrace{\int_{\Omega} \rho \mathbf{u} \cdot \nabla e^{T_l(v_l)} \psi_s}_{=: I} = \underbrace{\int_{\Omega} g(x, e^{T_l(v_l)}) \psi_s}_{=: II}. \quad (54)$$

• We will now evaluate the terms I and II.

The term I is simplified by writing one of its integrand factors, namely $\nabla e^{T_l(v_l)} \psi_s$ as a gradient. More precisely we have $\nabla e^{T_l(v_l)} \psi_s = \nabla \zeta_l$, with $\zeta_l \in H^1(\Omega) \cap L^\infty$ (see Lemma 14 in the Appendix). Hence by applying the divergence formula we see that I vanishes:

$$I = \int_{\Omega} \rho \mathbf{u} \cdot \nabla \zeta_l \stackrel{\text{div. formula}}{=} - \int_{\Omega} \zeta_l \underbrace{\text{div } \rho \mathbf{u}}_{=0} + \underbrace{\int_{\partial\Omega} \rho \zeta_l \mathbf{u} \cdot \mathbf{n} d\sigma}_{=0 \text{ by (14)}} = 0. \quad (55)$$

We next estimate the term II:

$$II = \int_{\Omega} g(x, e^{T_l(v_l)}) \psi_s \leq \underbrace{\int_{A_{s,l}^+} g^+(x, e^{T_l(v_l)}) \psi_s}_{=: II_1} + \underbrace{\int_{A_{s,l}^-} g^-(x, e^{T_l(v_l)}) \psi_s}_{=: II_2}.$$

Remark now that on $A_{s,l}^-$ we have $v_l \leq -s \leq 0$, which implies that $e^{T_l(v_l)} \leq 1$. Consequently by using the assumption (46) we obtain:

$$II_2 \leq \int_{\Omega} \gamma \psi_s \mathbf{1}_{A_{s,l}^-}. \quad (56)$$

The term II_1 is majorated as follows:

$$II_1 \stackrel{\text{by (47)}}{\leq} \int_{A_{s,l}^+} (\gamma_1 + \gamma_2 h(e^l)) \psi_s \stackrel{\text{by (53)}}{\leq} \int_{\Omega} (\gamma_1 + 1) \psi_s \mathbf{1}_{A_{s,l}^+} \quad (57)$$

• At this point by using the estimates (55), (56) and (57) together with (54) and (45), we obtain:

$$\kappa_0 \int_{\Omega} |\nabla \psi_s|^2 \leq \int_{\Omega} \underbrace{(\gamma \mathbf{1}_{A_{s,l}^-} + (\gamma_1 + 1) \mathbf{1}_{A_{s,l}^+})}_{=: E} \psi_s. \quad (58)$$

Note that $E \in L^r(\Omega)$ and: $\|E\|_{L^r} \leq C_0$, $C_0 = \|\gamma\|_{L^r} + \|\gamma_1 + 1\|_{L^r}$. On the other hand (see Lemma 5) there exists $\tilde{E} \in (L^\beta(\Omega))^3$ satisfying: $\|\tilde{E}\|_{(L^\beta)^N} \leq \|E\|_{L^r}$, and $\int_{\Omega} E \varphi = \int_{\Omega} \tilde{E} \nabla \varphi \quad \forall \varphi \in H_0^1$. Recall also that we have assumed in (8) that $r > \frac{N}{2}$ which implies $\beta > N$. By again using the Hölder inequality we obtain for $\varphi \in H_0^1$:

$$\int_{\Omega} \tilde{E} \nabla \varphi \leq \|\tilde{E}\|_{(L^\beta)^N} \|\varphi\|_{H_0^1} |\text{Supp } \varphi|^{\frac{\beta-2}{2\beta}}.$$

Consequently (58) leads to:

$$\begin{aligned} \kappa_0 \|\psi_s\|_{H_0^1}^2 &\leq \int_{\Omega} \tilde{E} \nabla \psi_s \leq C_0 |A_{s,l}|^{\frac{\beta-2}{2\beta}} \|\psi_s\|_{H_0^1} \\ &\stackrel{\text{Young ineq.}}{\leq} \frac{\kappa_0}{2} \|\psi_s\|_{H_0^1}^2 + \frac{C_0^2}{2\kappa_0} |A_{s,l}|^{\frac{\beta-2}{\beta}}. \end{aligned}$$

Let $\Phi := \frac{\beta-2}{\beta} > \frac{N-2}{N}$. We have obtained the estimate:

$$\|\psi_s\|_{H_0^1}^2 \leq \tilde{C}_0 |A_{s,l}|^\Phi, \quad \tilde{C}_0 = \frac{C_0^2}{4\kappa_0^2}.$$

By now using the Stampacchia estimates (see Subsection 3.2) we obtain the existence of a real Λ independent of l such that $|A_{\Lambda,l}| = 0$. Hence: $\|\tilde{v}_l\|_{L^\infty(\Omega)} \leq \Lambda$, $\Lambda = C(\|\zeta_0\|_{L^\infty(\partial\Omega)}, |\Omega|, N, r) \frac{C_0}{\kappa_0}$

□

Proof of Proposition 9

i) Existence and estimates.

Let $(v_l)_{l \geq 1}$ be the sequence given in Lemma 10. Let also $l \geq 1$ be given and $\tau_l := \frac{1}{h(e^l)}$. We assume that $\|\gamma_2\|_{L^r} \leq \tau_l$. It follows from Lemma 10.ii) that $\|v_l\|_{L^\infty} \leq K$, where $K = K(\kappa_0, DATA)$ (K independent of l) is the integer defined by $K = [\Lambda] + 1$.

Let now $\tau := \tau_K$ and assume that $\|\gamma_2\|_{L^r} \leq \tau$. Then we have: $\|v_K\|_{L^\infty} \leq K$. Hence $T_K(v_K) = v_K$. On the other hand v_K satisfies (S_K) , that is:

$$\begin{cases} \rho \mathbf{u} \nabla e^{v_K} - \text{div}((\nu e^{v_K} + \kappa) \nabla v_K) = g(x, e^{v_K}) & \text{in } \Omega \\ v_K = \ln \zeta_0 & \text{on } \partial\Omega \end{cases}$$

Let $\zeta := e^{v_K}$. We have (see Theorem 13 in the Appendix) $\zeta \in H^1(\Omega) \cap L^\infty(\Omega)$ and $\nabla \zeta = \zeta \nabla v_K$. Consequently ζ is a solution of problem (S).

