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1 HIGH ORDER HOMOGENIZATION OF THE STOKES SYSTEM IN 2 A PERIODIC POROUS MEDIUM*

FLORIAN FEPPON[†]

Abstract. We derive high order homogenized models for the incompressible Stokes system in 4 5 a cubic domain filled with periodic obstacles. These models have the potential to unify the three classical limit problems (namely the "unchanged" Stokes system, the Brinkman model, and the 6 Darcy's law) corresponding to various asymptotic regimes of the ratio $\eta \equiv a_{\varepsilon}/\varepsilon$ between the radius a_{ε} of the holes and the size ε of the periodic cell. What is more, a novel, rather surprising feature 8 9 of our higher order effective equations is the occurrence of odd order differential operators when the obstacles are not symmetric. Our derivation relies on the method of two-scale power series expansions and on the existence of a "criminal" ansatz, which allows to reconstruct the oscillating velocity and 11 pressure $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ as a linear combination of the derivatives of their formal average $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$ weighted 13 by suitable corrector tensors. The formal average $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$ is itself the solution to a formal, infinite 14 order homogenized equation, whose truncation at any finite order is in general ill-posed. Inspired by the variational truncation method of [53, 27], we derive, for any $K \in \mathbb{N}$, a well-posed model of order 2K + 2 which yields approximations of the original solutions with an error of order $O(\varepsilon^{K+3})$ 16in the L^2 norm. Furthermore, the error improves up to the order $O(\varepsilon^{2K+4})$ if a slight modification 17 of this model remains well-posed. Finally, we find asymptotics of all homogenized tensors in the low 18 volume fraction limit $\eta \to 0$ and in dimension $d \ge 3$. This allows us to obtain that our effective 19equations converge coefficient-wise to either of the Brinkman or Darcy regimes which arise when η is respectively equivalent, or greater than the critical scaling $\eta_{\rm crit} \sim \varepsilon^{2/(d-2)}$. 2021

22 Key words. Homogenization, higher order models, porous media, Stokes system, strange term.

23 AMS subject classifications. 35B27, 76M50, 35330

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1. Introduction. This article is concerned with the high order homogenization 24 of the Stokes system in a periodic porous medium. Let $D := (0, L)^d$ be a d-dimensional 25box filled with periodic obstacles $\omega_{\varepsilon} := \varepsilon(\mathbb{Z}^d + \eta T) \cap D$ (the setting is illustrated on 26Figure 1). The parameter ε denotes the size of the periodic cell, it is equal to $\varepsilon := L/N$ 27 where $N \in \mathbb{N}$ is a large integer and L is the length of the box. The parameter η is the 28 scaling ratio between the radius $a_{\varepsilon} := \eta \varepsilon$ of the obstacles and the length ε of the cells. 29The total fluid domain is denoted by $D_{\varepsilon} := D \setminus \overline{\omega_{\varepsilon}}$ and it is assumed to be connected. 30 $P = (0,1)^d$ is the unit cell and $Y = P \setminus \eta T$ denotes its fluid component. 31

32 We consider $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon}) \in H^1(D_{\varepsilon}, \mathbb{R}^d) \times L^2(D_{\varepsilon})/\mathbb{R}$ the solution to the Stokes system

33 (1.1) $\begin{cases} -\Delta \boldsymbol{u}_{\varepsilon} + \nabla p_{\varepsilon} = \boldsymbol{f} \text{ in } D_{\varepsilon} \\ \operatorname{div}(\boldsymbol{u}_{\varepsilon}) = 0 \\ \boldsymbol{u}_{\varepsilon} = 0 \text{ on } \partial \omega_{\varepsilon} \\ \boldsymbol{u}_{\varepsilon} \text{ is } D\text{-periodic,} \end{cases}$

where $\mathbf{f} \in \mathcal{C}_{per}^{\infty}(D, \mathbb{R}^d)$ (and all its derivatives) is a smooth, *D*-periodic right handside. The goal of this paper is to derive high order effective models for (1.1); i.e. a family of well-posed partial differential equations posed in the homogeneous domain *D* (without the holes) and whose solutions approximate the macroscopic behavior of $(\mathbf{u}_{\varepsilon}, p_{\varepsilon})$ at any desired order of accuracy in ε as $\varepsilon \to 0$.

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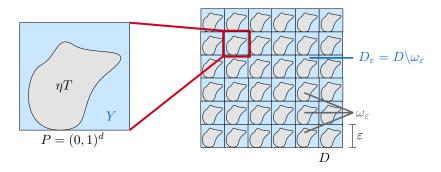


FIG. 1. The perforated domain $D_{\varepsilon} = D \setminus \omega_{\varepsilon}$ and the unit cell $Y = P \setminus (\eta T)$.

The literature [52, 47, 29, 5, 7, 4, 8] describes the occurrence of different asymp-39 totic regimes depending on how the size $a_{\varepsilon} = \eta \varepsilon$ of the holes compares to the critical size $\sigma_{\varepsilon} := \varepsilon^{d/(d-2)}$ in dimension $d \geq 3$ (if d = 2, then these regimes depend on 40 41 how $\log(a_{\varepsilon})$ compares to $-\varepsilon^{-2}$, see [7]). In loose mathematical terms, these can be 42 summarized as follows (see e.g. [5, 7] for the precise statements): 43

• if $a_{\varepsilon} = o(\sigma_{\varepsilon})$, then the holes have no effect and $(u_{\varepsilon}, p_{\varepsilon})$ converges as $\varepsilon \to 0$ to the solution (u, n) of the Stokes equation in the homogeneous domain D: 44 D45

to the solution
$$(\boldsymbol{u},p)$$
 of the Stokes equation in the homogeneous domain L

46 (1.2)
$$\begin{cases} -\Delta \boldsymbol{u} + \nabla \boldsymbol{p} = \boldsymbol{f} \text{ in } D \\ \operatorname{div}(\boldsymbol{u}) = 0 \\ \boldsymbol{u} \text{ is } D \text{-periodic.} \end{cases}$$

• if $a_{\varepsilon} = c\sigma_{\varepsilon}$ for a constant c > 0, then $(u_{\varepsilon}, p_{\varepsilon})$ converges as $\varepsilon \to 0$ to the 47 solution (\boldsymbol{u}, p) of the Brinkman equation 48

49 (1.3)
$$\begin{cases} -\Delta \boldsymbol{u} + cF\boldsymbol{u} + \nabla p = \boldsymbol{f} \text{ in } D \\ \operatorname{div}(\boldsymbol{u}) = 0 \\ \boldsymbol{u} \text{ is } D \text{-periodic,} \end{cases}$$

- where the so-called strange term cFu involves a symmetric positive definite 50 $d \times d$ matrix F which can be computed by means of an exterior problem in 51 $\mathbb{R}^d \setminus T$ (see [4] and section 5).
- if $\sigma_{\varepsilon} = o(a_{\varepsilon})$ and $a_{\varepsilon} = \eta \varepsilon$ with $\eta \to 0$ as $\varepsilon \to 0$, then the holes are "large" and $(a_{\varepsilon}^{d-2}\varepsilon^{-d}\boldsymbol{u}_{\varepsilon},p_{\varepsilon})$ converges to the solution (\boldsymbol{u},p) of the Darcy problem 54

55 (1.4)
$$\begin{cases} F\boldsymbol{u} + \nabla p = \boldsymbol{f} \text{ in } D \\ \operatorname{div}(\boldsymbol{u}) = 0 \text{ in } D \\ \boldsymbol{u} \text{ is } D \text{-periodic,} \end{cases}$$

- where F is the same symmetric positive definite $d \times d$ matrix as in (1.3). 56
- if $a_{\varepsilon} = \eta \varepsilon$ with the ratio η fixed, then $(\varepsilon^{-2} u_{\varepsilon}, p_{\varepsilon})$ converges to the solution 57 (\boldsymbol{u}, p) of the Darcy problem 58

59 (1.5)
$$\begin{cases} M^0 \boldsymbol{u} + \nabla p = \boldsymbol{f} \text{ in } D\\ \operatorname{div}(\boldsymbol{u}) = 0 \text{ in } D\\ \boldsymbol{u} \text{ is } D\text{-periodic,} \end{cases}$$

60 where M^0 is another positive symmetric $d \times d$ matrix (which depends on η). 61 Furthermore $M^0/|\log(\eta)| \to F$ if d = 2, and $M^0/\eta^{d-2} \to F$ (if $d \ge 3$) when 62 $\eta \to 0$, so that there is a continuous transition from (1.5) to (1.4), see [6].

One of the long-term motivations driving this work is the need to lay down theo-63 retical material that would allow to optimize the design of fluid systems by homog-64 enization methods similar to those available in the context of mechanical structures 65 [21, 20, 10, 50, 14]. To date, the Brinkman [24, 25, 30] and the Darcy models [56, 51] 66 are commonly used by topology optimization algorithms in order to conveniently in-67 terpolate the physics of the fluid at intermediate "gray" regions featuring locally a 68 mixture of fluid and solid. However, the above conclusions imply that these models 69 are consistent only in specific ranges of obstacle sizes a_{ε} : the Brinkman model (1.3) 70 71is relevant when there are none or tiny obstacles, while the Darcy models (1.4) and (1.5) should be used at locations where the obstacles are large enough. The arising 72of these different regimes (1.2)-(1.5) is consequently a major obstacle towards the de-73 velopment of 'de-homogenization' methods [14, 37, 50, 39, 40] for the optimal design 74of fluid systems, which would enable to interpret "gray" designs as locally periodic 75 "black and white" microstructures (featuring for instance many small tubes or thin 76 plates). 77

It turns out that there is a continuous transition between these regimes which can be captured by higher order homogenized equations, which is the object of the present article. These higher order models are obtained by adding corrective terms scaled by increasing powers of ε to the Darcy equation (1.5); they yield more accurate approximations of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ when ε is "not so small". For a desired order $K \in \mathbb{N}$, the homogenized model of order 2K + 2 reads

84 (1.6)
$$\begin{cases} \sum_{k=0}^{2K+2} \varepsilon^{k-2} \mathbb{D}_{K}^{k} \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon,K}^{*} + \nabla q_{\varepsilon,K}^{*} = \boldsymbol{f}, \\ \operatorname{div}(\boldsymbol{v}_{\varepsilon,K}^{*}) = 0, \\ \boldsymbol{v}_{\varepsilon,K}^{*} \text{ is } D\text{-periodic} \end{cases}$$

where $(\boldsymbol{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*)$ is a high order homogenized approximation of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$. The coefficient \mathbb{D}_{K}^{k} is a k-th order matrix valued tensor which can be computed by a procedure involving the resolution of cell problems; it makes $\mathbb{D}_{K}^{k} \cdot \nabla^{k}$ a differential operator of order k (the notation is defined in section 2 below). Finally, the high order equation (1.6) encompasses at least the Brinkman and the Darcy regimes in the sense that it converges coefficient-wise to either of (1.3) and (1.4) for the corresponding asymptotic regime of the scaling η (see Remarks 5.6 and 5.7) (the analysis of the subcritical case leading to the Stokes regime (1.2) requires more sophisticated arguments which are to be investigated in future works).

A rather striking feature of (1.6) is the arising of *odd* order differential operators 94 (these vanish, however, in case the obstacle ηT is symmetric with respect to the cell 95 axes; see Corollary 3.16). This fact is closely related to the vectorial nature of the 96 Stokes system (1.1): the tensors \mathbb{D}_{K}^{k} are symmetric and antisymmetric valued matrices 97 for respectively even and odd values of k. This property ensures that eventually, 98 $\mathbb{D}_{K}^{k} \cdot \nabla^{k}$ is a symmetric operator for any $0 \leq k \leq 2K + 2$ (see Remark 3.12). To our 99 knowledge, such terms have so far not been proposed in the literature seeking similar 100 higher order corrections for the Stokes system, although these have been observed 101 in other vectorial contexts [27, 28, 53]. Most of the available works have focused on 102103 situations with low regularity for f, T and D (see [52, 5]), where the homogenization

process can be justified only for the approximation at the leading order in ε . Error bounds for higher order approximations of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ (namely for the truncation of the ansatz (1.7) below) have been obtained in [46, 26], without relating these to effective models. A few additional works have sought corrector terms from physical modelling

108 considerations [35, 18, 17], without considering odd order operators.

Our derivation is inspired from the works [19, 53, 15]; it is based on (non stan-109 110 dard) two-scale asymptotic expansions and formal operations on related power series which give rise to several families of tensors and homogenized equations at any order. 111 We extend our previous works [34, 33] where we investigated the cases of the perfo-112 rated Poisson problem and of the perforated elasticity system. Expectedly, the major 113difficulty in extending the analysis to (1.6) is the treatment of the pressure variable 114 p_{ε} and of the incompressibility constraint div $(\boldsymbol{u}_{\varepsilon}) = 0$. Note that the D-periodicity 115assumption on f and u_{ε} is made in order to eliminate additional difficulties related 116to the arising of boundary layers (see [43, 22, 23, 11]). 117

118 The starting point of the method of two-scale expansions is to postulate an ansatz 119 for the velocity and pressure solution $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$:

120 (1.7)
$$\boldsymbol{u}_{\varepsilon}(x) = \sum_{i=0}^{+\infty} \varepsilon^{i+2} \boldsymbol{u}_i(x, x/\varepsilon), \quad p_{\varepsilon}(x) = \sum_{i=0}^{+\infty} \varepsilon^i(p_i^*(x) + \varepsilon p_i(x, x/\varepsilon)), \quad x \in D_{\varepsilon},$$

where the functions $u_i(x, y)$ and $p_i(x, y)$ are *P*-periodic with respect to $y \in P$, and *D*-periodic with respect to $x \in D$. In (1.7), the oscillating function $p_i(x, y)$ is required to be of zero average with respect to y:

$$\int_{Y} p_i(x, y) \mathrm{d}y = 0, \quad \forall i \ge 0.$$

- 121 The aim of the homogenization process is to obtain effective equations for the formal
- 122 "infinite order" homogenized averages $\boldsymbol{u}^*_{\varepsilon}$ and p^*_{ε} defined by

123 (1.8)
$$\boldsymbol{u}_{\varepsilon}^{*}(x) := \sum_{i=0}^{+\infty} \varepsilon^{i+2} \int_{Y} \boldsymbol{u}_{i}(x,y) \mathrm{d}y, \quad p_{\varepsilon}^{*}(x) := \sum_{i=0}^{+\infty} \varepsilon^{i} p_{i}^{*}(x), \quad x \in D.$$

In Proposition 3.7 below, we obtain that $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$ solves the following formal "infiniteorder" homogenized equation,

126 (1.9)
$$\begin{cases} \sum_{k=0}^{+\infty} \varepsilon^{k-2} M^k \cdot \nabla^k \boldsymbol{u}_{\varepsilon}^* + \nabla p_{\varepsilon}^* = \boldsymbol{f}, \\ \operatorname{div}(\boldsymbol{u}^*) = 0, \\ \boldsymbol{u}_{\varepsilon}^* \text{ is } D\text{-periodic}, \end{cases}$$

which involves a family of constant matrix-valued tensors $(M^k)_{k\in\mathbb{N}}$. Classically, truncating directly (1.9) yields, in general, an ill-posed model [12]. Several methods have been proposed to address this issue in order to obtain nonetheless well-posed higher order equations [16, 13, 1, 2, 15]. In our case, we adapt an idea from [53], whereby the coefficients \mathbb{D}_K^k are obtained thanks to a minimization principle (described in section 4) which makes indeed (1.6) well-posed. It is based on the existence of remarkable identities which relate the oscillating solution $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ to its formal average $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$:

134 (1.10)
$$\begin{cases} \boldsymbol{u}_{\varepsilon}(x) = \sum_{i=0}^{+\infty} \varepsilon^{i} N^{i}(x/\varepsilon) \cdot \nabla^{i} \boldsymbol{u}_{\varepsilon}^{*}(x) \\ p_{\varepsilon}(x) = p_{\varepsilon}^{*}(x) + \sum_{i=0}^{+\infty} \varepsilon^{i-1} \boldsymbol{\beta}^{i}(x/\varepsilon) \cdot \nabla^{i} \boldsymbol{u}_{\varepsilon}^{*}(x), \end{cases} \quad \forall x \in D_{\varepsilon},$$

where $(N^i(y))_{i\in\mathbb{N}}$ and $(\beta^i(y))_{i\in\mathbb{N}}$ are different families of respectively matrix valued and vector valued *P*-periodic tensors (of order *i*). The ansatz (1.10) is substantially different from (1.7); following [15], we call it "criminal" because the expansions of (1.10) depend on $\boldsymbol{u}_{\varepsilon}^*$ which is itself a formal power series in ε (eqn. (1.8)).

139 The order of accuracy at which the solution $(\boldsymbol{v}_{\varepsilon,K}^*, p_{\varepsilon,K}^*)$ yields an approximation of the original solution $(u_{\varepsilon}, p_{\varepsilon})$ is determined by how many leading coefficients of (1.6) 140and (1.9) coincide (Proposition 4.5). In Proposition 4.10, we show that $\mathbb{D}_{K}^{k} = M^{k}$ for $0 \leq k \leq K$, which allows to infer error estimates of order $O(\varepsilon^{K+3})$ in the $L^{2}(D)$ norm. 141 142It may seem disappointing that one needs to solve an equation of order 2K+2 in order 143 to obtain approximations of order $O(\varepsilon^{K+3})$ "only". This shortcoming is related to the 144 zero-divergence constraint: in the scalar and elasticity cases considered in [34, 33], it 145 turns out that K+1 extra coefficients coincide, namely $\mathbb{D}_{K}^{k} = M^{k}$ for $0 \leq k \leq 2K+1$, which yields error estimates of order $O(\varepsilon^{2K+4})$. In the present context devoted to the Stokes system (1.1), the equation obtained by substituting \mathbb{D}_{K}^{k} with M^{k} in (1.6) for 146147 148 $K+1 \le k \le 2K+1,$ 149

150 (1.11)
$$\begin{cases} \varepsilon^{2K} \mathbb{D}_{K}^{2K+2} \cdot \nabla^{2K+2} \widehat{\boldsymbol{v}}_{\varepsilon,K}^{*} + \sum_{k=0}^{2K+1} \varepsilon^{k-2} M^{k} \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon,K}^{*} + \nabla \widehat{q}_{\varepsilon,K}^{*} = \boldsymbol{f} \\ \operatorname{div}(\widehat{\boldsymbol{v}}_{\varepsilon,K}^{*}) = 0 \\ \widehat{\boldsymbol{v}}_{\varepsilon,K}^{*} \text{ is } D\text{-periodic,} \end{cases}$$

corresponds to applying the truncation method of [53] to the mixed variational formulation rather than to the minimization problem associated with (1.1) (see Remark 4.11). While the minimization principle ensures that (1.6) is well-posed, we do not know whether this is the case for (1.11). However if it is, then Proposition 4.5 implies that (1.11) improves the approximation accuracy up to the order $O(\varepsilon^{2K+4})$.

The article outlines as follows. Notation conventions related to tensors and technical assumptions are exposed in section 2.

