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Spectroscopic imaging of a dilute cell suspension

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Spectroscopic imaging of a dilute cell suspension*

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Abstract

The paper aims at analytically exhibiting for the first time the fundamental mechanisms underlying the fact that effective biological tissue electrical properties and their frequency dependence reflect the tissue composition and physiology. For doing so, a homogenization theory is derived to describe the effective admittivity of cell suspensions. A new formula is reported for dilute cases that gives the frequency-dependent effective admittivity with respect to the membrane polarization. Different microstructures are shown to be distinguishable via spectroscopic measurements of the overall admittivity using the spectral properties of the membrane polarization. The Debye relaxation times associated with the membrane polarization tensor are shown to be able to give the microscopic structure of the medium. A natural measure of the admittivity anisotropy is introduced and its dependence on the frequency of applied current is derived. A Maxwell-Wagner-Fricke formula is given for concentric circular cells, and the results can be extended to the random cases. A randomly deformed periodic medium is also considered and a new formula is derived for the overall admittivity of a dilute suspension of randomly deformed cells.

Mathematics Subject Classification (MSC2000): 35R30, 35B30.

Keywords: electrical impedance spectroscopy, stochastic homogenization, Maxwell-Wagner-Fricke formula, Debye relaxation time.

Contents

1 Introduction

3

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Problem settings and main results	5
2.1 Periodic domain	5
2.2 Electrical model of the cell	5
2.3 Governing equation	9
2.4 Main results in the periodic case	9
2.5 Description of the random cells and interfaces	12
2.6 Stationary ergodic setting	12
2.7 Main results in the random case	14
Analysis of the problem	16
3.1 Existence and uniqueness of a solution	16
3.2 Energy estimate	17
Homogenization	18
4.1 Two-scale asymptotic expansions	19
4.2 Convergence	22
4.2.1 A priori estimates	24
4.2.2 Two-scale convergence	26
Effective admittivity for a dilute suspension	30
5.1 Computation of the effective admittivity	30
5.2 Maxwell-Wagner-Fricke formula	33
5.3 Debye relaxation times	33
5.4 Properties of the membrane polarization tensor and the Debye relaxation times .	34
5.5 Anisotropy measure	35
Spectroscopic imaging of a dilute suspension	35
6.1 Spectroscopic conductivity imaging	35
6.2 Selective spectroscopic imaging	36
6.3 Spectroscopic measurement of anisotropy	38
	00
Stochastic homogenization of randomly deformed conductivity resistant membranes	38
7.1 Auxiliary problem: proof of Theorem 2.3	38
7.2 Proof of the homogenization theorem	42
7.2.1 Oscillating test functions	43
7.2.2 Proof of the homogenization theorem	44
7.3 Effective admittivity of a dilute suspension	48
Numerical simulations	50
	- •
Concluding remarks	50
A Extension lemmas 53	
Poincaré-Wirtinger inequality	56
Equivalence of the two norms on $W_{arepsilon}$	57
	Problem settings and main results 2.1 Periodic domain 2.2 Electrical model of the cell 2.3 Governing equation 2.4 Main results in the periodic case 2.5 Description of the random cells and interfaces 2.6 Stationary ergodic setting 2.7 Main results in the random case 2.8 Analysis of the problem 3.1 Existence and uniqueness of a solution 3.2 Energy estimate 4.1 Two-scale asymptotic expansions 4.2 Convergence 4.2.1 A priori estimates 4.2.2 Two-scale convergence 4.2.1 A priori estimates 4.2.2 Two-scale convergence 5.1 Computation of the effective admittivity 5.2 Maxwell-Wagner-Fricke formula 5.3 Debye relaxation times 5.4 Properties of the membrane polarization tensor and the Debye relaxation times 5.5 Anisotropy measure 5.6 Spectroscopic conductivity imaging 6.1 Spectroscopic maging of a dilute suspension 6.1 Spectroscopic measurem

D Technical lemma

1 Introduction

The electric behavior of biological tissue under the influence of an electric field at frequency ω can be characterized by its frequency-dependent effective admittivity $k_{ef} := \sigma_{ef}(\omega) + i\omega\epsilon_{ef}(\omega)$, where σ_{ef} and ε_{ef} are respectively its effective conductivity and permittivity. Electrical impedance spectroscopy assesses the frequency dependence of the effective admittivity by measuring it across a range of frequencies from a few Hz to hundreds of MHz. Effective admittivity of biological tissues and its frequency dependence vary with tissue composition, membrane characteristics, intra-and extra-cellular fluids and other factors. Hence, the admittance spectroscopy provides information about the microscopic structure of the medium and physiological and pathological conditions of the tissue.

The determination of the effective, or macroscopic, property of a suspension is an enduring problem in physics [44]. It has been studied by many distinguished scientists, including Maxwell, Poisson [52], Faraday, Rayleigh [54], Fricke [31], Lorentz, Debye, and Einstein [26]. Many studies have been conducted on approximate analytic expressions for overall admittivity of a cell suspension from the knowledge of pointwise conductivity distribution, and these studies were mostly restricted to the simplified model of a strongly dilute suspension of spherical or ellipsoidal cells.

In this paper, we consider a periodic suspension of identical cells of arbitrary shape. We apply at the boundary of the medium an electric field of frequency ω . The medium outside the cells has an admittivity of $k_0 := \sigma_0 + i\omega\epsilon_0$. Each cell is composed of an isotropic homogeneous core of admittivity k_0 and a thin membrane of constant thickness δ and admittivity $k_m := \sigma_m + i\omega\epsilon_m$. The thickness δ is considered to be very small relative to the typical cell size and the membrane is considered very resistive, *i.e.*, $\sigma_m \ll \sigma_0$. In this context, the potential in the medium passes an effective discontinuity over the cell boundary; the jump is proportional to its normal derivative with a coefficient of the effective thickness, given by $\delta k_0 / k_m$. The normal derivative of the potential is continuous across the cell boundaries.

We use homogenization techniques with asymptotic expansions to derive a homogenized problem and to define an effective admittivity of the medium. We prove a rigorous convergence of the original problem to the homogenized problem via two-scale convergence. For dilute cell suspensions, we use layer potential techniques to expand the effective admittivity in terms of cell volume fraction. Through the effective thickness, $\delta k_0/k_m$, the first-order term in this expansion can be expressed in terms of a membrane polarization tensor, *M*, that depends on the operating frequency ω . We retrieve the Maxwell-Wagner-Fricke formula for concentric circular-shaped cells. This explicit formula has been generalized in many directions: in three dimension for concentric spherical cells; to include higher power terms of the volume fraction for concentric, confocal ellipses and ellipsoids; see [14, 15, 28, 29, 30, 43, 55, 56, 57].

The imaginary part of *M* is positive for δ small enough. Its two eigenvalues are maximal for frequencies $1/\tau_i$, i = 1, 2, of order of a few MHz with physically plausible parameters values. This dispersion phenomenon well known by the biologists is referred to as the β -dispersion. The associated characteristic times τ_i correspond to Debye relaxation times. Given this, we show that different microscopic organizations of the medium can be distinguished via

 τ_i , i = 1, 2, alone. The relaxation times τ_i are computed numerically for different configurations: one circular or elliptic cell, two or three cells in close proximity. The obtained results illustrate the viability of imaging cell suspensions using the spectral properties of the membrane polarization. The Debye relaxation times are shown to be able to give the microscopic structure of the medium.

In the second part of this paper, we show that our results can be extended to the random case by considering a randomly deformed periodic medium. We also derive a rigorous homogenization theory for cells (and hence interfaces) that are randomly deformed from a periodic structure by random, ergodic, and stationary deformations. We prove a new formula for the overall conductivity of a dilute suspension of randomly deformed cells. Again, the spectral properties of the membrane polarization can be used to classify different microscopic structures of the medium through their Debye relaxation times. For recent works on effective properties of dilute random media, we refer to [7, 17].

Our results in this paper have potential applicability in cancer imaging, food sciences and biotechnology [41, 42], and applied and environmental geophysics. They can be used to model and improve the MarginProbe system for breast cancer [61], which emits an electric field and senses the returning signal from tissue under evaluation. The greater vascularization, differently polarized cell membranes, and other anatomical differences of tumors compared with healthy tissue cause them to show different electromagnetic signatures. The ability of the probe to detect signals characteristic of cancer helps surgeons ensure the removal of all unwanted tissue around tumor margins.

Another commercial medical system to which our results can be applied is ZedScan [62]. ZedScan is based on electrical impedance spectroscopy for detecting neoplasias in cervical disease [1, 20]. Malignant white blood cells can be also detected using induced membrane polarization [53]. In food quality inspection, spectroscopic conductivity imaging can be used to detect bacterial cells [12, 59]. In applied and environmental geophysics, induced membrane polarization can be used to probe up to subsurface depths of thousands of meters [58, 60].

The structure of the rest of this paper is as follows. Section 2 introduces the problem settings and state the main results of this work. Section 3 is devoted to the analysis of the problem. We prove existence and uniqueness results and establish useful *a priori* estimates. In section 4 we consider a periodic cell suspension and derive spectral properties of the overall conductivity. In section 5 we consider the problem of determining the effective property of a suspension of cells when the volume fraction goes to zero. Section 6 is devoted to spectroscopic imaging of a dilute suspension. We make use of the asymptotic expansion of the effective admittivity in terms of the volume fraction to image a permittivity inclusion. We also discuss selective spectroscopic imaging using a pulsed approach. Finally, we introduce a natural measure of the conductivity anisotropy and derive its dependence on the frequency of applied current. In section 7 we extend our results to the case of randomly deformed periodic media. In section 8 we provide numerical examples that support our findings. A few concluding remarks are given in the last section. For simplicity, we only treat the two-dimensional case. Our results can be extended into the three dimensional setting [35].

2 Problem settings and main results

The aim of this section is to introduce the problem settings and state the main results of this paper.

2.1 Periodic domain

We consider the probe domain Ω to be a bounded open set of \mathbb{R}^2 of class C^2 . The domain contains a periodic array of cells whose size is controlled by ε . Let *C* be a $C^{2,\eta}$ domain being contained in the unit square $Y = [0,1]^2$, see Figure 2.1. Here, $0 < \eta < 1$ and *C* represents a reference cell. We divide the domain Ω periodically in each direction in identical squares $(Y_{\varepsilon,n})_n$ of size ε , where

$$Y_{\varepsilon,n} = \varepsilon n + \varepsilon Y.$$

Here, $n \in N_{\varepsilon} := \{n \in \mathbb{Z}^2 | Y_{\varepsilon,n} \cap \Omega \neq \emptyset\}.$

We consider that a cell $C_{\varepsilon,n}$ lives in each small square $Y_{\varepsilon,n}$. As shown in Figure 2.4, all cells are identical, up to a translation and scaling of size ε , to the reference cell *C*:

$$\forall n \in N_{\varepsilon}, \quad C_{\varepsilon,n} = \varepsilon n + \varepsilon C.$$

So are their boundaries $(\Gamma_{\varepsilon,n})_{n \in N_{\varepsilon}}$ to the boundary Γ of *C*:

$$\forall n \in N_{\varepsilon}, \quad \Gamma_{\varepsilon,n} = \varepsilon n + \varepsilon \Gamma.$$

Let us also assume that all the cells are strictly contained in Ω , that is for every $n \in N_{\varepsilon}$, the boundary $\Gamma_{\varepsilon,n}$ of the cell $C_{\varepsilon,n}$ does not intersect the boundary $\partial\Omega$:

$$\partial \Omega \cap (\bigcup_{n \in N_{\varepsilon}} \Gamma_{\varepsilon,n}) = \emptyset.$$

2.2 Electrical model of the cell

Set for any open set *D* of \mathbb{R}^2 :

$$L_0^2(D) := \left\{ f \in L^2(D) \middle| \int_{\partial D} f(x) ds(x) = 0 \right\}$$

and

$$H^1(D) := \left\{ f \in L^2(D) \left| |\nabla f| \in L^2(D) \right\}.$$

We consider in this section the reference cell *C* immersed in a domain *D*. We apply a sinusoidal electrical current $g \in L^2_0(\partial D)$ with angular frequency ω at the boundary of *D*.

The medium outside the cell, $D \setminus C$, is a homogeneous isotropic medium with admittivity $k_0 := \sigma_0 + i\omega\epsilon_0$. The cell *C* is composed of an isotropic homogeneous core of admittivity k_0 and a thin membrane of constant thickness δ with admittivity $k_m := \sigma_m + i\omega\epsilon_m$. We make the following assumptions :

$$\sigma_0 > 0, \sigma_m > 0, \epsilon_0 > 0, \epsilon_m \ge 0.$$

If we apply a sinusoidal current $g(x)\sin(\omega t)$ on the boundary ∂D in the low frequency range below 10 MHz, the resulting complex-valued time harmonic potential \check{u} is governed by

$$\begin{cases} \nabla \cdot (k_0 + (k_m - k_0)\chi_{\Gamma^{\delta}}) \nabla \check{u}) = 0 \quad \text{in } D\\ k_0 \frac{\partial \check{u}}{\partial n}\Big|_{\partial D} = g, \end{cases}$$

where $\Gamma^{\delta} := \{x \in C : \operatorname{dist}(x, \Gamma) < \delta\}$ and $\chi_{\Gamma^{\delta}}$ is the characteristic function of the set Γ^{δ} .

The membrane thickness δ is considered to be very small compared to the typical size ρ of the cell, *i.e.*, $\delta/\rho \ll 1$. According to the transmission condition, the normal component of the current density $k_0 \frac{\partial \check{u}}{\partial n}$ can be approximately regarded as continuous across the thin membrane Γ .

We set $\beta := \frac{\delta}{k_m}$. Since the membrane is very resistive, *i.e.* $\sigma_m / \sigma_0 \ll 1$, the potential \check{u} in D undergoes a jump across the cell membrane Γ , which can be approximated at first order by $\beta k_0 \frac{\partial \check{u}}{\partial n}$. A rigorous proof of this result, based on asymptotic expansions of layer potentials, can be found in [37].

More precisely, we approximate \check{u} by u defined as the solution of the following equations [37, 50, 51]:

$$\begin{cases} \nabla \cdot k_0 \nabla u = 0 & \text{in } D \setminus \overline{C}, \\ \nabla \cdot k_0 \nabla u = 0 & \text{in } C, \\ k_0 \frac{\partial u}{\partial n}\Big|_+ = k_0 \frac{\partial u}{\partial n}\Big|_- & \text{on } \Gamma, \\ u\Big|_+ - u\Big|_- - \beta k_0 \frac{\partial u}{\partial n} = 0 & \text{on } \Gamma, \\ k_0 \frac{\partial u}{\partial n}\Big|_{\partial D} = g, \quad \int_{\partial D} g(x) ds(x) = 0, \quad \int_{D \setminus \overline{C}} u(x) dx = 0. \end{cases}$$
(2.1)

Here *n* is the outward unit normal vector and $u|_{\pm}(x)$ denotes $\lim_{t\to 0^+} u(x \pm tn(x))$ for *x* on the concerned boundary. Likewise, $\frac{\partial u}{\partial n}\Big|_{\pm} := \lim_{t\to 0^+} \nabla u(x \pm tn(x)) \cdot n(x)$.

Equation (2.1) is the starting point of our analysis.

For any open set *B* in \mathbb{R}^2 , we denote $H^1_{\mathbb{C}}(B)$ the Sobolev space $H^1(B)/\mathbb{C}$ which can be represented as :

$$H^1_{\mathbb{C}}(B) = \left\{ u \in H^1(B) \mid \int_B u(x) dx = 0 \right\}.$$

The following result holds.

Lemma 2.1. There exists a unique solution $u := (u^+, u^-)$ in $H^1_{\mathbb{C}}(D^+) \times H^1(D^-)$ to (2.1).



Figure 2.1: Schematic illustration of a unit period Y.

Proof. To prove the well-posedness of (2.1) we introduce the following Hilbert space: V := $H^1_{\mathbb{C}}(D) \times H^1(D)$ equipped with the following natural norm for our problem:

$$\forall u \in V \|u\|_{V} = \|\nabla u^{+}\|_{L^{2}(D^{+})} + \|\nabla u^{-}\|_{L^{2}(D^{-})} + \|u^{+} - u^{-}\|_{L^{2}(\Gamma)}.$$

We write the variational formulation of (2.1) as follows:

Find $u \in V$ such that for all $v := (v^+, v^-) \in V$:

$$\begin{cases} \int_{D^+} k_0 \nabla u^+(x) \cdot \nabla \overline{v}^-(x) \, dx + \int_{D^-} k_0 \nabla u^+(x) \cdot \nabla \overline{v}^-(x) \, dx \\ + \frac{1}{\beta k_0} \int_{\Gamma} (u^+ - u^-) \overline{(v^+ - v^-)} \, d\sigma(x) = \frac{1}{ko} \int_{\partial \Omega} g \overline{v} \, d\sigma(x) \end{cases}$$

Since $\Re(k_0) = \sigma_0 > 0$ and $\Re(\frac{1}{\beta k_0}) = \frac{\sigma_m \sigma_0 + \varepsilon_m \varepsilon_0}{\delta |k_0|} > 0$, we can apply Lax-Milgram theory to

obtain existence and uniqueness of a solution to problem (2.1).

We conclude this subsection with a few numerical simulations to illustrate the typical profile of the potential *u*. We consider an elliptic domain *D* in which lives an elliptic cell. We choose to virtually apply at the boundary of *D* an electrical current $g = e^{i*30r}$.

We use for the different parameters the following realistic values:

- the typical size of eukaryotes cells: $\rho \simeq 10 100 \,\mu\text{m}$;
- the ratio between the membrane thickness and the size of the cell: $\delta/\rho = 0.7 \cdot 10^{-3}$;



Figure 2.2: Real and imaginary parts of the potential u outside and inside the cell.



Figure 2.3: Gradient vector fields of the real and imaginary parts of u.

- the conductivity of the medium and the cell: $\sigma_0 = 0.5 \text{ S.m}^{-1}$;
- the membrane conductivity: $\sigma_m = 10^{-8} \text{ S.m}^{-1}$;
- the permittivity of the medium and the cell: $\epsilon_0 = 90 \times 8.85 \cdot 10^{-12} \text{ F.m}^{-1}$;
- the membrane permittivity: $\varepsilon_m = 3.5 \times 8.85 \cdot 10^{-12} \text{ F.m}^{-1}$;
- the frequency: $\omega = 10^6$ Hz.

Note that the assumptions of our model $\delta \ll \rho$ and $\sigma_m \ll \sigma_0$ are verified.

The real and imaginary parts of *u* outside and inside the cell are represented in Figure 2.2. We can observe that the potential jumps across the cell membrane. We plot the outside and inside gradient vector fields; see Figure 2.3.

2.3 Governing equation

We denote by Ω_{ε}^+ the medium outside the cells and Ω_{ε}^- the medium inside the cells:

$$\Omega_{\varepsilon}^{+} = \Omega \cap (\bigcup_{n \in N_{\varepsilon}} Y_{\varepsilon,n} \setminus \overline{C_{\varepsilon,n}}), \quad \Omega_{\varepsilon}^{-} = \bigcup_{n \in N_{\varepsilon}} C_{\varepsilon,n}.$$

Set $\Gamma_{\varepsilon} := \bigcup_{n \in N_{\varepsilon}} \Gamma_{\varepsilon,n}$. By definition, the boundaries $\partial \Omega_{\varepsilon}^+$ and $\partial \Omega_{\varepsilon}^-$ of respectively Ω_{ε}^+ and Ω_{ε}^- satisfy:

$$\partial \Omega^+_{\varepsilon} = \partial \Omega \cup \Gamma_{\varepsilon}, \quad \partial \Omega^-_{\varepsilon} = \Gamma_{\varepsilon},$$

We apply a sinusoidal current $g(x)\sin(\omega t)$ at $x \in \partial\Omega$, where $g \in L^2_0(\partial\Omega)$. The induced time-harmonic potential u_{ε} in Ω satisfies:

$$\begin{cases} \nabla \cdot k_0 \nabla u_{\varepsilon}^+ = 0 & \text{in } \Omega_{\varepsilon}^+, \\ \nabla \cdot k_0 \nabla u_{\varepsilon}^- = 0 & \text{in } \Omega_{\varepsilon}^-, \\ k_0 \frac{\partial u_{\varepsilon}^+}{\partial n} = k_0 \frac{\partial u_{\varepsilon}^-}{\partial n} & \text{on } \Gamma_{\varepsilon}, \\ u_{\varepsilon}^+ - u_{\varepsilon}^- - \varepsilon \beta k_0 \frac{\partial u_{\varepsilon}^+}{\partial n} = 0 & \text{on } \Gamma_{\varepsilon}, \\ k_0 \frac{\partial u_{\varepsilon}^+}{\partial n} \Big|_{\partial \Omega} = g, \quad \int_{\partial \Omega} g(x) ds(x) = 0, \quad \int_{\Omega_{\varepsilon}^+} u_{\varepsilon}^+(x) dx = 0, \end{cases}$$
where $u_{\varepsilon} = \begin{cases} u_{\varepsilon}^+ & \text{in } \Omega_{\varepsilon}^+, \\ u_{\varepsilon}^- & \text{in } \Omega_{\varepsilon}^-. \end{cases}$

Note that the previously introduced constant β , *i.e.*, the ratio between the thickness of the membrane of *C* and its admittivity, becomes $\varepsilon\beta$. Because the cells $(C_{\varepsilon,n})_{n\in N_{\varepsilon}}$ are in squares of size ε , the thickness of their membranes is given by $\varepsilon\delta$ and consequently, a factor ε appears.