Moreover $\|\zeta\|_{L^\infty(\Omega)} \leq e^K = C(\kappa_0, DATA)$.

On the other hand $v_K \geq -K$ almost everywhere, implies that $\zeta \geq e^{-K}$ a.e. in Ω . Hence we obtain (48) by setting $\zeta_{\min} = e^{-K}$ and $\zeta_{\max} = e^K$.

• The estimation (49) is obtained as follows.

By using the test function ψ_s^+ instead of ψ_s we obtain:

$$\zeta_{\max} \leq \Lambda_1 := e^{C \frac{\|\gamma_1\|_{L^r}}{\kappa_0}}. \quad (59)$$

This last estimation is only a first step in order to obtain the majoration for ζ_{\max} announced in (49).

In fact, let $\tilde{\kappa} := \nu + \frac{\kappa}{\zeta}$. We have $0 < \nu + \frac{\kappa_0}{\Lambda_1} =: \tilde{\kappa}_0 \leq \tilde{\kappa} \leq \nu + \frac{\kappa_0}{\zeta_{\min}} < \infty$, and

$$\begin{cases} \rho \mathbf{u} \cdot \nabla \zeta - \operatorname{div}(\tilde{\kappa} \nabla \zeta) = g^- + g^+ & \text{in } \Omega \\ \zeta = \zeta_0 & \text{on } \partial\Omega \end{cases}$$

we can then consider the decomposition $\zeta = \zeta_1 + \zeta_2$, where ζ_1 (resp. ζ_2) $\in H^1(\Omega) \cap L^\infty(\Omega)$ satisfies the following problem (S_1) (resp. (S_2)):

$$(S_1) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \zeta_1 - \operatorname{div}(\tilde{\kappa} \nabla \zeta_1) = g^- & \text{in } \Omega \\ \zeta_1 = \zeta_0 & \text{on } \partial\Omega \end{cases}$$

$$(S_2) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \zeta_2 - \operatorname{div}(\tilde{\kappa} \nabla \zeta_2) = g^+ & \text{in } \Omega \\ \zeta_2 = 0 & \text{on } \partial\Omega \end{cases}$$

Note that the second member in the PDE in (S_1) is negative. Hence, by the maximum principle (see [25] p.80 or [2] p.191 for a simplified situation) we obtain: $\zeta_1 \leq \|\zeta_0\|_{L^\infty(\partial\Omega)}$ a.e. in Ω .

By using next the Stampacchia technique (see again Subsection 3.2, Remark 6) we major the function ζ_2 as follows:

$$\zeta_2 \leq C' \frac{\|g^+\|_{L^r}}{\tilde{\kappa}_0} \underset{\text{by (59)}}{\leq} e^{C \frac{\|\gamma_1\|_{L^r}}{\kappa_0}} \frac{\|g^+\|_{L^r}}{k_0}, \quad C = C(DATA).$$

This leads to the majoration (49).

ii) Regularity results.

• If we assume that $\mathbf{u} \in (L^\infty(\Omega))^N$ then $\rho \mathbf{u} \in (L^\infty(\Omega))^N$. Moreover by using the estimates (48) the diffusion coefficient $\nu + \frac{\kappa}{\zeta}$ is bounded from above and below and $g(x, \zeta) \in L^r(\Omega)$ with $r > 3/2$ fixed. Hence (see Lemma 5): $g(x, \zeta) \in W^{-1, \beta}(\Omega)$ with $\beta > N$. Consequently, by using the De Giorgi-Nash Theorem (see [12] Th. 8.22) we obtain: $\zeta \in \mathcal{C}^{0, \alpha}(\overline{\Omega})$, for some $\alpha > 0$.

• Assume that in addition we have:

$$\begin{aligned} \partial\Omega, \zeta_0 \text{ are of class } \mathcal{C}^{2, \alpha}, \quad g \in \mathcal{C}^{0, \alpha}(\overline{\Omega} \times \mathbb{R}^+), \\ \rho \mathbf{u} \in (\mathcal{C}^{1, \alpha}(\overline{\Omega}))^N, \quad \kappa \in \mathcal{C}^{1, \alpha}(\overline{\Omega}). \end{aligned}$$

We have proved in the previous point that $\zeta \in \mathcal{C}^{0, \alpha}(\overline{\Omega})$. We now iterate the Schauder estimates as follows. In a first step we see that $g(x, \zeta)$ and $\nu + \frac{\kappa}{\zeta}$ are in $\mathcal{C}^{0, \alpha}(\overline{\Omega})$, and by applying Theorem 2.7 in [4] we obtain $\zeta \in \mathcal{C}^{1, \alpha}(\overline{\Omega})$. Consequently we now obtain (see Appendix B in [8]) $g(x, \zeta), \nu + \frac{\kappa}{\zeta} \in \mathcal{C}^{1, \alpha}(\overline{\Omega})$ and by using Theorem 2.8 in [4] we finally obtain $\zeta \in \mathcal{C}^{2, \alpha}(\overline{\Omega})$. Hence ζ is a classical solution of (S). □

4.2 Proof of Proposition 7

Let $n \in \mathbb{N}$, and θ_n, φ_n be given. We assume that condition (K_n) is satisfied. Recall that this implies in particular: $\varphi_n \leq \varphi_{\max}$, where $\varphi_{\max} > \|b\|_{L^\infty(\partial\Omega)}$ was fixed in (42). Hence, let $\varepsilon := \varphi_{\max} - \|b\|_{L^\infty(\partial\Omega)} > 0$.