In section 3, we introduce cell problems and their solution tensors $(\mathcal{X}^k, \boldsymbol{\alpha}^k)$ which 158allow to identify the functions u_i , p_i^* and p_i in the ansatz (1.7). We show that the 159formal average $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$ solves the infinite order homogenized equation (1.9) involving the tensors M^k . After defining the tensors $N^k(y)$ and $\beta^k(y)$, we derive the "criminal" 161ansatz (1.10) expressing $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ in terms of p_{ε}^* and of the derivatives of $\boldsymbol{u}_{\varepsilon}^*$. Through-162out this section, a number of algebraic properties are stated for the various tensors 163coming at play, such as the symmetry and the antisymmetry of the matrix valued 164tensors M^k for respectively even and odd values of k, and the simplifications taking 165place in case the obstacle ηT is symmetric with respect to the cell axes. 166

167 Section 4 details the truncation process of the infinite order equation (1.9) leading 168 to the well-posed model (1.6). We then provide an error analysis of the homogenized 169 approximations of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ generated by our procedure: our main result is stated in

Corollary 4.15 where we show that the solution $(\boldsymbol{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*)$ of (1.6) yield approxima-170 tions of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ in the $L^2(D_{\varepsilon})$ norm of order K+3 and K+1 for the velocity and 171 the pressure respectively. We establish explicit formulas relating the coefficients \mathbb{D}^k_K 172to the coefficients M^k and we briefly discuss the improvement provided by (1.11) in 173 174case it is well-posed. The last section 5 investigates asymptotics of the tensors M^k in the low volume 175fraction limit where the scaling of the obstacle η converges to zero. Our main result is 176Corollary 5.5 where we obtain the "coefficient-wise" convergence of the infinite order 177 homogenized equation as well as the one of (1.6) towards either of the Brinkman or 178

Darcy regimes (1.4) and (1.5) when η is respectively equivalent or greater than the critical size $\eta_{\rm crit} \sim \varepsilon^{2/(d-2)}$, and towards the Stokes regime (1.3) for $\eta = o(\varepsilon^{2/(d-2)})$ 179180 in the case K = 0. Although our error estimates for (1.6), are a priori not uniform 181 in η , this suggests that our higher order model (1.6) has the potential to yield valid 182approximations in any regime of size of holes (at least for K = 0 or above the critical 183scale). Note that our analysis is unfortunately unsufficient to establish the convergence 184of the high order coefficients $\varepsilon^{k-2}M^k$ with k>2 towards 0 as $\eta \to 0$. Future works 185 will investigate higher order asymptotics of the tensors M^k in the subcritical regime 186 $\eta = o(\varepsilon^{2/(d-2)})$ which are required to establish or invalidate such a claim. 187

188 **2.** Setting and notation conventions related to tensors. In the sequel, we 189 consider the following two classical assumptions for the distributions of the holes ω_{ε} 190 (we recall the schematic of Figure 1), following [5]:

191 (H1) $Y = P \setminus (\eta T) \subset P$, as a subset of the unit torus (opposite matching faces 192 of $(0, 1)^d$ are identified) is a smooth connected set with non-empty interior.

193 (H2) The fluid component $D_{\varepsilon} = D \setminus \omega_{\varepsilon}$ is a smooth connected set.

194 Remark 2.1. Assumption (H1) does not necessarily imply (H2), see [3] for a coun-195 terexample. Assumption (H1) is not very restrictive and can easily be generalized to 196 the case where the subset Y has m connected components with $m \in \mathbb{N}$ (see Appendix 197 7.5.6 in [33]). Assumption (H2) is stronger, but is also more connected to physical 198 applications. It forbids the existence of isolated fluid inclusions. Most of our deri-199 vations only assume (H1). However, we rely on both assumptions (H1) and (H2) in 200 order to obtain error bounds section 4, because we use some technical results of [5].

201 Below and further on, we consider scalar and vectorial functions such as

202 (2.1)
$$u : D \times P \rightarrow \mathbb{R}$$
, $u : D \times P \rightarrow \mathbb{R}^d$
 $(x,y) \mapsto u(x,y)$, $(x,y) \mapsto u(x,y)$

which are both D and P-periodic with respect to respectively the first and the second variable, and which vanish on the hole $D \times (\eta T)$. The arguments x and y of u(x, y) are respectively called the "slow" and the "fast" or "oscillating" variable. With a small abuse of notation, the partial derivative with respect to the variable y_j (respectively x_j) is simply written ∂_j instead of ∂_{y_j} (respectively ∂_{x_j}) when the context is clear.

The star—"*"– symbol is used to indicate that a quantity is "macroscopic" in the sense that it does not depend on the fast variable x/ε ; e.g. $(\boldsymbol{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*)$ or $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$ in (1.6) and (1.9). In the particular case where a two-variable quantity u(x, y) is given such as (2.1), $u^*(x)$ always denotes the average of $y \mapsto u(x, y)$ with respect to the y variable:

$$u^*(x) := \int_P u(x,y) \mathrm{d}y = \int_Y u(x,y) \mathrm{d}y, \qquad x \in D,$$

where the last equality is a consequence of u vanishing on $P \setminus Y = \overline{\eta T}$. When a function $\mathcal{X} : P \to \mathbb{R}$ depends only on the y variable, we find occasionally more convenient to write its cell average with the usual angle bracket symbols:

$$\langle \mathcal{X} \rangle := \int_P \mathcal{X}(y) \mathrm{d}y.$$

In all what follows, unless otherwise specified, the Einstein summation convention over repeated *subscript* indices is assumed (but never on *superscript* indices). Vectors $\boldsymbol{b} \in \mathbb{R}^d$ are written in bold face notation.

The notation conventions used for tensor related operations are summarized in the nomenclature below. Some of them are not standard; they allow to avoid to systematically write partial derivative indices (e.g. $1 \le i_1 \ldots i_k \le d$) and to distinguish them from spatial indices (e.g. $1 \le l, m \le d$) associated with vector or matrix components.

Scalar, vector, and matrix valued tensors and their coordinates 215Vector of \mathbb{R}^d 216 b $_{b^k}^{(b_j)_{1\leq j\leq d}}$ Coordinates of the vector \boldsymbol{b} 217 Scalar valued tensor of order k $(b_{i_1...i_k}^k \in \mathbb{R} \text{ for } 1 \leq i_1, \ldots, i_k \leq d)$ Vector valued tensor of order k $(b_{i_1...i_k}^k \in \mathbb{R}^d \text{ for } 1 \leq i_1, \ldots, i_k \leq d)$ Matrix valued tensor of order k $(B_{i_1...i_k}^k \in \mathbb{R}^{d \times d} \text{ for } 1 \leq i_1, \ldots, i_k \leq d)$ Coordinates of the vector valued tensor \boldsymbol{b}^k $(b_j^k \text{ is a scalar tensor of order } \boldsymbol{b}^k)$ 218 b^k 219 B^k 220 $(b_i^k)_{1 \le j \le d}$ 221 k). 2.2.2 $(B_{lm}^k)_{1 \le l,m \le d}$ Coefficients of the matrix valued tensor B^k $(B_{lm}^k)_{1 \le l,m \le d}$ is a scalar tensors of 223 2.2.4 order k). Coefficient of the vector valued tensor \boldsymbol{b}^k $(1 \leq i_1, \ldots i_k, j \leq d)$ $\begin{array}{c} b_{i_1\dots i_k,j}^k \\ B_{i_1\dots i_k,lm}^k \end{array}$ 225Coefficients of the matrix valued tensor B^k $(1 \le i_1, \ldots, i_k, l, m \le d)$ 226**Tensor products** 227 $b^p \otimes c^{k-p}$ Tensor product of scalar tensors b^p and c^{k-p} : 228 $(b^p \otimes c^{k-p})_{i_1 \dots i_k} := b^p_{i_1 \dots i_n} c^{k-p}_{i_{n+1} \dots i_k}.$ (2.2)229 Tensor product of a scalar tensors a^p and a vector valued tensor \boldsymbol{b}^{k-p} : $a^p \otimes \boldsymbol{b}^{k-p}$ 230 $(a^p \otimes \boldsymbol{b}^{k-p})_{i_1 \dots i_k} := a^p_{i_1 \dots i_p} \boldsymbol{b}^{k-p}_{i_{p+1} \dots i_k}$ (2.3)231 $B^p \otimes C^{k-p}$ Tensor product of matrix valued tensors B^p and C^{k-p} : 232 $(B^p \otimes C^{k-p})_{i_1...i_k,lm} := B^p_{i_1...i_n,l_j} C^{k-p}_{i_{n+1}...i_k,jm}.$ (2.4)233Hence a matrix product is implicitly assumed in the notation $B^p \otimes C^{k-p}$. 234 $B^p: C^{k-p}$ Tensor product and Frobenius product of matrix tensors B^p and C^{k-p} : 235 $(B^{p}:C^{k-p})_{i_{1}\ldots i_{k}}:=B^{p}_{i_{1}\ldots i_{p},lm}C^{k-p}_{i_{p+1}\ldots i_{k},lm}.$ (2.5)236Tensor product and inner product of vector valued tensors \boldsymbol{b}^p and \boldsymbol{c}^{k-p} : $\boldsymbol{b}^p \cdot \boldsymbol{c}^{k-p}$ 237 $(\boldsymbol{b}^{p} \cdot \boldsymbol{c}^{k-p})_{i_{1}...i_{k}} := b^{p}_{i_{1}...i_{n},m} c^{k-p}_{i_{n+1}...i_{k},m}.$ 238 (2.6)

239	$B^p \cdot c^{k-p}$	Tensor product of a matrix tensor B^p and a vector tensors c^{k-p} :		
240		(2.7)	$(B^{p} \cdot \boldsymbol{c}^{k-p})_{i_{1} \dots i_{k}, l} := B^{p}_{i_{1} \dots i_{p}, lm} c^{k-p}_{i_{p+1} \dots i_{k}, m}.$	
241		Hence a matri	ix-vector product is implicitly assumed in $B^p \cdot c^{k-p}$.	
242	Contraction	n with partial derivatives		
243	$b^k\cdot \nabla^k$	Differential operator of order k associated with a scalar tensor b^k :		
244		(2.8)	$b^k \cdot abla^k := b^k_{i_1 \dots i_k} \partial^k_{i_1 \dots i_k}.$	
245	$oldsymbol{b}^k \cdot abla^k$	Differential operator of order k associated with a vector tensor \boldsymbol{b}^k : for		
246		any smooth vector field $\boldsymbol{v} \in \mathcal{C}^{\infty}_{per}(D, \mathbb{R}^d),$		
247		(2.9)	$oldsymbol{b}^k \cdot abla^k oldsymbol{v} = b^k_{i_1 \ldots i_k, l} \partial^k_{i_1 \ldots i_k} v_l.$	
248 249	$B^k \cdot \nabla^k$	Differential operator of order k associated with a matrix valued tensor B^k : for any smooth vector field $\boldsymbol{v} \in \mathcal{C}^{\infty}_{per}(D, \mathbb{R}^d)$,		
250		(2.10)	$(B^k \cdot \nabla^k \boldsymbol{v})_l = B^k_{i_1 \dots i_k, lm} \partial^k_{i_1 \dots i_k} v_m.$	
251	Special tensors			
252	$(\hat{e}_j)_{1 < j < d}$	Vectors of the canonical basis of \mathbb{R}^d .		
253	e_j	Scalar valued tensor of order 1 given by $e_{j,i_1} := \delta_{i_1j}$ (with $1 \le j \le d$).		
	δ_{ij}	Kronecker symbol: $\delta_{ij} = 1$ if $i = j$ and $\delta_{ij} = 0$ if $i \neq j$.		
	I	Identity tensor of order 2:		
		$I_{i_1i_2} = \delta_{i_1i_2}.$		

255The identity tensor is another notation for the Kronecker tensor and it holds $I = e_j \otimes e_j$ with summation on the index $1 \leq j \leq d$. 256 J^{2k} Tensor of order 2k defined by:

$$J^{2k} := \overbrace{I \otimes I \otimes \cdots \otimes I}^{k \text{ times}}.$$

257

With a small abuse of notation, we consider zeroth order tensors b^0 to be constants 258(i.e. $b^0 \in \mathbb{R}$ if b^0 is scalar) and we still denote by $b^0 \otimes c^k := b^0 c^k$ the tensor product 259with a k-th order tensor c^k . The same convention also applies to vector valued and 260 matrix valued tensors. 261

In all what follows, a k-th order tensor b^k (scalar, vector or matrix valued) truly makes sense when contracted with k partial derivatives, as in (2.8)–(2.10). Therefore all the tensors considered throughout this work are identified to their symmetrization:

$$b_{i_1...i_k}^k \equiv \frac{1}{k!} \sum_{\sigma \in \mathfrak{S}_k} b_{i_{\sigma(1)}...i_{\sigma(k)}},$$

where \mathfrak{S}_k is the permutation group of order k. Consequently, the order in which the 262(derivative) indices i_1, \ldots, i_k are written in $b_{i_1...i_k}^k$ does not matter. 263

264 Finally, in the whole work, we write C, C_K or $C_K(f)$ to denote universal constants that do not depend on ε but whose values may change from lines to lines (and which 265may depend on η or on the obstacle T). 266

Remark 2.2. In a limited number of places, the superscript or subscript indices 267 $p,q \in \mathbb{N}$ are used. Naturally, these are not to be confused with the pressure variables 268 p_{ε} or q_{ε} introduced in (1.1). 269

3. Infinite order homogenized equation and criminal ansatz. We start 270by identifying the two-scale structure of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ which arise in the form of the ansatz 271272 (1.7). Because it helps emphasizing the arising of Cauchy products, we assume, in this section only, that the right-hand side f can be formally decomposed into a power 273series in ε : 274

275 (3.1)
$$\forall x \in D, \ \boldsymbol{f}(x) = \sum_{i=0}^{+\infty} \varepsilon^i \boldsymbol{f}^i(x).$$

3.1. Identification of the "classical" ansatz: tensors $(\mathcal{X}^k, \boldsymbol{\alpha}^k)$. Inserting 276(1.7) into the Stokes system (1.1) yields the following cascade of equations: 277

$$(3.2) \quad \begin{cases} -\Delta_{yy} \boldsymbol{u}_{i+2} + \nabla_y p_{i+2} = \boldsymbol{f}_{i+2} - \nabla_x p_{i+2}^* - \nabla_x p_{i+1} + \Delta_{xy} \boldsymbol{u}_{i+1} + \Delta_{xx} \boldsymbol{u}_i, \\ \operatorname{div}_y(\boldsymbol{u}_{i+2}) = -\operatorname{div}_x(\boldsymbol{u}_{i+1}), \\ \boldsymbol{u}_{-2} = \boldsymbol{u}_{-1} = 0, \ p_{-1} = 0, \\ \boldsymbol{u}_i(x, \cdot) = 0 \text{ on } \partial(\eta T) \\ \boldsymbol{u}_i(x, \cdot) \text{ is } P \text{-periodic for any } x \in D, \\ \boldsymbol{u}_i(\cdot, y) \text{ is } D \text{-periodic for any } y \in P, \end{cases}$$

for any $i \geq -2$, where the operators $-\Delta_{yy}$, $-\Delta_{xy}$, $-\Delta_{yy}$ are defined by

$$-\Delta_{xx} = -\operatorname{div}_x(\nabla_x \cdot), \quad -\Delta_{xy} = -\operatorname{div}_x(\nabla_y \cdot) - \operatorname{div}_y(\nabla_x \cdot), \quad -\Delta_{yy} := -\operatorname{div}_y(\nabla_y \cdot).$$

In order to solve (3.2), we introduce a family of respectively vector valued tensors 279

 $(\mathcal{X}_{j}^{k}(y))_{1 \leq j \leq d}$ and scalar valued tensors $(\alpha_{j}^{k}(y))_{1 \leq j \leq d}$ defined by induction as the unique solutions in $H_{per}^{1}(Y, \mathbb{R}^{d}) \times L^{2}(Y)/\mathbb{R}$ to the following cell problems: 280

281

282 (3.3)
$$\begin{cases} -\Delta_{yy} \mathcal{X}_j^0 + \nabla_y \alpha_j^0 = \boldsymbol{e}_j \text{ in } Y, \\ \operatorname{div}_y(\mathcal{X}_j^0) = 0 \text{ in } Y \end{cases}$$

283 (3.4)
$$\begin{cases} -\Delta_{yy} \mathcal{X}_j^1 + \nabla_y \alpha_j^1 = (2\partial_l \mathcal{X}_j^0 - \alpha_j^0 \boldsymbol{e}_l) \otimes \boldsymbol{e}_l \text{ in } Y \\ \operatorname{div}_y(\mathcal{X}_j^1) = -(\mathcal{X}_j^0 - \langle \mathcal{X}_j^0 \rangle) \cdot \boldsymbol{e}_l \otimes \boldsymbol{e}_l \text{ in } Y \end{cases}$$

$$\begin{cases} 284 \quad (3.5) \\ 285 \end{cases} \begin{cases} -\Delta_{yy} \boldsymbol{\mathcal{X}}_{j}^{k+2} + \nabla_{y} \alpha_{j}^{k+2} = (2\partial_{l} \boldsymbol{\mathcal{X}}_{j}^{k+1} - \alpha_{j}^{k+1} \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \boldsymbol{\mathcal{X}}_{j}^{k} \otimes \boldsymbol{I} \text{ in } \boldsymbol{Y} \\ \text{div}_{y} (\boldsymbol{\mathcal{X}}_{j}^{k+2}) = -(\boldsymbol{\mathcal{X}}_{j}^{k+1} - \langle \boldsymbol{\mathcal{X}}_{j}^{k+1} \rangle) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} \text{ in } \boldsymbol{Y} \end{cases} \forall k \ge 0.$$

Equations (3.3)-(3.5) are supplemented with the following boundary conditions: 286

287 (3.6)
$$\begin{cases} \int_{Y} \alpha_{j}^{k} dy = 0\\ \boldsymbol{\mathcal{X}}_{j}^{k} = 0 \text{ on } \partial(\eta T) \\ (\boldsymbol{\mathcal{X}}_{j}^{k}, \alpha_{j}^{k}) \text{ is } P\text{-periodic} \end{cases} \quad \forall k \ge 0.$$

288

Remark 3.1. In view of the notation conventions of section 2, the non bold symbols $\otimes e_I$ and $\otimes I$ indicate the arising of extra partial derivatives indices. For instance, the first line of (3.5) must be understood as

$$-\Delta_{yy} \boldsymbol{\mathcal{X}}_{j,i_1,\ldots,i_{k+2}}^{k+2} + \nabla \alpha_{j,i_1,\ldots,i_{k+2}}^{k+2} = 2\partial_{i_{k+2}} \boldsymbol{\mathcal{X}}_{j,i_1,\ldots,i_{k+1}}^{k+1} - \alpha_{j,i_1,\ldots,i_{k+1}}^{k+1} \boldsymbol{e}_{i_{k+2}} + \boldsymbol{\mathcal{X}}_{j,i_1,\ldots,i_k}^k \delta_{i_{k+1}i_{k+2}}.$$

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We introduce the k-th order matrix valued tensors \mathcal{X}^k whose columns are the vector valued tensors (\mathcal{X}_i^k) :

$$(\mathcal{X}_{ij}^k(y))_{1 \le i,j \le d} := \begin{bmatrix} \mathcal{X}_1^k(y) & \dots & \mathcal{X}_d^k(y) \end{bmatrix}, \forall y \in Y, \quad \forall k \ge 0.$$

We also denote by $\boldsymbol{\alpha}^k$ the k-th order vector valued tensor whose coordinates are the scalar tensors α_i^k :

$$\boldsymbol{\alpha}^k(y) := (\alpha_j^k(y))_{1 \le j \le d}, \, \forall y \in Y, \quad \forall k \ge 0.$$

Following the conventions of section 2, we use a star notation to denote the average of respectively the tensor \mathcal{X}^k and of the vector fields u_i :

291 (3.7)
$$\mathcal{X}^{k*} := \int_Y \mathcal{X}^k(y) \mathrm{d}y, \, \forall k \ge 0, \qquad \boldsymbol{u}_i^*(x) := \int_Y \boldsymbol{u}_i(x, y) \mathrm{d}y, \, \forall x \in D, \, \forall i \ge 0.$$

292 The tensors \mathcal{X}^k and $\boldsymbol{\alpha}^k$ enable to solve the cascade of equations (3.2):

i

293 PROPOSITION 3.2. Assume (H1). The solutions $u_i(x, y)$, $p_i(x, y)$ of the cascade 294 of equations (3.2) are given by

295 (3.8)

$$u_i(x,y) = \sum_{k=0} \mathcal{X}^k(y) \cdot \nabla^k (f_{i-k}(x) - \nabla p_{i-k}^*(x))$$
$$p_i(x,y) = \sum_{k=0}^i \alpha^k(y) \cdot \nabla^k (f_{i-k}(x) - \nabla p_{i-k}^*(x)),$$

where the functions p_i^* are uniquely determined recursively as the solutions to the following elliptic system: for any $i \ge 0$,

$$\begin{cases} -\operatorname{div}_{x}(\mathcal{X}^{0*}\nabla_{x}p_{i}^{*}) = -\operatorname{div}_{x}(\mathcal{X}^{0*}\boldsymbol{f}_{i}) \\ & -\sum_{k=1}^{i}\operatorname{div}(\mathcal{X}^{k*}\cdot\nabla^{k}(\boldsymbol{f}_{i-k}-\nabla_{x}p_{i-k}^{*})) \text{ in } D_{\varepsilon}, \\ & \int_{D}p_{i}^{*}\operatorname{d}x = 0 \\ & p_{i}^{*} \text{ is } D\text{-periodic.} \end{cases}$$

