

ON VON KARMAN MODELING FOR TURBULENT FLOW NEAR A WALL*

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Abstract. Mixing-length models are often used by engineers in order to take into account turbulence phenomena in a flow. This kind of model is obtained by adding a turbulent viscosity to the laminar one in Navier-Stokes equations. When the flow is confined between two close walls, von Karman’s model consists of adding a viscosity which depends on the rate of strain multiplied by the square of distance to the wall. In this short paper, we present a mathematical analysis of such modeling. In particular, we explain why von Karman’s model is numerically ill-conditioned when using a finite element method with a small laminar viscosity. Details of analysis can be found in [1], [2].

Key words. Stokes Equations, Weighted Sobolev Spaces, Finite Element Method.

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1. Introduction. Let $\Omega \subset \mathbb{R}^3$ be an open bounded domain with a Lipschitz boundary $\partial\Omega$. We assume that Ω is occupied by a fluid with velocity \mathbf{u} and pressure p , satisfying the incompressible Navier-Stokes equations:

$$\frac{\partial \mathbf{u}}{\partial t} + (\mathbf{u} \cdot \nabla) \mathbf{u} - \operatorname{div}(2\nu\epsilon(\mathbf{u})) + \nabla p = \mathbf{F}, \tag{1}$$

$$\operatorname{div}(\mathbf{u}) = 0, \tag{2}$$

where ν is the kinematic viscosity, \mathbf{F} is the external force and $\epsilon(\mathbf{u})$ is the strain tensor defined by

$$\epsilon(\mathbf{u})_{i,j} = \frac{1}{2} \left(\frac{\partial \mathbf{u}_i}{\partial x_j} + \frac{\partial \mathbf{u}_j}{\partial x_i} \right), \quad 1 \leq i, j \leq 3. \tag{3}$$

The boundary conditions can be of several types such as the adherent condition $\mathbf{u} = 0$ on a part of $\partial\Omega$, and slip conditions $\mathbf{u} \cdot \mathbf{n} = 0$ with $(\epsilon(\mathbf{u}) \cdot \mathbf{n}) \cdot \mathbf{t}_j = 0$, $j = 1, 2$ on another part of $\partial\Omega$. Here \mathbf{n} is the outward normal to $\partial\Omega$ and \mathbf{t}_j , $j = 1, 2$ are two corresponding unit tangent vectors.

When the domain is very flat, for instance $\Omega = (0, 1) \times (0, 1) \times (0, \varepsilon)$ with $0 < \varepsilon \ll 1$, the viscosity of von Karman turbulent model depends on the square of the distance to the wall multiplied by $|\epsilon(\mathbf{u})|$. More precisely

$$\nu = \nu(\mathbf{x}, \mathbf{u}) = \nu_0 + \beta_t |\epsilon(\mathbf{u})| d^2(\mathbf{x}, \partial\Omega). \tag{4}$$

Here, β_t is a positive parameter, $d(\mathbf{x}, \partial\Omega)$ is the distance from $\mathbf{x} \in \Omega$ to $\partial\Omega$, i.e. $d(\mathbf{x}, \partial\Omega) = \inf_{\mathbf{y} \in \partial\Omega} |\mathbf{x} - \mathbf{y}|$, $\nu_0 > 0$ is the constant kinematic laminar viscosity of the fluid and

$$|\epsilon(\mathbf{u})| = \left(\sum_{1 \leq i, j \leq 3} \epsilon(\mathbf{u})_{i,j}^2 \right)^{\frac{1}{2}}. \tag{5}$$

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We generalize this viscosity with ν to be

$$\nu = \nu_0 + \beta_t |\epsilon(\mathbf{u})| l^{2-\alpha} d^\alpha(\mathbf{x}, \partial\Omega), \tag{6}$$

where l is a characteristic length of the domain Ω and α is a constant parameter such that $0 \leq \alpha \leq 2$ ($\alpha = 2$ corresponds to von Karman modeling).

Our first consideration concerns the stationary Stokes problem related to (1)-(2) with a viscosity defined by (6).