2.4 Main results in the periodic case

We set $Y^+ := Y \setminus \overline{C}$ and $Y^- := C$. For any open set D in \mathbb{R}^2 , we denote $H^1_{\mathbb{C}}(D)$ the Sobolev space $H^1(D)/\mathbb{C}$ which can be represented as :

$$H^1_{\mathbb{C}}(D) = \left\{ u \in H^1(D) \mid \int_D u(x) dx = 0 \right\}.$$

Throughout this paper, we assume that $dist(Y^-, \partial Y) = O(1)$. We write the solution u_{ε} as

$$\forall x \in \Omega \ u_{\varepsilon}(x) = u_0(x) + \varepsilon u_1(x, \frac{x}{\varepsilon}) + o(\varepsilon),$$
(2.3)

with

$$y \mapsto u_1(x,y)$$
 Y-periodic and $u_1(x,y) = \begin{cases} u_1^+(x,y) \text{ in } \Omega \times Y^+, \\ u_1^-(x,y) \text{ in } \Omega \times Y^-. \end{cases}$

The following theorem holds.



Figure 2.4: Schematic illustration of the periodic medium Ω .

- **Theorem 2.1.** (i) The solution u_{ε} to (2.2) two-scale converges to u_0 and $\nabla u_{\varepsilon}(x)$ two-scale converges to $\nabla u_0(x) + \chi_{Y^+}(y)\nabla_y u_1^+(x,y) + \chi_{Y^-}(y)\nabla_y u_1^-(x,y)$, where $\chi_{Y^{\pm}}$ are the characteristic functions of Y^{\pm} .
 - (ii) The function u_0 in (2.3) is the solution in $H^1_{\mathbb{C}}(\Omega)$ to the following homogenized problem:

$$\begin{cases} \nabla \cdot K^* \, \nabla u_0(x) = 0 & \text{in } \Omega, \\ n \cdot K^* \nabla u_0 = g & \text{on } \partial \Omega, \end{cases}$$
(2.4)

where K^* , the effective admittivity of the medium, is given by

$$\forall (i,j) \in \{1,2\}^2, \qquad K_{i,j}^* = k_0 \left(\delta_{ij} + \int_Y (\chi_{Y^+} \nabla w_i^+ + \chi_{Y^-} \nabla w_i^-) \cdot e_j \right), \qquad (2.5)$$

and the function $(w_i)_{i=1,2}$ are the solutions of the following cell problems:

$$\begin{cases} \nabla \cdot k_0 \nabla (w_i^+(y) + y_i) = 0 & \text{in } Y^+, \\ \nabla \cdot k_0 \nabla (w_i^-(y) + y_i) = 0 & \text{in } Y^-, \\ k_0 \frac{\partial}{\partial n} (w_i^+(y) + y_i) = k_0 \frac{\partial}{\partial n} (w_i^-(y) + y_i) & \text{on } \Gamma, \\ w_i^+ - w_i^- - \beta k_0 \frac{\partial}{\partial n} (w_i^+(y) + y_i) = 0 & \text{on } \Gamma, \\ y \longmapsto w_i(y) Y\text{-periodic.} \end{cases}$$
(2.6)

(iii) Moreover, u_1 can be written as

$$\forall (x,y) \in \Omega \times Y, \ u_1(x,y) = \sum_{i=1}^2 \frac{\partial u_0}{\partial x_i}(x) w_i(y).$$
(2.7)

We define the integral operator $\mathcal{L}_{\Gamma} : \mathcal{C}^{2,\eta}(\Gamma) \to \mathcal{C}^{1,\eta}(\Gamma)$, with $0 < \eta < 1$ by

$$\mathcal{L}_{\Gamma}[\varphi](x) = \frac{1}{2\pi} \int_{\Gamma} \frac{\partial^2 \ln |x - y|}{\partial n(x) \partial n(y)} \varphi(y) ds(y), \quad x \in \Gamma.$$
(2.8)

 \mathcal{L}_{Γ} is the normal derivative of the double layer potential \mathcal{D}_{Γ} .

Since \mathcal{L}_{Γ} is positive, one can prove that the operator $I + \alpha \mathcal{L}_{\Gamma} : \mathcal{C}^{2,\eta}(\Gamma) \to \mathcal{C}^{1,\eta}(\Gamma)$ is a bounded operator and has a bounded inverse provided that $\Re \alpha > 0$ [23, 47].

As the fraction f of the volume occupied by the cells goes to zero, we derive an expansion of the effective admittivity for arbitrary shaped cells in terms of the volume fraction. We refer to the suspension, as periodic dilute. The following theorem holds.

Theorem 2.2. *The effective admittivity of a periodic dilute suspension admits the following asymptotic expansion:*

$$K^{*} = k_0 \left(I + fM \left(I - \frac{f}{2}M \right)^{-1} \right) + o(f^2),$$
(2.9)

where $ho = \sqrt{|Y^-|}$, $f =
ho^2$,

$$M = \left(M_{ij} = \beta k_0 \int_{\rho^{-1}\Gamma} n_j \psi_i^*(y) ds(y) \right)_{(i,j) \in \{1,2\}^2},$$
(2.10)

and ψ_i^* is defined by

$$\psi_i^* = -\left(I + \beta k_0 \mathcal{L}_{\rho^{-1}\Gamma}\right)^{-1} [n_i].$$
(2.11)

Note that $\rho^{-1}\Gamma$ is the rescaled membrane and therefore, *M* is independent of ρ .

2.5 Description of the random cells and interfaces

We describe the domains occupied by the cells. As mentioned earlier, they are formed by randomly deforming a periodic structure. We transform the aforementioned periodic structure by a random diffeomorphism $\Phi : \mathbb{R}^2 \to \mathbb{R}^2$. Let

$$\mathbb{R}_2^+ := \bigcup_{n \in \mathbb{Z}^2} (n + Y^+), \quad \mathbb{R}_2^- := \bigcup_{n \in \mathbb{Z}^2} (n + Y^-), \quad \Gamma := \bigcup_{n \in \mathbb{Z}^2} (n + \Gamma).$$
(2.12)

The cells, the environment and the interfaces are hence deformed to $\Phi(\mathbb{R}_2^-)$, $\Phi(\mathbb{R}_2^+)$ and $\Phi(\Gamma)$. We emphasize that the topology of these sets are the same as before. Finally, the deformed structure is scaled to size ε , where $0 < \varepsilon \ll 1$, by the dilation operator εI where I is the identity operator. The final sets $\varepsilon \Phi(\mathbb{R}_2^-)$, $\varepsilon \Phi(\Gamma)$ and $\varepsilon \Phi(\mathbb{R}_2^+)$ thus are realistic models for the random cells, membranes and the environment for the biological problem at hand.

To model the cells inside an arbitrary bounded domain Ω as in (2.2), we would like to set $\Omega_{\varepsilon}^{+} := \Omega \cap \varepsilon \Phi(\mathbb{R}_{2}^{-})$ and $\Gamma_{\varepsilon} := \Omega \cap \varepsilon \Phi(\Gamma)$. However, a technicality is encountered, precisely, the intersection of $\varepsilon \Phi(\Gamma)$ with the boundary $\partial \Omega$ may not be empty. In this case, some cells are cut by the boundary of the body, which is not physically admissible. Moreover, an arbitrary diffeomorphism Φ may allow some deformed cells in $\varepsilon \Phi(\mathbb{R}_{2}^{-})$ to get arbitrarily close to each other. This imposes difficulties for rigorous mathematical analysis. In order to resolve these issues, we will impose a few conditions on Φ and refine the above construction in the next subsection.

2.6 Stationary ergodic setting

Let $(\mathcal{O}, \mathcal{F}, \mathbb{P})$ be some probability space on which $\Phi(x, \gamma) : \mathbb{R}^2 \times \mathcal{O} \to \mathbb{R}^2$ is defined. Throughout this paper, we assume that \mathcal{F} is countably generated so that the space $L^2(\mathcal{O})$ is separable. For a random variable $X \in L^1(\mathcal{O}, d\mathbb{P})$, we denote its expectation by

$$\mathbb{E}X = \int_{\mathcal{O}} X(\gamma) d\mathbb{P}(\gamma)$$

Throughout this paper, we assume that the group $(\mathbb{Z}^2, +)$ acts on \mathcal{O} by some action $\{\tau_n : \mathcal{O} \to \mathcal{O}\}_{n \in \mathbb{Z}^2}$, and that for all $n \in \mathbb{Z}^2$, τ_n is \mathbb{P} -preserving, that is,

$$\mathbb{P}(A) = \mathbb{P}(\tau_n A), \text{ for all } A \in \mathcal{F}.$$

We assume further that the action is *ergodic*, which means that for any $A \in \mathcal{F}$, if $\tau_n A = A$ for all $n \in \mathbb{Z}^2$, then necessarily $\mathbb{P}(A) \in \{0, 1\}$.

Following [19], we say that a random process $F \in L^1_{loc}(\mathbb{R}^2, L^1(\mathcal{O}))$ is (discrete) *stationary* if

$$\forall n \in \mathbb{Z}^2$$
, $F(x+n,\gamma) = F(x,\tau_n\gamma)$ for almost every *x* and *\gamma*. (2.13)

Clearly, a deterministic periodic function is a special case of stationary process. However, we precise that the above notion of stationarity is different from the classical one, see for instance [49] and [39]. Throughout this paper, we presume stationarity in the sense of (2.13) if not stated otherwise. What makes this notion useful is the following version of ergodic theorem [25, 34].

Proposition 2.1. (*i*) Let $F \in L^{\infty}(\mathbb{R}^2, L^1(\mathcal{O}))$ be a stationary random process. Equip \mathbb{Z}^2 with the norm $|n|_{\infty} = \max_{1 \le i \le 2} |n_i|$ for all $n \in \mathbb{Z}^2$. Then

$$\frac{1}{(2N+1)^2} \sum_{|n|_{\infty} \le N} F(x, \tau_n \gamma) \xrightarrow{L^{\infty}} \mathbb{E}F(x, \cdot) \quad \text{for a.e. } \gamma \in \mathcal{O}.$$
(2.14)

This implies in particular that

$$F\left(\frac{x}{\varepsilon},\gamma\right) \xrightarrow[\varepsilon \to 0]{L^{\infty} weak-*}{\varepsilon \to 0} \mathbb{E}\left(\int_{Y} F(x,\cdot)dx\right) \quad \text{for a.e. } \gamma \in \mathcal{O}.$$
(2.15)

(*ii*) For $p \in (1, \infty)$, suppose $F \in L^p(\mathcal{O}, L^p_{loc}(\mathbb{R}^2))$ is a stationary random process, then the above convergence results still hold if we replace L^∞ by L^p_{loc} .

We assume that for every $\gamma \in \mathcal{O}$, $\Phi(\cdot, \gamma)$ is a diffeomorphism from \mathbb{R}^2 to \mathbb{R}^2 and that it satisfies

$$\nabla \Phi(x, \gamma)$$
 is stationary. (2.16)

$$\operatorname{ess\,inf}_{\gamma\in\mathcal{O},x\in\mathbb{R}^2}\det(\nabla\Phi(x,\gamma))=\kappa>0,\tag{2.17}$$

$$\operatorname{ess\,sup}_{\gamma \in \mathcal{O}, x \in \mathbb{R}^2} |\nabla \Phi(x, \gamma)|_F = \kappa' > 0, \tag{2.18}$$

where $|\cdot|_F$ is the Frobenius norm and ess inf and ess sup are the essential infimum and the essential supremum, respectively. We point out that Φ^{-1} automatically satisfies similar conditions with constants κ_1 and κ'_1 . By the uniform Lipschitz assumption on Φ and Φ^{-1} above, we have

$$|(\kappa_1')^{-1}|y_1-y_2| \le |\Phi(y_1,\gamma)-\Phi(y_2,\gamma)| \le \kappa'|y_1-y_2|.$$

So the cells remain well separated after the deformation. To avoid the intersection of $\partial\Omega$ and the random cells $\varepsilon \Phi(\mathbb{R}_2^-)$, we erase those intersecting the boundary. More precisely, given a bounded and simply connected open set Ω with smooth boundary and a small number $\varepsilon \ll 1$, we denote by $\Omega_{1/\varepsilon}$ the scaled set $\{x \in \mathbb{R}^2 \mid \varepsilon x \in \Omega\}$. Let $\widetilde{\Omega_{1/\varepsilon}}$ be the shrunk set

$$\widetilde{\Omega_{1/\varepsilon}} := \{ x \in \Omega_{1/\varepsilon} \mid \operatorname{dist}(x, \partial \Omega_{1/\varepsilon}) \ge \operatorname{dist}(Y^-, \partial Y) \}.$$

We introduce for $n \in \mathbb{Z}^2$, Y_n and Y_n^{\pm} the translated cubes, reference cells and reference environments: $Y_n := n + Y$, $Y_n^{\pm} := n + Y^{\pm}$. Let $\mathcal{I}_{\varepsilon} \subset \mathbb{Z}^2$ be the indices of cubes Y_n such that $Y_n \in \widetilde{\Omega_{1/\varepsilon}}$. Note that $\mathcal{I}_{\varepsilon}$ corresponds to N_{ε} in the periodic case. We set Ω_{ε}^- to be

$$\Omega_{\varepsilon}^{-} := \bigcup_{n \in \mathcal{I}_{\varepsilon}} \varepsilon \Phi(Y_{n}^{-})$$
(2.19)

and then $\Omega_{\varepsilon}^{+} = \Omega \setminus \overline{\Omega_{\varepsilon}^{-}}$. We also define the following two notations:

$$E_{\varepsilon} := \bigcup_{n \in \mathcal{I}_{\varepsilon}} \varepsilon \Phi(Y_n) \quad \text{and} \quad K_{\varepsilon} := \Omega \setminus \overline{E_{\varepsilon}}.$$
(2.20)

Clearly, E_{ε} encloses all the cells in $\varepsilon \Phi(Y_n^-)$, $n \in \mathcal{I}_{\varepsilon}$ and their immediate surroundings $\varepsilon \Phi(Y_n^+)$; K_{ε} is a cushion layer near the boundary that prevents the cells from touching the boundary. From the construction we see that

$$\inf_{x\in\Omega_{\varepsilon}^{-}}\operatorname{dist}(x,\partial\Omega)\geq C\varepsilon\quad\text{and}\quad\sup_{x\in K_{\varepsilon}}\operatorname{dist}(x,\partial\Omega)\leq C\varepsilon.$$
(2.21)

Furthermore, we can check that

$$\sup_{n,j\in\mathcal{I}_{\varepsilon},n\neq j} \inf_{x\in\varepsilon\Phi(Y_n^-),y\in\varepsilon\Phi(Y_j^-)} |x-y| \ge C\varepsilon.$$
(2.22)

Above, the constants *C* vary but all of them depend only on dist(Y^- , Y), κ and κ' and are uniform in ε . This shows that the cells in Ω are well separated, *i.e.*, with a distance comparable to (if not much larger than) the size of the cells.

2.7 Main results in the random case

The first important result in the random case concerns an auxiliary problem which produces oscillating test functions that are used in the stochastic homogenization procedure. In the following theorem, a function f^{ext} in $W^{1,s}_{\text{loc}}(\mathbb{R}^2)$ is said to be an extension of $f \in W^{1,s}_{\text{loc}}(\mathbb{R}^+_2)$ if $f^{\text{ext}} = f$ on \mathbb{R}^+_2 and $\|f^{\text{ext}}\|_{W^{1,s}(K)} \leq C(K, \mathbb{R}^+_2) \|f\|_{W^{1,s}(\mathbb{R}^+_2 \cap K)}$, for any compact subset *K*. The following theorem holds.

Theorem 2.3. Let $\Phi(\cdot, \gamma)$ be a random diffeomorphism from \mathbb{R}^2 to \mathbb{R}^2 defined on the probability space $(\mathcal{O}, \mathcal{F}, \mathbb{P})$, and assume that (2.16)(2.17)(2.18) hold. For a.e. $\gamma \in \mathcal{O}$ and for any fixed vector $p \in \mathbb{R}^2$, the system

$$\nabla \cdot k_{0}(\nabla w_{p}^{+}(y) + p) = 0 \qquad in \ \Phi(\mathbb{R}_{2}^{+}, \gamma),$$

$$\nabla \cdot k_{0}(\nabla w_{p}^{-}(y) + p) = 0 \qquad in \ \Phi(\mathbb{R}_{2}^{-}, \gamma),$$

$$k_{0} \frac{\partial w_{p}^{+}}{\partial n}(y) - k_{0} \frac{\partial w_{p}^{-}}{\partial n}(y) = 0 \qquad on \ \Phi(\Gamma, \gamma),$$

$$w_{p}^{+} - w_{p}^{-} = \beta k_{0} \left(\frac{\partial w_{p}^{+}}{\partial n}(y) + v_{y} \cdot p \right) \qquad on \ \Phi(\Gamma, \gamma),$$

$$w_{p}^{\pm}(y, \gamma) = \widetilde{w}_{p}^{\pm}(\Phi^{-1}(y, \gamma), \gamma), \qquad \nabla \widetilde{w}_{p}^{\pm} \ are \ stationary,$$

$$\exists \ \widetilde{w}_{p}^{\text{ext}} \in H_{\text{loc}}^{1}(\mathbb{R}^{2}) \ that \ extends \ \widetilde{w}_{p}^{+} \ s.t. \quad \nabla \widetilde{w}_{p}^{\text{ext}} = P\left(\nabla \widetilde{w}_{p}^{+}\right)$$

$$and \quad \mathbb{E}\left(\int_{Y} \nabla \widetilde{w}_{p}^{\text{ext}}(\widetilde{y}, \cdot) d\widetilde{y}\right) = 0,$$

admits a unique (up to an addition of a random variable) weak solution $w_p = w_p^+ \chi_{\Phi(\mathbb{R}_2^+)} + w_p^- \chi_{\Phi(\mathbb{R}_2^-)}$ in $H^1_{\text{loc}}(\Phi(\mathbb{R}_2^+)) \times H^1_{\text{loc}}(\Phi(\mathbb{R}_2^-))$ and the operator P above denotes the extension operator of Corollary A.1. The precise weak formulation of the system above is postponed to section 7, where the proof of this theorem is given; see (7.1). We remark that the non-unique additive random variable is not important and what matters is the fact that the gradient ∇w_p of the solution is unique. The second main result in the random case is the following homogenization theorem.

Theorem 2.4. Let Ω be a bounded and connected open subset of \mathbb{R}^2 with regular boundary. Let Φ be a random diffeomorphism on $(\mathcal{O}, \mathcal{F}, \mathbb{P})$ satisfying (2.16)(2.17)(2.18). Assume that the cells Ω_{ε}^- are constructed as in Section 2.6. Then for a.e. $\gamma \in \mathcal{O}$, the solution $u_{\varepsilon}(\cdot, \gamma) = (u_{\varepsilon}^+, u_{\varepsilon}^-)$ of (2.2) satisfies the following properties:

- (i) We can extend $u_{\varepsilon}^{+}(\cdot, \gamma)$ to $u_{\varepsilon}^{\text{ext}}(\cdot, \gamma) \in H^{1}(\Omega)$, where $u_{\varepsilon}^{\text{ext}}(\cdot, \gamma)$ converges weakly, as $\varepsilon \to 0$, to a deterministic function $u_{0} \in H^{1}(\Omega)$.
- (ii) The function $u_{\varepsilon}(\cdot, \gamma)$ converges strongly in $L^2(\Omega)$ to u_0 above. Further, let Q be the trivial extension operator setting Qf = 0 outside the domain of f, and define

$$\varrho := \det\left(\mathbb{E}\int_{Y} \nabla \Phi(z,\cdot)dz\right)^{-1}, \quad \theta := \varrho \mathbb{E}\int_{Y^{-}} \det \nabla \Phi(z,\cdot)dz < 1, \quad (2.24)$$

where det denotes the determinant. Then, Qu_{ε}^{-} converges weakly to θu_{0} in $L^{2}(\Omega)$.

(iii) The function u_0 is the unique weak solution in $H^1_{\mathbb{C}}(\Omega)$ to the homogenized equation

$$\begin{cases} \nabla \cdot K^* \nabla u_0(x) = 0, & x \in \Omega, \\ n(x) \cdot K^* \nabla u_0(x) = g, & x \in \partial \Omega, \end{cases}$$
(2.25)

The homogenized admittivity coefficient K^* *is given by* $\forall (i, j) \in \{1, 2\}^2$ *,*

$$K_{ij}^{*} = k_0 \left(\delta_{ij} + \varrho \mathbb{E} \int_{\Phi(Y)} e_j \cdot (\chi_{\Phi(Y^+)} \nabla w_{e_i}^+ + \chi_{\Phi(Y^-)} \nabla w_{e_i}^-)(y, \cdot) \, dy \right),$$
(2.26)

where $\{e_i\}_{i=1}^2$ is the Euclidean basis of \mathbb{R}^2 and for each $p \in \mathbb{R}^2$, the pair of functions (w_p^+, w_p^-) is the unique solution to the auxiliary system (2.23).