- Step 1: We introduce

$$\kappa^{(n)} := \frac{C_\theta^{(n)}}{\varphi_n}, \quad g_\theta(x, u) := C_5 - C_3Fu^2 + C_4Du.$$

Hence the first subproblem in (Q_n) reads as:

$$(Q_{n.1}) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \theta_{n+1} - \operatorname{div} \left(\left(\nu + \frac{\kappa^{(n)}}{\theta_{n+1}} \right) \nabla \theta_{n+1} \right) = g_\theta(x, \theta_{n+1}) & \text{in } \Omega \\ \theta_{n+1} = a & \text{on } \partial\Omega \end{cases}$$

Moreover, it is easy to verify that:

$$0 < \underbrace{\frac{\alpha_\theta}{\varphi_{\max}}}_{=: \kappa_0} \leq \kappa^{(n)} \leq \left\| \frac{C_\theta^{(n)}}{\varphi_n} \right\|_{L^\infty} < \infty$$

$$g_\theta^+(x, u) = C_5 + C_4D^+u, \quad g_\theta^-(x, u) = -C_3Fu^2 + C_4D^-u.$$

Note also that κ_0 is independent of n . Hence we can apply Proposition 9.i) (take $\zeta = \theta_{n+1}, \kappa = \kappa^{(n)}, g = g_\theta, \gamma_1 = C_5, \gamma_2 = C_4D^+, h(t) = t, \gamma = C_3F + C_4|D^-|, \zeta_0 = a$). We obtain the existence of $\tau_0 > 0$ independent of n such that if $\|D^+\|_{L^r} \leq \tau_0$ then problem $(Q_{n.1})$ admits at least one weak solution $\theta_{n+1} \in H^1(\Omega) \cap L^\infty(\Omega)$. Moreover we have

$$0 < e^{-C\varphi_{\max}} \leq \theta_{n+1} \leq e^{C\varphi_{\max}} < \infty, \quad C = C(\text{DATA}). \quad (60)$$

- Step 2: Let now

$$\tilde{\kappa}^{(n)} := \frac{C_\varphi^{(n)}}{\theta_{n+1}}, \quad g_\varphi^{(n)}(x, u) = g_\varphi^{(n)}(x) = -\varphi_n(C_8\theta_{n+1}^{-1} + C_6\theta_{n+1}F - C_7D).$$

With these notations, the second subproblem in (Q_n) reads as:

$$(Q_{n.2}) \quad \begin{cases} \rho \mathbf{u} \cdot \nabla \varphi_{n+1} - \operatorname{div} \left(\left(\nu + \frac{\tilde{\kappa}^{(n)}}{\varphi_{n+1}} \right) \nabla \varphi_{n+1} \right) = g_\varphi^{(n)} & \text{in } \Omega \\ \varphi_{n+1} = b & \text{on } \partial\Omega \end{cases}$$

We verify that:

$$0 < \underbrace{\frac{\alpha_\varphi}{\theta_{\max}}}_{=: \tilde{\kappa}_0} \leq \tilde{\kappa}^{(n)} \leq \frac{\|C_\varphi^{(n)}\|_{L^\infty}}{\theta_{\min}} < \infty, \quad g_\varphi^+ = \varphi_n C_7 D^+ \leq \varphi_{\max} C_7 D^+,$$

$$g_\varphi^- = -\varphi_n(C_8\theta_{n+1}^{-1} + C_6\theta_{n+1}F - C_7D^-), \quad |g_\varphi^-| \stackrel{\text{by(61)}}{\leq} \tilde{C}(\text{DATA}, \varphi_{\max}).$$

Hence we can apply the proposition 9.i) (take now $\zeta = \varphi_{n+1}, \kappa = \tilde{\kappa}^{(n)}, g = g_\varphi^{(n)}, \gamma_1 = \varphi_{\max} C_7 D^+, h = 0, \gamma = \tilde{C}, \zeta_0 = b$). Then we obtain the existence of

a weak solution⁴ $\varphi_{n+1} \in H^1(\Omega) \cap L^\infty(\Omega)$ for problem $(Q_n.2)$. Moreover we have

$$0 < e^{-C'e^{C\varphi_{\max}}} \leq \varphi_{n+1} \leq e^{C'e^{C\varphi_{\max}}} < \infty, \quad (61)$$

where C, C' depend on DATA but not on n . Moreover, by using (49) we have:

$$\begin{aligned} \|\varphi_{n+1}\|_{L^\infty(\Omega)} &\leq \|b\|_{L^\infty(\partial\Omega)} + C'' \|g_\varphi^+\|_{L^r} e^{C\varphi_{\max}} \|\varphi_{n+1}\|_{L^\infty(\Omega)}, \quad C'' = C''(\text{DATA}), \\ &\leq \|b\|_{L^\infty(\partial\Omega)} + \underbrace{C'' \|C_7\|_{L^\infty} e^{C\varphi_{\max}}}_{:=K(\varphi_{\max}, \text{DATA})} \|D^+\|_{L^r} \|\varphi_{n+1}\|_{L^\infty(\Omega)}. \end{aligned} \quad (62)$$

Assume now that

$$\|D^+\|_{L^r} \leq \tau := \min(\tau_0, \frac{\varepsilon}{\varphi_{\max} K}).$$

Then (62) leads to

$$\|\varphi_{n+1}\|_{L^\infty(\Omega)} \leq \|b\|_{L^\infty(\partial\Omega)} + \frac{\varepsilon}{\varphi_{\max}} \|\varphi_{n+1}\|_{L^\infty(\Omega)},$$

and it follows:

$$\|\varphi_{n+1}\|_{L^\infty(\Omega)} \leq \frac{\|b\|_{L^\infty(\partial\Omega)}}{1 - \frac{\varepsilon}{\varphi_{\max}}} \leq \varphi_{\max}.$$

- *Final Step*: if we assume that $\|D^+\|_{L^r} \leq \tau$ then by using the results established in the previous two steps, we conclude that there exists a solution $(\theta_{n+1}, \varphi_{n+1}) \in (H^1 \cap L^\infty)^2$ for problem (Q_n) . Moreover this solution satisfies (K_{n+1}) and (43),(44) hold. \square

5 Proofs of the theorems

We begin by a lemma:

Lemma 11 *Under the assumptions of Proposition 7, we can extract a subsequence (still denoted by (θ_n, φ_n)) such that*

$$\theta_n \overset{*}{\rightharpoonup} \theta, \varphi_n \overset{*}{\rightharpoonup} \varphi \text{ in } L^\infty(\Omega), \quad \theta_n \rightharpoonup \theta, \varphi_n \rightharpoonup \varphi \text{ in } H^1(\Omega), \quad (63)$$

$$\frac{1}{\theta_n \varphi_n} \rightarrow \frac{1}{\theta \varphi}, C_\theta^{(n)} \rightarrow C_\theta(x, \theta, \varphi), C_\varphi^{(n)} \rightarrow C_\varphi(x, \theta, \varphi) \text{ in } L^p(\Omega), \quad p < \infty \quad (64)$$

Proof

The first properties in (63) follow directly from Proposition 7. By next using $\theta_n - \theta_0$ as test function in $(Q_n.1)$ and $\varphi_n - \varphi_0$ as test function in $(Q_n.2)$ we obtain a uniform bound for (θ_n) and (φ_n) in the H^1 -norm. Hence the second properties in (63) follow. Finally property (64) is obtained by using the dominated convergence theorem. In fact we have:

$$\begin{aligned} \frac{1}{\theta_n \varphi_n} &\rightarrow \frac{1}{\theta \varphi}, \quad C_\theta^{(n)} \rightarrow C_\theta(x, \theta, \varphi), \quad C_\varphi^{(n)} \rightarrow C_\varphi(x, \theta, \varphi) \quad \text{a.e in } \Omega, \\ \left| \frac{1}{\theta_n \varphi_n} \right| &\leq \frac{1}{\theta_{\min} \varphi_{\min}}, \quad |C_\theta^{(n)}| \leq \sup_{(v,w) \in \mathcal{K}} C_\theta(x, v, w), \quad |C_\varphi^{(n)}| \leq \sup_{(v,w) \in \mathcal{K}} C_\varphi(x, v, w), \end{aligned}$$

where $\mathcal{K} = [0, \theta_{\max}] \times [0, \varphi_{\max}]$. \square

⁴At this stage there is anymore additional condition needed for $\|D^+\|_{L^r}$ because $h = 0$

5.1 Proof of Theorem 1

By using (63) together with (64) we obtain:

$$\begin{aligned} \rho \mathbf{u} \cdot \nabla \theta_{n+1} &\rightharpoonup \rho \mathbf{u} \cdot \nabla \theta \text{ in } L^1(\Omega) \\ \left(\nu + \frac{C_\theta^{(n)}}{\theta_{n+1} \varphi_n}\right) \nabla \theta_{n+1} &\rightharpoonup \left(\nu + \frac{C_\theta}{\theta \varphi}\right) \nabla \theta \text{ in } L^q(\Omega), \quad q < 2 \\ \rho \mathbf{u} \cdot \nabla \varphi_{n+1} &\rightharpoonup \rho \mathbf{u} \cdot \nabla \varphi \text{ in } L^1(\Omega) \\ \left(\nu + \frac{C_\varphi^{(n)}}{\theta_{n+1} \varphi_{n+1}}\right) \nabla \varphi_{n+1} &\rightharpoonup \left(\nu + \frac{C_\varphi}{\theta \varphi}\right) \nabla \varphi \text{ in } L^q(\Omega), \quad q < 2. \end{aligned}$$

Moreover by using (63) together with the property $\theta_{n+1} \geq \theta_{\min} > 0$ we obtain:

$$\theta_{n+1}^2 \rightarrow \theta^2, \quad \theta_{n+1} \rightarrow \theta, \quad \theta_{n+1}^{-1} \rightarrow \theta^{-1} \quad \text{in } L^p(\Omega), \quad p < \infty.$$

Hence we can pass to the limit in the approximate problems (Q_n) . We obtain a weak solution (θ, φ) for problem (Q). That is for all $\psi \in \mathcal{D}(\Omega)$:

$$\begin{aligned} \int_{\Omega} \rho \mathbf{u} \cdot \nabla \theta \psi + \int_{\Omega} \left(\nu + \frac{C_\theta}{\theta \varphi}\right) \nabla \theta \cdot \nabla \psi &= \int_{\Omega} (C_5 - C_3 F \theta^2 + C_4 \theta D) \psi \\ \int_{\Omega} \rho \mathbf{u} \cdot \nabla \varphi \psi + \int_{\Omega} \left(\nu + \frac{C_\varphi}{\theta \varphi}\right) \nabla \varphi \cdot \nabla \psi &= \int_{\Omega} -\varphi (C_8 \theta^{-1} + C_6 \theta F - C_7 D) \psi \\ \theta = a, \varphi = b &\quad \text{on } \partial \Omega \end{aligned}$$

Moreover this solution satisfies:

$$\theta, \varphi \in H^1(\Omega) \cap L^\infty(\Omega), \quad \theta, \varphi \geq \min(\theta_{\min}, \varphi_{\min}) > 0 \quad \text{a.e. in } \Omega.$$

5.2 Proof of Theorem 2

When $N = 2$ the function F has a stronger property of positivity (see lemma 12 in the Appendix):

$$F \geq \frac{D^2}{3}. \quad (65)$$

We will see that this last property allows one to obtain a weak solution for problem (Q) under the assumptions (A_0) and (A_1) but without assuming a low compressibility condition of the form (20).

In order to prove this result, we take over the proof of Proposition 7 with slight modifications: if (θ_n, φ_n) is given and satisfies (H_n) (it is not useful to consider (K_n) here) then problem (Q_n) has at least one solution $(\theta_{n+1}, \varphi_{n+1})$ satisfying in addition (H_{n+1}) and the estimates (43),(44).

- Step 1: By using property (65) we majore g_θ^+ as follows:

$$\begin{aligned} g_\theta^+(x, u) &= C_5 + C_4 D^+ u - C_3 F u^2 \leq C_5 + C_4 D^+ u - \frac{C_3}{3} (D^+ u)^2 \\ &\leq C_5 + \underbrace{D^+ u}_{\geq 0} \underbrace{\left(C_4 - \frac{C_3}{3} (D^+ u)\right)}_{\leq 0 \text{ if } D^+ u \geq \frac{3C_4}{C_3}} \leq C_5 + \frac{3C_4^2}{C_3}. \end{aligned}$$

Hence we have here estimated g_θ^+ independtly of the second variable. We then apply Proposition 9.i), but now we take $\gamma_1 = C_5 + \frac{3C_4^2}{C_3}$ instead of C_5

and $h(t) = 0$ instead of $h(t) = t$. It follow that there exists (without any condition on $\|D^+\|_{L^r}$ because $h = 0$) a weak solution $\theta_{n+1} \in H^1(\Omega) \cap L^\infty(\Omega)$ for problem $(Q_n.1)$, with the estimate:

$$0 < e^{-C\|\varphi_n\|_{L^\infty}} \leq \theta_{n+1} \leq e^{C\|\varphi_n\|_{L^\infty}} < \infty, \quad C = C(DATA). \quad (66)$$

- Step 2: By taking over the arguments presented in the proof of Proposition 7 we see that problem $(Q_n.2)$ has at least one positive solution $\varphi_{n+1} \in H^1(\Omega) \cap L^\infty(\Omega)$.