2. Analysis of a stationary Stokes problem. Let us consider the equations

$$-\operatorname{div}(2\nu\epsilon(\mathbf{u})) + \nabla p = \mathbf{F}, \text{ in } \Omega, \tag{7}$$

$$\operatorname{div}(\mathbf{u}) = 0, \text{ in } \Omega, \tag{8}$$

with the following boundary condition

$$\mathbf{u} = 0 \text{ on } \partial\Omega. \tag{9}$$

The details of analysis of problem defined by (7)-(9) can be found in [1] and [2]. In the following, we present results leading to our main theorem.

First we start by multiplying Equation (7) by a vectorial test function \mathbf{w} vanishing on the boundary $\partial\Omega$, and we proceed with an integration by part. We also multiply (8) by a test function q and we formally obtain the two following relations:

$$\int_{\Omega} 2\nu\epsilon(\mathbf{u}) * \epsilon(\mathbf{w})d\mathbf{x} - \int_{\Omega} p \operatorname{div}(\mathbf{w})d\mathbf{x} = \int_{\Omega} \mathbf{F} \cdot \mathbf{w}d\mathbf{x}, \tag{10}$$

$$\int_{\Omega} \operatorname{div}(\mathbf{u})qd\mathbf{x} = 0, \tag{11}$$

where $*$ denotes the tensorial product: $\epsilon(\mathbf{u}) * \epsilon(\mathbf{w}) = \sum_{1 \leq i,j \leq 3} \epsilon(\mathbf{u})_{i,j} \cdot \epsilon(\mathbf{w})_{i,j}$.

Let us remark that if ν is given by (6), then we have two kinds of integrals in (10) : $\int_{\Omega} \epsilon(\mathbf{u}) * \epsilon(\mathbf{w})d\mathbf{x}$ and $\int_{\Omega} d^\alpha(\mathbf{x}, \partial\Omega) |\epsilon(\mathbf{u})| \epsilon(\mathbf{u}) * \epsilon(\mathbf{w})d\mathbf{x}$. In order to give meaning to these two integrals, it is sufficient to take \mathbf{u} and \mathbf{w} in the Sobolev spaces $H_0^1(\Omega)^3$ for the first integral and $(W_{d^\alpha}^{1,3}(\Omega))^3$ for the second one, where $W_{d^\alpha}^{1,3}(\Omega)$ denotes the weighted Sobolev space given by

$$W_{d^\alpha}^{1,3}(\Omega) = \{w \in L_{d^\alpha}^3(\Omega) \text{ and } \nabla w \in (L_{d^\alpha}^3(\Omega))^3\}, \tag{12}$$

$$\text{with } L_{d^\alpha}^3(\Omega) = \left\{ w : \int |w(\mathbf{x})|^3 d^\alpha(\mathbf{x}, \partial\Omega)d\mathbf{x} < \infty \right\}. \tag{13}$$

Note that for $\alpha = 0$, $W_{d^0}^{1,3}(\Omega) = W^{1,3}(\Omega) \subset H^1(\Omega)$ and equation (10) can be posed in $W_0^{1,3}(\Omega)^3$. More precisely, when $\alpha = 0$ and for $\mathbf{F} \in L^{3/2}(\Omega)^3$, we can look for $\mathbf{u} \in W_0^{1,3}(\Omega)^3$ and $p \in L^{3/2}(\Omega)$ satisfying (10) and (11) for every $\mathbf{w} \in W_0^{1,3}(\Omega)^3$ and $q \in L^{3/2}(\Omega)$. In [3] and [4] one can find a proof of existence and uniqueness of solution to this problem, with p defined up to a constant.

For $\alpha < 2$ and close to 2, the situation is not so obvious. The space $W_{d^\alpha}^{1,3}(\Omega)$ is not embedded in $H^1(\Omega)$ (see [5]), and we have to take a Banach space of type $H_0^1(\Omega)^3 \cap (W_{d^\alpha}^{1,3}(\Omega))^3$ in order to analyse equation (10).