We mention the fact that K^* is uniformly elliptic, the proof of which is standard and is omitted. In the dilute limit $\rho := \sqrt{|Y^-|} \ll 1$, we obtain the following approximation of the effective permittivity for the dilute suspension:

$$K_{ij}^* = k_0 (I + f \mathbb{E} M_{ij}) + o(f), \qquad (2.27)$$

where ρ accounts for the averaged change of volume due to the random diffeomorphism and $f := \rho \rho^2$ is the volume fraction occupied by the cells ; the polarization matrix *M* is defined by

$$M_{ij} = \beta k_0 \int_{\rho^{-1} \Phi(\Gamma)} \tilde{\psi}_i n_j \, ds(\tilde{y}), \qquad (2.28)$$

where

$$\tilde{\psi}_i = -(I + \beta k_0 n \cdot \nabla \mathcal{D}_{\rho^{-1}\Phi(\Gamma)})^{-1}[n_i],$$

with $\mathcal{D}_{\rho^{-1}\Phi(\Gamma)}$ the double layer potential associated to the deformed inclusion scaled to the unit length scale.

3 Analysis of the problem

For a fixed ε , recall that $H^1_{\mathbb{C}}(\Omega^+_{\varepsilon})$ denotes the Sobolev space $H^1(\Omega^+_{\varepsilon})/\mathbb{C}$, which can be represented as

$$H^{1}_{\mathbb{C}}(\Omega^{+}_{\varepsilon}) = \left\{ u \in H^{1}(\Omega^{+}_{\varepsilon}) \mid \int_{\Omega^{+}_{\varepsilon}} u(x) dx = 0 \right\}.$$
(3.1)

The natural functional space for (2.2) is

$$W_{\varepsilon} := \left\{ u = u^+ \chi_{\varepsilon}^+ + u^- \chi_{\varepsilon}^- \mid u^+ \in H^1_{\mathbb{C}}(\Omega_{\varepsilon}^+), u^- \in H^1(\Omega_{\varepsilon}^-) \right\},$$
(3.2)

where χ^{\pm}_{ε} are the characteristic functions of the sets $\Omega^{\pm}_{\varepsilon}$. We can verify that

$$\|u\|_{W_{\varepsilon}} = \left(\|\nabla u^{+}\|_{L^{2}(\Omega_{\varepsilon}^{+})}^{2} + \|\nabla u^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2} + \varepsilon\|u^{+} - u^{-}\|_{L^{2}(\Gamma_{\varepsilon})}^{2}\right)^{\frac{1}{2}}$$
(3.3)

defines a norm on W_{ε} . In fact, as it will be seen in Proposition 3.2, this norm is equivalent to the standard norm on W_{ε} which is

$$\|u\|_{H^{1}_{\mathbb{C}}(\Omega^{+}_{\varepsilon})\times H^{1}(\Omega^{-}_{\varepsilon})} = \left(\|\nabla u^{+}\|^{2}_{L^{2}(\Omega^{+}_{\varepsilon})} + \|\nabla u^{-}\|^{2}_{L^{2}(\Omega^{-}_{\varepsilon})} + \|u^{-}\|^{2}_{L^{2}(\Omega^{-}_{\varepsilon})}\right)^{\frac{1}{2}}.$$
(3.4)

3.1 Existence and uniqueness of a solution

Problem (2.2) should be understood through its weak formulation as follows: For a fixed $\varepsilon > 0$, find $u_{\varepsilon} \in W_{\varepsilon}$ such that

$$\begin{split} \int_{\Omega_{\varepsilon}^{+}} k_{0} \nabla u_{\varepsilon}^{+}(x) \cdot \nabla \overline{v^{+}}(x) dx + \int_{\Omega_{\varepsilon}^{-}} k_{0} \nabla u_{\varepsilon}^{-}(x) \cdot \nabla \overline{v^{-}}(x) ds(x) \\ &+ \frac{1}{\varepsilon \beta} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+} - u_{\varepsilon}^{-})(x) \overline{(v^{+} - v^{-})}(x) ds(x) = \int_{\partial \Omega} g(x) \overline{v^{+}}(x) ds(x), \end{split}$$
(3.5)

for any function $v \in W_{\varepsilon}$.

Define the sesquilinear form $a_{\varepsilon}(\cdot, \cdot)$ on $W_{\varepsilon} \times W_{\varepsilon}$ by

$$a_{\varepsilon}(u,v) := \int_{\Omega_{\varepsilon}^{+}} k_{0} \nabla u^{+} \cdot \nabla \overline{v^{+}} dx + \int_{\Omega_{\varepsilon}^{-}} k_{0} \nabla u^{-} \cdot \nabla \overline{v^{-}} dx + \frac{1}{\varepsilon \beta} \int_{\Gamma_{\varepsilon}} (u^{+} - u^{-}) \overline{(v^{+} - v^{-})} ds.$$
(3.6)

Associate the following anti-linear form on W_{ε} to the boundary data *g*:

$$\ell(u) := \int_{\partial \Omega} g \overline{u^+} ds. \tag{3.7}$$

The forms a_{ε} and ℓ are bounded. Moreover, a_{ε} is coercive in the following sense

$$\Re k_0^{-1}a_{\varepsilon}(u,u) = \left(\int_{\Omega_{\varepsilon}^+} |\nabla u^+|^2 dx + \int_{\Omega_{\varepsilon}^-} |\nabla u^-|^2 dx\right) + \frac{1}{\varepsilon\beta'} \int_{\Gamma_{\varepsilon}} |u^+ - u^-|^2 ds \ge C \|u\|_{W_{\varepsilon}}^2, \quad (3.8)$$

where $\beta' := \delta(\sigma_0 \sigma_m + \omega^2 \epsilon_0 \epsilon_m) / (\sigma_m^2 + \omega^2 \epsilon_m^2)$. Consequently, due to the Lax–Milgram theorem we have existence and uniqueness for (2.2) for each fixed ϵ and for almost every $\gamma \in \mathcal{O}$. Note that *C* can be chosen independent of ϵ .

Proposition 3.1. Let $g \in H^{-1/2}(\partial\Omega)$. There exists a unique $u_{\varepsilon} \in W_{\varepsilon}$ so that

$$a_{\varepsilon}(u_{\varepsilon},\varphi) = \ell(\varphi), \quad \forall \varphi \in W_{\varepsilon}.$$
(3.9)

To end this subsection we remark that the two norms on W_{ε} are equivalent.

Proposition 3.2. The norm $\|\cdot\|_{W_{\varepsilon}}$ is equivalent with the standard norm on $H^1_{\mathbb{C}}(\Omega^+_{\varepsilon}) \times H^1(\Omega^-_{\varepsilon})$. Moreover, we can find two positive constants $C_1 < C_2$, independent of ε , so that

$$C_1 \|u\|_{W_{\varepsilon}} \le \|u\|_{H^1_{C} \times H^1} \le C_2 \|u\|_{W_{\varepsilon}}, \tag{3.10}$$

for any $u \in H^1_{\mathbb{C}}(\Omega^+_{\varepsilon}) \times H^1(\Omega^-_{\varepsilon})$.

Similar equivalence relation was established by Monsurrò [45], whose method can be adapted easily to the current case. For the sake of completeness, we present the details in Appendix C.

3.2 Energy estimate

For any fixed $\gamma \in \mathcal{O}$ and a sequence of $\varepsilon \to 0$, by solving (2.2) we obtain the sequence $u_{\varepsilon} = u_{\varepsilon}^{+} \chi_{\varepsilon}^{+} + u_{\varepsilon}^{-} \chi_{\varepsilon}^{-}$. We obtain some *a priori* estimates for u_{ε} .

We first recall that the extension theorem (Theorem A.2) yields a Poincaré–Wirtinger inequality in $H^1_{\mathbb{C}}(\Omega^+_{\varepsilon})$ with a constant independent of ε . Indeed, Corollary B.1 shows that for all $v^+ \in H^1_{\mathbb{C}}(\Omega^+_{\varepsilon})$, there exists a constant *C*, independent of ε , such that

$$\|v^+\|_{L^2(\Omega^+_{\epsilon})} \le C \|\nabla v^+\|_{L^2(\Omega^+_{\epsilon})}.$$

Similarly, we can find a constant, independent of ε , by applying the trace theorem in $H^1(\Omega_{\varepsilon}^+)$. Using Corollary B.2, the following result holds.

Proposition 3.3. Let $g \in H^{-\frac{1}{2}}(\partial\Omega)$. For almost any $\gamma \in \mathcal{O}$, let $\Omega = \Omega_{\varepsilon}^+ \cup \Gamma_{\varepsilon} \cup \Omega_{\varepsilon}^-$. Then there exist constants C's, independent of ε and γ , such that the solution u_{ε} to (2.2) satisfies the following estimates:

$$\|\nabla u_{\varepsilon}^{+}\|_{L^{2}(\Omega_{\varepsilon}^{+})} + \|\nabla u_{\varepsilon}^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})} \le C|k_{0}|^{-1}\|g\|_{H^{-\frac{1}{2}}(\partial\Omega)'}$$
(3.11)

$$\|u_{\varepsilon}^{+} - u_{\varepsilon}^{-}\|_{L^{2}(\Gamma_{\varepsilon})} \leq C|k_{0}|^{-1}\sqrt{\varepsilon\beta'}\|g\|_{H^{-\frac{1}{2}}(\partial\Omega)}.$$
(3.12)

Proof. By taking $\varphi = u_{\varepsilon}$ in (3.9), and taking the real part of resultant equality, we get

$$\|\nabla u_{\varepsilon}^{+}\|_{L^{2}(\Omega_{\varepsilon}^{+})}^{2} + \|\nabla u_{\varepsilon}^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2} + (\varepsilon\beta')^{-1}\|u_{\varepsilon}^{+} - u_{\varepsilon}^{-}\|_{L^{2}(\Gamma_{\varepsilon})}^{2} = \Re k_{0}^{-1} \langle g, u_{\varepsilon}^{+} \rangle.$$
(3.13)

Here $\langle g, u_{\varepsilon}^+ \rangle = \int_{\partial \Omega} g \overline{u_{\varepsilon}^+} ds$ is the pairing on $H^{-\frac{1}{2}}(\partial \Omega) \times H^{\frac{1}{2}}(\partial \Omega)$, for which we have the estimate

$$|\langle g, u_{\varepsilon}^{+} \rangle| \leq ||g||_{H^{-\frac{1}{2}}(\partial\Omega)} ||u_{\varepsilon}^{+}||_{H^{\frac{1}{2}}(\partial\Omega)} \leq C_{1} ||g||_{H^{-\frac{1}{2}}(\partial\Omega)} ||u_{\varepsilon}^{+}||_{H^{1}(\Omega_{\varepsilon}^{+})}.$$

thanks to the Cauchy–Schwarz inequality and Corollary (B.2). C_1 is here a constant which does not depend on ε .

Applying Proposition 3.2 yields

$$|\langle g, u_{\varepsilon}^{+} \rangle| \leq C_{2} ||g||_{H^{-\frac{1}{2}}(\partial \Omega)} ||u_{\varepsilon}||_{W_{\varepsilon}},$$

with a constant C_2 independent of ε .

Using this in (3.13) along with the coercivity of *a* we get

$$\|u_{\varepsilon}\|_{W_{\varepsilon}} \leq C_3 |k_0|^{-1} \|g\|_{H^{-\frac{1}{2}}(\partial\Omega)}$$

where C_3 is still independent of ε .

It follows also that

$$|\langle g, u_{\varepsilon}^{+}\rangle| \leq C_2 C_3 |k_0|^{-1} ||g||_{H^{-\frac{1}{2}}(\partial\Omega)}$$

Substituting this estimate into the right-hand side of (3.13), we get the desired estimates. \Box

Next, we apply the extension theorem (Theorem A.2) to obtain a bounded sequence in $H^1(\Omega)$ for which we can extract a converging subsequence.

Proposition 3.4. Suppose that the same conditions of the previous proposition hold. Let $P_{\gamma}^{\varepsilon} : H^1(\Omega_{\varepsilon}^+) \to H^1(\Omega)$ be the extension operator of Theorem A.2. Then we have

$$\|P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+}\|_{H^{1}(\Omega)} \leq C|k_{0}|^{-1}\|g\|_{H^{-\frac{1}{2}}(\partial\Omega)'}$$
(3.14)

and

$$\|P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+}-u_{\varepsilon}\|_{L^{2}(\Omega)} \leq C\varepsilon|k_{0}|^{-1}(1+\sqrt{\beta'})\|g\|_{H^{-\frac{1}{2}}(\partial\Omega)}.$$
(3.15)

Proof. The first inequality is a direct result of (A.11), (B.2) and (3.11). For the second inequality, we have

$$\begin{split} \|P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+}-u_{\varepsilon}\|_{L^{2}(\Omega)} &= \|P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+}-u_{\varepsilon}^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})} \\ &\leq C\sqrt{\varepsilon}\|P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+}-u_{\varepsilon}^{-}\|_{L^{2}(\Gamma_{\varepsilon})} + C\varepsilon\|\nabla(P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+}-u_{\varepsilon}^{-})\|_{L^{2}(\Omega_{\varepsilon}^{-})}. \end{split}$$

Here, we have used estimate (C.3). Now, $\|P_{\gamma}^{\varepsilon}u_{\varepsilon}^{+} - u_{\varepsilon}^{-}\|_{L^{2}(\Gamma_{\varepsilon})} = \|u_{\varepsilon}^{+} - u_{\varepsilon}^{-}\|_{L^{2}(\Gamma_{\varepsilon})}$ is bounded in (3.12). The second term is bounded from above by

$$C\varepsilon \|\nabla P^{\varepsilon}_{\gamma} u^{+}_{\varepsilon}\|_{L^{2}(\Omega^{-}_{\varepsilon})} + C\varepsilon \|\nabla u^{-}_{\varepsilon}\|_{L^{2}(\Omega^{-}_{\varepsilon})} \leq C\varepsilon (\|\nabla u^{+}_{\varepsilon}\|_{L^{2}(\Omega^{+}_{\varepsilon})} + \|\nabla u^{-}_{\varepsilon}\|_{L^{2}(\Omega^{-}_{\varepsilon})}),$$

where we have used again (A.11). This gives the desired estimates.

Remark 3.1. As a consequence of the previous proposition, we get a sequence in $H^1(\Omega)$, namely $P^{\varepsilon}_{\gamma}u^+_{\varepsilon}$, which is a good estimate of u_{ε} in $L^2(\Omega)$ and from which we can extract a subsequence weakly converging in $H^1(\Omega)$ and strongly in $L^2(\Omega)$.

4 Homogenization

We follow [5, 6] to derive a homogenized problem for the model with two-scale asymptotic expansions and to prove a rigorous two-scale convergence. In [45], the homogenization of an analogue problem is developed and proved with another method.

4.1 Two-scale asymptotic expansions

We assume that the solution u_{ε} admits the following two-scale asymptotic expansion

$$\forall x \in \Omega \ u_{\varepsilon}(x) = u_0(x) + \varepsilon u_1(x, \frac{x}{\varepsilon}) + o(\varepsilon),$$

with

$$y \mapsto u_1(x,y)$$
 Y-periodic and $u_1(x,y) = \begin{cases} u_1^+(x,y) \text{ in } \Omega \times Y^+, \\ u_1^-(x,y) \text{ in } \Omega \times Y^-. \end{cases}$

We choose a test function φ_{ε} of the same form as u_{ε} :

$$\forall x \in \Omega, \ \varphi_{\varepsilon}(x) = \varphi_0(x) + \varepsilon \varphi_1(x, \frac{x}{\varepsilon}),$$

with φ_0 smooth in Ω , $\varphi_1(x, .)$ *Y*-periodic,

$$arphi_1(x,y) = \left\{ egin{array}{c} arphi_1^+(x,y) \ {
m in}\,\Omega imes Y^+, \ arphi_1^-(x,y) \ {
m in}\,\Omega imes Y^-, \end{array}
ight.$$

and φ_1^{\pm} smooth in $\Omega \times Y^{\pm}$.

In order to prove items (ii) and (iii) in Theorem 2.1, we perform an asymptotic expansion of the variational formulation (3.9). We thus inject these ansatz in the variational formulation and only consider the order 0 of the different integrals.

At order 0,

$$\nabla u_{\varepsilon}(x) = \nabla u_0(x) + \nabla_y u_1(x, \frac{x}{\varepsilon}) + o(\varepsilon).$$

Thanks to Lemma 4.1, we then have for the two first integrals:

$$\begin{split} \int_{\Omega_{\varepsilon}^{+}} k_{0} \left(\nabla u_{0}(x) + \nabla_{y} u_{1}^{+}(x, \frac{x}{\varepsilon}) \right) \cdot \left(\nabla \overline{\varphi}_{0}(x) + \nabla_{y} \overline{\varphi}_{1}^{+}(x, \frac{x}{\varepsilon}) \right) dx \\ &= \int_{\Omega} \int_{Y^{+}} k_{0} \left(\nabla u_{0}(x) + \nabla_{y} u_{1}^{+}(x, y) \right) \cdot \left(\nabla \overline{\varphi}_{0}(x) + \nabla_{y} \overline{\varphi}_{1}^{+}(x, y) \right) dx dy + o(\varepsilon) \end{split}$$

and

$$\begin{split} \int_{\Omega_{\varepsilon}^{-}} k_{0} \left(\nabla u_{0}(x) + \nabla_{y} u_{1}^{-}(x, \frac{x}{\varepsilon}) \right) \cdot \left(\nabla \overline{\varphi}_{0}(x) + \nabla_{y} \overline{\varphi}_{1}^{-}(x, \frac{x}{\varepsilon}) \right) dx \\ &= \int_{\Omega} \int_{Y^{-}} k_{0} \left(\nabla u_{0}(x) + \nabla_{y} u_{1}^{-}(x, y) \right) \cdot \left(\nabla \overline{\varphi}_{0}(x) + \nabla_{y} \overline{\varphi}_{1}^{-}(x, y) \right) dx dy + o(\varepsilon). \end{split}$$

We write the third integral in (3.6) as the sum, over all squares $Y_{\varepsilon,n}$, of integrals on the boundaries $\Gamma_{\varepsilon,n}$. We have

$$\begin{split} \frac{1}{\beta\varepsilon} \int_{\Gamma_{\varepsilon}} \left(u_{\varepsilon}^{+}(x,\frac{x}{\varepsilon}) - u_{\varepsilon}^{-}(x,\frac{x}{\varepsilon}) \right) \left(\overline{\varphi}_{\varepsilon}^{+}(x,\frac{x}{\varepsilon}) - \overline{\varphi}_{\varepsilon}^{-}(x,\frac{x}{\varepsilon}) \right) ds(x) \\ &= \frac{1}{\beta\varepsilon} \sum_{n \in N_{\varepsilon}} \int_{\Gamma_{\varepsilon,n}} \left(u_{\varepsilon}^{+}(x,\frac{x}{\varepsilon}) - u_{\varepsilon}^{-}(x,\frac{x}{\varepsilon}) \right) \left(\overline{\varphi}_{\varepsilon}^{+}(x,\frac{x}{\varepsilon}) - \overline{\varphi}_{\varepsilon}^{-}(x,\frac{x}{\varepsilon}) \right) ds(x). \end{split}$$

Let $x_{0,n}$ be the center of $Y_{\varepsilon,n}$ for each $n \in N_{\varepsilon}$. We perform Taylor expansions with respect to the variable x around $x_{0,n}$ for all functions $(u_i)_{i \in \{1,2\}}$ and $(\varphi_i)_{i \in \{1,2\}}$ on $Y_{\varepsilon,n}$. After the change of variables $\varepsilon(y - y_{0,n}) = x - x_{0,n}$, we obtain that

$$u_{\varepsilon}(x) = u_0(x_{0,n}) + \varepsilon u_1(x,y) + \varepsilon(y - y_{0,n}) \cdot \nabla u_0(x_{0,n}) + o(\varepsilon),$$

$$\varphi_{\varepsilon}(x) = \varphi_0(x_{0,n}) + \varepsilon \varphi_1(x,y) + \varepsilon(y - y_{0,n}) \cdot \nabla \varphi_0(x_{0,n}) + o(\varepsilon).$$

Consequently, the third integral in the variational formulation (3.9) becomes

$$\frac{\varepsilon^2}{\beta} \sum_{n \in N_{\varepsilon}} \int_{\Gamma_n} \left(u_1^+(x_{0,n}, y) - u_1^-(x_{0,n}, y) \right) \left(\overline{\varphi}_1^+(x_{0,n}, y) - \overline{\varphi}_1^-(x_{0,n}, y) \right) ds(y).$$

Finally, Lemma 4.1 gives us that

$$\begin{split} \frac{1}{\varepsilon\beta} \int_{\Gamma_{\varepsilon}} \left(u_{\varepsilon}^{+} - u_{\varepsilon}^{-} \right) \left(\overline{\varphi}_{\varepsilon}^{+} - \overline{\varphi}_{\varepsilon}^{-} \right) ds \\ &= \frac{1}{\beta} \int_{\Omega} \int_{\Gamma} \left(u_{1}^{+}(x,y) - u_{1}^{-}(x,y) \right) \left(\overline{\varphi}_{1}^{+}(x,y) - \overline{\varphi}_{1}^{-}(x,y) \right) dx ds(y) + o(\varepsilon). \end{split}$$

Moreover, we can easily see that

$$\int_{\partial\Omega} g \overline{\varphi}_{\varepsilon}^+ ds = \int_{\partial\Omega} g \overline{\varphi}_0 ds + o(\varepsilon).$$

The order 0 of the variational formula is thus given by

$$\begin{split} &\int_{\Omega} \int_{Y^+} k_0 \left(\nabla u_0(x) + \nabla_y u_1^+(x,y) \right) \cdot \left(\nabla \overline{\varphi}_0(x) + \nabla_y \overline{\varphi}_1^+(x,y) \right) dx dy \\ &+ \int_{\Omega} \int_{Y^-} k_0 \left(\nabla u_0(x) + \nabla_y u_1^-(x,y) \right) \cdot \left(\nabla \overline{\varphi}_0(x) + \nabla_y \overline{\varphi}_1^-(x,y) \right) dx dy \\ &+ \frac{1}{\beta} \int_{\Omega} \int_{\Gamma} \left(u_1^+(x,y) - u_1^-(x,y) \right) \left(\overline{\varphi}_1^+(x,y) - \overline{\varphi}_1^-(x,y) \right) dx ds(y) \\ &- \int_{\partial \Omega} g(x) \overline{\varphi}_0(x) ds(x) = 0. \end{split}$$

By taking $\varphi_0 = 0$, it follows that

$$\begin{split} &\int_{\Omega} \int_{Y^+} k_0 \left(\nabla u_0(x) + \nabla_y u_1^+(x,y) \right) \cdot \nabla_y \overline{\varphi}_1^+(x,y) dx dy \\ &+ \int_{\Omega} \int_{Y^-} k_0 \left(\nabla u_0(x) + \nabla_y u_1^-(x,y) \right) \cdot \nabla_y \overline{\varphi}_1^-(x,y) dx dy \\ &+ \frac{1}{\beta} \int_{\Omega} \int_{\Gamma} \left(u_1^+(x,y) - u_1^-(x,y) \right) \left(\overline{\varphi}_1^+(x,y) - \overline{\varphi}_1^-(x,y) \right) dx ds(y) = 0, \end{split}$$

which is exactly the variational formulation of the cell problem (2.6) and definition (2.7) of u_1 .