Hence at this point we have obtained a weak solution $(\theta_{n+1}, \varphi_{n+1})$ for (Q_n) satisfying in addition (H_{n+1}) . It remains to prove that the estimates (43),(44) hold. We have made a first step in this direction by proving (66). We will now prove:

$$0 < e^{-C\|\theta_{n+1}\|_{L^\infty}} \leq \varphi_{n+1} \leq \|b\|_{L^\infty(\partial\Omega)}, \quad C = C(DATA). \quad (67)$$

In fact, by using the additional assumption (A_1) we majore the function $g_\varphi^{(n)}$ as follows:

$$\begin{aligned} g_\varphi^{(n)}(x) &= -\varphi_n(C_8\theta_{n+1}^{-1} + C_6\theta_{n+1}F - C_7D) \stackrel{\text{by (65)}}{\leq} -\varphi_n(C_8\theta_{n+1}^{-1} + \frac{C_6}{3}\theta_{n+1}D^2 - C_7D) \\ &\leq -\frac{\rho\varphi_n}{\theta_{n+1}}(c_8 + \frac{c_6}{3}(\theta_{n+1}D)^2 - c_7(\theta_{n+1}D)) = -\frac{\rho\varphi_n}{\theta_{n+1}}\mathcal{P}(\theta_{n+1}D), \end{aligned}$$

with $\mathcal{P}(X) := \frac{c_6}{3}X^2 - c_7X + c_8$. We remark that the discriminant Δ of \mathcal{P} is negative: $\Delta = c_7^2 - \frac{4}{3}c_6c_8 = -4.864 * 10^{-3} < 0$. It follows that \mathcal{P} is positive and consequently $g_\varphi^{(n)}$ is negative. Hence by applying Proposition 9.i) with now $g^+ \equiv 0$ we obtain (67).

- Final Step: By using (66) together with (67) we obtain the estimates (43) and (44). Hence we have recovered the conclusions of Proposition 7. The remaindee of the proof for Theorem 3 is exactly the same as for Theorem 1: we can extract a subsequence with the properties (63)-(63). These properties are sufficient to pass to the limit $n \rightarrow \infty$ in (Q_n) and Theorem 2 follows.

5.3 Proof of Theorem 3

Let (θ, φ) be a weak solution for (Q) in the class \mathcal{S} and consider the notations:

$$g_\theta(x, u) := C_5 - C_3Fu^2 + C_4Du, \quad g_\varphi(x, u) := -u(C_8\theta^{-1} + C_6\theta F - C_7D).$$

- i) It suffices to remark that the coefficients $\frac{C_\theta}{\varphi}$ and $\frac{C_\varphi}{\theta}$ are bounded from above and below, and g_θ, g_φ are Caratheodory functions satisfying (46) and (47). Hence we can apply the first point in Proposition 9-ii) in each equation of (Q). We obtain: $\theta, \varphi \in \mathcal{C}^{0,\alpha}(\bar{\Omega})$, for some $\alpha > 0$.
- ii) Assume that in addition we have:

$$\begin{aligned} \partial\Omega, a, b \text{ are of class } \mathcal{C}^{2,\alpha}, \quad F \in \mathcal{C}^{0,\alpha}(\bar{\Omega}), \\ \rho\mathbf{u} \in (\mathcal{C}^{1,\alpha}(\bar{\Omega}))^N, \quad C_\theta, C_\varphi \in \mathcal{C}^{1,\alpha}(\bar{\Omega} \times (\mathbb{R}^+)^2). \end{aligned}$$

We remark now that the conditions in the second part of Proposition 9-ii) are satisfied for each equation of (Q). Hence $\theta, \varphi \in \mathcal{C}^{2,\alpha}(\bar{\Omega})$ and it is a classical solution of (Q).

6 Appendix

6.1 Derivation of the $\theta - \varphi$ model.

The model is constructed from the $k - \varepsilon$ one which takes the form:

$$\partial_t k + \mathbf{u} \cdot \nabla k - \frac{c_\nu}{\rho} \operatorname{div} \left(\left(\nu + \rho \frac{k^2}{\varepsilon} \right) \nabla k \right) = c_\nu \frac{k^2}{\varepsilon} F - \frac{2}{3} k D - \varepsilon, \quad (68)$$

$$\partial_t \varepsilon + \mathbf{u} \cdot \nabla \varepsilon - \frac{c_\varepsilon}{\rho} \operatorname{div} \left(\left(\nu + \rho \frac{k^2}{\varepsilon} \right) \nabla \varepsilon \right) = c_1 k F - \frac{2c_1}{3c_\nu} \varepsilon D - c_2 \frac{\varepsilon^2}{k}, \quad (69)$$

where $D(x, t) := \operatorname{div} \mathbf{u}(x, t)$, $F(x, t) := \frac{1}{2} |\nabla \mathbf{u} + (\nabla \mathbf{u})^T|^2 - \frac{2}{3} D(x, t)^2 \geq 0$ (see the next subsection) and $c_\nu, c_\varepsilon, c_1, c_2$ are generally taken as positive constants. Their usual values are (see [23] p.122):

$$c_\nu = 0.09, \quad c_\varepsilon = 0.07, \quad c_1 = 0.128, \quad c_2 = 1.92. \quad (70)$$