Before considering problem (10) – (11), let us recall some results on weighted Sobolev spaces found in [5] and used in [1], [2]. For a 3×3 tensor κ depending on

$\mathbf{x} \in \Omega$, we adopt the notations:

$$\|\boldsymbol{\kappa}\|_{(L^2(\Omega))^{3 \times 3}} = \sum_{1 \leq i, j \leq 3} \left(\int_{\Omega} |\boldsymbol{\kappa}_{i,j}(\mathbf{x})|^2 d\mathbf{x} \right)^{1/2},$$

$$\|\boldsymbol{\kappa}\|_{(L^3_{d^\alpha}(\Omega))^{3 \times 3}} = \sum_{1 \leq i, j \leq 3} \left(\int_{\Omega} |\boldsymbol{\kappa}_{i,j}(\mathbf{x})|^3 d^\alpha(\mathbf{x}, \partial\Omega) d\mathbf{x} \right)^{1/3}.$$

- There exists $\alpha_0 \in]0, 2[$ such that for $\alpha \in [0, \alpha_0[$, $W_{d^\alpha}^{1,3}(\Omega) \subset H^1(\Omega)$, but for $\alpha > \alpha_0$, $W_{d^\alpha}^{1,3}(\Omega) \cap H^1(\Omega) \neq W_{d^\alpha}^{1,3}(\Omega)$.
- For $0 \leq \alpha < 2$, there is a trace operator $\Gamma : W_{d^\alpha}^{1,3}(\Omega) \rightarrow L^3(\partial\Omega)$ and $W_{d^\alpha,0}^{1,3}(\Omega) = \{w \in W_{d^\alpha}^{1,3}(\Omega) : \Gamma(w) = 0\}$.
- For $0 \leq \alpha < 2$, we have a Korn inequality in $W_{d^\alpha,0}^{1,3}(\Omega)$: there exists a positive constant β such that $\beta \|\nabla \mathbf{w}\|_{(L^3_{d^\alpha}(\Omega))^{3 \times 3}} \leq \|\epsilon(\mathbf{w})\|_{(L^3_{d^\alpha}(\Omega))^{3 \times 3}}$ for every $\mathbf{w} \in W_{d^\alpha,0}^{1,3}(\Omega)^3$ (see [6]).
- For $\alpha \geq 2$, a vanished trace of a function of $W_{d^\alpha}^{1,3}(\Omega)$ has no meaning and the problem of existence of a Korn inequality is open!

Starting from these considerations, we now consider only the case when $\alpha \in [0, 2[$. Let us define the Banach spaces

$$\mathbf{V} = H_0^1(\Omega)^3 \cap W_{d^\alpha,0}^{1,3}(\Omega)^3, \tag{14}$$

$$\mathbf{V}_{\text{div}} = \{\mathbf{w} \in \mathbf{V} : \text{div}(\mathbf{w}) = 0\}, \tag{15}$$

provided with the norm $\|\mathbf{w}\|_{\mathbf{V}} = \|\nabla \mathbf{w}\|_{(L^2(\Omega))^{3 \times 3}} + \|\nabla \mathbf{w}\|_{(L^3_{d^\alpha}(\Omega))^{3 \times 3}}$.

By taking $\mathbf{w} \in \mathbf{V}_{\text{div}}$ in (10), we have to look for $\mathbf{u} \in \mathbf{V}_{\text{div}}$ satisfying

$$\int_{\Omega} 2\nu\epsilon(\mathbf{u}) * \epsilon(\mathbf{w}) d\mathbf{x} = \int_{\Omega} \mathbf{F} \cdot \mathbf{w} d\mathbf{x} \quad \forall \mathbf{w} \in \mathbf{V}_{\text{div}}. \tag{16}$$

In order to give a meaning to the right side of equalities (10) or (16), we assume $\mathbf{F} \in L^2(\Omega)^3 \oplus L^{3/2}_{d^{-\alpha/2}}(\Omega)^3$.