By taking $\varphi_1 = 0$, we recover the variational formulation of the homogenized problem (2.4):

$$\begin{split} &\int_{\Omega} \int_{Y^{+}} k_0 \left(\nabla u_0(x) + \nabla_y u_1^{+}(x,y) \right) \cdot \nabla \overline{\varphi}_0(x) dx dy \\ &+ \int_{\Omega} \int_{Y^{-}} k_0 \left(\nabla u_0(x) + \nabla_y u_1^{-}(x,y) \right) \cdot \nabla \overline{\varphi}_0(x) dx dy \\ &- \int_{\partial \Omega} g(x) \overline{\varphi}_0(x) ds(x) = 0. \end{split}$$

We introduce some function spaces, which will be very useful in the following:

- $C^{\infty}_{\sharp}(D)$ is the space of functions, which are *Y* periodic and in $C^{\infty}(D)$,
- $L^2_{\sharp}(D)$ is the completion of $C^{\infty}_{\sharp}(D)$ in the L^2 -norm,
- $H^1_{\sharp}(D)$ is the completion of $C^{\infty}_{\sharp}(D)$ in the H^1 -norm,
- $L^2(\Omega, H^1_{\sharp}(D))$ is the space of square integrable functions on Ω with values in the space $H^1_{\sharp}(D)$,
- $\mathcal{D}(\Omega)$ is the space of infinitely smooth functions with compact support in Ω ,
- D(Ω, C[∞]_µ(D)) is the space of infinitely smooth functions with compact support in Ω and with values in the space C[∞]_µ,

where *D* is *Y*, Y^+ , Y^- or Γ .

The following lemma was used in the preceding proof. It follows from [5, Lemma 3.1].

Lemma 4.1. *Let f be a smooth function. We have*

(i)
$$\varepsilon^{2} \sum_{n \in N_{\varepsilon}} \int_{\Gamma_{\varepsilon,n}} f(x_{0,n}, y) ds(y) = \int_{\Omega} \int_{\Gamma} f(x, y) dx ds(y) + o(\varepsilon);$$

(ii) $\int_{\Omega_{\varepsilon}^{+}} f(x, \frac{x}{\varepsilon}) dx = \int_{\Omega} \int_{Y^{+}} f(x, y) dx dy + o(\varepsilon)$
and $\int_{\Omega_{\varepsilon}^{-}} f(x, \frac{x}{\varepsilon}) dx = \int_{\Omega} \int_{Y^{-}} f(x, y) dx dy + o(\varepsilon).$

We prove that the following lemmas hold.

Lemma 4.2. The homogenized problem admits a unique solution in $H^1_{\mathbb{C}}(\Omega)$.

Proof. The effective admittivity can be rewritten as

$$\begin{split} K_{i,j}^{*} &= k_{0} \int_{Y^{+}} (\nabla w_{i}^{+} + e_{i}) \cdot (\nabla \overline{w_{j}^{+}} + e_{j}) dx + k_{0} \int_{Y^{-}} (\nabla w_{i}^{-} + e_{i}) \cdot (\nabla \overline{w_{j}^{-}} + e_{j}) dx \\ &+ \frac{1}{\beta} \int_{\Gamma} (w_{i}^{+} - w_{i}^{-}) (\overline{w_{j}^{+}} - \overline{w_{j}^{-}}) ds, \quad i, j = 1, 2. \end{split}$$

Therefore, it follows that, for $\xi = (\xi_1, \xi_2) \in \mathbb{R}^2$,

$$K^{*}\xi\cdot\xi = k_{0}\int_{Y^{+}}|\nabla w^{+}+\xi|^{2}dx + k_{0}\int_{Y^{-}}|\nabla w^{-}+\xi|^{2}dx + \frac{1}{\beta}\int_{\Gamma}|w^{+}-w^{-}|^{2}ds,$$

where $w = \sum_i \xi_i w_i$. Since $\Re \beta \ge 0$,

$$K^*\xi\cdot\xi\geq k_0\int_{Y^+}|
abla w^++\xi|^2dx+k_0\int_{Y^-}|
abla w^-+\xi|^2dx.$$

Consequently, it follows from [3] that there exist two positive constants C_1 and C_2 such that

$$C_1|\xi|^2 \le \Re K^* \xi \cdot \xi \le C_2|\xi|^2$$

Standard elliptic theory yields existence and uniqueness of a solution to the homogenized problem in $H^1_{\mathbb{C}}(\Omega)$.

Lemma 4.3. The cell problem (2.6) admits a unique solution in $H^1_{\sharp}(Y^+)/\mathbb{C} \times H^1_{\sharp}(Y^-)$.

Proof. Let us introduce the Hilbert space

$$W_{\sharp} := \left\{ v := v^{+} \chi_{Y^{+}} + v^{-} \chi_{Y^{-}} | (v^{+}, v^{-}) \in H^{1}_{\mathbb{C}}(Y^{+}) \times H^{1}(Y^{-}) \right\},\$$

equipped with the norm

$$\|v\|_{W_{\sharp}}^{2} = \|\nabla v^{+}\|_{L^{2}(Y^{+})}^{2} + \|\nabla v^{-}\|_{L^{2}(Y^{-})}^{2} + \|v^{+} - v^{-}\|_{L^{2}(\Gamma)}^{2}.$$

We consider the following problem:

$$\begin{cases} \text{Find } w_{i} \in W_{\sharp} \text{ such that for all } \varphi \in W_{\sharp} \\ \int_{Y^{+}} k_{0} \nabla w_{i}^{+}(y) \cdot \nabla \overline{\varphi}^{+}(y) dy + \int_{Y^{-}} k_{0} \nabla w_{i}^{-}(y) \cdot \nabla \overline{\varphi}^{-}(y) dy \\ + \frac{1}{\beta} \int_{\Gamma} \left(w_{i}^{+} - w_{i}^{-} \right) (y) \left(\overline{\varphi}^{+} - \overline{\varphi}^{-} \right) (y) ds(y) = \\ - \int_{Y^{+}} k_{0} \nabla y_{i} \cdot \nabla \overline{\varphi}^{+}(y) dy - \int_{Y^{-}} k_{0} \nabla y_{i} \cdot \nabla \overline{\varphi}^{-}(y) dy. \end{cases}$$

$$(4.1)$$

Lax–Milgram theorem gives us existence and uniqueness of a solution. Moreover, one can show that this ensures the existence of a unique solution in $H^1_{\sharp}(Y^+)/\mathbb{C} \times H^1_{\sharp}(Y^-)$ for the cell problem (2.6).

We present in the following numerical examples (Figures 4.1–4.4) the real and imaginary parts of the solutions w_1 and w_2 of the cell problems.

4.2 Convergence

We present in this section a rigorous proof of the convergence of the original problem to the homogenized one. We use for this purpose the two-scale convergence technique and hence need first of all some bounds on u_{ε} to ensure the convergence.



Figure 4.1: *Real and imaginary parts of the cell problem solution* w_1 .



Figure 4.2: *Real and imaginary parts of the cell problem solution* w_2 .



Figure 4.3: Gradient vector field of the real and imaginary parts of w_1 .



Figure 4.4: Gradient vector field of the real and imaginary parts of w₂.

4.2.1 A priori estimates

Theorem 4.1. The function u_{ε}^+ is uniformly bounded with respect to ε in $H^1(\Omega_{\varepsilon}^+)$, *i.e.*, there exists a constant *C*, independent of ε , such that

$$\|u_{\varepsilon}^+\|_{H^1(\Omega_{\varepsilon}^+)} \leq C.$$

Proof. Combining (3.11) and Poincaré - Wirtinger inequality, we obtain immediately the wanted result.

The proof of the following result follows the one of Lemma 2.8 in [45].

Lemma 4.4. There exists a constant *C*, which does not depend on ε , such that for all $v \in W_{\varepsilon}$:

$$\|v^-\|_{L^2(\Omega_{\varepsilon}^-)} \leq C \|v\|_{W_{\varepsilon}}.$$

Proof. We write the norm $\|v^-\|_{L^2(\Omega_{\varepsilon}^-)}$ as a sum over all the cells.

$$\|v^-\|_{L^2(\Omega_{\varepsilon}^-)}^2 = \sum_{n \in N_{\varepsilon}} \|v^-\|_{L^2(Y_{\varepsilon,n}^-)}^2 = \sum_{n \in N_{\varepsilon}} \int_{Y_{\varepsilon,n}^-} |v^-(x)|^2 dx.$$

We perform the change of variable $y = \frac{x}{\varepsilon}$ and get

$$\|v^-\|_{L^2(\Omega_{\varepsilon}^-)}^2 = \varepsilon^2 \sum_{n \in N_{\varepsilon}} \int_{Y_n^-} |v_{\varepsilon}^-(y)|^2 dy,$$
(4.2)

where $v_{\varepsilon}^{-}(y) := v^{-}(\varepsilon y)$ for all $y \in Y^{-}$ and $Y_{n}^{-} = n + Y^{-}$ with $n \in N_{\varepsilon}$.

Let *W* denote the following Hilbert space:

$$W := \left\{ v := v^+ \chi_{Y^+} + v^- \chi_{Y^-} | (v^+, v^-) \in H^1_{\mathbb{C}}(Y^+) \times H^1(Y^-) \right\},$$

equipped with the norm:

$$\|v\|_{W}^{2} = \|\nabla v^{+}\|_{L^{2}(Y^{+})}^{2} + \|\nabla v^{-}\|_{L^{2}(Y^{-})}^{2} + \|v^{+} - v^{-}\|_{L^{2}(\Gamma)}^{2}.$$

We first prove that there exists a constant C_1 , independent of ε , such that for every $v \in W$:

$$\|v^{-}\|_{L^{2}(Y^{-})} \le C_{1} \|v\|_{W}.$$
(4.3)

We proceed by contradiction. Suppose that for any $n \in \mathbb{N}^*$, there exists $v_n \in W_{\varepsilon}$ such that

$$\|v_n^-\|_{L^2(Y^-)} = 1$$
 and $\|v_n\|_W \le \frac{1}{n}$.

Since $||v_n^-||_{L^2(Y^-)} = 1$ and $||\nabla v_n^-||_{L^2(Y^-)} \le ||v_n||_W \le \frac{1}{n}$, v_n^- is bounded in $H^1(Y^-)$. Therefore it converges weakly in $H^1(Y^-)$. By weak compactness, we can extract a subsequence, still denoted v_n^- , such that v_n^- converges strongly in $L^2(Y^-)$. We denote by l its limit.

Besides, ∇v_n^- converges strongly to 0 in $L^2(Y^-)$. We thus have $\nabla l = 0$ and l constant in Y^- . By applying in Y^+ the trace theorem and Poincaré–Wirtinger inequality to v_n^+ , one also gets

that

$$\|v_n^-\|_{L^2(\Gamma)} \le \|v_n^+ - v_n^-\|_{L^2(\Gamma)} + \|v_n^+\|_{L^2(\Gamma)} \le \|v_n^+ - v_n^-\|_{L^2(\Gamma)} + C\|v_n^+\|_{H^1(Y^+)} \le \frac{C'}{n}.$$

Consequently, v_n^- converges strongly to 0 in $L^2(\Gamma)$ and l = 0 on Γ . We have then l = 0 in Υ^- , which leads to a contradiction. This proves (4.3). We can now find an upper bound to (4.2):

$$\|v^-\|_{L^2(\Omega_{\varepsilon}^-)}^2 \leq \varepsilon^2 C_1 \sum_{n \in N_{\varepsilon}} \int_{Y_n^+} |\nabla v_{\varepsilon}^+(y)|^2 dy + \int_{Y_n^-} |\nabla v_{\varepsilon}^-(y)|^2 dy + \int_{\Gamma_n} |v_{\varepsilon}^+(y) - v_{\varepsilon}^-(y)|^2 ds(y).$$

After the change of variable $x = \varepsilon y$, one gets

$$\|v^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2} \leq \varepsilon C_{1}\left(\|\nabla v^{+}\|_{L^{2}(\Omega_{\varepsilon}^{+})}^{2} + \|\nabla v^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2} + \varepsilon \|v^{+} - v^{-}\|_{L^{2}(\Gamma_{\varepsilon})}^{2}\right).$$

Since $\varepsilon < 1$, there exists a constant C_2 , which does not depend on ε such that for every $v \in W_{\varepsilon}$,

$$\|v^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})} \leq C_{2}\|v\|_{W_{\varepsilon}}$$

which completes the proof.

Theorem 4.2. u_{ε}^{-} is uniformly bounded in ε in $H^{1}(\Omega_{\varepsilon}^{-})$, *i.e.*, there exists a constant C independent of ε , such that

 $\|u_{\varepsilon}^{-}\|_{H^{1}(\Omega_{\varepsilon}^{-})} \leq C.$

Proof. By definition of the norm on W_{ε} , $\|\nabla u_{\varepsilon}^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2} \leq \|u_{\varepsilon}\|_{W_{\varepsilon}}^{2}$.

We thus have with the result of Lemma 4.4:

$$\|u_{\varepsilon}^{-}\|_{H^{1}(\Omega_{\varepsilon}^{-})}^{2} \leq C_{1}\|u_{\varepsilon}\|_{W_{\varepsilon}'}^{2}$$
(4.4)

with a constant C_1 which does not depend on ε .

Furthermore, using the result of Theorem 4.1, there exists a constant C_2 independent of ε such that

$$|a(u_{\varepsilon},u_{\varepsilon})|\leq C_2.$$

We use the coercivity of *a* and get a uniform bound in ε of u_{ε} in W_{ε} . This bound and (4.4) complete the proof.

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4.2.2 Two-scale convergence

We first recall the definition of two-scale convergence and a few results of this theory [2, 48].

Definition 4.1. A sequence of functions u_{ε} in $L^2(\Omega)$ is said to two-scale converge to a limit u_0 belonging to $L^2(\Omega \times Y)$ if, for any function ψ in $L^2(\Omega, C_{\sharp}(Y))$, we have

$$\lim_{\varepsilon \to 0} \int_{\Omega} u_{\varepsilon}(x) \psi(x, \frac{x}{\varepsilon}) dx = \int_{\Omega} \int_{Y} u_0(x, y) \psi(x, y) dx dy$$

This notion of two-scale convergence makes sense because of the next compactness theorem.

Theorem 4.3. From each bounded sequence u_{ε} in $L^2(\Omega)$, we can extract a subsequence, and there exists a limit $u_0 \in L^2(\Omega \times Y)$ such that this subsequence two-scale converges to u_0 .

Two-scale convergence can be extended to sequences defined on periodic surfaces.

Proposition 4.1. For any sequence u_{ε} in $L^{2}(\Gamma_{\varepsilon})$ such that

$$\varepsilon \int_{\Gamma_{\varepsilon}} |u_{\varepsilon}|^2 dx \le C, \tag{4.5}$$

there exists a subsequence, still denoted u_{ε} , and a limit function $u_0 \in L^2(\Omega, L^2(\Gamma))$ such that u_{ε} two-scale converges to u_0 in the sense

$$\lim_{\varepsilon \to 0} \varepsilon \int_{\Gamma_{\varepsilon}} u_{\varepsilon}(x) \psi(x, \frac{x}{\varepsilon}) ds(x) = \int_{\Omega} \int_{\Gamma} u_0(x, y) \psi(x, y) dx ds(y),$$

for any function $\psi \in L^2(\Omega, C_{\sharp}(Y))$.

Remark 4.1. If u_{ε} and ∇u_{ε} are bounded in $L^{2}(\Omega)$, one can prove by using for example [4, Lemma 2.4.9] that u_{ε} verifies the uniform bound (4.5). The two-scale limit on the surface is then the trace on Γ of the two-scale limit of u_{ε} in Ω .

In order to prove item (i) in Theorem 2.1, we need the following results.

Lemma 4.5. Let the functions $(u_{\varepsilon})_{\varepsilon}$ be the sequence of solutions of (2.2). There exist functions $u(x) \in H^1(\Omega)$, $v^+(x,y) \in L^2(\Omega, H^1_{\sharp}(Y^+))$ and $v^-(x,y) \in L^2(\Omega, H^1_{\sharp}(Y^-))$ such that, up to a subsequence,

$$\begin{pmatrix} u_{\varepsilon} \\ \chi_{\varepsilon}^{+}(\frac{x}{\varepsilon})\nabla u_{\varepsilon}^{+} \\ \chi_{\varepsilon}^{-}(\frac{x}{\varepsilon})\nabla u_{\varepsilon}^{-} \end{pmatrix} \quad two-scale \ converge \ to \quad \begin{pmatrix} u(x) \\ \chi_{Y^{+}}(y) \left(\nabla u(x) + \nabla_{y}v^{+}(x,y)\right) \\ \chi_{Y^{-}}(y) \left(\nabla u(x) + \nabla_{y}v^{-}(x,y)\right) \end{pmatrix}$$

Proof. We denote by $\tilde{\cdot}$ the extension by zero of functions on Ω_{ε}^+ and Ω_{ε}^- in the respective domains Ω_{ε}^- and Ω_{ε}^+ .

From the previous estimates, $\tilde{u}_{\varepsilon}^{\pm}$ and $\widetilde{\nabla u}_{\varepsilon}^{\pm}$ are bounded sequences in $L^{2}(\Omega)$. Up to a subsequence, they two-scale converge to $\tau^{\pm}(x, y)$ and $\xi^{\pm}(x, y)$. Since $\tilde{u}_{\varepsilon}^{\pm}$ and $\widetilde{\nabla u}_{\varepsilon}^{\pm}$ vanish in $\Omega_{\varepsilon}^{\mp}$, so do τ^{\pm} and ξ^{\pm} .

Consider $\varphi \in \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y))^2$ such that $\varphi = 0$ for $y \in \overline{Y^-}$. By integrating by parts, it follows that

$$\varepsilon \int_{\Omega_{\varepsilon}^{+}} \nabla u_{\varepsilon}^{+}(x) \cdot \overline{\varphi}(x, \frac{x}{\varepsilon}) dx = -\int_{\Omega_{\varepsilon}^{+}} u_{\varepsilon}^{+}(x) \left(\operatorname{div}_{y} \overline{\varphi}(x, \frac{x}{\varepsilon}) + \varepsilon \operatorname{div}_{x} \overline{\varphi}(x, \frac{x}{\varepsilon}) \right) dx.$$

We take the limit of this equality as $\varepsilon \to 0$:

$$0 = -\int_{\Omega} \int_{Y^+} \tau^+(x,y) \operatorname{div}_y \overline{\varphi}(x,y) dx dy.$$

Therefore, τ^+ does not depend on y in Y^+ , *i.e.*, there exists a function $u^+ \in L^2(\Omega)$ such that $\tau^+(x, y) = \chi_{Y^+}(y)u^+(x)$ for all $(x, y) \in \Omega \times Y$.