We then consider the new variables

$$\theta = \frac{k}{\varepsilon}, \quad \varphi = k^\alpha \varepsilon^\beta, \quad (71)$$

with α and β to be chosen appropriately. Let D_t denote the total derivative operator. By using (68) together with (69) we obtain an equation for θ :

$$\begin{aligned} D_t \theta &= \frac{\partial \theta}{\partial t} + \mathbf{u} \cdot \nabla \theta = \frac{1}{\varepsilon} D_t k - \frac{k}{\varepsilon^2} D_t \varepsilon = -c_3 \theta^2 F + c_4 \theta D + c_5 + \operatorname{Diff}_\theta, \\ c_3 &= c_1 - c_\nu, \quad c_4 = \frac{2}{3} \left(\frac{c_1}{c_\nu} - 1 \right), \quad c_5 = c_2 - 1, \end{aligned} \quad (72)$$

where $\operatorname{Diff}_\theta$ denotes the collected terms coming from the viscous one in the k and ε equations. The equation for φ is obtained in the same way:

$$\begin{aligned} D_t \varphi &= \alpha k^{\alpha-1} \varepsilon^\beta D_t k + \beta k^\alpha \varepsilon^{\beta-1} D_t \varepsilon = \alpha k^{\alpha-1} \varepsilon^\beta \left(c_\nu \frac{k^2}{\varepsilon} F - \frac{2}{3} k D - \varepsilon \right) + \operatorname{Diff}1_\varphi \\ &+ \beta k^\alpha \varepsilon^{\beta-1} \left(c_1 k F - \frac{2c_1}{3c_\nu} \varepsilon D - c_2 \frac{\varepsilon^2}{k} \right) + \operatorname{Diff}2_\varphi = F k^{\alpha+1} \varepsilon^{\beta-1} (\alpha c_\nu + \beta c_1) \\ &- k^\alpha \varepsilon^\beta D \frac{2}{3} \left(\alpha + \beta \frac{c_1}{c_\nu} \right) - k^{\alpha-1} \varepsilon^{\beta+1} (\alpha + \beta c_2) + \operatorname{Diff}_\varphi, \\ &= -\varphi (c_6 \theta F - c_7 D + c_8 \theta^{-1}) + \operatorname{Diff}_\varphi, \\ c_6 &= -\alpha c_\nu - \beta c_1, \quad c_7 = -\frac{2}{3} \left(\alpha + \beta \frac{c_1}{c_\nu} \right), \quad c_8 = \alpha + \beta c_2, \end{aligned} \quad (73)$$

where $\operatorname{Diff}_\varphi = \operatorname{Diff}1_\varphi + \operatorname{Diff}2_\varphi$ is the sum of the terms coming from the viscous one in the k and ε equations.

The usual constant values for the parameters c_3, c_4, c_5 appearing in the model are obtained by replacing the values (70) in the expressions (72). This leads to:

$$c_3 = 0.038, \quad c_4 = 0.2815, \quad c_5 = 0.92.$$

At this stage it remains to choose appropriately α and β in (71), and to model the terms $\operatorname{Diff}_\theta$ and $\operatorname{Diff}_\varphi$.

It is shown in [23], page 67, that a good choice in the incompressible situation

(i.e. when $D = 0$) is for instance $\alpha = -3, \beta = 2$. This leads to the following constant values:

$$c_6 = 0.014, \quad c_7 = 0.104, \quad c_8 = 0.84, \quad (74)$$

and this makes the dynamic stable for the equation in φ . I.e. in the absence of the viscous part Diff_φ , we have: $D_t\varphi \leq 0$.

In the compressible situation the authors suggest in [23], page 125 to consider another choice: $\alpha = -2\frac{c_1}{c_\nu} \approx -2.88$ and $\beta = 2$ which makes again the dynamic stable.

Nevertheless in this last situation, the variable φ does not have a clear physical meaning (whereas when $\alpha = -3, \beta = 2$ we have $\varphi = \frac{\varepsilon^2}{k^3}$ and $L := \varphi^{-1/2}$ represents a length scale of turbulence (see [26])). Moreover, a carefully estimation schows (see Lemma 12 in the next subsection) that when $N = 2$ we have $F \geq \frac{1}{3}D^2$. This leads to:

$$-\varphi(c_6\theta F - c_7D + c_8\theta^{-1}) \leq -\frac{\varphi}{\theta}\left(\frac{c_6}{3}(\theta D)^2 - c_7(\theta D) + c_8\right) = -\frac{\varphi}{\theta}\mathcal{P}(\theta D),$$

with $\mathcal{P}(X) := \frac{c_6}{3}X^2 - c_7X + c_8$. Hence the choice $\alpha = -3, \beta = 2$ makes again the dynamic stable when $N = 2$. In fact, in this case c_6, c_7, c_8 take the values (74) and the discriminant Δ of \mathcal{P} is $\Delta = c_7^2 - \frac{4}{3}c_6c_8 = -4.864 * 10^{-3} < 0$. Consequently $\mathcal{P}(X) \geq 0$ and $D_t\varphi \leq 0$ in the absence of the viscous terms.

In consequence we point out that the choice $\alpha = -3, \beta = 2$ is also interesting in the compressible situation. We shall make this choice in all the situations. Our analysis (see Theorem 1) shows that this leads to a well posed model even when $N = 3$ under an additional assumption of low compressibility of the flow.

The terms Diff_θ and Diff_φ are modelled (see [16]) by:

$$\text{Diff}_\theta = \frac{1}{\rho} \text{div}((\nu + c_\theta\nu_t)\nabla\theta), \quad \text{Diff}_\varphi = \frac{1}{\rho} \text{div}((\nu + c_\varphi\nu_t)\nabla\theta), \quad (75)$$

where $\nu_t := \rho c_\nu \frac{k^2}{\varepsilon} = \rho \frac{c_\nu}{\theta\varphi}$ is the turbulent viscosity coming from the equation of k , and c_θ, c_φ are two new parameters for the model.

The determination of the parameters c_θ and c_φ can be realized in the same way as for the determination of the coefficients arising in the $k - \varepsilon$ model (see [23]). In [22] a constant value for both c_θ and c_φ was numerically tuned from a simulation of a Poiseuil flow. However better results are obtained if we allow c_θ, c_φ to be some positive functions (see [24]). In our analysis we allow the coefficients to be of a very general form, in particular c_θ, c_φ may depend on x, θ and φ . We only assume that they are Caratheodory functions and that they satisfy some positivity and boundedness properties (see (17)-(19), where $C_\theta = \rho c_\theta$ and $C_\varphi = \rho c_\varphi$).