Recognizing Cauchy products, the identities (3.8) and (3.9) rewrite formally in terms of equality of formal power series:

301 (3.10)
$$\boldsymbol{u}_{\varepsilon}(x) = \sum_{i=0}^{+\infty} \varepsilon^{i+2} \mathcal{X}^{i}(x/\varepsilon) \cdot \nabla^{i}(\boldsymbol{f}(x) - \nabla p_{\varepsilon}^{*}(x)),$$

302 (3.11)
$$p_{\varepsilon}(x) = p_{\varepsilon}^{*}(x) + \sum_{i=0}^{+\infty} \varepsilon^{i+1} \boldsymbol{\alpha}^{i}(x/\varepsilon) \cdot \nabla^{i}(\boldsymbol{f}(x) - \nabla p_{\varepsilon}^{*}(x)),$$

303 (3.12)
$$\operatorname{div}(\boldsymbol{u}_{\varepsilon}^{*}(x)) = 0 \text{ where } \boldsymbol{u}_{\varepsilon}^{*}(x) = \sum_{i=0}^{+\infty} \varepsilon^{i+2} \mathcal{X}^{i*} \cdot \nabla^{i}(\boldsymbol{f}(x) - \nabla p_{\varepsilon}^{*}(x))$$
304

Proof. The result is proved by induction. The case i = -1 is straightforward thanks to the convention $u_{-1} = p_{-1} = 0$. In this proof we use the short-hand

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notation $h_i(x) = f_i(x) - \nabla p_i^*(x)$. Assuming (3.8) and (3.9) hold till rank i + 1 with $i \ge -2$, we compute, substituting (3.8) into (3.2):

$$(3.13) \quad \begin{cases} (-\Delta_{yy}\boldsymbol{u}_{i+2} + \nabla_{y}p_{i+2})(x,y) \\ = h_{i+2,j}(x)\boldsymbol{e}_{j} + (2\partial_{l}\boldsymbol{\mathcal{X}}_{j}^{0}(y) - \alpha_{j}^{0}(y)\boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l}) \cdot \nabla h_{i+1,j}(x) \\ + \sum_{k=0}^{i} ((2\partial_{l}\boldsymbol{\mathcal{X}}_{j}^{k+1}(y) - \alpha_{j}^{k+1}(y)\boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \boldsymbol{\mathcal{X}}_{j}^{k}(y) \otimes I) \cdot \nabla^{k+2}h_{i-k,j}(x) \\ \operatorname{div}_{y}(\boldsymbol{u}_{i+2})(x,y) = -\sum_{k=0}^{i+1} (\boldsymbol{\mathcal{X}}_{j}^{k}(y) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}) \cdot \nabla^{k+1}h_{i+1-k,j}(x). \end{cases}$$

The system (3.13) admits a unique solution (u_{i+2}, p_{i+2}) with $\int_Y p_{i+2}(x, y) dy = 0$ if and only if the following compatibility condition (the so-called "Fredholm alternative") holds (for any $i \geq -1$):

$$\int_{Y} \operatorname{div}_{y}(\boldsymbol{u}_{i+2})(x,y) \mathrm{d}y = -\sum_{k=0}^{i+1} [\langle \boldsymbol{\mathcal{X}}_{j}^{k} \rangle \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}] \cdot \nabla^{k+1} h_{i+1-k,j}(x) = 0.$$

The above equation determines p_{i+1}^* given the values of p_k^* for $0 \le k \le i$:

$$(\langle \boldsymbol{\mathcal{X}}_{j}^{0} \rangle \cdot \boldsymbol{e}_{l}) \partial_{l} (f_{i+1,j} - \partial_{j} p_{i+1}^{*}) = -\sum_{k=1}^{i+1} [\langle \boldsymbol{\mathcal{X}}_{j}^{k} \rangle \cdot \boldsymbol{e}_{l} \otimes e_{l}] \cdot \nabla^{k+1} (f_{i+1-k,j} - \partial_{j} p_{i+1-k}^{*}),$$

which is (3.9) at order i + 1. This identity allows to rewrite $\operatorname{div}_{u}(u_{i+2})$ as

311 (3.14)
$$\operatorname{div}_{y}(\boldsymbol{u}_{i+2})(x,y) = -\sum_{k=0}^{i+1} [(\boldsymbol{\mathcal{X}}_{j}^{k}(y) - \langle \boldsymbol{\mathcal{X}}_{j}^{k} \rangle) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}] \cdot \nabla^{k+1} h_{i+1-k,j}(x).$$

By linearity, (3.13) and (3.14) and the definitions of $(\mathcal{X}_j^k, \alpha_j^k)$ through the cell problems (3.3)–(3.5) imply the result at rank i + 2.

Remark 3.3. The truncation of the series (3.12) at first order yields the wellknown Darcy's law [52]. The next terms of the series have been obtained in [46, 26], at least up to the order i = 1.

317 Remark 3.4. The ansatz (3.10) is already non standard (when compared to (1.7)) 318 because it features p_{ε}^* which is a formal power series in ε (recall (1.8)).

The next proposition establishes the symmetry and antisymmetry of the matrices \mathcal{X}^{k*} (eqn. (3.7)) for respectively odd and even values of k. We note that similar identities have been found for the Poisson [34] or the wave equation [1].

PROPOSITION 3.5. For any $k \ge 0$ and $0 \le p \le k$, $1 \le i, j \le d$, the following identity holds for the matrix valued tensor \mathcal{X}^{k*} :

324 (3.15)
$$\mathcal{X}_{ij}^{k*} = (-1)^p \int_Y ((-\Delta_{yy} \mathcal{X}_i^p + \nabla \alpha_i^p) \cdot \mathcal{X}_j^{k-p} + \nabla \alpha_j^{k-p} \cdot \mathcal{X}_i^p - \mathcal{X}_j^{k-p-1} \cdot \mathcal{X}_i^{p-1} \otimes I) \mathrm{d}y$$

with $\boldsymbol{\mathcal{X}}_{i}^{-1} = 0$ by convention. In particular, for any $k \geq 0$, $\boldsymbol{\mathcal{X}}^{2k*}$ and $\boldsymbol{\mathcal{X}}^{2k+1*}$ take values respectively in the set of $d \times d$ symmetric and antisymmetric matrices:

327 (3.16)
$$\mathcal{X}_{ij}^{2k*} = (-1)^k \int_Y (\nabla \mathcal{X}_i^k : \nabla \mathcal{X}_j^k + \nabla \alpha_i^k \cdot \mathcal{X}_j^k + \nabla \alpha_j^k \cdot \mathcal{X}_i^k - \mathcal{X}_i^{k-1} \cdot \mathcal{X}_j^{k-1} \otimes I) \mathrm{d}y$$

328

$$\mathcal{X}_{ij}^{2k+1*} = (-1)^k \int_Y (\mathcal{X}_i^k \cdot \nabla \mathcal{X}_j^k - \mathcal{X}_j^k \cdot \nabla \mathcal{X}_i^k + \alpha_i^k \mathcal{X}_j^k - \alpha_j^k \mathcal{X}_i^k) \cdot \boldsymbol{e}_l \otimes \boldsymbol{e}_l \mathrm{d}y \\ + (-1)^k \int_Y (\mathcal{X}_j^{k-1} \cdot \mathcal{X}_j^k - \mathcal{X}_i^{k-1} \cdot \mathcal{X}_j^k) \mathrm{d}y.$$

Proof. The result holds for p = 0 because

$$\mathcal{X}_{ij}^{k*} = \int_{Y} \mathcal{X}_{j}^{k} \cdot \boldsymbol{e}_{i} \mathrm{d}y = \int_{Y} \mathcal{X}_{j}^{k} \cdot (-\Delta_{yy} \mathcal{X}_{i}^{0} + \nabla \alpha_{i}^{0}) \mathrm{d}y.$$

Assuming now that (3.15) holds till rank p with $k > p \ge 0$, we prove the result at 330 rank p + 1. We write, after an integration by parts and by using (3.3)-(3.5): 331

332
$$\mathcal{X}_{ij}^{k*} = (-1)^p \int_Y [-\mathcal{X}_i^p \cdot \Delta \mathcal{X}_j^{k-p} - \alpha_i^p \operatorname{div}(\mathcal{X}_j^{k-p}) - \alpha_j^{k-p} \operatorname{div}(\mathcal{X}_i^p) - \mathcal{X}_i^{k-p-1} \cdot \mathcal{X}_i^{p-1} \otimes I] dy$$

333
$$= (-1)^{p} \int_{Y} \left[(\boldsymbol{\mathcal{X}}_{i}^{p} \cdot (2\partial_{l}\boldsymbol{\mathcal{X}}_{j}^{k-p-1} - \alpha_{j}^{k-p-1}\boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \boldsymbol{\mathcal{X}}_{j}^{k-p-2} \otimes I - \nabla \alpha_{j}^{k-p}) \cdot \boldsymbol{\mathcal{X}}_{i}^{p} \right]$$
334
$$+ \alpha_{i}^{p} \boldsymbol{\mathcal{X}}_{j}^{k-p-1} \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} + \alpha_{j}^{k-p} \boldsymbol{\mathcal{X}}_{i}^{p-1} \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} - \boldsymbol{\mathcal{X}}_{j}^{k-p-1} \cdot \boldsymbol{\mathcal{X}}_{i}^{p-1} \otimes I \right] \mathrm{d}\boldsymbol{y}$$

335
$$= (-1)^p \int_Y \left[-\boldsymbol{\mathcal{X}}_j^{k-p-1} \cdot \left((2\partial_l \boldsymbol{\mathcal{X}}_i^p - \alpha_i^p \boldsymbol{e}_l) \otimes \boldsymbol{e}_l + \boldsymbol{\mathcal{X}}_i^{p-1} \otimes \boldsymbol{I} \right) + \alpha_j^{k-p-1} \operatorname{div}(\boldsymbol{\mathcal{X}}_i^{p+1}) \right]$$

336
$$-\nabla \alpha_j^{k-p} \cdot \boldsymbol{\mathcal{X}}_i^p - \alpha_j^{k-p} \operatorname{div}(\boldsymbol{\mathcal{X}}_i^p) + \boldsymbol{\mathcal{X}}_j^{k-p-2} \cdot \boldsymbol{\mathcal{X}}_i^p \otimes I] \mathrm{d}y$$

337
$$= (-1)^p \int_Y \left[-\mathcal{X}_j^{k-p-1} \cdot \left(-\Delta_{yy} \mathcal{X}_i^{p+1} + \nabla \alpha_i^{p+1} \right) - \nabla \alpha_j^{k-p-1} \cdot \mathcal{X}_i^{p+1} + \mathcal{X}_j^{k-p-2} \cdot \mathcal{X}_i^p \otimes I \right] \mathrm{d}y,$$

whence (3.15) at rank
$$p + 1$$
. Finally, the expression (3.16) for \mathcal{X}_{ij}^{2k*} is obtained by
setting $k \leftarrow 2k$ and $p \leftarrow k$ in (3.15). The expression for \mathcal{X}_{ij}^{2k+1*} is obtained by setting
 $k \leftarrow 2k + 1$ and $p \leftarrow k$ and performing an integration by parts.

3.2. Derivation of the infinite order homogenized equation and of the 343 criminal ansatz. We now proceed on the derivation of the infinite order homogenized 344 equation (1.9). Let us recall the classical positive definiteness of the Darcy tensor \mathcal{X}^{0*} . 345

COROLLARY 3.6. Assume (H1). The matrix $\mathcal{X}^{0*} = (\mathcal{X}^{0*}_{ij})_{1 \leq i,j \leq d}$ (defined in 346 (3.7)) is positive symmetric definite. 347

Proof. See [52] or Corollary 7.8 in [33]. 348

Hence, the following definition of the tensors $(M^k)_{k \in \mathbb{N}}$ makes sense. 349

PROPOSITION 3.7. Let M^k be the tensor of order k defined by induction as fol-350 lows: 351

352 (3.18)
$$\begin{cases} M^0 = (\mathcal{X}^{0*})^{-1} \\ M^k = -(\mathcal{X}^{0*})^{-1} \sum_{p=0}^{k-1} \mathcal{X}^{k-p*} \otimes M^p, \quad \forall k \ge 1. \end{cases}$$

Then the source terms f_i (eqn. (3.1)) can be expressed in terms of the averaged summands $u_i^*(x)$ and $p_i^*(x)$ ((1.8) and (3.7)) through the following identity:

355 (3.19)
$$\forall i \ge 0, \ \boldsymbol{f}_i(x) - \nabla p_i^*(x) = \sum_{k=0}^i M^k \cdot \nabla^k \boldsymbol{u}_{i-k}^*(x)$$

Recognizing a Cauchy product, (3.19) and (3.12) rewrite formally as the "infinite order" homogenized system (1.9) for the formal average $(\boldsymbol{u}_{\varepsilon}^*, \boldsymbol{p}_{\varepsilon}^*)$ defined in (1.8).

Proof. The proof is identical to the one of Proposition 5 in [33], it amounts to average the first line of (3.8) with respect to y and to solve the resulting triangular system determining $f_{i-k} - \nabla p_{i-k}^*$ in terms of u_i^* .

The definition (3.18) essentially states that $\sum_{k=0}^{+\infty} \varepsilon^{k-2} M^k \cdot \nabla^k$ is the inverse of the formal power series $\sum_{k=0}^{+\infty} \varepsilon^{k+2} \mathcal{X}^{k*} \cdot \nabla^k$. In this spirit, it is even possible to write a fully explicit formula (see [34], Proposition 6 and Remark 2 for the proof):

364 PROPOSITION 3.8. For any $k \ge 1$, the tensor M^k is explicitly given by

365 (3.20)
$$M^k = \sum_{p=1}^k (-1)^p \sum_{\substack{i_1 + \dots + i_p = k \\ 1 \le i_1, \dots, i_p \le k}} (\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{i_1*} \otimes \dots \otimes (\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{i_p*} \otimes (\mathcal{X}^{0*})^{-1}.$$

We now introduce matrix valued tensors N^k and vector valued tensors $\boldsymbol{\beta}^k$ which allow

to obtain the "criminal ansatz" (1.10) expressing the velocity and pressure $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ in terms of their formal average $(\boldsymbol{u}_{\varepsilon}^*, p_{\varepsilon}^*)$.

PROPOSITION 3.9. Let N^k and β^k be respectively the k-th order matrix valued and vector valued tensors defined for any $k \in \mathbb{N}$ by

$$N^{k}(y) := \sum_{p=0}^{k} \mathcal{X}^{k-p}(y) \otimes M^{p}, \qquad \beta^{k}(y) := \sum_{p=0}^{k} (-1)^{p} M^{p} \cdot \boldsymbol{\alpha}^{k-p}(y), \qquad \forall y \in Y.$$

Then the summands $u_i(x, y)$ and $p_i(x, y)$ of (3.10) and (3.11) are given for any $i \ge 0$ in terms of the averages u_i^* (eqn. (3.7)) and p_i^* as follows:

371 (3.21)
$$\boldsymbol{u}_{i}(x,y) = \sum_{k=0}^{i} N^{k}(y) \cdot \nabla^{k} \boldsymbol{u}_{i-k}^{*}(x), \qquad p_{i}(x,y) = \sum_{k=0}^{i} \boldsymbol{\beta}^{k}(y) \cdot \nabla^{k} \boldsymbol{u}_{i-k}^{*}(x).$$

Recognizing Cauchy products, the identities (3.21) can be rewritten formally as the "criminal ansatz" (1.10).

Proof. The result is obtained by substituting (3.19) into (3.8) which yields

$$\begin{aligned} \boldsymbol{u}_{i}(x,y) &= \sum_{p=0}^{i} \sum_{q=0}^{i-p} \mathcal{X}^{p}(y) \otimes M^{q} \cdot \nabla^{p+q} \boldsymbol{u}_{i-p-q}^{*}(x) \\ &= \sum_{k=0}^{i} \sum_{p=0}^{k} (\mathcal{X}^{p}(y) \otimes M^{p-k}) \cdot \nabla^{k} \boldsymbol{u}_{i-k}^{*}(x) \text{ (change of indices } k = p+q) \end{aligned}$$

from where the identity (3.21) for $u_i(x, y)$ follows by inverting the summation. Similarly, we obtain

$$p_i(x,y) = \sum_{k=0}^{i} \sum_{p=0}^{k} ((M^{p-k})^T \cdot \alpha^p(y)) \cdot \nabla^k u_{i-k}^*(x),$$

hence (3.21) by using $(M^{p-k})^T = (-1)^{p-k} M^{p-k}$ (see Corollary 3.11 below). 374 In what follows, we denote by $(N_j^k)_{1 \le j \le d}$ and by $(\beta_j^k)_{1 \le j \le d}$ respectively the column vectors and the coefficients of $N^k(y)$ and $\beta^k(y)$:

$$\forall 1 \leq i, j \leq d, \ \mathbf{N}_j^k := N^k \mathbf{e}_j \text{ and } \beta_j^k := \boldsymbol{\beta}^k \cdot \mathbf{e}_j$$

- In addition, the convention $N_j^{-1} = 0$ is assumed. We shall in the sequel use several 375
- times the following properties of (N_j^k, β_j^k) which are dual to those of $(\mathcal{X}_j^k, \alpha_j^k)$. 376
- PROPOSITION 3.10. The k-th order tensors N^k , $(N_j^k)_{1 \le j \le d}$, β^k and $(\beta_j^k)_{1 \le j \le d}$ 377 378 satisfy:
- $(i) \int_Y N^0(y) dy = I \text{ and } \int_Y N^k(y) dy = 0 \text{ for any } k \ge 1;$ 379(ii) $\int_{Y} \boldsymbol{\beta}^{k}(y) dy = 0$ for any $k \ge 0$; 380
- (iii) For any $k \ge -2$ and $1 \le j \le d$,
- 381

382 (3.22)
$$\begin{cases} -\Delta_{yy} \boldsymbol{N}_{j}^{k+2} + \nabla \beta_{j}^{k+2} = (2\partial_{l} \boldsymbol{N}_{j}^{k+1} - \beta_{j}^{k+1} \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \boldsymbol{N}_{j}^{k} \otimes \boldsymbol{I} + M^{k+2} \boldsymbol{e}_{j}, \\ \operatorname{div}(\boldsymbol{N}_{j}^{k+2}) = -(\boldsymbol{N}_{j}^{k+1} - \langle \boldsymbol{N}_{j}^{k+1} \rangle) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}; \end{cases}$$

Proof. (i) and (ii) are straightforward consequences of (3.18). 383

(iii) is obtained by writing, for $k \ge 0$ (implicit summation on the repeated index 384 j assumed): 385

$$-\Delta_{yy} \mathbf{N}_{j}^{k+2} + \nabla \beta_{j}^{k+2} = -\Delta_{yy} \left(\sum_{p=0}^{k+2} \mathbf{\mathcal{X}}_{i}^{k+2-p}(y) \otimes M_{ij}^{p} \right) + \nabla \left(\sum_{p=0}^{k+2} \alpha_{i}^{k+2-p}(y) \otimes M_{ij}^{p} \right)$$

$$= \sum_{p=0}^{k} \left[(2\partial_{l} \mathbf{\mathcal{X}}_{i}^{k+1-p} - \alpha_{i}^{k+1-p} \mathbf{e}_{l}) \otimes e_{l} + \mathbf{\mathcal{X}}_{i}^{k-p} \otimes I \right] M_{ij}^{p}$$

$$+ (2\partial_{l} \mathbf{\mathcal{X}}_{i}^{0} - \alpha_{i}^{0} \mathbf{e}_{l}) M_{ij}^{k+1} + M_{ij}^{k+2} \mathbf{e}_{i}$$

$$= (2\partial_{l} \mathbf{N}_{j}^{k+1} - \beta_{i}^{k+1} \mathbf{e}_{l}) \otimes e_{l} + \mathbf{N}_{j}^{k} \otimes I + M^{k+2} \mathbf{e}_{j}.$$

$$\operatorname{div}(\mathbf{N}_{j}^{k+2}) = \sum_{p=0}^{k+2} \operatorname{div}(\mathbf{\mathcal{X}}_{i}^{k+2-p}) M_{ij}^{p} = -\sum_{p=0}^{k+1} M_{ij}^{p} (\mathbf{\mathcal{X}}_{i}^{k+1-p} - \langle \mathbf{\mathcal{X}}_{i}^{k+1-p} \rangle) \cdot \mathbf{e}_{l} \otimes e_{l}.$$
The proof is identical for $k = -1$ and $k = -2$.