Take now $\varphi \in \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y))^2$ such that $\varphi = 0$ for $y \in \overline{Y^-}$ and $\operatorname{div}_y \varphi = 0$. Similarly, we have

$$\int_{\Omega_{\varepsilon}^{+}} \nabla u_{\varepsilon}^{+}(x) \cdot \overline{\varphi}(x, \frac{x}{\varepsilon}) dx = -\int_{\Omega_{\varepsilon}^{+}} u_{\varepsilon}^{+}(x) \operatorname{div}_{x} \overline{\varphi}(x, \frac{x}{\varepsilon}) dx,$$

and thus

$$\int_{\Omega} \int_{Y^+} \xi^+(x,y) \cdot \overline{\varphi}(x,y) dx dy = -\int_{\Omega} \int_{Y^+} u^+(x) \operatorname{div}_x \overline{\varphi}(x,y) dx dy.$$
(4.6)

For φ independent of y, this implies that $u^+ \in H^1(\Omega)$. Furthermore, if we integrate by parts the right-hand side of (4.6), we get

$$\int_{\Omega} \int_{Y^+} \xi^+(x,y) \cdot \overline{\varphi}(x,y) dx dy = \int_{\Omega} \int_{Y^+} \nabla u^+(x) \cdot \overline{\varphi}(x,y) dx dy,$$

for all $\varphi \in \mathcal{D}(\Omega, C^{\infty}_{\#}(\Upsilon^{+}))^{2}$ such that $\operatorname{div}_{y} \varphi = 0$ and $\varphi(x, y) \cdot n(y) = 0$ for y on Γ .

Since the orthogonal of the divergence-free functions are exactly the gradients, there exists a function $v^+ \in L^2(\Omega, H^1_{\sharp}(Y^+))$ such that

$$\xi^+(x,y) = \chi_{Y^+}(y) \left(\nabla u^+(x) + \nabla_y v^+(x,y) \right),$$

for all $(x, y) \in \Omega \times Y$.

Likewise, there exist functions $u^- \in H^1(\Omega)$ and $v^- \in L^2(\Omega, H^1_{\sharp}(Y^-))$ such that

$$\tau^{-}(x,y) = \chi_{Y^{-}}(y)u^{-}(x), \text{ and } \xi^{-}(x,y) = \chi_{Y^{-}}(y)\left(\nabla u^{-}(x) + \nabla_{y}v^{-}(x,y)\right),$$

for all $(x, y) \in \Omega \times Y$.

Furthermore, thanks to Remark 4.1, we have also

$$\varepsilon \int_{\Gamma_{\varepsilon}} u_{\varepsilon}^{\pm}(x) \overline{\varphi}(x, \frac{x}{\varepsilon}) dx \xrightarrow[\varepsilon \to 0]{} \int_{\Omega} \int_{\Gamma} u^{\pm}(x, y) \overline{\varphi}(x, y) dx dy,$$

for all $\varphi \in L^2(\Omega, C^{\infty}_{\sharp}(\Gamma))$.

Recall that u_{ε} is a solution to the following variational form:

$$\begin{split} \int_{\Omega_{\varepsilon}^{+}} k_{0} \nabla u_{\varepsilon}^{+}(x) \cdot \nabla \overline{\varphi}_{\varepsilon}^{+}(x) dx + \int_{\Omega_{\varepsilon}^{-}} k_{0} \nabla u_{\varepsilon}^{-}(x) \cdot \nabla \overline{\varphi}_{\varepsilon}^{-}(x) dx \\ &+ \frac{1}{\varepsilon \beta} \int_{\Gamma_{\varepsilon}} \left(u_{\varepsilon}^{+} - u_{\varepsilon}^{-} \right) \left(\overline{\varphi}_{\varepsilon}^{+} - \overline{\varphi}_{\varepsilon}^{-} \right) ds - k_{0} \int_{\partial \Omega} g \overline{\varphi}_{\varepsilon}^{+} ds = 0, \end{split}$$

for all $(\varphi_{\varepsilon}^+, \varphi_{\varepsilon}^-) \in (H^1(\Omega_{\varepsilon}^+), H^1(\Omega_{\varepsilon}^-)).$

We multiply this equality by ε^2 and take the limit when ε goes to 0. The first two terms disappear and we obtain, for all $(\varphi^+, \varphi^-) \in \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^-))$:

$$\frac{1}{\beta}\int_{\Omega}\int_{\Gamma}(u^{+}(x)-u^{-}(x))(\overline{\varphi}^{+}(x,y)-\overline{\varphi}^{-}(x,y))dxdy=0.$$

Thus $u^+(x) = u^-(x)$ for all $x \in \Omega$, and u_{ε} two-scale converges to $u = u^+ = u^- \in H^1(\Omega)$. This completes the proof.

Now, we are ready to prove Theorem 2.1. For this, we need to show that $u = u_0$, that $v^+ - u_1^+$ is constant, and that $v^- - u_1^-$ is constant, where u_1^\pm is defined in (2.7). The uniqueness of a solution for the homogenized problem and the cell problems will then allow us to conclude the convergence, not only up to a subsequence.

Proof. We first want to retrieve the expression of u_1 as a test function of the derivatives of u_0 and the cell problem solutions w_i .

We choose in the variational formulation (3.5) a function φ_{ε} of the form

$$\varphi_{\varepsilon}(x) = \varepsilon \varphi_1(x, \frac{x}{\varepsilon}),$$

where $\varphi_1 \in \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^-)).$

Lemma 4.5 shows the two-scale convergence of the following three terms:

$$\begin{split} &\int_{\Omega_{\varepsilon}^{+}} k_{0} \nabla u_{\varepsilon}^{+}(x) \cdot \nabla \overline{\varphi}_{\varepsilon}^{+}(x) dx \quad \xrightarrow{\epsilon \to 0} \quad \int_{\Omega} \int_{Y^{+}} k_{0} \left(\nabla u(x) + \nabla_{y} v^{+}(x,y) \right) \cdot \nabla_{y} \overline{\varphi}_{1}^{+}(x,y) dx dy \\ &\int_{\Omega_{\varepsilon}^{-}} k_{0} \nabla u_{\varepsilon}^{-}(x) \cdot \nabla \overline{\varphi}_{\varepsilon}^{-}(x) dx \quad \xrightarrow{\epsilon \to 0} \quad \int_{\Omega} \int_{Y^{-}} k_{0} \left(\nabla u(x) + \nabla_{y} v^{-}(x,y) \right) \cdot \nabla_{y} \overline{\varphi}_{1}^{-}(x,y) dx dy \\ &\int_{\partial \Omega} g(x) \overline{\varphi}_{\varepsilon}^{+}(x) ds(x) \quad \xrightarrow{\epsilon \to 0} \qquad 0. \end{split}$$

We can not take directly the limit as $\varepsilon \rightarrow 0$ in the last term:

$$\begin{split} \frac{1}{\varepsilon\beta} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+}(x) - u_{\varepsilon}^{-}(x))(\overline{\varphi}_{\varepsilon}^{+}(x) - \overline{\varphi}_{\varepsilon}^{-}(x))ds(x) \\ &= \frac{1}{\beta} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+}(x) - u_{\varepsilon}^{-}(x))\left(\overline{\varphi}_{1}^{+}(x,\frac{x}{\varepsilon}) - \overline{\varphi}_{1}^{-}(x,\frac{x}{\varepsilon})\right)ds(x). \end{split}$$

Lemma D.1 ensures the existence of a function $\theta \in (\mathcal{D}(\Omega, H^1_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, H^1_{\sharp}(Y^-)))^2$ such that for all $\psi \in H^1_{\sharp}(Y^+)/\mathbb{C} \times H^1_{\sharp}(Y^-)$:

$$\int_{Y^{+}} \nabla \psi^{+}(y) \cdot \overline{\theta}^{+}(x,y) dy + \int_{Y^{-}} \nabla \psi^{-}(y) \cdot \overline{\theta}^{-}(x,y) dy + \int_{\Gamma} \left(\psi^{+}(y) - \psi^{-}(y)\right) \left(\overline{\varphi}_{1}^{+}(x,y) - \overline{\varphi}_{1}^{-}(x,y)\right) ds(y) = 0.$$
(4.7)

We make the change of variables $y = \frac{x}{\varepsilon}$, sum over all $(Y_{\varepsilon,n})_{n \in N_{\varepsilon}}$, and choose $\psi = u_{\varepsilon}$ to get

$$\begin{split} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+}(x) - u_{\varepsilon}^{-}(x)) \left(\overline{\varphi}_{1}^{+}(x, \frac{x}{\varepsilon}) - \overline{\varphi}_{1}^{-}(x, \frac{x}{\varepsilon}) \right) ds(x) = \\ &- \int_{\Omega_{\varepsilon}^{+}} \nabla u_{\varepsilon}^{+}(x, \frac{x}{\varepsilon}) \cdot \theta^{+}(x, \frac{x}{\varepsilon}) dx - \int_{\Omega_{\varepsilon}^{-}} \nabla u_{\varepsilon}^{-}(x, \frac{x}{\varepsilon}) \cdot \theta^{-}(x, \frac{x}{\varepsilon}) dx. \end{split}$$

We can now take the limit as ε goes to 0:

$$\lim_{\varepsilon \to 0} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+}(x) - u_{\varepsilon}^{-}(x)) \left(\overline{\varphi}_{1}^{+}(x, \frac{x}{\varepsilon}) - \overline{\varphi}_{1}^{-}(x, \frac{x}{\varepsilon})\right) ds(x) = -\int_{Y^{+}} \left(\nabla u(x) + \nabla_{y} v^{+}(x, y)\right) \cdot \overline{\theta}^{+}(x, y) dx dy - \int_{Y^{-}} \left(\nabla u(x) + \nabla_{y} v^{-}(x, y)\right) \cdot \overline{\theta}^{-}(x, y) dx dy.$$

Finally, the variational formula (4.7) gives us

$$\begin{split} \lim_{\varepsilon \to 0} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+}(x) - u_{\varepsilon}^{-}(x)) \left(\overline{\varphi}_{1}^{+}(x, \frac{x}{\varepsilon}) - \overline{\varphi}_{1}^{-}(x, \frac{x}{\varepsilon}) \right) ds(x) &= \\ \int_{\Omega} \int_{\Gamma} \left(v^{+}(y) - v^{-}(y) \right) \left(\overline{\varphi}_{1}^{+}(x, y) - \overline{\varphi}_{1}^{-}(x, y) \right) ds(y). \end{split}$$

For $\varphi_{\varepsilon}(x) = \varepsilon \varphi_1(x, \frac{x}{\varepsilon})$, with $\varphi_1 \in \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^-))$, the two-scale limit of the variational formula is

$$\begin{split} &\int_{\Omega} \int_{Y^+} k_0 \left(\nabla u(x) + \nabla_y v^+(x,y) \right) \cdot \nabla_y \overline{\varphi}_1^+(x,y) dx dy \\ &+ \int_{\Omega} \int_{Y^-} k_0 \left(\nabla u(x) + \nabla_y v^-(x,y) \right) \cdot \nabla_y \overline{\varphi}_1^-(x,y) dx dy \\ &+ \frac{1}{\beta} \int_{\Omega} \int_{\Gamma} \left(v^+(y) - v^-(y) \right) \left(\overline{\varphi}_1^+(x,y) - \overline{\varphi}_1^-(x,y) \right) ds(y) = 0. \end{split}$$

By density, this formula holds true for $\varphi_1 \in L^2(\Omega, H^1_{\sharp}(Y^+)) \times L^2(\Omega, H^1_{\sharp}(Y^-))$. One can recognize the formula verified by u_1^{\pm} as the definition of the cell problems. Hence, separation of variables and uniqueness of the solutions of the cell problems in *W* give

$$v^{-}(x,y) = u_{1}^{-} = \sum_{i=1,2} \frac{\partial u_{0}}{\partial x_{i}}(x)w_{i}^{-}(y)$$

and, up to a constant:

$$v^+(x,y) = u_1^+ = \sum_{i=1,2} \frac{\partial u_0}{\partial x_i}(x) w_i^+(y).$$

We now choose in the variational formula verified by u_{ε} a test function $\varphi_{\varepsilon}(x) = \varphi(x)$, with $\varphi \in C^{\infty}(\overline{\Omega})$.

The limit of (3.5) as ε goes to 0 is then given by

$$\begin{split} &\int_{\Omega} \int_{Y^{+}} k_{0} \left(\nabla u(x) + \nabla_{y} v^{+}(x, y) \right) \cdot \nabla \overline{\varphi}(x) dx dy \\ &+ \int_{\Omega} \int_{Y^{-}} k_{0} \left(\nabla u(x) + \nabla_{y} v^{-}(x, y) \right) \cdot \nabla \overline{\varphi}(x) dx dy \\ &+ \int_{\partial \Omega} g(x) \overline{\phi}(x) ds(x) = 0. \end{split}$$

By density, this formula hold true for $\varphi \in H^1(\Omega)$, which leads exactly to the variational formula of the homogenized problem (2.4). Since the solution of this problem is unique in $H^1_{\mathbb{C}}(\Omega)$, u_{ε} converges to u_0 , not only up to a subsequence. Likewise, ∇u_{ε} two-scale converges to $\nabla u_0 + \chi_{Y^+} \nabla_y u_1^+ + \chi_{Y^-} \nabla_y u_1^-$.

5 Effective admittivity for a dilute suspension

In general, the effective admittivity given by formula (2.5) can not be computed exactly except for a few configurations. In this section, we consider the problem of determining the effective property of a suspension of cells when the volume fraction $|Y^-|$ goes to zero. In other words, the cells have much less volume than the medium surrounding them. This kind of suspension is called dilute. Many approximations for the effective properties of composites are based on the solution for dilute suspension.

5.1 Computation of the effective admittivity

We investigate the periodic double-layer potential used in calculating effective permittivity of a suspension of cells. We introduce the periodic Green function G_{\sharp} , for the Laplace equation in Y, given by

$$orall x \in Y, \qquad G_{\sharp}(x) = \sum_{n \in \mathbb{Z}^2 \setminus \{0\}} rac{e^{i2\pi n \cdot x}}{4\pi^2 |n|^2}.$$

The following lemma from [11, 9] plays an essential role in deriving the effective properties of a suspension in the dilute limit.

Lemma 5.1. The periodic Green function G_{\sharp} admits the following decomposition:

$$\forall x \in Y, \qquad G_{\sharp}(x) = \frac{1}{2\pi} \ln |x| + R_2(x),$$
 (5.1)

where R_2 is a smooth function with the following Taylor expansion at 0:

$$R_2(x) = R_2(0) - \frac{1}{4}(x_1^2 + x_2^2) + O(|x|^4).$$
(5.2)

Let $L_0^2(\Gamma) := \left\{ \varphi \in L^2(\Gamma) \Big| \int_{\Gamma} \varphi(x) ds(x) = 0 \right\}.$

We define the periodic double-layer potential \mathcal{D}_{Γ} of the density function $\varphi \in L^2_0(\Gamma)$:

$$\widetilde{\mathcal{D}}_{\Gamma}[\varphi](x) = \int_{\Gamma} \frac{\partial}{\partial n_y} G_{\sharp}(x-y) \varphi(y) ds(y).$$

The double-layer potential has the following properties [9].

Lemma 5.2. Let $\varphi \in L^2_0(\Gamma)$. $\widetilde{\mathcal{D}}_{\Gamma}[\varphi]$ verifies:

(i)
$$\Delta \widetilde{\mathcal{D}}_{\Gamma}[\varphi] = 0$$
 in Y^+ ,
 $\Delta \widetilde{\mathcal{D}}_{\Gamma}[\varphi] = 0$ in Y^- ,

(*ii*)
$$\frac{\partial}{\partial n} \widetilde{\mathcal{D}}_{\Gamma}[\varphi] \Big|_{+} = \frac{\partial}{\partial n} \widetilde{\mathcal{D}}_{\Gamma}[\varphi] \Big|_{-}$$
 on Γ ,
(*iii*) $\widetilde{\mathcal{D}}_{\Gamma}[\varphi] \Big|_{\pm} = \left(\mp \frac{1}{2}I + \widetilde{\mathcal{K}}_{\Gamma} \right) [\varphi]$ on Γ ,

where $\widetilde{\mathcal{K}}_{\Gamma}: L^2_0(\Gamma) \mapsto L^2_0(\Gamma)$ is the Neumann–Poincaré operator defined by

$$orall x \in \Gamma$$
, $\widetilde{\mathcal{K}}_{\Gamma}[arphi](x) = \int_{\Gamma} rac{\partial}{\partial n_y} G_{\sharp}(x-y) arphi(y) ds(y).$

The following integral representation formula holds.

Theorem 5.1. Let w_i be the unique solution in W of (2.6) for i = 1, 2. w_i admits the following integral representation in Y:

$$w_{i} = -\beta k_{0} \widetilde{\mathcal{D}}_{\Gamma} \left(I + \beta k_{0} \widetilde{\mathcal{L}}_{\Gamma} \right)^{-1} [n_{i}], \qquad (5.3)$$

where $\widetilde{\mathcal{L}}_{\Gamma} = \frac{\partial \widetilde{D}_{\Gamma}}{\partial n}$ and $n = (n_i)_{i=1,2}$ is the outward unit normal to Γ .

Proof. Let $\varphi := -\beta k_0 \left(I + \beta k_0 \widetilde{\mathcal{L}}\right)^{-1} [n_i]. \ \varphi$ verifies :

$$\int_{\Gamma} \varphi(y) ds(y) = -\beta k_0 \int_{\Gamma} \frac{\partial}{\partial n} (\widetilde{\mathcal{D}}_{\Gamma}[\varphi](y) + y_i) ds(y) = 0.$$

The first equality comes from the definition of φ and the second from an integration by parts and the fact that $\widetilde{\mathcal{D}}_{\Gamma}[\varphi]$ and *I* are harmonic. Consequently, $\varphi \in L^2_0(\Gamma)$.

We now introduce $V_i := \widetilde{\mathcal{D}}_{\Gamma}[\varphi]$. V_i is solution to the following problem:

$$\begin{cases} \nabla \cdot k_0 \nabla V_i = 0 & \text{in } Y^+, \\ \nabla \cdot k_0 \nabla V_i = 0 & \text{in } Y^-, \\ k_0 \frac{\partial V_i}{\partial n}\Big|_+ = k_0 \frac{\partial V_i}{\partial n}\Big|_- & \text{on } \Gamma, \\ V_i|_+ - V_i|_- = \varphi & \text{on } \Gamma, \\ y \longmapsto V_i(y) \text{ } Y\text{-periodic.} \end{cases}$$

We use the definitions of φ and V_i and recognize that the last problem is exactly problem (2.6). The uniqueness of the solution in *W* gives us the wanted result.

From Theorem 2.1, the effective admittivity of the medium K^* is given by

$$\forall (i,j) \in \{1,2\}^2, \qquad K_{i,j}^* = k_0 \left(\delta_{ij} + \int_Y \nabla w_i \cdot e_j\right).$$

After an integration by parts, we get

$$\forall (i,j) \in \{1,2\}^2, \qquad K_{i,j}^* = k_0 \left(\delta_{ij} + \int_{\partial Y} w_i(y) n_j(y) \, ds(y) - \int_{\Gamma} \left(w_i^+ - w_i^- \right) n_j(y) \, ds(y) \right).$$

Because of the *Y*-periodicity of w_i , we have: $\int_{\partial Y} w_i(y) n_j ds(y) = 0$. Finally, the integral representation 5.3 gives us that

$$\forall (i,j) \in \{1,2\}^2, \qquad K_{i,j}^* = k_0 \left(\delta_{ij} - (\beta k_0) \int_{\Gamma} \left(I + \beta k_0 \widetilde{\mathcal{L}}_{\Gamma} \right)^{-1} [n_i] n_j \, ds(y) \right).$$

We consider that we are in the context of a dilute suspension, *i.e.*, the size of the cell is small compared to the square: $|Y^-| \ll |Y| = 1$. We perform the change of variable: $z = \rho^{-1}y$ with $\rho = |Y^-|^{\frac{1}{2}}$ and obtain that

$$\forall (i,j) \in \{1,2\}^2, \qquad K_{i,j}^* = k_0 \left(\delta_{ij} - \rho^2(\beta k_0) \int_{\rho^{-1}\Gamma} \left(I + \rho\beta k_0 \widetilde{\mathcal{L}}_{\Gamma}\right)^{-1} [n_i](\rho z) n_j(z) \, ds(z)\right),$$

where *n* is the outward unit normal to Γ . Note that, in the same way as before, β becomes $\rho\beta$ when we rescale the cell.

Let us introduce $\varphi_i = -(I + \rho \beta k_0 \tilde{\mathcal{L}}_{\Gamma})^{-1}[n_i]$ and $\psi_i(z) = \varphi_i(\rho z)$ for all $z \in \rho^{-1}\Gamma$. From (5.1), we get, for any $z \in \rho^{-1}\Gamma$, after changes of variable in the integrals:

$$\widetilde{\mathcal{L}}_{\Gamma}[\varphi_i](\rho z) = \frac{\partial}{\partial n} \widetilde{\mathcal{D}}_{\Gamma}[\varphi_i](\rho z) = \rho^{-1} \frac{\partial}{\partial n} \mathcal{D}_{\rho^{-1}\Gamma}[\psi_i](z) + \frac{\partial}{\partial n(z)} \int_{\rho^{-1}\Gamma} \frac{\partial}{\partial n(y)} R_2(\rho z - \rho y) \varphi(\rho y) ds(y).$$

Besides, the expansion (5.2) gives us that the estimate

$$\nabla R_2(\rho(z-y)) \cdot n(y) = -\frac{\rho}{2}(z-y) \cdot n(y) + O(\rho^3),$$

holds uniformly in $z, y \in \rho^{-1}\Gamma$.