6.2 Positivity of the function F

In this paragraph we will establish some properties of positivity for the function F appearing in the models.

Let $\mathcal{M}_N(\mathbb{R})$ denote the vector space of the N -square matrix with real coefficients, equipped with the scalar product:

$$A : B = \sum_i \sum_j a_{ij} b_{ij}, \quad \forall A = (a_{ij}), B = (b_{ij}) \in \mathcal{M}_N(\mathbb{R}).$$

Hence $|A|^2 := \sqrt{A : A}$ defines a norm on $\mathcal{M}_N(\mathbb{R})$.

For a vector field $\mathbf{u} : \Omega \rightarrow \mathbb{R}^N$ we classically defines the gradient $\nabla \mathbf{u}$ and the divergence $\text{div } \mathbf{u}$ ($=:D$) by:

$$\begin{aligned} \nabla \mathbf{u} : \Omega &\rightarrow \mathcal{M}_N(\mathbb{R}), & (\nabla \mathbf{u}(x))_{ij} &= \frac{\partial u_i(x)}{\partial x_j}, \\ D : \Omega &\rightarrow \mathbb{R}, & D(x) &= \sum_i \frac{\partial u_i(x)}{\partial x_i} = \text{Tr}(\nabla \mathbf{u}). \end{aligned}$$

Recall that the function F was defined by the formula:

$$F(x) := \frac{1}{2} |\nabla \mathbf{u} + (\nabla \mathbf{u})^T|^2 - \frac{2}{3} D(x)^2, \quad (76)$$

and by an easy calculation we obtain:

$$F(x) = (\nabla \mathbf{u} + (\nabla \mathbf{u})^t) : \nabla \mathbf{u} - \frac{2}{3} D(x)^2. \quad (77)$$

This last expression is sometimes chosen (for instance in [23]) to equivalently define F .

The important fact is that we always have $F \geq 0$ but moreover, when $N = 2$ the stronger estimate: $F \geq \frac{1}{3} D^2$ holds. These properties are established in the following lemma:

Lemma 12 *The function F satisfies the estimates:*

$$\begin{aligned} F &= \frac{2}{3} ((\partial_1 u_1 - \partial_2 u_2)^2 + (\partial_1 u_1 - \partial_3 u_3)^2 + (\partial_2 u_2 - \partial_3 u_3)^2 \\ &\quad + (\partial_2 u_1 + \partial_1 u_2)^2 + (\partial_3 u_1 + \partial_1 u_3)^2 + (\partial_3 u_2 + \partial_2 u_3)^2) \geq 0, \quad \text{when } N = 3. \\ F &= (\partial_1 u_1 - \partial_2 u_2)^2 + (\partial_2 u_1 + \partial_1 u_2)^2 + \frac{1}{3} D^2 \geq \frac{1}{3} D^2, \quad \text{when } N = 2. \end{aligned}$$

Proof

Let $N = 2$ or 3 , and $M := \nabla \mathbf{u} + (\nabla \mathbf{u})^t$. Then $M_{ij} = \partial_j u_i + \partial_i u_j$ and we obtain:

$$\begin{aligned} |M|^2 &= \sum_i ((2\partial_i u_i)^2 + \sum_{j \neq i} (\partial_j u_i + \partial_i u_j)^2) = 4 \sum_i (\partial_i u_i)^2 + 2 \sum_i \sum_{j > i} (\partial_j u_i + \partial_i u_j)^2, \\ F &= \frac{1}{2} |M|^2 - \frac{2}{3} D^2 = 2 \underbrace{\sum_i (\partial_i u_i)^2}_{:=A} - \frac{2}{3} (\sum_i \partial_i u_i)^2 + \sum_i \sum_{j > i} (\partial_j u_i + \partial_i u_j)^2 \end{aligned}$$

The term A is evaluated separately in the cases $N = 2$ and $N = 3$.

We remark that $2(a^2 + b^2 + c^2) - \frac{2}{3}(a + b + c)^2 = \frac{2}{3}((a - b)^2 + (a - c)^2 + (b - c)^2)$.

Hence, when $N = 3$ we have

$$A = \frac{2}{3} ((\partial_1 u_1 - \partial_2 u_2)^2 + (\partial_1 u_1 - \partial_3 u_3)^2 + (\partial_2 u_2 - \partial_3 u_3)^2),$$

and we obtain the expression announced for F .

In the same way, we remark that

$$\begin{aligned} 2(a^2 + b^2) - \frac{2}{3}(a + b)^2 &= \frac{2}{3}(a - b)^2 + \frac{2}{3}(a^2 + b^2) \\ &= \frac{2}{3}(a - b)^2 + \frac{1}{3}(a + b)^2 + \frac{1}{3}(a - b)^2 = (a - b)^2 + \frac{1}{3}(a + b)^2. \end{aligned}$$

Hence, when $N = 2$ we obtain:

$$A = (\partial_1 u_1 - \partial_2 u_2)^2 + \frac{1}{3} \underbrace{(\partial_1 u_1 + \partial_2 u_2)^2}_{=D^2},$$

and the expression for F follows. \square

6.3 On the low compressibility assumption

We will show here that without any assumption of low compressibility of the form (20), problem (Q) may be very hard to analyze when $N=3$, and singular solutions or non existence of weak solution may occur.

When the dimension equals two we have seen in Theorem 2 that a condition of low compressibility is not necessary. The reason is related to the fact that a stronger property of positivity for F holds in this case, i.e. we have: $F \geq \frac{D^2}{3}$. When the dimension equals 3 such a property does not hold in general.