- The proof is identical for k = -1 and k = -2. 386
- The identity (3.22) allows to infer important properties characterizing the tensors M^k 387 which are similar to those of Proposition 3.5. 388

COROLLARY 3.11. For any $1 \le p \le k - 1$, it holds

$$M_{ij}^{k} = (-1)^{p+1} \int_{Y} ((-\Delta_{yy} \mathbf{N}_{i}^{p} + \nabla \beta_{i}^{p}) \cdot \mathbf{N}_{j}^{k-p} + \nabla \beta_{j}^{k-p} \cdot \mathbf{N}_{i}^{p} - \mathbf{N}_{i}^{p-1} \cdot \mathbf{N}_{j}^{k-p-1} \otimes I) \mathrm{d}y.$$

Consequently, for any $k \geq 0$, 389

• M^{2k} is a symmetric matrix valued tensor, and the following identities hold:

$$M_{ij}^0 = \int_Y \nabla \boldsymbol{N}_i^0 : \nabla \boldsymbol{N}_j^0 \mathrm{d}y,$$

$$\forall k \ge 1, \ M_{ij}^{2k} = (-1)^{k+1} \int_{Y} (\nabla \mathbf{N}_{i}^{k} : \nabla \mathbf{N}_{j}^{k} + \nabla \beta_{i}^{k} \cdot \mathbf{N}_{j}^{k} + \nabla \beta_{j}^{k} \cdot \mathbf{N}_{i}^{k} - \mathbf{N}_{i}^{k-1} \cdot \mathbf{N}_{j}^{k-1} \otimes I) \mathrm{d}y.$$

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• M^{2k+1} is an antisymmetric matrix valued tensor and it holds:

392
$$M_{ij}^{2k+1} = (-1)^{k+1} \int_{Y} (\boldsymbol{N}_{i}^{k} \cdot \nabla \boldsymbol{N}_{j}^{k} - \boldsymbol{N}_{j}^{k} \cdot \nabla \boldsymbol{N}_{i}^{k} + \beta_{i}^{k} \boldsymbol{N}_{j}^{k} - \beta_{j}^{k} \boldsymbol{N}_{i}^{k}) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} \mathrm{d}\boldsymbol{y}$$

$$+ (-1)^{k+1} \int_{Y} (\boldsymbol{N}_{j}^{k-1} \cdot \boldsymbol{N}_{i}^{k} - \boldsymbol{N}_{i}^{k-1} \cdot \boldsymbol{N}_{j}^{k}) \otimes I \mathrm{d}y.$$

Proof. The proof is very similar to the one of Proposition 3.5 and is omitted, see 395 also Proposition 7.34 in [33]. Π 396

Remark 3.12. The antisymmetry of odd order tensors M^{2k+1} ensures that the 397 associated differential operators $\varepsilon^{2k-1}M^{2k+1} \cdot \nabla^{2k+1}$ arising in the "infinite order" 398 homogenized equation (1.9) are symmetric. Indeed, the antisymmetry of M^{2k+1} "com-399 pensates" the one induced by odd order derivatives which makes $M^{2k+1} \cdot \nabla^{2k+1}$ be a 400symmetric operator: for two vector fields $\boldsymbol{u} := (u_i)_{1 \le i \le d}, \boldsymbol{v} = (v_i)_{1 \le i \le d}$, it holds 401 402

390 391

Remark 3.13. It is not completely straightforward to exhibit an instance of hole 406 ∂T and $k \in \mathbb{N}$ for which we can actually prove that M^{2k+1} is not zero. However 407 simple numerical evidences tend to confirm this conjecture, see section 7.4.5 in [33] 408 for an example featuring $M^1 \neq 0$ in the case of the elasticity system. 409

3.3. Simplifications for the tensors \mathcal{X}^{k*} and M^k in case of symmetries. 410 In the final part of this section, we examine how the symmetries of the obstacle ηT 411 with respect to the cell axes reflect into the coefficients of the matrix valued tensors 412 \mathcal{X}^{k*} and M^k . Our final result is stated in Corollary 3.16, which implies that odd 413 order tensors \mathcal{X}^{2k+1} and M^{2k+1} vanish in case ηT is symmetric with respect to the 414 cell axes. It is based on the following elementary lemma: 415

LEMMA 3.14. Let $S \in \mathbb{R}^{d \times d}$ an orthogonal symmetry, i.e. $S = S^T$ and SS = I. 416 The following identities hold for any smooth vector field $\boldsymbol{\mathcal{X}}$ and scalar field $\boldsymbol{\alpha}$: 417

 $-\Delta(S\mathcal{X} \circ S) + \nabla(\alpha \circ S) = S(-\Delta\mathcal{X} + \nabla\alpha) \circ S,$ 418 (3.23)

419 (3.24)
$$\operatorname{div}(S\mathcal{X} \circ S) = \operatorname{div}(\mathcal{X}) \circ S,$$

$$439 \quad (3.25) \qquad \qquad \partial_i (S\mathcal{X} \circ S) = S_{ij} S(\partial_j \mathcal{X}) \circ S$$

Proof. The first two identities are obtained by writing

$$-\Delta(S\boldsymbol{\mathcal{X}}\circ S) + \nabla(\alpha\circ S) = -S\partial_{ij}\boldsymbol{\mathcal{X}}\circ SS_{il}S_{jl} + S(\nabla\alpha)\circ S$$
$$= -S(\Delta\boldsymbol{\mathcal{X}} + \nabla\alpha)\circ S,$$

$$\operatorname{div}(S\boldsymbol{\mathcal{X}}\circ S)=\operatorname{Tr}(\nabla(S\boldsymbol{\mathcal{X}}\circ S))=\operatorname{Tr}(S(\nabla\boldsymbol{\mathcal{X}})\circ SS)=\operatorname{Tr}((\nabla\boldsymbol{\mathcal{X}})\circ S)=\operatorname{div}(\boldsymbol{\mathcal{X}})\circ S.$$

Identity (3.25) is an elementary consequence of the chain rule. 422

PROPOSITION 3.15. If the cell $Y = P \setminus (\eta T)$ is invariant with respect to a sym-423

metry S, i.e. S(Y) = Y, then the following identity holds for the tensors $(\mathcal{X}_{I}^{k}, \alpha_{I}^{k})$ 424 (defined in (3.3) - (3.5)):425

426 (3.26)
$$S\boldsymbol{\mathcal{X}}_{i_1\dots i_k,l}^k \circ S = S_{i_1j_1}\dots S_{i_kj_k}S_{lm}\boldsymbol{\mathcal{X}}_{j_1\dots j_k,m}^k,$$

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$$\begin{array}{l} 427\\ 428 \end{array} (3.27) \qquad \qquad \alpha_{i_1\dots i_k,l}^k \circ S = S_{i_1j_1}\dots S_{i_kj_k} S_{lm} \alpha_{j_1\dots j_k,m}^k \end{array}$$

429 with implicit summation over the repeated indices j_1, \ldots, j_k and m. As a consequence,

430 the following identities hold for the constant matrix valued tensors \mathcal{X}^{k*} and M^k :

431 (3.28)
$$\mathcal{X}_{i_1...i_k,lm}^{k*} = S_{i_1j_1}...S_{i_kj_k}S_{lp}S_{mq}\mathcal{X}_{j_1...j_k,pq}^{k*}$$

$$M_{i_1...i_k,lm}^k = S_{i_1j_1}...S_{i_kj_k}S_{lp}S_{mq}M_{j_1...j_k,pq}^k$$

Proof. We prove (3.26) and (3.27) by induction. Applying Proposition 3.15 yields

$$\begin{cases} -\Delta_{yy}(S\boldsymbol{\mathcal{X}}_{l}^{0}\circ S)+\nabla_{y}(\alpha_{l}^{0}\circ S)=S\boldsymbol{e}_{l}\circ S=S\boldsymbol{e}_{l}=S_{mj}\boldsymbol{e}_{m},\\ \operatorname{div}(S\boldsymbol{\mathcal{X}}_{l}^{0}\circ S)=0. \end{cases}$$

Since the cell is symmetric with respect to S, $(S\mathcal{X}_l^0 \circ S, \alpha_l^0 \circ S)$ satisfies the same boundary conditions (3.6) than $S_{mj}(\mathcal{X}_m^0, \alpha_m^0)$. Therefore these vector fields are equal and we infer (3.26) and (3.27) at rank k = 0. We then write, for a given $1 \le i_1 \le d$:

$$\begin{cases} -\Delta_{yy}(S\boldsymbol{\mathcal{X}}_{i_{1},l}^{1}\circ S) + \nabla_{y}(\alpha_{i_{1},l}^{1}\circ S) = S(2\partial_{i_{1}}\boldsymbol{\mathcal{X}}_{l}^{0} - \alpha_{l}^{0}\boldsymbol{e}_{i_{1}})\circ S\\ = S_{i_{1}j_{1}}(2\partial_{j_{1}}(S\boldsymbol{\mathcal{X}}_{l}^{0}\circ S) - \alpha_{l}^{0}\circ S\boldsymbol{e}_{j_{1}}) = S_{i_{1}j_{1}}S_{lm}(2\partial_{j_{1}}\boldsymbol{\mathcal{X}}_{m}^{0} - \alpha_{m}^{0}\boldsymbol{e}_{j_{1}}),\\ \operatorname{div}_{y}(S\boldsymbol{\mathcal{X}}_{i_{1},l}^{1}\circ S) = -(\boldsymbol{\mathcal{X}}_{l}^{0}\circ S - \langle \boldsymbol{\mathcal{X}}_{l}^{0}\rangle) \cdot \boldsymbol{e}_{i_{1}}\\ = -S_{lm}S(\boldsymbol{\mathcal{X}}_{m}^{0} - \langle \boldsymbol{\mathcal{X}}_{m}^{0}\rangle) \cdot \boldsymbol{e}_{i_{1}} = -S_{i_{1}j_{1}}S_{lm}(\boldsymbol{\mathcal{X}}_{m}^{0} - \langle \boldsymbol{\mathcal{X}}_{m}^{0}\rangle) \cdot \boldsymbol{e}_{j_{1}}, \end{cases}$$

where we have used $\langle \boldsymbol{\mathcal{X}}_{l}^{0} \rangle = \langle \boldsymbol{\mathcal{X}}_{l}^{0} \circ S \rangle$. This implies similarly (3.26) and (3.27) at rank k = 1. Assuming now the result holds till rank k + 1 with $k \geq 0$, it holds:

$$\begin{cases} -\Delta_{yy}(S\boldsymbol{\mathcal{X}}_{i_{1}...i_{k+2},l}^{k+2}\circ S) + \nabla_{y}(\alpha_{i_{1}...i_{k+2},l}^{k+2}\circ S) \\ &= S(2\partial_{i_{k+2}}\boldsymbol{\mathcal{X}}_{i_{1}...i_{k+1},l}^{k+1} - \alpha_{i_{1}...i_{k+1},l}^{k+1}\boldsymbol{e}_{i_{k+2}})\circ S + S\boldsymbol{\mathcal{X}}_{i_{1}...i_{k},l}^{k}\circ S\delta_{i_{k+1}i_{k+2}} \\ &= S_{i_{k+2}j_{k+2}}(2\partial_{j_{k+2}}(S\boldsymbol{\mathcal{X}}_{i_{1}...i_{k+1},l}^{k+1}\circ S) - \alpha_{i_{1}...i_{k+1},l}^{k+1}\circ S\boldsymbol{e}_{j_{k+2}}) \\ &+ S_{i_{k+1}j_{k+1}}S_{i_{k+2}j_{k+2}}\delta_{j_{k+1}j_{k+2}}S\boldsymbol{\mathcal{X}}_{i_{1}...i_{k},l}^{k}\circ S \\ &= S_{i_{1}j_{1}}\ldots S_{i_{k+2}i_{k+2}}S_{lm}[(2\partial_{j_{k+2}}\boldsymbol{\mathcal{X}}_{j_{1}...j_{k+1},m}^{k+1} - \alpha_{j_{1}...j_{k+1},m}^{k+1}\boldsymbol{e}_{j_{k+2}}) \\ &+ \delta_{j_{k+1}j_{k+2}}\boldsymbol{\mathcal{X}}_{j_{1}...j_{k},m}^{k}] \\ \operatorname{div}_{y}(S\boldsymbol{\mathcal{X}}^{k+2}\circ S) = -(\boldsymbol{\mathcal{X}}_{i_{1}...i_{k+1},l}^{k+1}\circ S - <\boldsymbol{\mathcal{X}}_{i_{1}...i_{k+1},l}^{k+1} >) \dots \boldsymbol{e}_{i_{k+2}} \\ &= -S_{i_{1}j_{1}}\ldots S_{i_{k+1}j_{k+1}}S_{lm}S(\boldsymbol{\mathcal{X}}_{j_{1}...j_{k+1},m}^{k+1} - <\boldsymbol{\mathcal{X}}_{j_{1}...j_{k+1},m}^{k+1} >) \cdot \boldsymbol{e}_{i_{k+2}} \\ &= -S_{i_{1}j_{1}}\ldots S_{i_{k+2}j_{k+2}}S_{lm}(\boldsymbol{\mathcal{X}}_{j_{1}...j_{k+1},m}^{k+1} - <\boldsymbol{\mathcal{X}}_{j_{1}...j_{k+1},m}^{k+1} >) \cdot \boldsymbol{e}_{j_{k+2}}, \end{cases}$$

hence (3.26) and (3.27) at rank k + 2. A change of variable then yields:

$$\mathcal{X}_{i_1\dots i_k,lm}^{k*} = \int_Y \boldsymbol{e}_l \cdot \boldsymbol{\mathcal{X}}_{i_1\dots i_k,m}^k \mathrm{d}y = \int_Y (S\boldsymbol{e}_l) \cdot (S\boldsymbol{\mathcal{X}}_{i_1\dots i_k,m}^k \circ S) \mathrm{d}y.$$

434 This implies (3.28), and then (3.29) by using (3.20).

435 We apply the above result to two possible families of symmetries:

• for $1 \leq l \leq d$, the symmetry S^l with respect to the cell axis e_l :

$$S^l = 1 - 2\boldsymbol{e}_l \boldsymbol{e}_l^T;$$

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• for $1 \le m, l \le d$ with $m \ne l$, the symmetry S^{lm} with respect to the diagonal axis $\boldsymbol{e}_l - \boldsymbol{e}_m$:

$$S^{lm} = I - \boldsymbol{e}_l \boldsymbol{e}_l^T - \boldsymbol{e}_m \boldsymbol{e}_m^T + \boldsymbol{e}_l \boldsymbol{e}_m^T + \boldsymbol{e}_l \boldsymbol{e}_m^T + \boldsymbol{e}_m \boldsymbol{e}_l^T$$

1. If the cell Y is symmetric with respect to all cell axes Corollary 3.16. $(e_l)_{1 \leq l \leq d}$, i.e. $S^l(Y) = Y$ for any $1 \leq l \leq d$, then

$$\mathcal{X}^{k*}_{i_1...i_k,pq} = 0 \ and \ M^k_{i_1...i_k,pq} = 0$$

whenever any given integer $1 \leq l \leq d$ occurs an odd number of times in the 436indices $i_1 \dots i_k, p, q$. In particular, this implies $\mathcal{X}^{2k+1*} = 0$ and $M^{2k+1} = 0$.

> 2. If the cell Y is symmetric with respect to all diagonal axes $e_l - e_m$, i.e. $S^{l,m}(Y) = Y$ for any $1 \leq l < m \leq d$, then for any permutation $\sigma \in \mathfrak{S}_d$,

$$\mathcal{X}_{\sigma(i_1)\ldots\sigma(i_k),\sigma(p)\sigma(q)}^{k*} = \mathcal{X}_{i_1\ldots i_k,pq}^{k*} \text{ and } M_{\sigma(i_1)\ldots\sigma(i_k),\sigma(p)\sigma(q)}^k = M_{i_1\ldots i_k,pq}^k.$$

Proof. The result is obtained by applying (3.28) and (3.29) to the particular 438 symmetries S^l and S^{lm} . See also Corollary 3 in [34]. 439

- Let us illustrate how the previous properties translate for the tensors M^0 , M^2 and 440 M^4 : 441
 - if the cell Y is symmetric with respect to all cell axes $(e_l)_{1 \le l \le d}$, only the coefficients of the form $M_{i,i}^0$ are non zero. For M^2 , only

$$M_{ii,jj}^2, M_{ij,ij}^2, M_{ii,ii}^2$$

with $i \neq j$ are non zero. For M^4 , only the coefficients of the form

$$M_{iijj,kk}^4, M_{iijk,jk}^4, M_{iiii,jj}^4, M_{iijj,ii}^4, M_{iiij,ij}^4, M_{iiii,ii}^4$$

are non zero with distinct integers i, j, k.

• If in addition the obstacle is symmetric with respect to all diagonal axes, then the values of the above coefficients do not depend on the choice of the distinct integers i, j, k. As a result, M^0 is proportional to the identity tensor, M^2 reduces to at most three coefficients (the material is said to be orthotropic), and M^4 reduces to at most 6 coefficients for $d \ge 3$, and to 4 coefficients for d=2. For instance there are three constants α, β, γ such that $M^2 \cdot \nabla$ is the operator

$$M^2 \cdot \nabla \boldsymbol{v} = \alpha \Delta \boldsymbol{v} + \beta \nabla \operatorname{div}(\boldsymbol{v}) + \gamma \sum_{i=1}^d \partial_{ii} v_i \boldsymbol{e}_i$$

4. Homogenized equations of order 2K + 2: tensors \mathbb{D}_K^k . In this section, 443we derive the well-posed high order homogenized system (1.6) and we justify the 444445 homogenization process by means of quantitative error estimates.