We thus get the following expansion:

$$\widetilde{\mathcal{L}}_{\Gamma}[\varphi_i](\rho z) = \rho^{-1} \mathcal{L}_{\rho^{-1}\Gamma}[\psi_i](z) - \frac{\rho}{2} \sum_{j=1,2} n_j \int_{\rho^{-1}\Gamma} n_j \psi_i(y) ds(y) + O(\rho^4).$$

Using ψ_i^* defined by (2.11) we get on $\rho^{-1}\Gamma$:

$$\psi_i = \psi_i^* + \beta k_0 \frac{\rho^2}{2} \sum_{j=1,2} \psi_j^* \int_{\rho^{-1}\Gamma} n_j(y) \psi_i(y) ds(y) + O(\rho^4).$$
(5.4)

By iterating the formula (5.4), we obtain on $\rho^{-1}\Gamma$ that

$$\psi_i = \psi_i^* + \beta k_0 \frac{\rho^2}{2} \sum_{j=1,2} \psi_j^* \int_{\rho^{-1}\Gamma} n_j(y) \psi_i^*(y) ds(y) + O(\rho^4).$$

Therefore, one can easily see that Theorem 2.2 holds.

5.2 Case of concentric circular-shaped cells: the Maxwell-Wagner-Fricke formula

We consider in this section that the cells are disks of radius r_0 . $\rho^{-1}\Gamma$ becomes a circle of radius r_0 .

For all $g \in L^2((0, 2\pi))$, we introduce the Fourier coefficients:

$$\forall m \in \mathbb{Z}, \quad \hat{g}(m) = \frac{1}{2\pi} \int_0^{2\pi} g(\varphi) e^{-im\varphi} d\varphi,$$

and have then for all $\varphi \in (0, 2\pi)$:

$$g(\varphi) = \sum_{m=-\infty}^{\infty} \hat{g}(m) e^{im\varphi}.$$

For $f \in C^{2,\eta}(\rho^{-1}\Gamma)$, we obtain after a few computations:

$$\forall \theta \in]0, 2\pi[, (I + \beta k_0 \mathcal{L}_{\rho^{-1}\Gamma})^{-1}[f](\theta) = \sum_{m \in \mathbb{Z}^*} \left(1 + \beta k_0 \frac{|m|}{2r_0}\right)^{-1} \hat{f}(m) e^{im\theta}.$$

For $p = 1, 2, \psi_p^* = -(I + \beta k_0 \mathcal{L}_{\rho^{-1}\Gamma})^{-1}[n_p]$ then have the following expression:

$$\forall \theta \in (0, 2\pi), \quad \psi_p^* = -\left(1 + \frac{\beta k_0}{2r_0}\right)^{-1} n_p.$$

Consequently, we get for $(p,q) \in \{1,2\}^2$:

$$M_{p,q}=-\delta_{pq}rac{eta k_0 \pi r_0}{1+rac{eta k_0}{2r_0}},$$

and hence,

$$\Im M_{p,q} = \delta_{p,q} \frac{\pi r_0 \delta \omega (\epsilon_m \sigma_0 - \epsilon_0 \sigma_m)}{(\sigma_m + \sigma_0 \frac{\delta}{2r_0})^2 + \omega^2 (\epsilon_m + \epsilon_0 \frac{\delta}{2r_0})^2}.$$
(5.5)

Formula (5.5) is the two-dimensional version of the Maxwell-Wagner-Fricke formula, which gives the effective admittivity of a dilute suspension of spherical cells covered by a thin membrane.

An explicit formula for the case of elliptic cells can be derived by using the spectrum of the integral operator $\mathcal{L}_{\rho^{-1}\Gamma}$, which can be identified by standard Fourier methods [36].

5.3 Debye relaxation times

From (5.5), it follows that the imaginary part of the membrane polarization attains its maximum with respect to the frequency at \int_{1}^{2}

$$rac{1}{ au} = rac{\sigma_m + \sigma_0 \, rac{o}{2r_0}}{arepsilon_m + arepsilon_0 \, rac{\delta}{2r_0}}.$$

This dispersion phenomenon due to the membrane polarization is well known and referred to as the β -dispersion. The associated characteristic time τ corresponds to a Debye relaxation time.

For arbitrary-shaped cells, we define the first and second Debye relaxation times, τ_i , i = 1, 2, by

$$\frac{1}{\tau_i} := \arg \max_{\omega} |\lambda_i(\omega)|, \tag{5.6}$$

where $\lambda_1 \leq \lambda_2$ are the eigenvalues of the imaginary part of the membrane polarization tensor $M(\omega)$. Note that if the cell is of circular shape, $\lambda_1 = \lambda_2$.

As it will be shown later, the Debye relaxation times can be used for identifying the microstructure.

5.4 Properties of the membrane polarization tensor and the Debye relaxation times

In this subsection, we derive important properties of the membrane polarization tensor and the Debye relaxation times defined respectively by (2.10) and (5.6). In particular, we prove that the Debye relaxation times are invariant with respect to translation, scaling, and rotation of the cell.

First, since the kernel of $\mathcal{L}_{\rho^{-1}\Gamma}$ is invariant with respect to translation, it follows that $M(C, \beta k_0)$ is invariant with respect to translation of the cell *C*.

Next, from the scaling properties of the kernel of $\mathcal{L}_{\rho^{-1}\Gamma}$ we have

$$M(sC,\beta k_0) = s^2 M(C,\frac{\beta k_0}{s})$$

for any scaling parameter s > 0.

Finally, we have

$$M(\mathcal{R}C,\beta k_0) = \mathcal{R}M(C,\beta k_0)\mathcal{R}^t$$
 for any rotation \mathcal{R} ,

where *t* denotes the transpose.

Therefore, the Debye relaxation times are translation and rotation invariant. Moreover, for scaling, we have

$$au_i(hC,eta k_0)= au_i(C,rac{eta k_0}{h}), \qquad i=1,2, \quad h>0.$$

Since β is proportional to the thickness of the cell membrane, β/h is nothing else than the real rescaled coefficient β for the cell *C*. The Debye relaxation times (τ_i) are therefore invariant by scaling.

Since $\mathcal{L}_{\rho^{-1}\Gamma}$ is self-adjoint, it follows that *M* is symmetric. Finally, we show positivity of the imaginary part of the matrix *M* for δ small enough.

We consider that the cell contour Γ can be parametrized by polar coordinates. We have, up to $O(\delta^3)$,

$$M + \beta \rho^{-1} |\Gamma| = -\beta^2 \int_{\rho^{-1}\Gamma} n \mathcal{L}_{\rho^{-1}\Gamma}[n] \, ds, \qquad (5.7)$$

where again we have assumed that $\sigma_0 = 1$ and $\epsilon_0 = 0$.

Recall that

$$\beta = \frac{\delta \sigma_m}{\sigma_m^2 + \omega^2 \varepsilon_m^2} - i \frac{\delta \omega \varepsilon_m}{\sigma_m^2 + \omega^2 \varepsilon_m^2}.$$

Hence, the positivity of $\mathcal{L}_{\rho^{-1}\Gamma}$ yields

$$\Im M \geq rac{\delta \omega arepsilon_m}{2
ho (\sigma_m^2 + \omega^2 arepsilon_m^2)} |\Gamma| I$$

for δ small enough, where *I* is the identity matrix.

Finally, by using (5.7) one can see that the eigenvalues of $\Im M$ have one maximum each with respect to the frequency. Let l_i , i = 1, 2, $l_1 \ge l_2$, be the eigenvalues of $\int_{\rho^{-1}\Gamma} n \mathcal{L}_{\rho^{-1}\Gamma}[n] ds$. We have

$$\lambda_{i} = \frac{\delta\omega\varepsilon_{m}}{\rho(\sigma_{m}^{2} + \omega^{2}\varepsilon_{m}^{2})}|\Gamma| - \frac{2\delta^{2}\omega\varepsilon_{m}\sigma_{m}}{(\sigma_{m}^{2} + \omega^{2}\varepsilon_{m}^{2})^{2}}l_{i}, \quad i = 1, 2.$$
(5.8)

Therefore, τ_i is the inverse of the positive root of the following polynomial in ω :

$$-\varepsilon_m^4 |\Gamma| \omega^4 + 6\delta \varepsilon_m^2 \sigma_m l_i \rho \omega^2 + \sigma_m^4 |\Gamma|$$

5.5 Anisotropy measure

Anisotropic electrical properties can be found in biological tissues such as muscles and nerves. In this subsection, based on formula (2.9), we introduce a natural measure of the conductivity anisotropy and derive its dependence on the frequency of applied current. Assessment of electrical anisotropy of muscle may have useful clinical application. Because neuromuscular diseases produce substantial pathological changes, the anisotropic pattern of the muscle is likely to be highly disturbed [21, 32]. Neuromuscular diseases could lead to a reduction in anisotropy for a range of frequencies as the muscle fibers are replaced by isotropic tissue.

Let $\lambda_1 \leq \lambda_2$ be the eigenvalues of the imaginary part of the membrane polarization tensor $M(\omega)$. The function

$$\omega \mapsto \frac{\lambda_1(\omega)}{\lambda_2(\omega)}$$

can be used as a measure of the anisotropy of the conductivity of a dilute suspension. Assume $\epsilon_0 = 0$. As frequency ω increases, the factor βk_0 decreases. Therefore, for large ω , using the expansions in (5.8) we obtain that

$$\frac{\lambda_1(\omega)}{\lambda_2(\omega)} = 1 + (l_1 - l_2) \frac{2\delta\sigma_m \rho}{(\sigma_m^2 + \omega^2 \varepsilon_m^2)|\Gamma|} + O(\delta^2),$$
(5.9)

where $l_1 \leq l_2$ are the eigenvalues of $\int_{\rho^{-1}\Gamma} n \mathcal{L}_{\rho^{-1}\Gamma}[n] ds$.

Formula (5.9) shows that as the frequency increases, the conductivity anisotropy decreases. The anisotropic information can not be captured for

$$\omega \gg \frac{1}{\varepsilon_m} ((l_1 - l_2) \frac{2\delta \sigma_m \rho}{|\Gamma|} - \sigma_m^2)^{1/2}.$$

6 Spectroscopic imaging of a dilute suspension

6.1 Spectroscopic conductivity imaging

We now make use of the asymptotic expansion of the effective admittivity in terms of the volume fraction $f = \rho^2$ to image a permittivity inclusion. Consider *D* to be a bounded domain

in Ω with admittivity $1 + fM(\omega)$, where $M(\omega)$ is a membrane polarization tensor and f is the volume fraction of the suspension in D. The inclusion D models a suspension of cells in the background Ω . For simplicity, we neglect the permittivity ϵ_0 of Ω and assume that its conductivity $\sigma_0 = 1$. We also assume that $M(\omega)$ is isotropic. At the macroscopic scale, if we inject a current g on $\partial\Omega$, then the electric potential satisfies:

$$\begin{cases} \nabla \cdot (1 + fM(\omega)\chi_D)\nabla u = 0 & \text{in } \Omega, \\ \frac{\partial u}{\partial n}\Big|_{\partial\Omega} = g, \quad \int_{\partial\Omega} g(x)ds(x) = 0, \quad \int_{\Omega} u(x)dx = 0. \end{cases}$$
(6.1)

The imaging problem is to detect and characterize *D* from measurements of *u* on $\partial \Omega$.

Integrating by parts and using the trace theorem for the double-layer potential [23, 47], we obtain, $\forall x \in \partial \Omega$,

$$\frac{1}{2}u(x) + \frac{1}{2\pi}\int_{\partial\Omega}\frac{(x-y)\cdot n(x)}{|x-y|^2}u(y)ds(y) + \frac{1}{2\pi}\int_{\partial\Omega}g(y)\ln|x-y|ds(y)$$

$$= \frac{f}{2\pi}M(\omega)\int_D\nabla u(y)\cdot\frac{(x-y)}{|x-y|^2}dy.$$
(6.2)

Since *f* is small,

$$\int_D \nabla u(y) \cdot \frac{(x-y)}{|x-y|^2} dy \simeq \int_D \nabla U(y) \cdot \frac{(x-y)}{|x-y|^2} dy$$

holds uniformly for $x \in \partial \Omega$, where *U* is the background solution, that is,

$$\begin{cases} \Delta U = 0 & \text{in } \Omega, \\ \frac{\partial U}{\partial n}\Big|_{\partial \Omega} = g, \quad \int_{\Omega} U(x) dx = 0, \end{cases}$$

Therefore, taking the imaginary part of (6.2) yields

$$\frac{1}{2}\Im u(x) + \frac{1}{2\pi} \int_{\partial\Omega} \frac{(x-y) \cdot n(x)}{|x-y|^2} \Im u(y) ds(y) \simeq \frac{f}{2\pi} \Im M(\omega) \int_D \nabla U(y) \cdot \frac{(x-y)}{|x-y|^2} dy, \tag{6.3}$$

uniformly for $x \in \partial \Omega$, provided that *g* is real. Finally, taking the argument of the maximum of the right-hand side in (6.3) with respect to the frequency ω gives the Debye relaxation time of the suspension in *D*.

6.2 Selective spectroscopic imaging

A challenging applied problem is to design a selective spectroscopic imaging approach for suspensions of cells. Using a pulsed imaging approach [33, 38], we propose a simple way to selectively image dilute suspensions. Again, we assume for the sake of simplicity that $\epsilon_0 = 0$ and $\sigma_0 = 1$.

In the time-dependant regime, the electrical model for the cell (2.1) is replaced with

$$u(x,t) = \int \hat{h}(\omega)\hat{u}(x,\omega)e^{i\omega t}d\omega,$$

where $\hat{u}(x, \omega)$ is the solution to

$$\Delta \hat{u}(\cdot,\omega) = 0 \qquad \qquad \text{in } D \setminus \overline{C},$$

$$\Delta \hat{u}(\cdot,\omega) = 0 \qquad \qquad \text{in } C,$$

$$\frac{\partial \hat{u}(\cdot,\omega)}{\partial n}\Big|_{+} = \frac{\partial \hat{u}(\cdot,\omega)}{\partial n}\Big|_{-} \qquad \text{on } \Gamma,$$
(6.4)

$$\begin{cases} \frac{\partial \hat{u}(\cdot,\omega)}{\partial n}\Big|_{+} = \frac{\partial \hat{u}(\cdot,\omega)}{\partial n}\Big|_{-} & \text{on } \Gamma, \\ \hat{u}(\cdot,\omega)\Big|_{+} - \hat{u}(\cdot,\omega)\Big|_{-} - \beta(\omega)\frac{\partial \hat{u}(\cdot,\omega)}{\partial n} = 0 & \text{on } \Gamma, \\ \frac{\partial \hat{u}(\cdot,\omega)}{\partial n}\Big|_{\partial D} = f, \quad \int_{\partial D} \hat{u}(\cdot,\omega)ds = 0, \end{cases}$$

and

$$h(t) = \int \hat{h}(\omega) e^{i\omega t} d\omega$$

is the pulse shape. The support of *h* is assumed to be compact.

At the macroscopic scale, if we inject a pulsed current, g(x)h(t), on $\partial\Omega$, then the electric potential u(x, t) in the presence of a suspension occupying *D* is given by

$$u(x,t) = \int \hat{h}(\omega)\hat{u}(x,\omega)e^{i\omega t}d\omega,$$

where

$$\begin{cases} \nabla \cdot (1 + fM(\omega)\chi_D)\nabla \hat{u}(\cdot, \omega) = 0 & \text{in } \Omega, \\ \frac{\partial \hat{u}(\cdot, \omega)}{\partial n}\Big|_{\partial \Omega} = g, \quad \int_{\partial \Omega} \hat{u}(\cdot, \omega) ds = 0. \end{cases}$$

Assume that we are in the presence of two suspensions occupying the domains D_1 and D_2 inside Ω . From (6.2) it follows that

$$\frac{1}{2}\hat{u}(x,\omega) + \frac{1}{2\pi} \int_{\partial\Omega} \frac{(x-y) \cdot n(x)}{|x-y|^2} \hat{u}(y,\omega) ds(y) + \frac{1}{2\pi} \int_{\partial\Omega} g(y) \ln|x-y| ds(y) \simeq \frac{f_1}{2\pi} M_1(\omega) \int_{D_1} \nabla U(y) \cdot \frac{(x-y)}{|x-y|^2} dy + \frac{f_2}{2\pi} M_2(\omega) \int_{D_2} \nabla U(y) \cdot \frac{(x-y)}{|x-y|^2} dy,$$
(6.5)

and therefore,

$$\frac{1}{2}u(x,t) + \frac{1}{2\pi} \int_{\partial\Omega} \frac{(x-y) \cdot n(x)}{|x-y|^2} u(y,t) ds(y) + \frac{1}{2\pi} h(t) \int_{\partial\Omega} g(y) \ln |x-y| ds(y)$$

$$\simeq \frac{f_1}{2\pi} \mathcal{M}_1(t) \int_{D_1} \nabla U(y) \cdot \frac{(x-y)}{|x-y|^2} dy + \frac{f_2}{2\pi} \mathcal{M}_2(t) \int_{D_2} \nabla U(y) \cdot \frac{(x-y)}{|x-y|^2} dy,$$
(6.6)

uniformly in $x \in \partial \Omega$ and $t \in \operatorname{supp} h$, where

$$\mathcal{M}_i(t) := \int \hat{h}(\omega) M_i(\omega) e^{i\omega t} d\omega, \quad i = 1, 2.$$

As it will be shown in section 8, by comparing the Debye relaxation times associated to M_1 and M_2 , one can design the pulse shape h in order to image selectively D_1 or D_2 . For example, one can selectively image D_1 by taking $\hat{h}(\omega)$ close to zero around the Debye relaxation time of M_2 and close to one around the Debye relaxation time of M_1 .

6.3 Spectroscopic measurement of anisotropy

In this subsection we assume that *M* is anisotropic and consider the solution *u* to (6.1). We want to assess the anisotropy of the inclusion *D* of admittivity $1 + fM(\omega)$ from measurements of *u* on the boundary $\partial\Omega$.

From (6.3) it follows that

$$\int_{\partial\Omega} g(x) \left[\frac{1}{2} \Im u(x) + \frac{1}{2\pi} \int_{\partial\Omega} \frac{(x-y) \cdot n(x)}{|x-y|^2} \Im u(y) ds(y) \right] ds(x)$$

$$\simeq \frac{f}{2\pi} \int_D \Im M(\omega) \nabla U(y) \cdot \nabla U(y) dy,$$
(6.7)

provided that *g* is real. Now, taking constant current sources corresponding to $g = a \cdot n$, where $a \in \mathbb{R}^2$ is a unit vector, yields

$$\mathcal{S}[a] := \int_{\partial\Omega} g(x) \left[\frac{1}{2} \Im u(x) + \frac{1}{2\pi} \int_{\partial\Omega} \frac{(x-y) \cdot n(x)}{|x-y|^2} \Im u(y) ds(y) \right] ds(x) \simeq \frac{f}{2\pi} \Im M(\omega) |a|^2 |D|.$$

Since

$$\frac{\min_{a} \mathcal{S}[a]}{\max_{a} \mathcal{S}[a]} \simeq \frac{\lambda_{1}(\omega)}{\lambda_{2}(\omega)},$$

where λ_1 and λ_2 (with $\lambda_1 \leq \lambda_2$) are the eigenvalues of $\Im M$, it follows from subsection 5.5 that

$$\omega \mapsto \frac{\min_a \mathcal{S}[a]}{\max_a \mathcal{S}[a]}$$

is a natural measure of conductivity anisotropy. This measure may be used for the detection and classification of neuromuscular diseases via measurement of muscle anisotropy [21, 32].

7 Stochastic homogenization of randomly deformed conductivity resistant membranes

The first main result of this section is to show that a rigorous homogenization theory can be derived when the cells (and hence interfaces) are randomly deformed from a periodic structure, and the random deformation is ergodic and stationary in the sense of (2.13).

7.1 Auxiliary problem: proof of Theorem 2.3

In this subsection, we prove Theorem 2.3 about the existence and uniqueness of the auxiliary problem. This is the key step in stochastic homogenization. The main difficulty is due to the fact that one does not have compactness in the general stationary ergodic setting.

We first make the weak formulation of the system (2.23) precise. To this end, we introduce the space $\widetilde{\mathcal{H}} := L^2(\mathcal{O}, H^1_{\text{loc}}(\mathbb{R}_2^+) \times H^1_{\text{loc}}(\mathbb{R}_2^-))$ and the space $\widetilde{\mathcal{H}}_S$ which is a subspace of $\widetilde{\mathcal{H}}$ where the elements are stationary. Define also the space $\mathcal{H} := \{w = \widetilde{w} \circ \Phi^{-1} \mid \widetilde{w} \in \widetilde{\mathcal{H}}\}$ and the space $\mathcal{H}_S := \{w = \widetilde{w} \circ \Phi^{-1} \mid \widetilde{w} \in \widetilde{\mathcal{H}}_S\}.$ We say that $w_p = w_p^+ \chi_{\Phi(\mathbb{R}_2^+)} + w_p^- \chi_{\Phi(\mathbb{R}_2^-)} \in \mathcal{H}$ is a weak solution to (2.23) if ∇w_p is stationary and for all $\varphi \in \mathcal{H}$ with compact support $K \subset \mathbb{R}^2$, it holds that

$$\mathbb{E}\int_{K\cap\Phi(\mathbb{R}_{2}^{+},\gamma)}k_{0}(p+\nabla w_{p}^{+})\cdot\nabla\overline{\varphi}dx + \mathbb{E}\int_{K\cap\Phi(\mathbb{R}_{2}^{-},\gamma)}k_{0}(p+\nabla w_{p}^{-})\cdot\nabla\overline{\varphi}dx + \mathbb{E}\int_{K\cap\Phi(\Gamma,\gamma)}\frac{1}{\beta}(w_{p}^{+}-w_{p}^{-})(\overline{\varphi}^{+}-\overline{\varphi}^{-})ds(x) = 0.$$
(7.1)

Since the integrals above do not control $\|w_p^{\pm}\|_{L^2(\Omega, L^2_{loc}(\Phi(\mathbb{R}^{\pm}_2)))}$ and the space \mathcal{H} does not possess Poincaré inequality, the existence of weak solutions is not immediate.