In fact, let for instance $\Omega = B_{\mathbb{R}^3}(0, 1)$ and

$$\rho = \left(\prod_{i=1}^3 (x_i + 4) \right)^{-1}, \quad \mathbf{u} = (u_1, u_2, u_3)^t, \quad u_i = x_i + 4.$$

Then a simple calculation gives: $D = 3$ and $F = 0$. Hence (Q) reads as:

$$\begin{cases} \rho \mathbf{u} \cdot \nabla \theta - \operatorname{div} \left(\left(\nu + \frac{C_\theta}{\theta \varphi} \right) \nabla \theta \right) = 3C_4 \theta + C_5 & \text{in } \Omega \\ \rho \mathbf{u} \cdot \nabla \varphi - \operatorname{div} \left(\left(\nu + \frac{C_\varphi}{\theta \varphi} \right) \nabla \varphi \right) = \varphi (3C_7 - C_8 \theta^{-1}) & \text{in } \Omega \\ \theta = a, \varphi = b & \text{on } \partial\Omega \end{cases}$$

In this situation the problem becomes hard to analyze. Assume however that we have obtained a solution (θ, φ) in the class \mathcal{S} . Then the equation satisfied by θ is closely related to:

$$(R) \quad \begin{cases} -\operatorname{div}(\eta \nabla \theta) = 3C_4 e^\theta + C_5 & \text{in } \Omega, \\ \theta = a & \text{on } \partial\Omega, \end{cases}$$

with η bounded from above and below. Hence a contradiction can occur because the problem (R) may not have any weak solution (see for instance [3, 19]). Note that in the considered example ρ and \mathbf{u} satisfy all the conditions needed in (A_0) , except $\mathbf{u} \cdot \mathbf{n} = 0$ on $\partial\Omega$ but this is not restrictive for the purpose here. In fact we can consider the domain $\Omega_1 = B_{\mathbb{R}^3}(0, 2)$ which contains Ω and we can extend ρ, \mathbf{u} in Ω_1 in such a way that all the conditions in (A_0) are satisfied. Hence we obtain an example within the main situation of the study but the evocated problems remain the same.

6.4 A generalized chain rule

Let $u \in W^{1,p}(\Omega)$ and $G : \mathbb{R} \rightarrow \mathbb{R}$ be a Lipschitz function. We recall here some useful properties of the composed function $G(u)$. In particular we shall see that $G(u) \in W^{1,p}(\Omega)$. Moreover in some situations we also have $\frac{\partial}{\partial x_i} G(u) = G'(u) \frac{\partial u}{\partial x_i}$.

The main result we have in mind is Theorem 13 which is due to Stampacchia. In particular we point out that the additional assumption $G(0) = 0$ for the Lipschitz function G is only necessary if Ω is unbounded and $p \neq \infty$, or if we want that $G(u)$ has a vanishing trace on $\partial\Omega$ when u has it (this last situation was in fact the case of interest of Stampacchia).

Theorem 13 *Let G be a Lipschitz real function, $\Omega \subset \mathbb{R}^N$ be a bounded open Lipschitz domain, and $u \in W^{1,p}(\Omega)$, with $p \in [1; \infty]$. We have:*

- i) $G(u) \in W^{1,p}(\Omega)$. Moreover if $u \in W_0^{1,p}(\Omega)$ and $G(0) = 0$ then $G(u) \in W_0^{1,p}(\Omega)$.
- ii) If G' has a finite number of discontinuity⁵. Then the weak derivatives of $G(u)$ are given by the formula:

$$\frac{\partial}{\partial x_i} G(u) = G'(u) \frac{\partial u}{\partial x_i} \quad \text{a.e. in } \Omega. \quad (78)$$

Proof

See the appendix in [25] for the original proof or [12] Theorem 7.8 and [9] Theorem 4 for alternative proofs and additional comments.

We also recall that the formula (78) may be interpreted in some critical points. In fact, let $(t_i)_{i=1, \dots, n}$ denote the points of discontinuity of G' and let $E_i := \{x \in \Omega : u(x) = t_i\}$ be the associated level sets for the function u . Let $1 \leq i \leq n$ be a fixed integer. If $|E_i| > 0$ then the formula (78) has a priori no sense in this last set which is not negligible. Nevertheless it can be shown (see [25]) that $\frac{\partial u}{\partial x_i} = 0$ on such a set. Hence we interpret the right hand side of (78) as zero in the critical set E_i . \square

We now establish some technical results used in the proof of Lemma 10. Let $v_l \in H^1(\Omega)$ be given and consider the function $h_{\pm} \in L_{loc}^1(\mathbb{R})$ defined by:

$$h_{\pm}(y) := e^y(y - T_{s,\pm}(y)). \quad (79)$$

We introduce the functions

$$\begin{aligned} g_{\pm}(t) &:= \int_0^{T_l(t)} h_{\pm}(y) dy, \quad t \in \mathbb{R} \\ \zeta_{l,\pm}(x) &:= g_{\pm}(v_l(x)) \quad \text{a.e. } x \in \Omega. \end{aligned}$$

We have:

Lemma 14 *The function $\zeta_{l,\pm}$ has the properties:*

$$\begin{aligned} \zeta_{l,\pm} &\in H^1(\Omega) \cap L^\infty(\Omega), \\ \nabla \zeta_{l,\pm} &= \nabla e^{T_l(v_l)} \psi_s^{\pm}. \end{aligned}$$

⁵The derivative G' of G takes here the classical sense

Proof

A simple majoration gives

$$|\zeta_{l,\pm}(x)| \leq \int_0^l |h_{\pm}(y)| dy, \quad \text{a.e. in } \Omega,$$

hence $\zeta_{l,\pm} \in L^\infty(\Omega)$.

We next remark that g_{\pm} is a Lipschitz function and its classical derivative is given by:

$$g'_{\pm}(t) = h_{\pm}(T_l(t)) \mathbf{1}_{\{|t| \leq l\}} \quad \forall t \neq l, -l.$$

Hence by using Theorem 13 we obtain $\zeta_{l,\pm} \in H^1(\Omega)$ and

$$\begin{aligned} \nabla \zeta_{l,\pm} &= h_{\pm}(T_l(v_l)) \mathbf{1}_{\{|v_l| \leq l\}} \nabla v_l = h_{\pm}(v_l) \mathbf{1}_{\{|v_l| \leq l\}} \nabla v_l \\ &= e^{v_l} \nabla v_l \mathbf{1}_{\{|v_l| \leq l\}} \psi_s^{\pm}(v_l) = \nabla e^{T_l(v_l)} \psi_s^{\pm}. \end{aligned}$$

□

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