The formal identities (1.10) lead us to introduce, for any order $K \in \mathbb{N}$, the 446 truncated ansatz $W_{\varepsilon,K}(v)$ and $Q_{\varepsilon,K}(v,\phi)$ for the reconstructed velocity and pressure: 447

448 (4.1)
$$\boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v})(x,y) := \sum_{k=0}^{K} \varepsilon^k N^k(y) \cdot \nabla^k \boldsymbol{v}(x), \qquad x \in D, \ y \in Y$$

449 (4.2)
$$Q_{\varepsilon,K}(\boldsymbol{v},\phi)(x,y) := \phi(x) + \sum_{k=0}^{K} \varepsilon^{k-1} \boldsymbol{\beta}^{k}(y) \cdot \nabla^{k} \boldsymbol{v}(x), \qquad x \in D, \ y \in Y$$

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for any $\boldsymbol{v} \in H^{K+1}(D,\mathbb{R}^d)$ and $\phi \in L^2(D)$ which are sought to approximate the 451 homogenized averages $\boldsymbol{u}_{\varepsilon}^*$ and p_{ε}^* respectively. Similarly we denote by $\boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v})$ and 452 $\widetilde{Q}_{\varepsilon,K}(\boldsymbol{v},\phi)$ the reconstructed *oscillating* functions defined for any $x \in D_{\varepsilon}$ by 453

454 (4.3)
$$\widetilde{\boldsymbol{W}}_{\varepsilon,K}(\boldsymbol{v})(x) := \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v})(x,x/\varepsilon), \quad \widetilde{Q}_{\varepsilon,K}(\boldsymbol{v},\phi)(x) := Q_{\varepsilon,K}(\boldsymbol{v},\phi)(x,x/\varepsilon).$$

Most of the results of this section are consequences of the following observation: 455

LEMMA 4.1. For any $K' \in \mathbb{N}$, $(\boldsymbol{v}, \phi) \in H^1(D, \mathbb{R}^d) \times L^2(D)$, the reconstructed 456 velocity and pressure $(\widetilde{W}_{\varepsilon,K'}(\boldsymbol{v}),\widetilde{Q}_{\varepsilon,K'}(\boldsymbol{v},\phi))$ of (4.3) satisfy 457

$$-\Delta \widetilde{\boldsymbol{W}}_{\varepsilon,K'}(\boldsymbol{v}) + \nabla \widetilde{\boldsymbol{Q}}_{\varepsilon,K'}(\boldsymbol{v},\phi)$$

$$= \sum_{k=0}^{K'} \varepsilon^{k-2} M^k \nabla^k \boldsymbol{v} - \varepsilon^{K'-1} ((2\partial_l \boldsymbol{N}_j^{K'} - \beta_j^{K'} \boldsymbol{e}_l) \otimes \boldsymbol{e}_l) (\cdot/\varepsilon) \cdot \nabla^{K'+1} v_j.$$

$$+ (\boldsymbol{N}_j^{K'-1} \otimes I) (\cdot/\varepsilon) \cdot \nabla^{K'+1} v_j - \varepsilon^{K'} \boldsymbol{N}_j^{K'} (\cdot/\varepsilon) \otimes I \cdot \nabla^{K'+2} v_j,$$

$$(459)$$

460 (4.5)
$$\operatorname{div}(\widetilde{\boldsymbol{W}}_{\varepsilon,K'}(\boldsymbol{v})) = \operatorname{div}(\boldsymbol{v}) + \varepsilon^{K'} \boldsymbol{N}_{j}^{K'}(\cdot/\varepsilon) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} \cdot \nabla^{K'+1} \boldsymbol{v}_{j}.$$

Proof. (4.4) and (4.5) are obtained by applying the Laplace and gradient opera-461 tors on (4.1) and (4.2) and by using the identity (3.22). Г 462

4.1. Sufficient conditions leading to error estimates. The purpose of this 463 part is to demonstrate that a sequence of functions $(\boldsymbol{v}_{\varepsilon}^*, q_{\varepsilon}^*)_{\varepsilon>0}$ yields an approxima-464tion of $(\boldsymbol{u}_{\varepsilon}, p_{\varepsilon})$ at the order $O(\varepsilon^{K'})$ provided it solves the infinite order homogenized 465equation (1.9) up to a remainder of order $O(\varepsilon^{K'+1})$. The derivation of a *finite-order* 466 homogenized equation such as (1.6) reduces then to determine 2K + 2 - K' tensors 467 \mathbb{D}_{K}^{k} for $K' + 1 \leq k \leq 2K + 2$ such that the equation 468

469 (4.6)
$$\sum_{k=K'+1}^{2K+2} \varepsilon^{k-2} \mathbb{D}_{K}^{k} \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon}^{*} + \sum_{k=0}^{K'} \varepsilon^{k-2} M^{k} \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon}^{*} + \nabla q_{\varepsilon}^{*} = \boldsymbol{f}$$

is well-posed. The proof is based on the next three technical results. 470

LEMMA 4.2. There exists a constant C independent of $\varepsilon > 0$ such that for any $\boldsymbol{v} \in H^1(D_{\varepsilon}, \mathbb{R}^d)$ with $\boldsymbol{v} = 0$ on $\partial \omega_{\varepsilon}$, the following Poincaré inequality holds:

$$||\boldsymbol{v}||_{L^2(D_{\varepsilon},\mathbb{R}^d)} \leq C\varepsilon ||\nabla \boldsymbol{v}||_{L^2(D,\mathbb{R}^{d\times d})}$$

Proof. See e.g. [44] or the appendix of [49]. 471

The next lemma states the existence of a continuous right inverse for the divergence 472

 B_{ε} —so-called a Bogovskii's operator—with a bound explicit in ε on the uniform con-473tinuity constant. 474

LEMMA 4.3. Assume (H1) and (H2). Then there exists a linear operator B_{ε} : $L^{2}(D) \rightarrow H^{1}(D, \mathbb{P}^{d})$ satisfying for any $\phi \in L^{2}(D)$ with $\int \phi dx = 0$: 475

476
$$L^{2}(D_{\varepsilon}) \rightarrow H^{2}(D_{\varepsilon}, \mathbb{R}^{n})$$
 satisfying, for any $\phi \in L^{2}(D_{\varepsilon})$ with $\int_{D_{\varepsilon}} \phi dx = 0$.

477 (i)
$$\operatorname{div}(B_{\varepsilon}\phi) = \phi$$
 in D_{ε} ,

(ii) $B_{\varepsilon}\phi = 0$ on $\partial \omega_{\varepsilon}$ and $B_{\varepsilon}\phi$ is D-periodic, 478

(iii) $||\nabla(B_{\varepsilon}\phi)||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} \leq C\varepsilon^{-1}||\phi||_{L^{2}(D_{\varepsilon})}$, for a constant C > 0 independent of ϕ 479and ε . 480

481 *Proof.* See [32], Lemma 2.1, or [33], Lemma 7.9.

COROLLARY 4.4. Assume (H1) and (H2). For any $\mathbf{h} \in L^2(D_{\varepsilon}, \mathbb{R}^d)$ and $g \in$ 482 $L^2(D_{\varepsilon})$ satisfying $\int_{D_{\varepsilon}} g dx = 0$, let $(\boldsymbol{v}, \phi) \in H^1(D_{\varepsilon}, \mathbb{R}^d) \times L^2(D_{\varepsilon})$ be the unique solution 483to the Stokes problem 484

485 (4.7)
$$\begin{cases} -\Delta \boldsymbol{v} + \nabla \phi = \boldsymbol{h} \text{ in } D_{\varepsilon} \\ \operatorname{div}(\boldsymbol{v}) = g \text{ in } D_{\varepsilon} \\ \int_{D_{\varepsilon}} \phi \mathrm{d}x = 0 \\ \boldsymbol{v} = 0 \text{ on } \partial \omega_{\varepsilon} \\ \boldsymbol{v} \text{ is } D\text{-periodic.} \end{cases}$$

There exists a constant C independent of ε , **h** and g such that 486

$$487 \quad (4.8) \qquad ||\nabla \boldsymbol{v}||_{L^2(D_{\varepsilon},\mathbb{R}^{d\times d})} + \varepsilon||\phi||_{L^2(D_{\varepsilon})} \le C(\varepsilon||\boldsymbol{h}||_{L^2(D_{\varepsilon},\mathbb{R}^d)} + \varepsilon^{-1}||g||_{L^2(D_{\varepsilon})}),$$

Proof. We use the operator B_{ε} of Lemma 4.3 to lift the divergence of v. Let us define the vector field $\boldsymbol{w} := \boldsymbol{v} - B_{\varepsilon}g \in H^1_{per}(D_{\varepsilon}, \mathbb{R}^d)$ which satisfies

$$\begin{cases} \operatorname{div}(\boldsymbol{w}) = 0 \text{ in } D_{\varepsilon}, \\ \boldsymbol{w} = 0 \text{ on } \partial \omega_{\varepsilon}. \end{cases}$$

After an integration by part, we obtain:

$$\begin{aligned} ||\nabla \boldsymbol{w}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})}^{2} &= \int_{D_{\varepsilon}} \boldsymbol{h} \cdot \boldsymbol{w} \mathrm{d}x - \int_{D_{\varepsilon}} \nabla(B_{\varepsilon}g) : \nabla \boldsymbol{w} \mathrm{d}x \\ &\leq ||\boldsymbol{h}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} ||\boldsymbol{w}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} + ||\nabla(B_{\varepsilon}g)||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} ||\nabla \boldsymbol{w}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} \\ &\leq C(\varepsilon ||\boldsymbol{h}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} + ||\nabla(B_{\varepsilon}g)||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})}) ||\nabla \boldsymbol{w}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})}, \end{aligned}$$

where the last inequality is a consequence of Lemma 4.2. Therefore, simplifying by $||\nabla w||_{L^2(D_{\varepsilon},\mathbb{R}^{d\times d})}$ and using the point (iii) of Lemma 4.3 yields

$$\begin{aligned} ||\nabla \boldsymbol{v}||_{L^{2}(D,\mathbb{R}^{d\times d})} &\leq ||\nabla \boldsymbol{w}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} + ||\nabla (B_{\varepsilon}g)||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} \\ &\leq C(\varepsilon||\boldsymbol{h}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} + \varepsilon^{-1}||g||_{L^{2}(D_{\varepsilon})}), \end{aligned}$$

which proves the first part of the bound (4.8) on ∇v . The bound on the pressure is 488 489 then obtained by using $B_{\varepsilon}\phi$ as a test function: we write

490
$$||\phi||^{2}_{L^{2}(D_{\varepsilon})} = \int_{D_{\varepsilon}} \phi \operatorname{div}(B_{\varepsilon}\phi) \mathrm{d}x = -\int_{D_{\varepsilon}} \nabla \phi \cdot B_{\varepsilon}\phi \mathrm{d}x$$

491 $= \int_{D_{\varepsilon}} (-\Delta \boldsymbol{v} - \boldsymbol{h}) \cdot B_{\varepsilon}\phi \mathrm{d}x = \int_{D_{\varepsilon}} (\nabla \boldsymbol{v} \cdot \nabla (B_{\varepsilon}\phi) - \boldsymbol{h} \cdot B_{\varepsilon}\phi) \mathrm{d}x$

$$< ||\nabla v||_{L^2(D)}$$

492
493
$$\leq ||\nabla \boldsymbol{v}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})}||\nabla (B_{\varepsilon}\phi)||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} + ||\boldsymbol{h}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})}||B_{\varepsilon}\phi||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})}$$
493
$$\leq C(\varepsilon||\boldsymbol{h}||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} + \varepsilon^{-1}||g||_{L^{2}(D_{\varepsilon})})||\nabla (B_{\varepsilon}\phi)||_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})}$$

$$494 \\ 495 \\ \leq C\varepsilon^{-1}(\varepsilon||\boldsymbol{h}||_{L^2(D_\varepsilon,\mathbb{R}^d)} + \varepsilon^{-1}||g||_{L^2(D_\varepsilon)})||\phi||_{L^2(D_\varepsilon)},$$

which concludes the proof. 496

497 We are now in position to state the main result of this section.

498 PROPOSITION 4.5. Let $(\boldsymbol{v}_{\varepsilon}^*, q_{\varepsilon}^*)_{\varepsilon>0}$ be a sequence of functions of $H^1(D_{\varepsilon}, \mathbb{R}^d) \times$ 499 $L^2(D_{\varepsilon}), D$ -periodic, depending on ε (and possibly on K' and \boldsymbol{f}) satisfying the follow-500 ing conditions:

501 1. $(\boldsymbol{v}_{\varepsilon}, q_{\varepsilon})$ solves the infinite order homogenized equation (1.9) up to an error of 502 order $O(\varepsilon^{K'+1})$:

503 (4.9)
$$\left\| \sum_{k=0}^{K'} \varepsilon^{k-2} M^k \cdot \nabla^k \boldsymbol{v}_{\varepsilon}^* + \nabla q_{\varepsilon}^* - \boldsymbol{f} \right\|_{L^2(D,\mathbb{R}^d)} \leq C_{K'}(\boldsymbol{f}) \varepsilon^{K'+1}$$

504

505 (4.10)
$$\operatorname{div}(\boldsymbol{v}_{\varepsilon}^*) = 0 \ in \ D$$

507 (4.11)
$$(\boldsymbol{v}_{\varepsilon}^*, p_{\varepsilon}^*)$$
 is *D*-periodic.

508 2. For any $m \in \mathbb{N}$, there exists a constant C_m independent of ε such that

509 (4.12)
$$||\boldsymbol{v}_{\varepsilon}^{*}||_{H^{m}(D,\mathbb{R}^{d})} \leq C_{m}\varepsilon^{2}$$

Then the reconstructed functions $(\widetilde{W}_{\varepsilon,K'}(v_{\varepsilon}^*), \widetilde{Q}_{\varepsilon,K'-1}(q_{\varepsilon}^*, v_{\varepsilon}^*))$ (eqn. (4.1) and (4.2)) yield approximations of $(u_{\varepsilon}, p_{\varepsilon})$ of order O(K' + 2) in the $H^1(D_{\varepsilon}, \mathbb{R}^d)$ norm and $O(\varepsilon^{K'+3})$ in the $L^2(D_{\varepsilon}, \mathbb{R}^d)$ norm:

$$\begin{aligned} \left\| \nabla (\boldsymbol{u}_{\varepsilon} - \widetilde{\boldsymbol{W}}_{\varepsilon,K'}(\boldsymbol{v}_{\varepsilon}^{*})) \right\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} + \varepsilon \left\| p_{\varepsilon} - \widetilde{Q}_{\varepsilon,K'-1}(q_{\varepsilon}^{*},\boldsymbol{v}_{\varepsilon}^{*}) \right\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} \leq C_{K'}(\boldsymbol{f})\varepsilon^{K'+2}, \\ \left\| \boldsymbol{u}_{\varepsilon} - \widetilde{\boldsymbol{W}}_{\varepsilon,K'}(\boldsymbol{v}_{\varepsilon}^{*}) \right\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} \leq C_{K'}(\boldsymbol{f})\varepsilon^{K'+3}. \end{aligned}$$

Proof. According to Lemma 4.1 and (4.10), it holds

$$-\Delta \widetilde{\boldsymbol{W}}_{\varepsilon,K'+1}(\boldsymbol{v}_{\varepsilon}^{*}) + \nabla \widetilde{Q}_{\varepsilon,K'+1}(q_{\varepsilon}^{*},\boldsymbol{v}_{\varepsilon}^{*}) = \sum_{k=0}^{K'} \varepsilon^{k-2} M^{k} \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon}^{*} + \nabla q_{\varepsilon}^{*} + O_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})}(\varepsilon^{K'+1})$$
$$\operatorname{div}(\widetilde{\boldsymbol{W}}_{\varepsilon,K'+1}) = O_{L^{2}(D_{\varepsilon})}(\varepsilon^{K'+3}),$$

where we have used (4.12) to estimate the right-hand side terms. Applying now Corollary 4.4 to $(\boldsymbol{v}, \phi) \equiv (\boldsymbol{u}_{\varepsilon} - \widetilde{\boldsymbol{W}}_{\varepsilon,K'+1}(\boldsymbol{v}_{\varepsilon}^*), p_{\varepsilon} - \widetilde{Q}_{\varepsilon,K'}(\boldsymbol{q}_{\varepsilon}^*, \boldsymbol{v}_{\varepsilon}^*))$ yields the error estimate $\left\| \nabla (\boldsymbol{u}_{\varepsilon} - \widetilde{\boldsymbol{W}}_{\varepsilon,K'+1}(\boldsymbol{v}_{\varepsilon}^*)) \right\| + \varepsilon \left\| p_{\varepsilon} - \widetilde{Q}_{\varepsilon,K'+1}(\boldsymbol{q}_{\varepsilon}^*, \boldsymbol{v}_{\varepsilon}^*) \right\| \leq C_{K'}(\boldsymbol{f})\varepsilon^{K'+2}.$

$$\left\| \nabla (\boldsymbol{u}_{\varepsilon} - \boldsymbol{W}_{\varepsilon,K'+1}(\boldsymbol{v}_{\varepsilon}^{*})) \right\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} + \varepsilon \left\| p_{\varepsilon} - Q_{\varepsilon,K'+1}(q_{\varepsilon}^{*},\boldsymbol{v}_{\varepsilon}^{*}) \right\|_{L^{2}(D_{\varepsilon})} \leq C_{K'}(\boldsymbol{f})\varepsilon^{K+2}$$

Finally, remarking that the highest order terms are already of order $O(\varepsilon^{K'+2})$, i.e.

$$\begin{aligned} \|\nabla(\widetilde{\boldsymbol{W}}_{\varepsilon,K'+1}(\boldsymbol{v}_{\varepsilon}^{*}) - \widetilde{\boldsymbol{W}}_{\varepsilon,K'}(\boldsymbol{v}_{\varepsilon}^{*}))\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} &\leq C_{K'}\varepsilon^{K'+2}, \\ \varepsilon\|\widetilde{Q}_{\varepsilon,K'+1}(\boldsymbol{q}_{\varepsilon}^{*},\boldsymbol{v}_{\varepsilon}^{*}) - \widetilde{Q}_{\varepsilon,K'-1}(\boldsymbol{q}_{\varepsilon}^{*},\boldsymbol{v}_{\varepsilon}^{*})\|_{L^{2}(D_{\varepsilon})} &\leq C_{K'}\varepsilon^{K'+2}, \end{aligned}$$

510 we obtain the result by using the triangle's inequality.

511 Remark 4.6. We need only K-1 derivatives in the truncated criminal ansatz 512 $Q_{\varepsilon,K-1}(\boldsymbol{v}_{\varepsilon}^*, q_{\varepsilon}^*)$ for the pressure (eqn. (4.2)), because the term of highest order has a 513 norm of order ε^K while $\boldsymbol{v}_{\varepsilon}^*$ is of order ε^2 by the assumption (4.12).

514 Remark 4.7. As a result of the scaling ε^{-1} in Corollary 4.4, we pay a factor ε^{-1} 515 in the error induced by the non zero divergence constraint. However we are able to 516 obtain the right order of ε in the error estimates of Proposition 4.5 thanks to the use 517 of higher order terms of the ansatz (3.21) which are removed at the end of the proof. 518 This strategy is quite classical in the truncation analysis of two-scale expansions, see 519 e.g. [26, 11].

4.2. Construction of a well-posed higher order effective models by a minimization principle. We now derive the well-posed homogenized equation (1.6) of (finite) order 2K + 2 by following the variational method introduced by Smyshlyaev and Cherednychenko in [53] and used in the further works [28, 27, 34]. In the present context, of the Stokes system (1.1), we shall see that (1.6) can be obtained as (4.6) with K' = K, which yields error estimates of order $O(\varepsilon^{K+3})$ in the $L^2(D_{\varepsilon})$ norm.

Recall that the velocity u_{ε} solution to the Stokes system (1.1) is the unique minimizer of the constrained minimization problem

$$u_{\varepsilon} = \arg \min_{\boldsymbol{w} \in H^{1}(D_{\varepsilon}, \mathbb{R}^{d})} \quad J(\boldsymbol{w}, \boldsymbol{f}) := \int_{D} \left(\frac{1}{2} \nabla \boldsymbol{w} : \nabla \boldsymbol{w} - \boldsymbol{f} \cdot \boldsymbol{w}\right) dy$$
528 (4.13)
$$s.t. \begin{cases} \operatorname{div}(\boldsymbol{w}) = 0 \text{ in } D_{\varepsilon} \\ \boldsymbol{w} = 0 \text{ on } \partial \omega_{\varepsilon} \\ \boldsymbol{w} \text{ is } D\text{-periodic.} \end{cases}$$

In the context of the homogenization of linearized elasticity, the main idea of the method of [53] is to restrict (4.13) to functions of the form $\boldsymbol{w} = \widetilde{\boldsymbol{W}}_{\varepsilon,K}(\boldsymbol{v})$ given by (4.1), where $\boldsymbol{v} \in H^{K+1}(D, \mathbb{R}^d)$ is an unknown function sought to approximate $\boldsymbol{u}_{\varepsilon}^*$. In the present setting, we consider the following approximation of (4.13):

533 (4.14)
$$\min_{\boldsymbol{v}\in H^{K+1}(D,\mathbb{R}^d)} \quad J(\boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v}),\boldsymbol{f})$$
$$s.t. \quad \begin{cases} \operatorname{div}(\boldsymbol{v}) = 0 \text{ in } D, \\ \boldsymbol{v} \text{ is } D\text{-periodic.} \end{cases}$$

534 Note that (4.14) is not exactly the restriction of (4.13) to such functions $W_{\varepsilon,K}(v)$

because div $(\widetilde{W}_{\varepsilon,K}(\boldsymbol{v})) \neq 0$ (it is of order ε^K , see (4.5)). The next step of the process is to eliminate the oscillating variable x/ε in $J(\widetilde{W}_{\varepsilon,K}(\boldsymbol{v}), \boldsymbol{f})$ so as to obtain an effective energy $J_K^*(\boldsymbol{v}, \boldsymbol{f}, \varepsilon) \simeq J(\widetilde{W}_{\varepsilon,K}(\boldsymbol{v}), \boldsymbol{f})$ which does not involve oscillating functions. Such is achieved thanks to the classical lemma of two-scale convergence [9].