Remark 7.1. Due to the separability of $L^2(\mathcal{O})$, one can show that for almost all $\gamma \in \mathcal{O}$, the solution $w_p(\cdot, \gamma)$ satisfying (7.1) is also a weak solution in the usual sense. That is, for all $\phi = \tilde{\phi} \circ \Phi^{-1}$ where $\tilde{\phi} \in H^1_{loc}(\mathbb{R}^+_2) \times H^1_{loc}(\mathbb{R}^-_2)$ with compact support $K \subset \mathbb{R}^2$

$$\int_{\Phi(K \cap \mathbb{R}_{2}^{+}, \gamma)} k_{0}(p + \nabla w_{p}^{+}) \cdot \nabla \overline{\phi} dx + \int_{\Phi(K \cap \mathbb{R}_{2}^{-}, \gamma)} k_{0}(p + \nabla w_{p}^{-}) \cdot \nabla \overline{\phi} dx + \int_{\Phi(K \cap \Gamma, \gamma)} \frac{1}{\beta} (w_{p}^{+} - w_{p}^{-}) (\overline{\phi}^{+} - \overline{\phi}^{-}) ds(x) = 0.$$
(7.2)

We refer to [35, Proposition 4.1] for the proof.

Our strategy is as follows: First, an absorption term is added to regularize the problem. The sequence of regularized solutions, which corresponds to a sequence of vanishing regularization, have a converging gradient. Secondly, the potential field corresponding to the limiting gradient is shown to be a solution to the auxiliary problem. Finally, using regularity results and sub-linear growth of potential field with stationary gradient, we prove that the gradient of the solution to the auxiliary problem is unique.

Proof of Theorem 2.3. Step 1: The regularized auxiliary problem. Fix $p \in \mathbb{R}^2$. Consider the following regularized problem where an absorption $\alpha > 0$ is added.

$$\begin{cases} -\nabla \cdot k_0 (\nabla w_{p,\alpha}^{\pm}(y) + p) + \alpha w_{p,\alpha}^{\pm} = 0 & \text{in} \quad \Phi(\mathbb{R}_2^{\pm}, \gamma), \\ n \cdot k_0 \nabla w_{p,\alpha}^{-}(y) = n \cdot k_0 \nabla w_{p,\alpha}^{-}(y), & \text{in} \quad \Phi(\Gamma, \gamma), \\ w_{p,\alpha}^{+} - w_{p,\alpha}^{-} = \beta k_0 n \cdot (\nabla w_{p,\alpha}^{-} + p) & \text{in} \quad \Phi(\Gamma, \gamma), \\ w_{p,\alpha}^{\pm}(y, \gamma) = \widetilde{w}_{p,\alpha}^{\pm}(\Phi^{-1}(y, \gamma), \gamma), \text{ and} \quad \widetilde{w}_{p,\alpha}^{\pm} \text{ are stationary.} \end{cases}$$
(7.3)

We first construct a solution for the above equation in \mathcal{H}_S in a sense that seems weaker than (7.1). It can be verified that \mathcal{H}_S equipped with the inner product

$$(u,v)_{\mathcal{H}_{\mathcal{S}}} = \mathbb{E}\left(\int_{Y^{+}} \nabla \tilde{u} \cdot \nabla \overline{\tilde{v}} dx + \int_{Y^{-}} \nabla \tilde{u} \cdot \nabla \overline{\tilde{v}} dx + \int_{Y} \tilde{u} \overline{\tilde{v}} dx\right).$$
(7.4)

is a Hilbert space. For any fixed $\alpha > 0$, define the bilinear form $A_{\alpha} : \mathcal{H}_S \times \mathcal{H}_S \to \mathbb{R}$ by

$$\begin{aligned} A_{\alpha}(u,v) &= \mathbb{E}\left(\int_{\Phi(Y^{+})} k_{0} \nabla u^{+} \cdot \overline{\nabla v^{+}} dx + \int_{\Phi(Y^{-})} k_{0} \nabla u^{-} \cdot \overline{\nabla v^{-}} dx \right. \\ &+ \alpha \int_{\Phi(Y)} u \overline{v} dx + \frac{1}{\beta} \int_{\Phi(\Gamma_{0})} (u^{+} - u^{-}) \overline{(v^{+} - v^{-})} ds \right), \end{aligned}$$

and the linear functional $b_p : \mathcal{H}_S \to \mathbb{R}$ by

$$b_p(v) = -k_0 \mathbb{E}\left(\int_{\Phi(Y^+)} p \cdot \overline{\nabla v^+} dx + \int_{\Phi(Y^-)} p \cdot \overline{\nabla v^-} dx\right).$$

We verify that A_{α} is bounded and coercive, and b_p is bounded. By the Lax–Milgram theorem, there exists a unique element $w_{p,\alpha} \in \mathcal{H}_S$ satisfying

$$A_{\alpha}(w_{p,\alpha},\varphi) = b_p(\varphi), \quad \forall \varphi \in \mathcal{H}_S.$$
(7.5)

By choosing φ to be $w_{p,\alpha}$, we obtain the estimates:

$$\mathbb{E}\int_{Y^{\pm}} |\nabla \tilde{w}_{p,\alpha}^{\pm}|^2 \le C, \quad \mathbb{E}\int_{\Gamma_0} |\tilde{w}_{p,\alpha}^{+} - \tilde{w}_{p,\alpha}^{-}|^2 \le C, \quad \mathbb{E}\int_{Y^{\pm}} |\tilde{w}_{p,\alpha}^{\pm}|^2 \le \frac{C}{\alpha}.$$
 (7.6)

Next we argue that for almost all $\gamma \in O$, the solution $w_{p,\alpha}(\cdot, \gamma)$ above satisfies (7.3) in the usual distributional sense. That is, for any $\phi(x) \in C^{\infty}(\mathbb{R}^+_2) \cap C^{\infty}(\mathbb{R}^-_2)$, whose support *K* is a compactly contained in \mathbb{R}^2 , we have

$$\int_{K\cap\Phi(\mathbb{R}_{2}^{+},\gamma)} k_{0}(p+\nabla w_{p,\alpha}^{+}) \cdot \nabla \overline{\phi} dx + \int_{K\cap\Phi(\mathbb{R}_{2}^{-},\gamma)} k_{0}(p+\nabla w_{p,\alpha}^{-}) \cdot \nabla \overline{\phi} dx + \alpha \int_{K} w_{p,\alpha} \overline{\phi} dx + \int_{K\cap\Phi(\mathbb{F},\gamma)} \frac{1}{\beta} (w_{p,\alpha}^{+}-w_{p,\alpha}^{-}) (\overline{\phi}^{+}-\overline{\phi}^{-}) ds(x) = 0.$$
(7.7)

Indeed, due to the regularization, the above problem (with any fixed $\gamma \in \mathcal{O}$ and any fixed δ) admits a unique solution in the space $H^1_{\text{uloc}}(\Phi(\mathbb{R}^+_2, \gamma)) \times H^1_{\text{uloc}}(\Phi(\mathbb{R}^-_2, \gamma))$. This result is nontrivial and we refer to [35, Lemma 4.3] for the proof. Then one can verify that the solution $w_{p,\alpha}(\cdot, \gamma)$ is stationary and satisfies (7.5); therefore, it must agree with the solution provided by the Lax-Milgram theorem. As a consequence, $w_{p,\alpha}(x, \gamma)$ is also a weak solution in \mathcal{H} to (2.23) in the sense of (7.1).

Applying Corollary A.1 and Corollary A.2 to the family $\{\tilde{w}_{p,\alpha}\}_{\alpha}$, we obtain a family $\{\tilde{w}_{p,\alpha}^{ext} = P\tilde{w}_{p,\alpha}^{+}\}_{\alpha} \subset L^{2}(\mathcal{O}, H_{\text{loc}}^{1}(\mathbb{R}^{2}))$ and a family $\{w_{p,\alpha}^{ext} = P_{\gamma}w_{p,\alpha}^{+}\}_{\alpha}$. Further, $\{\tilde{w}_{p,\alpha}^{ext}\}_{\alpha}$ are stationary. They satisfy that $w_{p,\alpha}^{ext} = \tilde{w}_{p,\alpha}^{ext} \circ \Phi^{-1}$ and that

$$\mathbb{E}\int_{Y} |\nabla \tilde{w}_{p,\alpha}^{\text{ext}}|^2 \le C, \quad \mathbb{E}\int_{\Gamma_0} |\tilde{w}_{p,\alpha}^{\text{ext}} - \tilde{w}_{p,\alpha}^-|^2 \le C, \quad \mathbb{E}\int_{Y} |\tilde{w}_{p,\alpha}^{\text{ext}}|^2 \le \frac{C}{\alpha}.$$
(7.8)

Step 2: Extraction of a converging subsequence. The family $\{\tilde{w}_{p,\alpha}^{\text{ext}}\}_{\alpha}$ may be studied from two view points. Firstly, they form a bounded family in $\tilde{\mathcal{H}}_S$. Secondly, they belong to $\tilde{\mathcal{H}}$ and for any compact set $K \subset \mathbb{R}^2$, the estimates (7.8) imply that

$$\mathbb{E}\int_{K} |\nabla \tilde{w}_{p,\alpha}^{\text{ext}}|^{2} \leq C(K), \quad \mathbb{E}\int_{\mathbf{\Gamma}\cap K} |\tilde{w}_{p,\alpha}^{\text{ext}} - \tilde{w}_{p,\alpha}^{-}|^{2} \leq C(K), \quad \alpha \mathbb{E}\int_{K} |\tilde{w}_{p,\alpha}^{\text{ext}}|^{2} \leq C(K).$$
(7.9)

From the first point of view, there exists a subsequence, still denoted by $\nabla \tilde{w}_{p,\alpha}^{\text{ext}}$, which converges weakly as $\alpha \downarrow 0$ to a function $\tilde{\eta}_p^{\text{ext}} \in [L_S^2(\mathcal{O}, L_{\text{loc}}^2(\mathbb{R}^2)]^2$ where the subscript *S* indicates stationary. By a change of variable, we also have that $\nabla w_{p,\alpha}^{\text{ext}}$ converges in $[L^2(\mathcal{O}, L_{\text{loc}}^2(\mathbb{R}^2)]^2$ to η_p^{ext} and

$$\eta_p^{\text{ext}}(y,\gamma) = \nabla_y \Psi(y,\gamma) \tilde{\eta}_p^{\text{ext}}(\tilde{y},\gamma), \qquad (7.10)$$

where $\Psi = \Phi^{-1}$ and $\tilde{y} = \Psi(y)$. Moreover, as gradients, $\nabla_{\tilde{y}} \tilde{w}_{p,\alpha}^{\text{ext}}$ and $\nabla_{y} w_{p,\alpha}^{\text{ext}}$ are curl free. This property is preserved by their limits:

$$\partial_{y_i}(\eta_p^{\text{ext}})_j = \partial_{y_j}(\eta_p^{\text{ext}})_i, \quad \partial_{\tilde{y}_i}(\tilde{\eta}_p^{\text{ext}})_j = \partial_{\tilde{y}_j}(\tilde{\eta}_p^{\text{ext}})_i, \quad i, j \in \{1, 2\}.$$
(7.11)

That is to say, η_p^{ext} and $\tilde{\eta}_p^{\text{ext}}$ are also gradient functions. Consequently, there exist w_p^{ext} and \tilde{w}_p^{ext} such that $\eta_p^{\text{ext}} = \nabla_y w_p^{\text{ext}}$ and $\tilde{\eta}_p^{\text{ext}} = \nabla_{\tilde{y}} \tilde{w}_p^{\text{ext}}$. The relation (7.10) implies that $w_p^{\text{ext}}(y) = \tilde{w}_p^{\text{ext}}(\Psi(y,\gamma),\gamma) + C_p(\gamma)$ where $C_p(\gamma)$ is a random constant. We hence re-define \tilde{w}_p^{ext} by adding to it the random variable C_p so that $w_p^{\text{ext}} = \tilde{w}_p^{\text{ext}} \circ \Psi$. By the same token, we have that $\nabla \tilde{w}_{p,\alpha}^-$ and $\nabla w_{p,\alpha}^-$ converge (along the above subsequence) to $\tilde{\eta}_p^- \in [L_S^2(\mathcal{O}, L_{\text{loc}}^2(\mathbb{R}_2^-))]^2$ and $\eta_p^- \in [L^2(\mathcal{O}, L_{\text{loc}}^2(\Phi(\mathbb{R}_2^-)))]^2$ respectively. In addition, for some $\tilde{w}_p^- \in L^2(\mathcal{O}, H_{\text{loc}}^1(\mathbb{R}_2^-))$ and $w_p^- \in L^2(\mathcal{O}, H_{\text{loc}}^1(\mathbb{R}_2^-))$ satisfying that $w_p^- = \tilde{w}_p^- \circ \Psi$, we have $\tilde{\eta}_p^- = \nabla \tilde{w}_p^-$ and $\eta_p^- = \nabla w_p^-$. Similarly, due to the second bound in (7.8), one observes that $\{(\tilde{w}_{p,\alpha}^{\text{ext}} - \tilde{w}_{p,\alpha}^-)|_{\Gamma}\}_{\alpha}$ converges (through a subsequence) to some $\tilde{\zeta}_p \in L_S^2(\mathcal{O}, L_{\text{loc}}^2(\Gamma))$. Again, by a change of variable, $\{(w_{p,\alpha}^{\text{ext}} - w_{p,\alpha}^-)_{\Phi(\Gamma)}\}_{\alpha}$ converges to certain $\zeta_p \in L^2(\mathcal{O}, L_{\text{loc}}^2(\Gamma))$ and it holds that $\zeta_p = \tilde{\zeta}_p \circ \Psi$. Finally, since $\tilde{w}_{p,\alpha}^{\text{ext}}$ is stationary, one has $\mathbb{E} \int_{Y} \nabla_{\tilde{y}} \tilde{w}_{p,\alpha}^{\text{ext}} d\tilde{y} = 0$. Passing to the limit, we get

$$\mathbb{E}\int_{Y} \nabla_{\tilde{y}} \tilde{w}_{p}^{\text{ext}}(\tilde{y}) d\tilde{y} = 0.$$
(7.12)

Now, from the second point of view and the estimate (7.9), we can choose a further subsequence of the converging subsequence obtained from the first view point, still denoted by $\{\tilde{w}_{p,\alpha}\}_{\alpha}$ and so on, such that the family $\nabla \tilde{w}_{p,\alpha}^{ext}$ converges in $L^2(\mathcal{O}, [L_{loc}^2(\mathbb{R}^2)]^2)$ to $\tilde{\eta}_p^{ext}, \{\nabla \tilde{w}_{p,\alpha}^-\}_{\alpha}$ converges to $\tilde{\eta}_p^-$ and $\{(\tilde{w}_{p,\alpha}^{ext} - \tilde{w}_{p,\alpha}^-)|_{\Gamma}\}_{\alpha}$ converges in $L^2(\mathcal{O}, L_{loc}^2(\Gamma))$ to $\tilde{\zeta}_p$. We then verify that these functions are stationary, and by the ergodic theorem they agree with the limits $\tilde{\eta}_p^{ext}, \tilde{\eta}_p^$ and $\tilde{\zeta}_p$ obtained from the first point of view. As a result, we take expectation on the weak formulation (7.7), and then pass to the limit and obtain: for all $\varphi \in \mathcal{H}$ with compact support $K \subset \mathbb{R}^2$, we have

$$\mathbb{E}\int_{K\cap\Phi(\mathbb{R}_{2}^{+},\gamma)}k_{0}(p+\nabla w_{p}^{+})\cdot\nabla\overline{\varphi}dx + \mathbb{E}\int_{K\cap\Phi(\mathbb{R}_{2}^{-},\gamma)}k_{0}(p+\nabla w_{p}^{-})\cdot\nabla\overline{\varphi}dx + \mathbb{E}\int_{K\cap\Phi(\Gamma,\gamma)}\frac{1}{\beta}\zeta_{p}(\overline{\varphi}^{+}-\overline{\varphi}^{-})ds(x) = 0.$$
(7.13)

This is almost (7.1) except for the interface term. By choosing φ compactly supported in $\Phi(\mathbb{R}_2^+)$ (respectively in $\Phi(\mathbb{R}_2^-)$) we verify that the first (respectively the second) equation of (2.23) is satisfied. By choosing $\varphi^- = 0$ first and then $\varphi^+ = 0$, and applying the divergence theorem we check that the third line of (2.23) is satisfied and further

$$\beta k_0 \left(\frac{\partial w_p^{\pm}}{\partial n} (y) + v_y \cdot p \right) = \zeta_p.$$

To relate ζ_p with $w_p^+ - w_p^-$, for any $h \in C^{\infty}(\Phi(\Gamma))$ such that $\int_{\Phi(\Gamma_k)} h = 0$ for each $k \in \mathbb{Z}^d$, we can find φ^{\pm} such that for some fixed ball $B^{\delta} \subset \subset Y$ which contains Y^- and its boundary is well separated from ∂Y^-

$$\begin{cases} -\nabla \cdot k_0 \nabla \varphi^{\pm} = 0, & \text{in} \qquad \Phi(Y_k^-) \text{ and } \Phi(k + B^{\delta}) \setminus \Phi(Y_k^-), \\ n \cdot k_0 \nabla \varphi^+ = n \cdot k_0 \nabla \varphi^- = h, & \text{on} \qquad \Phi(\Gamma_k), \\ \varphi^+ = 0, & \text{on} \qquad \Phi(k + \partial B^{\delta}). \end{cases}$$

Use $\varphi = (\varphi^+, \varphi^-)$ in (7.13) and apply again the divergence theorem. We verify that $w_p^+ - w_p^- = \zeta_p + C(\omega, k)$ on each $\Phi(\Gamma_k)$. By adding the constant to w_p^- , we may set $\zeta = w_p^+ - w_p^-$; so the fourth line in (2.23) is proved. From the construction, it is easy to see that the last two lines of that equation also hold. To summarize $w_p^{\text{ext}}\chi_{\Phi(\mathbb{R}^+_2,\gamma)} + w_p^-\chi_{\Phi(\mathbb{R}^-_2)} \in \mathcal{H}$ provides a weak solution to (2.23).