Now we define the functional \mathbf{J} on \mathbf{V}_{div} by

$$\mathbf{J}(\mathbf{v}) = \int_{\Omega} [2A(\mathbf{x}, |\epsilon(\mathbf{v})|) - \mathbf{F} \cdot \mathbf{v}] d\mathbf{x}, \tag{17}$$

where

$$A(\mathbf{x}, s) = \nu_0 \frac{s^2}{2} + \beta_t \frac{s^3}{3} l^{2-\alpha} d^\alpha(\mathbf{x}, \partial\Omega). \tag{18}$$

It is easy to show that \mathbf{J} is Gâteaux-differentiable on \mathbf{V}_{div} , and its derivative at \mathbf{u} in direction \mathbf{w} is given by

$$D\mathbf{J}(\mathbf{u})(\mathbf{w}) = \int_{\Omega} 2\nu\epsilon(\mathbf{u}) * \epsilon(\mathbf{w}) d\mathbf{x} - \int_{\Omega} \mathbf{F} \cdot \mathbf{w} d\mathbf{x}, \tag{19}$$

with ν depending on \mathbf{u} via (6). It follows that if $\mathbf{u} \in \mathbf{V}_{\text{div}}$ is solution of (16), then $D\mathbf{J}(\mathbf{u})(\mathbf{w}) = 0$ for every $\mathbf{w} \in \mathbf{V}_{\text{div}}$. In [2], it is proven that the functional \mathbf{J} is continuous, coercive and strictly convex on \mathbf{V}_{div} . Then, there exists a unique \mathbf{u} which

minimizes \mathbf{J} on \mathbf{V}_{div} and $D\mathbf{J}(\mathbf{u}) = 0$. Moreover, \mathbf{u} is the unique solution of Problem (16). The main result of this analysis is

PROPOSITION 1. *When $0 \leq \alpha < 2$, and ν is defined by (6), there exists a unique solution of Problem (16) (see [2]).*

Concerning the pressure p which appears in Problem (10) – (11), it is well known that an inf – sup condition allows to obtain its existence and uniqueness by knowing the velocity $\mathbf{u} \in \mathbf{V}_{\text{div}}$. In [2], it is shown that the dual space of $L^3_{d^\alpha}(\Omega)$ can be identified to $L^{3/2}_{d^{-\alpha/2}}(\Omega)$ by means of an integral on Ω . Since $\text{div}(\mathbf{w}) \in L^2(\Omega) \cap L^3_{d^\alpha}(\Omega)$, the pressure p will be considered in the space

$$Q = L^2_0(\Omega) \oplus L^{3/2}_{d^{-\alpha/2,0}}(\Omega),$$

in order to give meaning to the second integral of (10). The index 0 in $L^2_0(\Omega)$ and $L^{3/2}_{d^{-\alpha/2,0}}(\Omega)$ signifies that the mean values of functions in these spaces are zero.

When α is close to 2 (recall that $\alpha = 2$ in von Karman modeling), the space $L^2(\Omega) \cap L^{3/2}_{d^{-\alpha/2}}(\Omega)$ is not reduced to the null space, and we are not able to verify the inequality $\inf_{q \in Q} \sup_{\mathbf{w} \in \mathbf{V}} \frac{\int_\Omega q \text{div}(\mathbf{w}) dx}{\|q\|_Q \|\mathbf{w}\|_{\mathbf{V}}} > 0$. As a consequence, even if we are looking for a pressure p defined up to a constant in Problem (10) – (11), we are not able to directly verify its existence and uniqueness.

However, the existence of a pressure satisfying (10) – (11) with the velocity \mathbf{u} given above, can be deduced from two inf – sup conditions. The first inf – sup is the pairing $H^1_0(\Omega)^3 - L^2_0(\Omega)$ and the other one is $W^{1,3}_{d^\alpha,0}(\Omega)^3 - L^{3/2}_{d^{-\alpha/2,0}}(\Omega)$, (see [2]). Due to these two inf-sup conditions, we can prove the existence and uniqueness of $\bar{p}_1 \in L^2_0(\Omega)$ and $\bar{p}_2 \in L^{3/2}_{d^{-\alpha/2,0}}(\Omega)$ satisfying:

$$\int_\Omega \bar{p}_1 \text{div}(\mathbf{w}) dx = \int_\Omega 2\nu_0 \epsilon(\mathbf{u}) * \epsilon(\mathbf{w}) dx - \int_\Omega \mathbf{F} \cdot \mathbf{w} dx, \quad \forall \mathbf{w} \in H^1_0(\Omega)^3,$$

$$\int_\Omega \bar{p}_2 \text{div}(\mathbf{w}) dx = \int_\Omega 2\beta_t |t|^{2-\alpha} d^\alpha(x, \partial\Omega) |\epsilon(\mathbf{u})| \epsilon(\mathbf{u}) * \epsilon(\mathbf{w}) dx, \quad \forall \mathbf{w} \in W^{1,3}_{d^\alpha,0}(\Omega)^3.$$

By setting $p = \bar{p}_1 + \bar{p}_2$, then $(\mathbf{u}, p) \in \mathbf{V} \times Q$ satisfies relations (10)-(11) for every $(\mathbf{w}, q) \in \mathbf{V} \times Q$.