LEMMA 4.8. Let ϕ be a $P = (0,1)^d$ -periodic function and $f \in \mathcal{C}^{\infty}_{per}(D)$ be a smooth D-periodic function. Then for any $p \in \mathbb{N}$, there exists a constant $C_p(f,\phi)$ independent of ε such that:

$$\left|\int_D f(x)\phi(x/\varepsilon)\mathrm{d}x - \int_D \int_P f(x)\phi(y)\mathrm{d}y\mathrm{d}x\right| \le C_p(f,\phi)\varepsilon^p.$$

539 *Proof.* See Appendix C. of [53] or Lemma 7.3 in [33].

Applying Lemma 4.8 to (4.14) in order to pass to the limit in the terms of $J(\widetilde{W}_{\varepsilon,K}(\boldsymbol{v}), \boldsymbol{f})$ which depends on the oscillating variable x/ε , we obtain the existence of a functional J_K^* such that for any $\boldsymbol{v} \in \mathcal{C}_{per}^{\infty}(D, \mathbb{R}^d)$, it holds

$$J(\boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v}),\boldsymbol{f}) = J_K^*(\boldsymbol{v},\boldsymbol{f},\varepsilon) + o(\varepsilon^p)$$

540 with $p \in \mathbb{N}$ arbitrarily large. The functional J_K^* is given explicitly by

541 (4.15)
$$J_K^*(\boldsymbol{v}, \boldsymbol{f}, \varepsilon) := \int_D \int_P \frac{1}{2} \left\| (\nabla_x + \varepsilon^{-1} \nabla_y) \left(\boldsymbol{W}_{\varepsilon, K}(\boldsymbol{v})(x, y) \right) \right\|^2 \mathrm{d}y \mathrm{d}x - \int_D \boldsymbol{f} \cdot \boldsymbol{v} \mathrm{d}x,$$

542 where we have used the point 1 of Proposition 3.10 to simplify the linear part of

the energy. Replacing $J(\widetilde{W}_{\varepsilon,K}(\boldsymbol{v}), \boldsymbol{f})$ by $J_K^*(\boldsymbol{v}, \boldsymbol{f}, \varepsilon)$ in (4.14) allows to obtain the homogenized equation (1.6) of order 2K + 2:

DEFINITION 4.9. For any $K \in \mathbb{N}$, we call homogenized equation of order 2K + 2545associated with the Stokes system (1.1) the Euler-Lagrange equation of the minimiza-546tion problem 547

548 (4.16)
$$\min \begin{array}{c} J_K^*(\boldsymbol{v}, \boldsymbol{f}, \varepsilon) \\ \boldsymbol{v} \in H^{K+1}(D, \mathbb{R}^d), \\ \operatorname{div}(\boldsymbol{v}) = 0 \ in \ D, \\ \boldsymbol{v} \ is \ D-periodic. \end{array}$$

This Euler-Lagrange equation can be written as (1.6) where the constant (matrix val-549ued) tensors \mathbb{D}_{K}^{k} are inferred from (4.15) and where $(\boldsymbol{v}_{\varepsilon,K}^{*}, q_{\varepsilon,K}^{*}) \in H^{K+1}(D, \mathbb{R}^{d}) \times$ 550 $L^{2}(D)$ defines the higher order homogenized solution. 551

The next two propositions verify that (1.6) is indeed a "good" candidate effective 552model, by relating the coefficients \mathbb{D}_{K}^{k} to the tensors M^{k} (in view of (4.6)), and by 553establishing the well-posedness of (1.6). 554

PROPOSITION 4.10. The coefficients of the matrix valued tensor \mathbb{D}_{K}^{k} are explicitly 555given for any $1 \leq i, j \leq d$ by: 556

557 (4.17)
$$\mathbb{D}_{K,ij}^{k} = \begin{cases} M^{k} \text{ if } 0 \le k \le K \\ M^{k} + \mathbb{A}_{K}^{k} \text{ if } K + 1 \le k \le 2K + 1 \\ (-1)^{K+1} \int_{Y} \mathbf{N}_{i}^{K} \cdot \mathbf{N}_{j}^{K} \otimes I \mathrm{d}y \text{ if } k = 2K + 2. \end{cases}$$

where the matrix valued tensor \mathbb{A}_{K}^{k} is given for any $K+1 \leq k \leq 2K+1$ by 558

559 (4.18)
$$\mathbb{A}_{K,ij}^k := (-1)^{K+1} \int_Y (\nabla \beta_j^{k-K-1} \cdot \mathbf{N}_i^{K+1} + (-1)^k \nabla \beta_i^{k-K-1} \cdot \mathbf{N}_j^{K+1}) \mathrm{d}y.$$

Proof. Let us denote by V_K the space 560

561 (4.19)
$$V_K := \{ \boldsymbol{v} \in H^{K+1}(D, \mathbb{R}^d) \, | \, \operatorname{div}(\boldsymbol{v}) = 0 \text{ and } \boldsymbol{v} \text{ is } D\text{-periodic} \}.$$

We identify the coefficients \mathbb{D}_{K}^{k} by computing the Euler-Lagrange equation associated 562with (4.15). For any $(\boldsymbol{v}, \boldsymbol{w}) \in V_K$, it holds, in a distributional sense: 563

564
$$\int_{D} \int_{P} (\nabla_{x} + \varepsilon^{-1} \nabla_{y}) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v}) : (\nabla_{x} + \varepsilon^{-1} \nabla_{y}) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w}) \mathrm{d}x \mathrm{d}y$$

565
$$= \int \int [(-\Delta_{xx} - \varepsilon^{-1} \Delta_{xy} - \varepsilon^{-2} \Delta_{yy}) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{v})]$$

565

566
$$\int_D \int_P^{1} (\nabla_x + \varepsilon^{-1} \nabla_y) Q_{\varepsilon,K}(\boldsymbol{v}, 0)] \cdot \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w}) dx dy$$

567
568 +
$$\int_D \int_Y Q_{\varepsilon,K}(\boldsymbol{v},0) [(\operatorname{div}_x + \varepsilon^{-1} \operatorname{div}_y) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w})] \mathrm{d}x \mathrm{d}y.$$

By using (4.1) and the point (i) of Proposition 3.10, the above quantity is equal to 569

570
$$\int_{D} \int_{Y} \left[\sum_{k=0}^{K} \varepsilon^{k-2} M^{k} \nabla^{k} \boldsymbol{v}(x) \right] \cdot \sum_{k=0}^{K} \varepsilon^{k} N^{k}(y) \cdot \nabla^{k} \boldsymbol{w}(x) \mathrm{d}x \mathrm{d}y$$

571
$$- \int_{D} \int_{Y} \varepsilon^{K-1} \left[(2\partial_{l} \boldsymbol{N}_{j}^{K}(y) - \beta_{j}^{K}(y) \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l}) \cdot \nabla^{K+1} v_{j}(x) \right] \cdot \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w}) \mathrm{d}x \mathrm{d}y$$

572
$$-\int_{D}\int_{Y} [\varepsilon^{K-1}(\boldsymbol{N}_{j}^{K-1}(y)\otimes I)\cdot\nabla^{K+1}\boldsymbol{v}(x)]\cdot\boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w})\mathrm{d}x\mathrm{d}y$$

573
$$-\int_{D}\int_{Y} [\varepsilon^{K}(\boldsymbol{N}_{j}^{K}(y)\otimes I)\cdot\nabla^{K+2}\boldsymbol{v}(x)]\cdot\boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w})\mathrm{d}x\mathrm{d}y$$

573

574
$$+ \int_D \int_Y \sum_{k=0}^K \varepsilon^{k+K-1} (\beta_j^k(y) \cdot \nabla^k v_j(x)) (\boldsymbol{N}_i^K(y) \cdot \boldsymbol{e}_l \otimes \boldsymbol{e}_l \cdot \nabla^{K+1} w_i(x)) \mathrm{d}y \mathrm{d}x$$

575
$$= \int_{D} \left(\sum_{k=0}^{K} \varepsilon^{k-2} M^{k} \cdot \nabla^{k} \boldsymbol{v} + \sum_{k=K+1}^{2K+2} \varepsilon^{k-2} \mathbb{D}_{K}^{k} \cdot \nabla^{k} \boldsymbol{v} \right) \cdot \boldsymbol{w} \mathrm{d}x$$

where we identify (by integration by parts) $\mathbb{D}_{K,ij}^{2K+2} := -(-1)^K \int_Y \mathbf{N}_i^K \cdot \mathbf{N}_j^K \otimes I dy$ as claimed. The coefficients of the tensor \mathbb{D}_K^k are given for $K + 1 \leq k \leq 2K + 1$ by

$$\begin{split} \mathbb{D}_{K,ij}^{k} &= -(-1)^{k-K-1} \int_{Y} ((2\partial_{l} \boldsymbol{N}_{j}^{K} - \beta_{j}^{K} \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \boldsymbol{N}_{j}^{K-1} \otimes \boldsymbol{I}) \cdot \boldsymbol{N}_{i}^{k-K-1} \mathrm{d}\boldsymbol{y} \\ &- (-1)^{k-K-2} \int_{Y} \boldsymbol{N}_{j}^{K} \cdot \boldsymbol{N}_{i}^{k-K-2} \otimes \boldsymbol{I} \mathrm{d}\boldsymbol{y} + (-1)^{K+1} \int_{Y} \beta_{j}^{k-K-1} \boldsymbol{N}_{i}^{K} \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} \mathrm{d}\boldsymbol{y} \\ &= -(-1)^{K+1} \int_{Y} (-1)^{k} (-\Delta_{yy} \boldsymbol{N}_{j}^{K+1} + \nabla \beta_{j}^{K+1} - \boldsymbol{M}^{K+1} \boldsymbol{e}_{j}) \cdot \boldsymbol{N}_{i}^{k-K-1} \mathrm{d}\boldsymbol{y} \\ &+ (-1)^{K+1} \int_{Y} [(-1)^{k} \boldsymbol{N}_{j}^{K} \cdot \boldsymbol{N}_{i}^{k-K-2} \otimes \boldsymbol{I} + \beta_{j}^{k-K-1} \boldsymbol{N}_{i}^{K} \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}] \mathrm{d}\boldsymbol{y}, \end{split}$$

where we have used extensively Proposition 3.10. We now distinguish two cases: 577 1. if k = K + 1, then the above expression reads

$$\begin{split} \mathbb{D}_{K,ij}^{K+1} &= (M^{K+1}\boldsymbol{e}_j) \cdot \boldsymbol{e}_i - \int_Y \boldsymbol{N}_j^{K+1} \cdot (-\Delta_{yy}\boldsymbol{N}_i^0 + \nabla\beta_i^0) \mathrm{d}y + \int_Y \nabla\beta_i^0 \cdot \boldsymbol{N}_j^{K+1} \\ &+ (-1)^{K+1} \int_Y \nabla\beta_j^{k-K-1} \cdot \boldsymbol{N}_i^{K+1} \mathrm{d}y = M_{ij}^{K+1} + \mathbb{A}_{ij}^{K+1}, \end{split}$$

578whence the result for k = K + 1;

2. if $K + 2 \le k \le 2K + 1$, then we read instead

$$\begin{split} \mathbb{D}_{K,ij}^{k} &= (-1)^{k+K} \int_{Y} ((-\Delta_{yy} \boldsymbol{N}_{j}^{K+1} + \nabla \beta_{j}^{K+1}) \cdot \boldsymbol{N}_{i}^{k-K-1} + \nabla \beta_{i}^{k-K-1} \cdot \boldsymbol{N}_{j}^{K+1}) \mathrm{d}y \\ &- (-1)^{k+K} \int_{Y} \boldsymbol{N}_{j}^{K} \cdot \boldsymbol{N}_{i}^{k-K-2} \otimes I \mathrm{d}y + \mathbb{A}_{ij}^{k}. \end{split}$$

Applying finally Corollary 3.11 with p = K + 1, we obtain that the first two term of 579the above equation are equal to $(-1)^k M_{ii}^k = M_{ij}^k$, which yields the final result. 580

Remark 4.11. In view of the proof of Proposition 4.10, it is possible to show that 581(1.11) (with the same definition for the leading coefficient \mathbb{D}_{K}^{2K+2}) is the strong form of the following "mixed" variational formulation: find $(\hat{v}_{\varepsilon,K}^{*}, \hat{q}_{\varepsilon,K}^{*}) \in V_{K} \times L^{2}(D)$ such 582583 that for any $(\boldsymbol{w}, \phi) \in V_K \times L^2(D)$, 584

586 (4.20)
$$\int_{D} \int_{Y} \left[(\varepsilon^{-1} \nabla_{y} + \nabla_{x}) \boldsymbol{W}_{\varepsilon,K}(\widehat{\boldsymbol{v}}_{\varepsilon,K}^{*}) \right] : \left[(\varepsilon^{-1} \nabla_{y} + \nabla_{x}) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w}) \right] dy dx$$

587
$$- \int_{D} \int_{Y} Q_{\varepsilon,K}(\widehat{\boldsymbol{v}}_{\varepsilon,K}^{*}, \widehat{q}_{\varepsilon,K}^{*}) \left[(\varepsilon^{-1} \operatorname{div}_{y} + \operatorname{div}_{x}) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{w}) \right] dy dx$$

588
589
$$-\int_D \int_Y Q_{\varepsilon,K}(\boldsymbol{w},\phi) \big[(\varepsilon^{-1} \mathrm{div}_y + \mathrm{div}_x) \boldsymbol{W}_{\varepsilon,K}(\boldsymbol{\hat{v}}_{\varepsilon,K}^*) \big] \mathrm{d}y \mathrm{d}x = \int_D \boldsymbol{f} \cdot \boldsymbol{w} \mathrm{d}x.$$

This result is quite surprising. Indeed, (4.20) is built from the truncated ansatz 590 $(W_{\varepsilon,K}(\widehat{v}_{\varepsilon,K}^*), Q_{\varepsilon,K}(\widehat{v}_{\varepsilon,K}^*, \widehat{q}_{\varepsilon,K}^*))$ which is expected to yield approximations of $(u_{\varepsilon}, p_{\varepsilon})$ 591 at order $O(\varepsilon^{K+3})$ only (eqn. (4.1) and (4.2)). However, the strong form (1.11) turns out to exhibit 2K + 1 "correct" coefficients M^k . As a result, if (1.11) is well-593 posed, the reconstructed oscillating functions $W_{\varepsilon,2K+1}(\widehat{v}_{\varepsilon,K}^*)$ and $Q_{\varepsilon,2K}(\widehat{v}_{\varepsilon,K}^*,\widehat{q}_{\varepsilon,K}^*)$ 594approximate u_{ε} and p_{ε} with an error rate as good as $O(\varepsilon^{2K+4})$ in the $L^2(D_{\varepsilon})$ norm (Proposition 4.5). This improvement (which had not been noticed in the original 596597paper [53]) actually holds in the context of the Poisson or elasticity equations for which there is no difference between (1.6) and (1.11) (see [34, 33]). Unfortunately in 598 the case of the Stokes system, we do not know whether the mixed formulation (4.20)599with the "velocity-dependent" pressure $Q_{\varepsilon,K}(\boldsymbol{w},\phi)$ yields a well-posed problem, hence 600 our commitment to consider (1.6) instead of (1.11). 601

The leading tensor \mathbb{D}_{K}^{2K+2} is nonnegative according to (4.17). Under a rather unrestrictive additional non-degeneracy assumption, we obtain that the minimization principle (4.16) makes (1.6) be a well posed problem.

PROPOSITION 4.12. Assume the dominant tensor $\mathbb{D}_{K}^{2K+2} = (-1)^{K+1} \mathbb{B}_{K}^{K+1,K+1}$ is non-degenerate, that is there exists a constant $\nu > 0$ such that for any constant vector tensor $\boldsymbol{\xi}^{K+1} = \boldsymbol{\xi}_{i_{1}...i_{K+1},p}^{K+1} \in \mathbb{R}^{d^{K+1}} \times \mathbb{R}^{d}$ of order K + 1, it holds

608 (4.21)
$$\int_{Y} [(N^{K} \otimes e_{l}) \cdot \boldsymbol{\xi}^{K+1}] \cdot [(N^{K} \otimes e_{l}) \cdot \boldsymbol{\xi}^{K+1}] dy \ge \nu \, \boldsymbol{\xi}_{i_{1}...i_{K+1}}^{K+1} \cdot \boldsymbol{\xi}_{i_{1}...i_{K+1}}^{K+1}$$

Then there exists a unique velocity and pressure couple $(\mathbf{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*) \in H^{K+1}(D, \mathbb{R}^d) \times L^2(D)/\mathbb{R}$ solving the higher order homogenized equation (1.6).

611 *Proof.* The proof relies on the positivity of the quadratic part of the energy 612 $J_K^*(\boldsymbol{v}, \boldsymbol{f}, \varepsilon)$. By adapting the arguments of the proof of Proposition 12 in [34], we 613 obtain indeed that the bilinear form associated with the energy (4.15) is coercive on 614 the space V_K defined in (4.19). This is enough to apply standard theory for sad-615 dle point problems involving the zero divergence constraint (see e.g. the textbooks 616 [55, 54, 38, 31]) which ensures the existence and uniqueness of a solution for (1.6).

617 Remark 4.13. The assumption (4.21) could fail for $K \ge 1$ in case the obstacle ηT 618 is invariant along some of the directions e_i of the cell P, however it is not restrictive. 619 Indeed, since the leading order tensor \mathbb{D}_K^{2K+2} has no influence on the error estimates 620 of Proposition 4.5, it is always possible to add to \mathbb{D}_K^{2K+2} a small non-negative tensor 621 making it non-degenerate.

4.3. Error estimates: justification of the homogenization process. We conclude this section by stating error estimates holding for the solution $(\boldsymbol{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*)$ to the high order homogenized model (1.6). We know from Proposition 4.10 that $\mathbb{D}_{K}^{k} = M^{k}$ for $0 \leq k \leq K$, therefore the assumptions of Proposition 4.5 are satisfied provided we verify the uniform regularity estimate (4.12).

LEMMA 4.14. The solution $(\boldsymbol{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*)$ of (1.6) is smooth and for any $m \in \mathbb{N}$, there exists a constant $C_m(\boldsymbol{f})$ depending only on m and \boldsymbol{f} such that

$$||\boldsymbol{v}_{\varepsilon,K}^*||_{H^m(D,\mathbb{R}^d)} \leq C_m(\boldsymbol{f})\varepsilon^2.$$

627 Proof. This result can be obtained by solving (1.6) explicitly with Fourier series 628 in the periodic domain D and by adapting the proof of Lemma 5 in [34].

Since we have verified that all the assumptions of Proposition 4.5 hold with K' = K, we are finally in position to state the following error bounds.

COROLLARY 4.15. Let $(\boldsymbol{v}_{\varepsilon,K}^*, q_{\varepsilon,K}^*)$ be the unique solution to the high order homogenized equation (1.6). There exists a constant $C_K(\boldsymbol{f})$ independent of ε (but depending on K, \boldsymbol{f} , and a priori on the shape of the hole (ηT)) such that the following error estimates hold:

$$\begin{split} \left\| \boldsymbol{u}_{\varepsilon} - \sum_{k=0}^{K} \varepsilon^{k} N^{k}(\cdot/\varepsilon) \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon,K}^{*} \right\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d})} &\leq C_{K}(\boldsymbol{f}) \varepsilon^{K+3}, \\ \left\| \nabla \left(\boldsymbol{u}_{\varepsilon} - \sum_{k=0}^{K} \varepsilon^{k} N^{k}(\cdot/\varepsilon) \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon,K}^{*} \right) \right\|_{L^{2}(D_{\varepsilon},\mathbb{R}^{d\times d})} &\leq C_{K}(\boldsymbol{f}) \varepsilon^{K+2}, \\ \left| p_{\varepsilon} - \left(q_{\varepsilon,K}^{*} + \sum_{k=0}^{K-1} \varepsilon^{k-1} \boldsymbol{\beta}^{k}(\cdot/\varepsilon) \cdot \nabla^{k} \boldsymbol{v}_{\varepsilon,K}^{*} \right) \right\|_{L^{2}(D_{\varepsilon})} &\leq C_{K}(\boldsymbol{f}) \varepsilon^{K+1}. \end{split}$$

631 Remark 4.16. As the reader may expect, error bounds with the same order of 632 convergence hold for the truncation at order K of the "classical" ansatz (1.7), see 633 [26, 46] up to the order K = 1, and in Proposition 7.37 of [33] at all orders.

5. Low volume fraction limits when the scaling η of the obstacle vanishes. In this section, we provide evidences that (1.6) is "well-behaved" in the sense that it has the potential to capture the homogenized regimes (1.2)–(1.4) in the low volume fraction limit where the size of the obstacles vanish. Our results supporting this claim are obtained by analyzing the asymptotics of the tensors \mathcal{X}^{k*} , M^k and \mathbb{D}_K^{2K+2} as the scaling ratio η converges to 0.

In this whole subsection, we assume for simplicity, that the space dimension is greater than 3:

 $d \geq 3.$

640 Similar results are expected to hold in dimension d = 2 but would require a different

treatment, as e.g. in [6, 41]. The hole ηT is assumed to be an non-empty open subset strictly included in the unit cell for any $\eta \leq 1$ (it does not touch the boundary): $\eta T \subset \subset P$.

Let us recall the definition of the Deny-Lions (or Beppo-Levi) space denoted by $\mathcal{D}^{1,2}(\mathbb{R}^d \setminus T, \mathbb{R}^d)$ (the reader is referred to [6, 4, 8] and also [48], p.59. for more details).

DEFINITION 5.1 (Deny-Lions space). The Deny-Lions space $\mathcal{D}^{1,2}(\mathbb{R}^d \setminus T, \mathbb{R}^d)$ is the completion of the space of smooth vector fields by the L^2 norm of their gradients:

$$\mathcal{D}^{1,2}(\mathbb{R}^d \backslash T, \mathbb{R}^d) := \overline{\mathcal{D}(\mathbb{R}^d \backslash T, \mathbb{R}^d)}^{||\nabla \cdot ||_{L^2(\mathbb{R}^d \backslash T, \mathbb{R}^d)}}$$

646 where $\mathcal{D}(\mathbb{R}^d \setminus T, \mathbb{R}^d)$ is the space of compactly supported smooth vector fields. When 647 $d \geq 3$, it is admits the following characterization:

649 $\mathcal{D}^{1,2}(\mathbb{R}^d \setminus T, \mathbb{R}^d)$

648

 $\underset{n \in \mathbb{Z}}{\text{for } measurable } |||v||_{L^{2d/(d-2)}(\mathbb{R}^d \setminus T, \mathbb{R}^d)} < +\infty \text{ and } ||\nabla v||_{L^2(\mathbb{R}^d \setminus T, \mathbb{R}^d \times d)} < +\infty \}.$

For any $1 \leq j \leq d$, we consider the unique solution (Ψ_j, σ_j) to the exterior Stokes problem

654 (5.1)
$$\begin{cases} -\Delta \Psi_j + \nabla \sigma_j = 0 \text{ in } \mathbb{R}^d \setminus T \\ \operatorname{div}(\Psi_j) = 0 \text{ in } \mathbb{R}^d \setminus T \\ \Psi_j = 0 \text{ on } \partial T \\ \Psi_j \to e_j \text{ at } \infty \\ \sigma_i \in L^2(\mathbb{R}^d \setminus T). \end{cases}$$

The convergence condition $\Psi_j \to e_j$ at infinity must be understood in the sense that $\Psi_j - e_j$ belongs to $\mathcal{D}^{1,2}(\mathbb{R}^d \setminus T, \mathbb{R}^d)$. Similarly, the pressures $(\sigma_j)_{1 \leq j \leq d}$ are uniquely determined by the condition $\sigma_j \in L^2(\mathbb{R}^d \setminus T)$ (see e.g. Lemma 1.1, article V. of [36]).

658 We denote by $F := (F_{ij})_{1 \le i,j \le d}$ the matrix collecting the drag force components:

659 (5.2)
$$F_{ij} := \int_{\mathbb{R}^d \setminus T} \nabla \Psi_i : \nabla \Psi_j \, \mathrm{d}x = -\int_{\partial T} \boldsymbol{e}_j \cdot (\nabla \Psi_i - \sigma_i I) \cdot \boldsymbol{n} \mathrm{d}s,$$

660 where the normal \boldsymbol{n} is pointing *inward* T.

661 **5.1. Technical estimates in the growing periodic domain** $\eta^{-1}P \setminus T$. In all 662 this section, vector fields of the rescaled cell $\eta^{-1}P$ are indicated by a tilde \tilde{v} notation. 663 For a given vector field $\tilde{v} \in L^2(\eta^{-1}P, \mathbb{R}^d)$, we denote by $\langle \tilde{v} \rangle$ the average $\langle \tilde{v} \rangle :=$ 664 $\eta^d \int_{n^{-1}P} \tilde{v}(y) dy$.

Let us recall that for any $\boldsymbol{v} \in H^1(P \setminus (\eta T), \mathbb{R}^d)$, if $\tilde{\boldsymbol{v}}$ is the rescaled function defined by $\tilde{\boldsymbol{v}}(y) := v(\eta y)$ in the rescaled cell $\eta^{-1}P \setminus T$, then the L^2 norms of $\boldsymbol{v}, \tilde{\boldsymbol{v}}$ and of their gradients are related by the following identities:

668
$$||\boldsymbol{v}||_{L^2(P\setminus(\eta T),\mathbb{R}^d)} = \eta^{d/2} ||\widetilde{\boldsymbol{v}}||_{L^2(\eta^{-1}P\setminus T,\mathbb{R}^d)}$$

$$\|\nabla \boldsymbol{v}\|_{L^2(P\setminus(\eta T),\mathbb{R}^{d\times d})} = \eta^{d/2-1} \|\nabla \widetilde{\boldsymbol{v}}\|_{L^2(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})}$$

The asymptotic behaviors of the tensors \mathcal{X}^{k*} , M^k are obtained by following the methodology of [6, 41, 34], which relies on several technical results stated in this part.

673 LEMMA 5.2. Assume $d \ge 3$. There exists a constant C > 0 independent of $\eta > 0$ 674 such that for any $\tilde{v} \in H^1(\eta^{-1}P \setminus T, \mathbb{R}^d)$ which vanishes on the hole ∂T and which is 675 $\eta^{-1}P$ periodic, the following inequalities hold:

676 (5.3) $||\widetilde{\boldsymbol{v}}||_{L^2(\eta^{-1}P\setminus T,\mathbb{R}^d)} \leq C\eta^{-d/2} ||\nabla\widetilde{\boldsymbol{v}}||_{L^2(\eta^{-1}P\setminus T,\mathbb{R}^d\times d)},$

677 (5.4)
$$|\langle \widetilde{\boldsymbol{v}} \rangle| \leq C ||\nabla \widetilde{\boldsymbol{v}}||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})}$$

678 (5.5) $||\widetilde{\boldsymbol{v}} - \langle \widetilde{\boldsymbol{v}} \rangle||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^n d)} \le C\eta^{-1} ||\nabla \widetilde{\boldsymbol{v}}||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})},$

 $\begin{cases} 870 \\ 880 \end{cases} (5.6) \qquad \qquad ||\widetilde{\boldsymbol{v}} - \langle \widetilde{\boldsymbol{v}} \rangle||_{L^{2d/(d-2)}(\eta^{-1}P \setminus T, \mathbb{R}^d)} \le C ||\nabla \widetilde{\boldsymbol{v}}||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^d \times d)}. \end{cases}$

681 *Proof.* See [41, 45, 6, 42].

EEMMA 5.3. Consider $\mathbf{h} \in L^2(\eta^{-1}P \setminus T, \mathbb{R}^d)$ and $g \in L^2(\eta^{-1}P \setminus T)$ a function satisfying $\int_{\eta^{-1}P \setminus T} g dx = 0$. Let $(\mathbf{v}, \phi) \in H^1(\eta^{-1}P \setminus T, \mathbb{R}^d) \times L^2(\eta^{-1}P \setminus T)$ be the

unique solution to the following Stokes system: 684

685 (5.7)
$$\begin{cases} -\Delta \boldsymbol{v} + \nabla \phi = \boldsymbol{h} \text{ in } \eta^{-1} P \setminus T \\ \operatorname{div}(\boldsymbol{v}) = g \text{ in } \eta^{-1} P \setminus T \\ \int_{\eta^{-1} P \setminus T} \phi \mathrm{d}x = 0 \\ \boldsymbol{v} = 0 \text{ on } \partial T \\ \boldsymbol{v} \text{ is } \eta^{-1} P \text{-periodic.} \end{cases}$$

There exists a constant C > 0 independent of η , h and g such that 686 687

688 (5.8)
$$||\nabla \boldsymbol{v}||_{L^2(\eta^{-1}P\setminus T, \mathbb{R}^{d\times d})} + ||\phi||_{L^2(\eta^{-1}P\setminus T)}$$

$$\leq C(\eta^{-1}||\boldsymbol{h} - \langle \boldsymbol{h} \rangle||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^d)} + \eta^{-d}|\langle \boldsymbol{h} \rangle| + ||g||_{L^2(\eta^{-1}P \setminus T)})$$

Proof. From Lemma 2.2.4 in [7], for any $\eta > 0$, there exists a linear "Bogovskii's" 691 operator B_{η} : $L^{2}(P \setminus (\eta T)) \to H^{1}(P \setminus (\eta T), \mathbb{R}^{d})$ satisfying for any $\phi \in L^{2}(P \setminus (\eta Y))$ 692 such that $\int_{P \setminus (\eta T)} \phi dy = 0$: 693

(i) $\operatorname{div}(B_\eta \phi) = \phi$, 694 (ii) $B_{\eta}\phi = 0$ on $\partial(\eta T)$, 695

(iii) $B_{\eta}\phi$ is *P*-periodic, 696

(iv) $||\nabla(B_\eta\phi)||_{L^2(P\setminus(\eta T),\mathbb{R}^{d\times d})} \leq C||\phi||_{L^2(P\setminus(\eta T))}$ for a constant C independent of η 697 698 and ϕ .

For any $\widetilde{\phi} \in L^2(\eta^{-1}P \setminus T)$ such that $\int_{\eta^{-1}P \setminus T} \widetilde{\phi} dy = 0$, we define

$$\widetilde{B}_{\eta}(\widetilde{\phi}) := \eta^{-1} \big[B_{\eta}(\widetilde{\phi}(\eta^{-1} \cdot))(\eta \cdot) \big].$$

The operator \widetilde{B}_{η} : $L^{2}(\eta^{-1}P \setminus T) \to H^{1}(\eta^{-1}P \setminus T, \mathbb{R}^{d})$ satisfies the following properties: 699 for any $\widetilde{\phi} \in L^2(\eta^{-1}P \setminus T)$ such that $\int_{\eta^{-1}P \setminus T} \widetilde{\phi} dx = 0$, 700

(i) $\operatorname{div}(\widetilde{B}_{\eta}\widetilde{\phi}) = \widetilde{\phi} \text{ in } \eta^{-1}P \backslash T$, 701

(ii) $\widetilde{B}_{\eta}\widetilde{\phi} = 0$ on ∂T , 702

(iii) $\widetilde{B}_{\eta}\widetilde{\phi}$ is $\eta^{-1}P$ -periodic, 703

(iv) $||\nabla(\tilde{B}_{\eta}\dot{\phi})||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})} \leq C||\widetilde{\phi}||_{L^{2}(\eta^{-1}P\setminus T)}$ for a constant C independent of η 704 and ϕ . 705

The proof follows then classically along the lines of Corollary 4.4. Upon an integration 706by parts and by using Lemma 5.2, it is readily obtained with $\boldsymbol{w} := \boldsymbol{v} - B_{\eta}g$: 707

708
$$||\nabla \boldsymbol{w}||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})}^{2} = \int_{\eta^{-1}P\setminus T} \boldsymbol{h} \cdot \boldsymbol{w} \mathrm{d}\boldsymbol{y}$$

$$= \int_{\eta^{-1}P \setminus T} (\boldsymbol{h} - \langle \boldsymbol{h} \rangle) \cdot (\boldsymbol{w} - \langle \boldsymbol{w} \rangle) \mathrm{d}y + \int_{\eta^{-1}P \setminus T} \langle \boldsymbol{h} \rangle \cdot \langle \boldsymbol{w} \rangle \mathrm{d}y$$

710
$$\leq C(||\boldsymbol{h} - \langle \boldsymbol{h} \rangle||_{L^{2}(\eta^{-1}P \setminus T, \mathbb{R}^{d})}||\boldsymbol{w} - \langle \boldsymbol{w} \rangle||_{L^{2}(\eta^{-1}P \setminus T, \mathbb{R}^{d})} + \eta^{-d}|\langle \boldsymbol{h} \rangle||\langle \boldsymbol{w} \rangle|)$$

$$\mathbb{F}_{112}^{11} \quad (5.9) \qquad \leq C(\eta^{-1} || \boldsymbol{h} - \langle \boldsymbol{h} \rangle ||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^d)} + \eta^{-d} |\langle \boldsymbol{h} \rangle |) || \nabla \boldsymbol{w} ||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^d \times d)}$$

for a constant C > 0 independent of η and h. This implies

$$\begin{aligned} ||\nabla \boldsymbol{v}||_{L^{2}(\eta^{-1}P\setminus T)} &\leq ||\nabla \boldsymbol{w}||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})} + ||\nabla(\widetilde{B}_{\eta}g)||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})} \\ &\leq C(||\nabla \boldsymbol{w}||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})} + ||g||_{L^{2}(\eta^{-1}P\setminus T)}), \end{aligned}$$

whence the bound on $||\nabla v||_{L^2(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})}$ by using (5.9). The bound for the pressure is obtained by writing

715

28

716
$$\int_{\eta^{-1}P\setminus T} \phi^2 \mathrm{d}x = \int_{\eta^{-1}P\setminus T} \phi \mathrm{div}(B_\eta \phi) \mathrm{d}x$$

$$71'$$

 718

$$= -\int_{\eta^{-1}P\setminus T} \nabla\phi \cdot B_{\eta}\phi dx = \int_{\eta^{-1}P\setminus T} (\nabla \boldsymbol{v}: \nabla(B_{\eta}\phi) - \boldsymbol{h}\cdot(B_{\eta}\phi)) dx$$

from where (5.8) follows analogously.

5.2. Asymptotic convergences of homogenized tensors in the low volume fraction limit $\eta \to 0$. The asymptotics of the corrector tensors $(\mathcal{X}_j^0, \alpha_j^0)$ and of \mathcal{X}^{0*} (defined in (3.3) and (3.7)) have been obtained in of Theorem 3.1 in [6]. The following proposition extends this result to the whole family of tensors $(\mathcal{X}^k, \alpha_j^k)_{k \in \mathbb{N}}$ and $(\mathcal{X}^{k*})_{k \in \mathbb{N}}$.

PROPOSITION 5.4. Assume $d \geq 3$. For any $k \geq 0$ and $1 \leq j \leq d$, denote by $(\widetilde{\boldsymbol{\mathcal{X}}}_{j}^{2k}, \widetilde{\alpha}_{j}^{2k})$ and $(\widetilde{\boldsymbol{\mathcal{X}}}_{j}^{2k+1}, \widetilde{\alpha}_{j}^{2k+1})$ the rescaled tensors in $\eta^{-1}P \setminus T$ defined by

$$\begin{cases} \widetilde{\boldsymbol{\mathcal{X}}}_{j}^{2k}(x) := \eta^{(d-2)(k+1)} \boldsymbol{\mathcal{X}}_{j}^{2k}(\eta x) \\ \widetilde{\alpha}_{j}^{2k}(x) := \eta^{(d-2)(k+1)-1} \alpha_{j}^{2k}(\eta x) \end{cases} \begin{cases} \widetilde{\boldsymbol{\mathcal{X}}}_{j}^{2k+1}(x) := \eta^{(d-2)(k+1)} \boldsymbol{\mathcal{X}}_{j}^{2k+1}(\eta x) \\ \widetilde{\alpha}_{j}^{2k+1}(x) := \eta^{(d-2)(k+1)-1} \alpha_{j}^{2k+1}(\eta x) \end{cases}$$

725 for any $x \in \eta^{-1}P \setminus T$. Then:

1. there exists a constant C independent of $\eta > 0$ such that

$$\begin{aligned} \forall \eta > 0, \ ||\nabla \widetilde{\boldsymbol{\mathcal{X}}}_{j}^{2k}||_{L^{2}(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})} + ||\widetilde{\alpha}_{j}^{2k}||_{L^{2}(\eta^{-1}P \setminus T)} \leq C, \\ \forall \eta > 0, \ ||\nabla \widetilde{\boldsymbol{\mathcal{X}}}_{j}^{2k+1}||_{L^{2}(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})} + ||\widetilde{\alpha}_{j}^{2k+1}||_{L^{2}(\eta^{-1}P \setminus T)} \leq C; \end{aligned}$$

726 2. the following convergences hold as $\eta \to 0$:

727 (5.10)
$$(\widetilde{\boldsymbol{\chi}}_{i}^{2k}, \widetilde{\alpha}_{i}^{2k}) \rightarrow (c_{ij}^{2k} \boldsymbol{\Psi}_{j}, c_{ij}^{2k} \sigma_{j})$$
 weakly in $H^{1}_{loc}(\mathbb{R}^{d} \setminus T, \mathbb{R}^{d}) \times L^{2}_{loc}(\mathbb{R}^{d} \setminus T),$
728 (5.11) $(\widetilde{\boldsymbol{\chi}}_{i}^{2k+1}, \widetilde{\alpha}_{i}^{2k+1}) \rightarrow (0, 0)$ weakly in $H^{1}_{loc}(\mathbb{R}^{d} \setminus T, \mathbb{R}^{d}) \times L^{2}_{loc}(\mathbb{R}^{d} \setminus T),$

728 (5.11)
$$(\mathcal{X}_i, \alpha_i^{2n+1}) \rightarrow (0,0)$$
 weakly in $H^1_{loc}(\mathbb{R}^n \setminus T, \mathbb{R}^n) \times 1$

729 (5.12)
$$\mathcal{X}^{2k*} \sim \frac{1}{\eta^{(d-2)(k+1)}} c^{2k},$$

730 (5.13)
$$\mathcal{X}^{2k+1*} = o\left(\frac{1}{\eta^{(d-2)(k+1)}}\right),$$

where c_{ij}^{2k} denotes the coefficients of the matrix valued tensor $c^{2k} := (c_{ij}^{2k})_{1 \le i,j \le d}$ of order 2k given by

$$c^{2k} := F^{-(k+1)}J^{2k}$$
 with $J^{2k} = \overbrace{I \otimes I \otimes \cdots \otimes I}^{k \text{ times}}$.

732 *Proof.* The result is proved by induction on k.