These considerations lead to:

PROPOSITION 2 (see [2]). *When $0 \leq \alpha < 2$, and ν is defined by (6), there exists a solution $(\mathbf{u}, p) \in \mathbf{V} \times Q$ satisfying (10) – (11) for every $(\mathbf{w}, q) \in \mathbf{V} \times Q$. Moreover, if (\mathbf{u}_1, p_1) and (\mathbf{u}_2, p_2) are two solutions, then $\mathbf{u}_1 = \mathbf{u}_2$. When $\alpha \in [0, \alpha_0[$, i.e when $W^{1,3}_{d^\alpha}(\Omega) \subset H^1(\Omega)$, then $p_1 = p_2$.*

3. Conclusions and remarks.

- When $\nu_0 = 0$ in (6), then $H^1_0(\Omega)^3$ can be dropped in the definition of \mathbf{V} in order to set Problem (10) – (11). It follows that for $0 \leq \alpha < 2$, the functional J can be minimized on $X_{\text{div}} = \{\mathbf{v} \in W^{1,3}_{d^\alpha,0}(\Omega)^3 : \text{div}(\mathbf{v}) = 0\}$. In the case $\nu_0 = 0$, for every $\alpha \in [0, 2[$, Problem (10) – (11) possesses a unique solution $(\mathbf{u}, p) \in W^{1,3}_{d^\alpha,0}(\Omega)^3 \times L^{3/2}_{d^{-\alpha/2,0}}(\Omega)$.

- For $\alpha = 2$, there is no trace of \mathbf{u} on $\partial\Omega$ when $\mathbf{u} \in W_{d^2}^{1,3}(\Omega)^3$. In this case it is not possible to take $\nu_0 = 0$ when $\mathbf{u} = 0$ is imposed on the boundary $\partial\Omega$. Moreover Korn's inequality is probably wrong in $W_{d^2}^{1,3}(\Omega)^3$ (Kalamajska's conjecture [6]). The analysis of von Karman model ($\alpha = 2$) remains open.
- A direct consequence of previous remarks is: when using a finite element approximation on von Karman model ($\alpha = 2$) and when the laminar viscosity ν_0 is "small" with respect to the numerical viscosity of the method, **the obtained results can strongly depend on the mesh!** In this case, it is necessary to take a very thin meshing close to the walls in order to remedy to this situation. Numerical examples are given in [1] and [2].

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REFERENCES

- [1] J. RAPPAZ, J. ROCHAT, *On non-linear Stokes problems with viscosity depending on the distance to the wall*, C.R. Acad. Sci. Paris, Ser.I 354 (2016), pp. 499–502.
- [2] J. RAPPAZ, J. ROCHAT, *On some weighted Stokes problems. Application on Smagorinsky models*, Contributions to Partial Differential Equations and Applications, Computational Methods in Applied Sciences, vol. 47, pp. 395–410, Springer (2018).
- [3] J. BARANGER AND K. NAJIB, *Analyse numérique des écoulements quasi-Newtoniens dont la viscosité obéit à la loi de puissance ou la loi de Carreau*, Numer. Math., 58 (1990), pp. 35–49.
- [4] D. SANDRI, *Sur l'approximation numérique des écoulements quasi-newtoniens dont la viscosité suit la loi de puissance ou la loi de Carreau*, Modélisation mathématique et analyse numérique, tome 27:2 (1993), pp. 131–155.
- [5] F. KUFNER, *Weighted Sobolev Spaces*, Elsevier, Paris, 1983.
- [6] A. KALAMAJSKA, *Coercive inequalities on weighted Sobolev spaces*, Colloquium Mathematicum, Vol. LXVI, 1993, FASC.