733 1. Case 2k with k = 0. The tensor $(\widetilde{\boldsymbol{\chi}}_{i}^{0}, \widetilde{\alpha}_{i}^{0})$ satisfies

734 (5.14)
$$\begin{cases} -\Delta \widetilde{\boldsymbol{\chi}}_{i}^{0} + \nabla \widetilde{\alpha}_{i}^{0} = \eta^{d} \boldsymbol{e}_{i} \text{ in } \mathbb{R}^{d} \backslash T \\ \operatorname{div}(\widetilde{\boldsymbol{\chi}}_{i}^{0}) = 0 \text{ in } \mathbb{R}^{d} \backslash T, \end{cases}$$

as well as the other boundary conditions of (5.7). Lemma 5.3 implies then

$$\|\nabla \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{0}\|_{L^{2}(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})} + \|\widetilde{\alpha}_{i}^{0}\|_{L^{2}(\eta^{-1}P \setminus T)} \leq C\eta^{-d}\eta^{d} |\langle \boldsymbol{e}_{i} \rangle| \leq C.$$

From (5.4), we also obtain that $\langle \tilde{\boldsymbol{\chi}}_{i}^{0} \rangle$ is bounded. Hence, up to extracting a subsequence, there exists a constant matrix $c^{0} := (c_{ij}^{0})_{1 \leq i,j \leq d}$, and fields $(\hat{\boldsymbol{\Psi}}_{i}^{0}, \hat{\sigma}_{i}^{0})_{1 \leq i \leq d}$ such that

$$\langle \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{0} \rangle \cdot \boldsymbol{e}_{j} \rightarrow c_{ij}^{0},$$

$$(\widetilde{\boldsymbol{\mathcal{X}}}_{i}^{0},\widetilde{\alpha}_{i}^{0}) \rightharpoonup (\widehat{\boldsymbol{\Psi}}_{i}^{0},\widehat{\sigma}_{i}^{0})$$
 weakly in $H^{1}_{loc}(\eta^{-1}P\backslash T,\mathbb{R}^{d}) \times L^{2}_{loc}(\eta^{-1}P\backslash T).$

Multiplying (5.14) by a compactly supported test function $\Phi \in \mathcal{C}^{\infty}_{c}(\mathbb{R}^{d} \setminus T)$ and integrating by parts yields

$$\int_{\eta^{-1}P\setminus T} (\nabla \widetilde{\boldsymbol{\mathcal{X}}}_i^0 : \nabla \boldsymbol{\Phi} - \widetilde{\alpha}_i^0 \operatorname{div}(\boldsymbol{\Phi})) \mathrm{d}x = \int_{\eta^{-1}P\setminus T} \eta^d \boldsymbol{\Phi} \cdot \boldsymbol{e}_i \mathrm{d}x.$$

Passing to the limit as $\eta \to 0$ implies then

$$\begin{cases} -\Delta \widehat{\Psi}_i^0 + \nabla \widehat{\sigma}_i^0 = 0 \text{ in } \mathbb{R}^d \backslash T \\ \operatorname{div}(\Psi_i^0) = 0 \text{ in } \mathbb{R}^d \backslash T \\ \Psi_i^0 = 0 \text{ on } \partial T. \end{cases}$$

By applying the point (5.6) of Lemma 5.2 and by using the lower semi-continuity of the Lebesgue space norms, we infer $(\widehat{\Phi}_i^0 - c_{ij}^0 e_j, \widehat{\sigma}_i^0) \in \mathcal{D}^{1,2}(\mathbb{R}^d \setminus T) \times L^2(\mathbb{R}^d \setminus T)$ (see the proof of Theorem 3.1 in [6] for a detailed justification). By linearity, it is then necessary that $(\widehat{\Phi}_i^0, \widehat{\sigma}_i^0) = (c_{ij}^0 \Psi_j, c_{ij}^0 \sigma_j)$ where (Ψ_j, σ_j) are the solution to the exterior problem (5.1). In order to identify the coefficient c_{ij}^0 , we integrate (5.1) by parts against the test function $\Phi = e_j$ then yields

$$0 = \eta^d \int_{\eta^{-1}P \setminus T} \delta_{ij} \mathrm{d}x + \int_{\partial T} \boldsymbol{e}_j \cdot (\nabla \widetilde{\boldsymbol{\mathcal{X}}}_i^0 - \widetilde{\alpha}_i^0 I) \cdot \boldsymbol{n} \mathrm{d}x.$$

Passing to the limit as $\eta \to 0$ by using the continuity of the drag force with respect to the weak convergence and (5.2) yields then

$$0 = \delta_{ij} + \int_{\partial T} \boldsymbol{e}_j \cdot (\nabla \widehat{\boldsymbol{\Phi}}_i^0 - \widehat{\sigma}_i^0) \cdot \boldsymbol{n} \mathrm{d}x = \delta_{ij} - c_{ip}^0 F_{pj}.$$

This implies $c^0 = F^{-1}$ as claimed and the convergence of the whole sequence by uniqueness of the limit. The asymptotic for \mathcal{X}^{0*} as $\eta \to 0$ follows by the change of variable $y = \eta x$:

$$\mathcal{X}_{ij}^{0*} = \boldsymbol{e}_i \cdot \int_{P \setminus (\eta T)} \boldsymbol{\mathcal{X}}_j^0 \mathrm{d}y = \eta^{2-d} \eta^d \boldsymbol{e}_i \cdot \int_{\eta^{-1} P \setminus T} \boldsymbol{\widetilde{\mathcal{X}}}_j^0 \mathrm{d}y \sim \eta^{2-d} \langle \boldsymbol{\widetilde{\mathcal{X}}}_j^0 \rangle \cdot \boldsymbol{e}_i \sim \eta^{2-d} c_{ji}^0.$$

735 2. Case 2k + 1 with k = 0. The tensor $(\widetilde{\boldsymbol{\chi}}_{i}^{1}, \widetilde{\alpha}_{i}^{1})$ satisfies

736 (5.15)
$$\begin{cases} -\Delta \widetilde{\boldsymbol{\chi}}_{i}^{1} + \nabla \widetilde{\alpha}_{i}^{1} = \eta (2\partial_{l} \widetilde{\boldsymbol{\chi}}_{i}^{0} - \widetilde{\alpha}_{j}^{0} \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} \text{ in } \eta^{-1} \boldsymbol{P} \backslash \boldsymbol{T} \\ \operatorname{div}(\widetilde{\boldsymbol{\chi}}_{i}^{1}) = -\eta (\widetilde{\boldsymbol{\chi}}_{j}^{0} - \langle \widetilde{\boldsymbol{\chi}}_{j}^{0} \rangle) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l} \text{ in } \eta^{-1} \boldsymbol{P} \backslash \boldsymbol{T}. \end{cases}$$

Applying Lemma 5.3 and remarking that $\langle 2\partial_l \tilde{\boldsymbol{\chi}}_i^0 - \tilde{\alpha}_i^0 \boldsymbol{e}_l \rangle = 0$, we obtain

$$||\nabla \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{1}||_{L^{2}(\eta^{-1}P \setminus T, \mathbb{R}^{d \times d})} + ||\widetilde{\alpha}_{i}^{1}||_{L^{2}(\eta^{-1}P \setminus T)} \leq C.$$

Integrating (5.15) by parts against a compactly supported test function $\Phi \in C^c(\mathbb{R}^d \setminus T)$ and passing to the limit as $\eta \to 0$, we obtain similarly the existence of a matrix valued tensor $c^1 := (c_{ij}^1)_{1 \le i,j \le d}$ (of order 1) such that, up to the extraction of a subsequence:

$$\langle \widetilde{\boldsymbol{\mathcal{X}}}_i^1 \rangle \cdot \boldsymbol{e}_j \to c_{ij}^1,$$

$$(\widetilde{\boldsymbol{\mathcal{X}}}_{i}^{1}, \widetilde{\alpha}_{i}^{1}) \rightharpoonup (c_{ij}^{1} \Psi_{j}, c_{ij}^{1} \sigma_{j})$$
 weakly in $H_{loc}^{1}(\eta^{-1}P \setminus T, \mathbb{R}^{d}) \times L_{loc}^{2}(\eta^{-1}P \setminus T).$

Integrating (5.15) by parts against the test function e_j and passing to the limit as $\eta \to 0$ yields in this situation $0 = c_{ij}^1 F_{pj}$ whence $c^1 = 0$.

3. General case. Assuming that the result holds till rank k, the differential equations satisfied by the rescaled tensors in $\eta^{-1}P \setminus T$ read:

741 (5.16)
$$\begin{cases} -\Delta \widetilde{\boldsymbol{\chi}}_{i}^{2k+2} + \nabla \widetilde{\alpha}_{i}^{2k+2} = \eta^{d-1} (\partial_{l} \widetilde{\boldsymbol{\chi}}_{i}^{2k+1} - \widetilde{\alpha}_{i}^{2k+1} \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \eta^{d} \widetilde{\boldsymbol{\chi}}_{i}^{2k} \otimes \boldsymbol{I} \\ \operatorname{div}(\widetilde{\boldsymbol{\chi}}_{i}^{2k+2}) = -\eta^{d-1} (\widetilde{\boldsymbol{\chi}}_{i}^{2k+1} - \langle \widetilde{\boldsymbol{\chi}}_{i}^{2k+1} \rangle) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}. \end{cases}$$

742

743 (5.17)
$$\begin{cases} -\Delta \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+3} + \nabla \widetilde{\alpha}_{i}^{2k+3} = \eta(\partial_{l} \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+2} - \widetilde{\alpha}_{i}^{2k+2} \boldsymbol{e}_{l}) \otimes \boldsymbol{e}_{l} + \eta^{d} \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+1} \otimes \boldsymbol{I} \\ \operatorname{div}(\widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+3}) = -\eta(\widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+2} - \langle \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+2} \rangle) \cdot \boldsymbol{e}_{l} \otimes \boldsymbol{e}_{l}. \end{cases}$$

Using Lemma 5.2, Lemma 5.3 and the point 1. of the proposition at rank k, we readily obtain

$$\begin{split} ||\nabla \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+2}||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})} + ||\widetilde{\alpha}_{i}^{2k+2}||_{L^{2}(\eta^{-1}P\setminus T)} \leq C, \\ ||\nabla \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+3}||_{L^{2}(\eta^{-1}P\setminus T,\mathbb{R}^{d\times d})} + ||\widetilde{\alpha}_{i}^{2k+3}||_{L^{2}(\eta^{-1}P\setminus T)} \leq C. \end{split}$$

Repeating the above arguments, we obtain, up to the extraction of a subsequence, the existence of matrix valued tensors c^{2k+2} and c^{2k+3} such that

$$\begin{split} &\langle \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+2} \rangle \cdot \boldsymbol{e}_{j} \rightharpoonup c_{ij}^{2k+2} \text{ and } \langle \widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+3} \rangle \cdot \boldsymbol{e}_{j} \rightharpoonup c_{ij}^{2k+3}, \\ &(\widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+2}, \widetilde{\alpha}_{i}^{2k+2}) \rightharpoonup (c_{ij}^{2k+2} \Psi_{j}, c_{ij}^{2k+2} \sigma_{j}) \text{ weakly in } H^{1}_{loc}(\mathbb{R}^{d} \backslash T, \mathbb{R}^{d}) \times L^{2}_{loc}(\mathbb{R}^{d} \backslash T), \\ &(\widetilde{\boldsymbol{\mathcal{X}}}_{i}^{2k+3}, \widetilde{\alpha}_{i}^{2k+3}) \rightharpoonup (c_{ij}^{2k+3} \Psi_{j}, c_{ij}^{2k+3} \sigma_{j}) \text{ weakly in } H^{1}_{loc}(\mathbb{R}^{d} \backslash T, \mathbb{R}^{d}) \times L^{2}_{loc}(\mathbb{R}^{d} \backslash T). \end{split}$$

The last step consists in integrating (5.16) and (5.17) by part against the test function e_j and to pass to the limit as $\eta \to 0$ in order to identify c_{ij}^{2k+2} and c_{ij}^{2k+3} . Performing this computation as above yields

$$0 = c_{ij}^{2k} \otimes I - c_{ip}^{2k+2} F_{pj},$$

$$0 = c_{ij}^{2k+1} \otimes I - c_{ip}^{2k+3} F_{pj}$$

from where we infer $c^{2k+2} = c^{2k}F^{-1} \otimes I$, $c^{2k+3} = c^{2k+1}F^{-1} \otimes I$, whence the result (recall $c^1 = 0$ from the point 2. of the proof).

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Using the identity of (3.20), we obtain the asymptotics for the coefficients M^k of the infinite order homogenized equation (1.9).

748 COROLLARY 5.5. Assume $d \ge 3$. The following convergences hold for the matrix 749 valued tensors M^k as $\eta \to 0$:

750 (5.18)
$$M^0 \sim \eta^{d-2} F,$$

751 (5.19)
$$M^1 = o\left(\eta^{d-2}\right),$$

752 (5.20)
$$M^2 \to -I,$$

753 (5.21)
$$\forall k \ge 1, \ M^{2k} = o\left(\frac{1}{\eta^{(d-2)(k-1)}}\right),$$
754 (5.22)
$$\forall k \ge 1, \ M^{2k+1} = o\left(\frac{1}{\eta^{(d-2)(k-1)}}\right)$$

754 (5.22)
$$\forall k \ge 1, M^{2k+1} = o\left(\frac{1}{\eta^{(d-2)(k-1)}}\right)$$

Proof. We replace the asymptotics of (5.4) in the explicit formula for the tensors M^k given in (3.20). (5.18) is an immediate consequence of $M^0 = (\mathcal{X}^{0*})^{-1}$. The convergence (5.20) is obtained by writing, according to (3.20):

$$\begin{split} M^2 &= -(\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{2*} \otimes (\mathcal{X}^{0*})^{-1} + (\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{1*} \otimes (\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{1*} \otimes (\mathcal{X}^{0*})^{-1} \\ &= -\frac{\eta^{2(d-2)}}{\eta^{2(d-2)}} F \otimes c^2 \otimes F + o\left(\frac{\eta^{3(d-2)}}{\eta^{2(d-2)}}\right) = -(FF^{-2}F) \otimes I + o(\eta^{d-2}) \\ &= -I + o(\eta^{d-2}). \end{split}$$

For M^{2k+1} with $k \ge 0$, we use (3.20) and we observe that, for any $0 \le p \le 2k+1$ and indices $1 \le i_1 \ldots i_p \le 2k+1$ such that $i_1 + \cdots + i_p = 2k+1$, there exists at least one odd index i_q with $1 \le q \le p$. Using (5.12) and (5.13), we arrive at

$$\begin{array}{l} & (\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{i_1*} \otimes \dots \otimes (\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{i_p*} \otimes (\mathcal{X}^{0*})^{-1} \\ & = o\left(\frac{\eta^{(p+1)(d-2)}}{\eta^{(p+\lfloor i_1/2 \rfloor + \dots + \lfloor i_p/2 \rfloor)(d-2)}}\right) = o\left(\frac{1}{\eta^{(d-2)(k-1)}}\right) \end{array}$$

which implies (5.19) and (5.22). For
$$M^{2k}$$
 with $k > 1$, we separate the summands
of (3.20) into two categories. For a given p indices such that $1 \le p \le 2k$ and
 $i_1 + \cdots + i_p = 2k$, there are only two possibilities:

1. either there exists at least one odd index i_q , in that case the above reasoning implies as well

$$(\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{i_1*} \otimes \cdots \otimes (\mathcal{X}^{0*})^{-1} \mathcal{X}^{i_p*} \otimes (\mathcal{X}^{0*})^{-1} = o\left(\frac{1}{\eta^{(d-2)(k-1)}}\right).$$

766 2. or all indices $i_1 + \cdots + i_p$ are even, in that case we may write, as $\eta \to 0$:

767
$$(\mathcal{X}^{0*})^{-1} \otimes \mathcal{X}^{i_1*}(\mathcal{X}^{0*})^{-1} \otimes \cdots \otimes (\mathcal{X}^{0*})^{-1} \mathcal{X}^{i_p*} \otimes (\mathcal{X}^{0*})^{-1}$$

768
$$\sim \frac{\eta^{(d-2)(p+1)}}{\eta^{(d-2)(p+k)}} F \otimes c^{i_1} \otimes \cdots \otimes F \otimes c^{i_p} \otimes F$$

769
$$\sim \frac{1}{\eta^{(d-2)(k-1)}} (FF^{-(i_1/2+1)}) \times \dots \times (FF^{-(i_p/2+1)}) FJ^{2k}$$

770
$$\sim \frac{1}{\eta^{(d-2)(k-1)}} F^{-(k-1)} \otimes J^{2k}.$$

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Note that in the latter case, the asymptotic does not depend on the choice of indices $i_1 + \cdots + i_p = 2k$. Therefore, by isolating the terms featuring only even indices $i_1 := 2j_1, \ldots, i_p := 2j_p$ in (3.20), we obtain finally

$$M^{2k} = \frac{1}{\eta^{(d-2)(k-1)}} F^{-(k-1)} \otimes J^{2k} \left(\sum_{\substack{p=1\\p=1}}^{2k} (-1)^p \sum_{\substack{2j_1 + \dots + 2j_p = 2k\\1 \le j_1, \dots, j_p \le k}} 1 \right) + o\left(\frac{1}{\eta^{(d-2)(k-1)}}\right).$$

The asymptotic (5.21) follows from the fact that the summation over p in the above 772 773 expression is zero (see e.g. the end of the proof of Corollary 6 in [34]). Π

Remark 5.6. We have therefore obtained the following asymptotic estimates for 774 the coefficients $\varepsilon^{k-2}M^k$ of the infinite order homogenized equation (1.9) as $\eta \to 0$: 775

$$\varepsilon^{-2}M^0 \sim (\eta^{d-2}/\varepsilon^2)F,$$

777
$$\varepsilon^{-1}M^1 = o\left(\varepsilon(\eta^{d-2}/\varepsilon^2)\right)$$
778
$$\varepsilon^0 M^2 \to -I,$$

778

779
$$\varepsilon^{2k-2}M^{2k} = o\left(\left(\varepsilon^2/\eta^{d-2}\right)^{k-1}\right) \text{ for } k \ge 1,$$

$$\varepsilon^{2k-1}M^{2k+1} = o\left(\varepsilon\left(\varepsilon^2/\eta^{d-2}\right)^{k-1}\right) \text{ for } k \ge 1$$

These asymptotics bring into play the ratio ε^2/η^{d-2} and so the critical scaling $\eta_{\rm crit} \sim$ 782 $\varepsilon^{2/(d-2)}$. They imply thus the "coefficient-wise" convergence of (1.9) to the Brinkman 783 regime (1.3) at the critical rate $\eta \sim \varepsilon^{2/(d-2)}$, in which case $\varepsilon^{-2}M^0 \to F$ and 784 $\varepsilon^{k-2}M^k \to 0$ for any k > 2. Note that $\varepsilon^0 M^2 \to -I$ whatever the rate of conver-785 gence at which $\eta \to 0$. The Darcy regimes (1.4) and (1.5) correspond to the situation 786 where $\eta^{d-2}/\varepsilon^2 \to +\infty$; in that case the zeroth order term $\varepsilon^{-2} \hat{M}^0$ is dominant. 787

Finally, the leading coefficients of the Stokes regime (1.2) are retrieved for $\eta =$ 788 $o(\varepsilon^{2/(d-2)})$, since in this case, $\varepsilon^{-2}M^0 \to 0$, $\varepsilon^{-1}M^1 \to 0$ and $\varepsilon^0 M^0 \to -I$. However the present analysis is not sufficient to conclude that the coefficients $\varepsilon^{k-2}M^k$ of order 789 790 k > 2 converge to zero in the subcritical regime $\eta = o(\varepsilon^{2/(d-2)})$. Indeed $\varepsilon^{k-2}M^k$ is just 791 bounded by $(\varepsilon^2/\eta^{d-2})^k$, a quantity which can potentially blow up for too small values 792 of η . This matter is to be addressed in a future work through a more accurate analysis 793 of the rate of convergence of the coefficients \mathcal{X}^{2k*} and \mathcal{X}^{2k+1*} in the asymptotic (5.12) 794 and (5.13). To date, let us note that Jing obtained recently $\mathcal{X}^{0*} = F/\eta^{d-2} + O(1)$ in 795 the scalar case (proof of the Lemma 5.1 in [41]) by using layer potential techniques. 796

Remark 5.7. From the estimates of (5.4), we obtain that the coefficients (\mathbb{D}_K^k) of (1.6) satisfy the same asymptotic convergences of Corollary 5.5. Indeed, by using the definition (3.21) we find that there exists a constant C > 0 independent of η such that for any $K \ge 0$, $K + 1 \le k \le 2K + 1$ and $1 \le j \le d$:

$$\begin{aligned} \|\nabla \boldsymbol{N}_{j}^{K}\|_{L^{2}(Y,\mathbb{R}^{d\times d})} &\leq C\eta^{-(d-2)\lfloor K/2 \rfloor} \eta^{d/2-1} \\ \|\beta_{j}^{k-K-1}\|_{L^{2}(Y,\mathbb{R}^{d})} &\leq C\eta^{-(d-2)\lfloor \frac{k-K-1}{2} \rfloor+1} \eta^{d/2}. \end{aligned}$$

Applying the inequality $\lfloor x \rfloor + \lfloor y \rfloor \leq \lfloor x + y \rfloor$, we obtain that the coefficients \mathbb{A}_{ij}^k of Proposition 4.10 satisfy for $K + 1 \le k \le 2K + 1$,

$$|\mathbb{A}_{ij}^k| \le C\eta^d \eta^{-\lfloor \frac{k-1}{2} \rfloor (d-2)} \le C\eta^2 M_{ij}^k.$$

⁷⁹⁷ Hence the discrepancy induced by the coefficients \mathbb{A}_{ij}^k is small and $\mathbb{D}_{K,ij}^k$ also satisfies ⁷⁹⁸ (5.21) and (5.22) for $K+1 \leq k \leq 2K+1$. Similarly (see Corollary 7 of [34]), we may ⁷⁹⁹ show that $\mathbb{D}_0^2 \to -I$. The remaining coefficients \mathbb{D}_K^k are equal to M^k .

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