Step 3: Uniqueness of ∇w_p . For any two solutions w_p^1 and w_p^2 of (2.23), let v_p be their difference. Then it satisfies (7.1) with p = 0. In addition, there is an extension of \tilde{v}_0^+ denoted by \tilde{v}_0^{ext} , such that $\nabla \tilde{v}_0^{\text{ext}} = P(\nabla \tilde{v}_0^+)$ where *P* is the extension operator of Corollary A.1, and \tilde{v}_0^{ext} satisfies

$$\nabla \tilde{v}_0^{\text{ext}} \text{ is stationary, and } \mathbb{E} \int_Y \nabla \tilde{v}_0^{\text{ext}} dx = 0.$$
 (7.14)

In the usual weak formulation of the equations satisfied by (v_0^+, v_0^-) , take v_0 itself as the test function and integrate over $\Phi(NY)$ for a large integer *N*. We get

$$\begin{split} \int_{\Phi(NY\cap\mathbb{R}_{2}^{+})} k_{0} |\nabla v_{0}^{+}|^{2} dx + \int_{\Phi(NY\cap\mathbb{R}_{2}^{-})} k_{0} |\nabla v_{0}^{-}|^{2} dx \\ &+ \beta^{-1} \int_{\Phi(NY\cap\Gamma)} |v_{0}^{+} - v_{0}^{-}|^{2} ds = \int_{\partial\Phi(NY)} k_{0} (n \cdot \nabla v_{0}^{+}) \overline{v_{0}^{+}} ds. \end{split}$$

Invoking elliptic regularity, we can show that $|\nabla v_0^+| \leq C \|\nabla v_0\|_{L^2(\Phi(Y_k))}$ along the boundary $\partial \Phi(NY)$. Let $\mathcal{K}_{\partial(NY)}$ denote the indices of those cubes. By a change of variables and extension, we also have such bound of $\nabla \tilde{v}_0^{\text{ext}}$. On the other hand, due to the assumption that $\nabla \tilde{v}_0^{\text{ext}}$ is stationary and satisfies $\mathbb{E} \int_Y \nabla \tilde{v}_0^{\text{ext}} dy = 0$, in light of Lemma A.5 of [13], we conclude that \tilde{v}_0^{ext} grows sublinearly almost surely. Consequently, the integral on the right is of order $o(N) \sum_{k \in \mathcal{K}_{\partial(NY)}} \|\nabla v_0\|_{L^2(\Phi(Y_k))}$. Divide the above equality by N^2 , and change variable in the integrals, we have, with $\mathcal{I}(N)$ being the indices of cubes $\{Y_n \subset NY\}$.

$$\frac{1}{N^2} \sum_{n \in \mathcal{I}(N)} \left[\int_{\Phi(Y_n^+)} \sigma_0 |\nabla v_0^+|^2 dx + \int_{\Phi(Y_n^-)} \sigma_0 |\nabla v_0^-|^2 dx + \Re \beta^{-1} \int_{\Phi(\Gamma_n)} |v_0^+ - v_0^-|^2 ds \right]$$

converges to the limit of $\frac{o(N)}{N^2} \sum_{k \in \mathcal{K}_{\partial(NY)}} \|\nabla \tilde{v}_0^{\text{ext}}\|_{L^2(Y_k)}$, which is zero by the fact that the cardinality of $\mathcal{K}_{\partial(NY)}$ is O(N) and by using ergodic theorem for this sum. By a change of variable with bounds (2.17) and (2.18), the above implies

$$\frac{1}{N^2} \sum_{n \in \mathcal{I}(N)} \left[\int_{Y_n^+} \sigma_0 |\nabla \tilde{v}_0^+|^2 d\tilde{x} + \int_{Y_n^-} \sigma_0 |\nabla \tilde{v}_0^-|^2 d\tilde{x} + \beta^{-1} \int_{\Gamma_n} |\tilde{v}_0^+ - \tilde{v}_0^-|^2(\tilde{x}) ds(\tilde{x}) \right] \longrightarrow 0$$

By the stationarity of the integrands and again the ergodic theorem, we have

$$\mathbb{E}\int_{Y^+}|\nabla \tilde{v}_0^+|^2d\tilde{x} + \mathbb{E}\int_{Y^-}|\nabla \tilde{v}_0^-|^2d\tilde{x} + \mathbb{E}\int_{\Gamma_0}|\tilde{v}_0^+ - \tilde{v}_0^-|^2(\tilde{x})ds(\tilde{x}) = 0.$$

This implies that $\tilde{v}_0^+ = \tilde{v}_0^- = C(\gamma)$ for some random constant. This proves the uniqueness of ∇w_p .

7.2 **Proof of the homogenization theorem**

In this section, we prove the homogenization theorem using the energy method, *i.e.*, the method of oscillating test functions [46].

7.2.1 Oscillating test functions

We first build the oscillating test functions. For a fixed vector $p \in \mathbb{R}^2$. Let $(w_p^+, w_p^-) \in \mathcal{H}$ be the unique solution (up to the addition of a random constant) of the auxiliary problem (2.23). In particular, w_p^+ has an extension w_p^{ext} . We define

$$\begin{cases} w_{1p}^{\varepsilon}(x,\gamma) = x \cdot p + \varepsilon w_p^{\text{ext}}(\frac{x}{\varepsilon},\gamma), & x \in \mathbb{R}^2, \\ w_{2p}^{\varepsilon}(x,\gamma) = x \cdot p + \varepsilon Q w_p^{-}(\frac{x}{\varepsilon},\gamma), & x \in \mathbb{R}^2. \end{cases}$$
(7.15)

Here and in the sequel, Q denotes the trivial extension operator which sets Qf = 0 outside the spatial support of f. By scaling the auxiliary problem, we verify that $(w_p^{\varepsilon_+}, w_p^{\varepsilon_-})$, where $w_p^{\varepsilon_+}$ is the restriction of w_{1p}^{ε} in $\varepsilon \Phi(\mathbb{R}_2^+)$ and $w_p^{\varepsilon_-}$ is the restriction of w_{2p}^{ε} in $\varepsilon \Phi(\mathbb{R}_2^-)$, satisfies

$$\begin{cases} \nabla \cdot k_0 \nabla w_p^{\varepsilon +} = 0 & \text{and} \quad \nabla \cdot k_0 \nabla w_p^{\varepsilon -} = 0 & \text{in } \varepsilon \Phi(\mathbb{R}_2^{\pm}), \\ k_0 n \cdot \nabla w_p^{\varepsilon +} = k_0 n \cdot \nabla w_p^{\varepsilon -} & \text{and} \quad w_p^{\varepsilon +} - w_p^{\varepsilon -} = \varepsilon \beta k_0 n \cdot \nabla w_p^{\varepsilon -} & \text{on } \varepsilon \Phi(\Gamma). \end{cases}$$

This means that for any test function $\varphi = (\varphi^+, \varphi^-) \in L^2(\mathcal{O}, H^1_{loc}(\varepsilon \Phi(\mathbb{R}^+_2) \times \varepsilon \Phi(\mathbb{R}^-_2)))$ compactly supported on a bounded open set $\mathscr{O} \subset \mathbb{R}^2$, we have that

$$\mathbb{E}\int_{\mathscr{O}\cap\varepsilon\Phi(\mathbb{R}_{2}^{+})}k_{0}\nabla w_{p}^{\varepsilon+}\cdot\overline{\nabla\varphi^{+}}dx + \mathbb{E}\int_{\mathscr{O}\cap\varepsilon\Phi(\mathbb{R}_{2}^{-})}k_{0}\nabla w_{p}^{\varepsilon-}\cdot\overline{\nabla\varphi^{-}}dx + (\varepsilon\beta)^{-1}\mathbb{E}\int_{\mathscr{O}\cap\varepsilon\Phi(\Gamma)}(w_{p}^{\varepsilon-}-w_{p}^{\varepsilon-})\overline{(\varphi^{+}-\varphi^{-})}ds = 0.$$
(7.16)

Clearly, this is the scaled version of (7.1). Remark 7.1 applies here also, so w_p^{ε} solves the above equations in the usual weak sense as well. We define the vector fields $\eta_p^{\varepsilon\pm} = k_0 \nabla w_p^{\varepsilon\pm}$. They satisfy the following convergence results.

Lemma 7.1. Let $w_p^{\varepsilon\pm}$ and the vector fields $\eta_p^{\varepsilon\pm}$ be defined as above and let $\mathscr{O} \subset \mathbb{R}^2$ be a bounded open set. Then as $\varepsilon \to 0$, we have the following:

$$w_{1p}^{\varepsilon} \to x \cdot p$$
, uniformly in \mathcal{O} a.s. in \mathcal{O} ; (7.17)

$$w_{2p}^{\varepsilon} \to x \cdot p,$$
 in $L^{2}(\mathcal{O})$ a.s. in \mathcal{O} . (7.18)

$$Q\eta_p^{\varepsilon\pm} \rightharpoonup \varrho \mathbb{E} \int_{\Phi(Y^{\pm})} k_0 \left(\nabla w_p^{\pm}(x, \cdot) + p \right) dx \qquad \text{in } [L^2(\mathscr{O})]^2 \qquad \text{a.s. in } \mathcal{O}.$$
(7.19)

Proof. To prove the first result, we recall that (w_p^+, w_p^-) solves (7.1) and by the elliptic regularity theorem adapted to the space \mathcal{H} we have

$$\mathbb{E}\int_{\Phi(Y^+,\gamma)} |\nabla w_p^+(x,\gamma)|^s dx < \infty, \quad \text{which implies} \quad \mathbb{E}\int_Y |\nabla \tilde{w}_p^{\text{ext}}(y,\gamma)|^s dy < \infty$$

for some s > 2. In addition, $\nabla \tilde{w}_p^{\text{ext}}$ is stationary and its integral over Y has mean zero. By a version of Birkhoff's ergodic theorem, see e.g. Theorem 9 of [40], we have that

$$\lim_{\varepsilon \to 0} \sup_{x \in K} \left| \varepsilon \tilde{w}^{\text{ext}} \left(\frac{x}{\varepsilon}, \gamma \right) \right| = 0 \quad \mathbb{P}\text{-a.s}$$

for any compact set $K \subset \mathbb{R}^2$. The desired convergence result follows from the relation between w_p^{ext} and \tilde{w}_p^{ext} .

For the second convergence result, we first observe the following decomposition

$$w_{2p}^{\varepsilon} - x \cdot p = \varepsilon \left(w_p^-(\frac{x}{\varepsilon}) - w_p^{\text{ext}}\left(\frac{x}{\varepsilon}\right) \right) \chi_{\varepsilon \Phi(\mathbb{R}_2^-)} + \varepsilon w_p^{\text{ext}}\left(\frac{x}{\varepsilon}\right) \chi_{\varepsilon \Phi(\mathbb{R}_2^-)}.$$

By the proof of the first result, the second item on the right above converges uniformly in \mathscr{O} to zero and it suffices to show that $J_{\varepsilon} := \|\varepsilon w_p^-(\varepsilon^{-1}x) - \varepsilon w_p^{\text{ext}}(\varepsilon^{-1}x)\|_{L^2(\varepsilon \Phi(\mathbb{R}^-_2) \cap \mathscr{O})}$ converges to zero. Given \mathscr{O} and ε , we can find $\mathcal{I}_{\varepsilon}(\mathscr{O}) \subset \mathbb{Z}^2$ such that $\mathscr{O} \subset \bigcup_{k \in \mathcal{I}_{\varepsilon}} \varepsilon \Phi(Y_n)$ and $|\mathcal{I}_{\varepsilon}| \leq C(\mathscr{O})\varepsilon^{-2}$. Then a.s. in \mathscr{O} we verify that

$$J_{\varepsilon} \leq \sum_{n \in \mathcal{I}_{\varepsilon}} \int_{\varepsilon \Phi(Y_n^-)} \varepsilon^2 \left| w_p^{\text{ext}}\left(\frac{x}{\varepsilon}\right) - w_p^-\left(\frac{x}{\varepsilon}\right) \right|^2 dx = \varepsilon^4 \sum_{n \in \mathcal{I}_{\varepsilon}} \int_{\Phi(Y_n^-)} \left| w_p^{\text{ext}}(x) - w_p^-(x) \right|^2 dx$$
$$\leq C \varepsilon^4 \sum_{n \in \mathcal{I}_{\varepsilon}} \int_{Y_n^-} \left| \tilde{w}_p^{\text{ext}}(y) - \tilde{w}_p^-(y) \right|^2 dy.$$

In the last inequality, we used the change of variable $y = \Phi^{-1}(x)$ and the bounds (2.17) and (2.18). Using the estimate (C.4), we have

$$J_{\varepsilon} \leq C\varepsilon^{2} \left[\frac{1}{|\mathcal{I}_{\varepsilon}|} \sum_{n \in \mathcal{I}_{\varepsilon}} \left(\int_{\Gamma_{n}} \left| \tilde{w}_{p}^{+}(y) - \tilde{w}_{p}^{-}(y) \right|^{2} ds(y) + \int_{Y_{n}^{-}} \left| \nabla \tilde{w}_{p}^{\text{ext}}(y) - \nabla \tilde{w}_{p}^{-}(y) \right|^{2} dy \right) \right].$$

Note that the integrands above are stationary and the item inside the bracket is ready for applying ergodic theorem. This item converges to

$$\mathbb{E}\int_{\Gamma_0}|\tilde{w}_p^+-\tilde{w}_p^-|^2(y)ds(y)+\mathbb{E}\int_{Y^-}|\nabla\tilde{w}_p^{\text{ext}}-\nabla\tilde{w}_p^-|^2dy,$$

which is bounded for example by (7.6) and (7.8). Consequently, $J_{\varepsilon} \rightarrow 0$, proving (7.18).

For the third convergence result, we set first

$$Q\tilde{\eta}_p^{\pm} = (k_0[p + (\nabla\Phi)^{-1}\nabla\tilde{w}_p^{\pm}])\chi_{\mathbb{R}_2^{\pm}}.$$

These functions are stationary and we have the relation $Q\eta_p^{\varepsilon\pm} = (Q\tilde{\eta}_p^{\pm}) \left(\Phi^{-1}(\frac{x}{\varepsilon}, \gamma) \right)$ holds. By an ergodic theorem adapted to the stationary ergodic setting of this paper given in Lemma 2.2. of [19], we obtain (7.19).

7.2.2 **Proof of the homogenization theorem**

In this subsection we prove the homogenization theorem using Tartar's energy method. Here is the strategy: In the first step, we recall the energy estimates for the solution u_{ε} to the problem (2.2) and extract a subsequence along which $u_{\varepsilon}^{\text{ext}}$ converges weakly in $H^1(\Omega)$ to some u_0 , and the trivially extended gradient functions $Q\nabla u_{\varepsilon}^+$ and $Q\nabla u_{\varepsilon}^-$ has weak limits in $[L^2(\Omega)]^d$. Passing to limits in the weak formulation of (2.2), we obtain equations for these limits and the proper boundary conditions. In step three we identify u_0 as the unique solution to a homogenized equation. This is done by choosing the oscillating test functions $(\varphi \overline{w_{1p}^{\varepsilon}}, \varphi \overline{w_{2p}^{\varepsilon}})$ for the u_{ε} -equation and the oscillating test functions $(\varphi \overline{u_{\varepsilon}^+}, \varphi \overline{u_{\varepsilon}^-})$ for the w_p^{ε} -equation. The uniqueness of the solution to the weak formulation of u_0 relies on the fact that the trivial extension of u_{ε}^- converges weakly in $L^2(\Omega)$ to θu_0 for some constant $\theta < 1$. This fact is proved in step two.

In view of Proposition 3.3, Lemma 7.1 and Theorem 2.3, there exisit a faithful subset \mathcal{O}_1 of \mathcal{O} such that the conclusions in these results hold. We henceforth fix a $\gamma \in \mathcal{O}_1$ and ignore the dependence on γ .

Proof of Theorem 2.4. Step 1: Extraction of converging subsequences. Let $(u_{\varepsilon}^+, u_{\varepsilon}^-)$ be the solution to (2.2). In particular, u_{ε}^+ has an extension $u^{\text{ext}} \in H^1(\Omega)$. Let the vector fields ξ_{ε}^{\pm} be $k_0 \nabla u_{\varepsilon}^{\pm}$. Then the estimates (3.14) and (3.11) show that

$$\|u_{\varepsilon}^{\text{ext}}\|_{H^{1}(\Omega)} + \|Q\xi_{\varepsilon}^{+}\|_{[L^{2}(\Omega)]^{2}} + \|Q\xi_{\varepsilon}^{-}\|_{[L^{2}(\Omega)]^{2}} \leq C.$$

Consequently, there exists a subsequence and functions $u_0 \in H^1(\Omega)$ and $\xi_1, \xi_2 \in [L^2(\Omega)]^2$, such that

$$u_{\varepsilon}^{\text{ext}} \to u_{0} \text{ weakly in } H^{1}(\Omega), \qquad u_{\varepsilon}^{\text{ext}} \to u_{0} \text{ strongly in } L^{2}(\Omega);$$

$$Q\xi_{\varepsilon}^{+} \to \xi_{1} \text{ weakly in } [L^{2}(\Omega)]^{2}, \qquad Q\xi_{\varepsilon}^{-} \to \xi_{2} \text{ weakly in } [L^{2}(\Omega)]^{2}.$$
(7.20)

In the proof of Proposition 3.4, we also proved that

$$u_{\varepsilon}^{\text{ext}}\chi_{\varepsilon}^{-} - Qu_{\varepsilon}^{-} \to 0 \text{ strongly in } L^{2}(\Omega).$$
 (7.21)

Now fix an arbitrary test function $\varphi \in C_0^{\infty}(\Omega)$. Take $(\varphi \chi_{\varepsilon}^+, \varphi \chi_{\varepsilon}^-)$ as a test function in (3.5). Then the interface term disappears and we get

$$\int_{\Omega} k_0(Q\xi_{\varepsilon}^+) \cdot \overline{\nabla \varphi} dx + \int_{\Omega} k_0(Q\xi_{\varepsilon}^-) \cdot \overline{\nabla \varphi} dx = 0.$$

Passing to the limit $\varepsilon \to 0$ along the subsequence above, one finds

$$\int_{\Omega} (\xi_1 + \xi_2) \cdot \overline{\nabla \varphi} dx = 0, \quad \forall \varphi \in C_0^{\infty}(\Omega).$$
(7.22)

Therefore, the limiting vector field $\xi_1 + \xi_2$ satisfies that

$$\nabla \cdot (\xi_1 + \xi_2) = 0, \quad \text{in } \mathcal{D}'(\Omega), \tag{7.23}$$

where $\mathcal{D}'(\Omega)$ denotes the space of tempered distributions on Ω . Now for any $\phi \in C^{\infty}(\partial\Omega)$, we may lift it to a smooth function $\varphi \in C^{\infty}(\overline{\Omega})$ such that $\varphi = \phi$ on $\partial\Omega$. Take $(\varphi\chi_{\varepsilon}^{+}, \varphi\chi_{\varepsilon}^{-})$ as the test function in (3.5) and pass to the limit; we get

$$\int_{\Omega} (\xi_1 + \xi_2) \cdot \overline{\nabla \varphi} dx = \int_{\partial \Omega} g \overline{\phi} \, ds.$$
(7.24)

Since $\xi_1 + \xi_2 \in L^2(\Omega)$ and $\nabla \cdot (\xi_1 + \xi_2) \in L^2(\Omega)$, the trace $n \cdot (\xi_1 + \xi_2)$ on the boundary $\partial \Omega$ is well defined. Applying the divergence theorem and (7.23) we get

$$\int_{\partial\Omega} n \cdot (\xi_1 + \xi_2) \overline{\phi} \, ds = \int_{\partial\Omega} g \overline{\phi} \, ds, \quad \forall \phi \in C^{\infty}(\overline{\Omega}).$$

This shows that, $n \cdot (\xi_1 + \xi_2) = g$ at $\partial \Omega$. Further, since the trace of $Q\xi_{\varepsilon}^{-}$ is zero for all ε , the same argument above shows that $n \cdot \xi_2 = 0$ at $\partial \Omega$. We hence get

$$n \cdot \xi_1 = g$$
 at $\partial \Omega$.

Step 2: Weak convergence of Qu_{ε}^- . We can write Qu_{ε}^- as $u_{\varepsilon}^{\text{ext}}\chi_{\varepsilon}^- + (Qu_{\varepsilon}^- - u_{\varepsilon}^{\text{ext}}\chi_{\varepsilon}^-)$. Due to (7.21) and the fact that $u_{\varepsilon}^{\text{ext}}$ converges strongly to u_0 , we only need to verify that χ_{ε}^- converges weakly to θ . To this purpose, fix an arbitrary open set *K* compactly supported in Ω , and observe that for sufficiently small ε , *K* is compactly supported in E_{ε} defined in (2.20). Then we have

$$\int_{K} \chi_{\Omega_{\varepsilon}^{-}} dx = \int_{K \cap \varepsilon \Phi(\mathbb{R}_{2}^{-})} dx = \int_{\varepsilon \Phi^{-1}\left(\frac{K}{\varepsilon}\right)} \chi_{\mathbb{R}_{2}^{-}}\left(\frac{z}{\varepsilon}\right) \det \nabla \Phi\left(\frac{z}{\varepsilon}, \gamma\right) dz.$$

In [19, 18], it is shown that the characteristic function $\varepsilon \Phi^{-1}\left(\frac{K}{\varepsilon}\right)$ converges strongly in $L^1(\mathbb{R}^2)$ to that of the set $[\mathbb{E} \int_Y \nabla \Phi(y, \cdot) dy]^{-1} K$. On the other hand, since the function $\chi_{\mathbb{R}^2} \det \nabla \Phi$ is stationary, by ergodic theorem, we have

$$\chi_{\mathbb{R}_{2}^{-}}(\frac{z}{\varepsilon})\det\nabla\Phi(\frac{z}{\varepsilon},\gamma)\stackrel{*}{\rightharpoonup}\mathbb{E}\int_{Y}\chi_{\mathbb{R}_{2}^{-}}\det\nabla\Phi(z,\gamma)dz=\theta\varrho^{-1},\quad\text{in }L^{\infty}(\mathbb{R}^{2}).$$

Here, θ is defined as in (2.24). Consequently, we observe that for any open set *K* compactly supported in Ω , we have

$$\int \chi_K \chi_{\Omega_{\varepsilon}^-} dx \to \theta \varrho^{-1} \int_{[\mathbb{E} \int_Y \nabla \Phi(y, \cdot) dy]^{-1} K} dx = \theta \varrho^{-1} \det \left(\mathbb{E} \int_Y \nabla \Phi(y, \cdot) dy \right)^{-1} |K| = \theta |K|.$$

Here, we used the fact that det $(\mathbb{E} \int_Y \nabla \Phi(y, \cdot) dy) = \varrho^{-1}$, a fact also proved in [19, 18]. Since linear combinations of characteristic functions of compact sets in Ω are dense in $L^2(\Omega)$, we get the desired result. The fact that $\theta < 1$ is due to the fact that Φ does not expand the cells dramatically. This completes the proof of item two of the theorem up to a subsequence.

Step 3: Identifying the limit. Fix an arbitrary test function $\varphi \in C_0^{\infty}(\Omega)$. Recall the definitions of Ω_{ε}^- , K_{ε} and E_{ε} in (2.19) and (2.20). For sufficiently small ε , the function φ is compactly supported in E_{ε} .

Choose $p = e_k, k = 1, 2$ where $e_1 = (1, 0)$ and $e_2 = (0, 1)$. Let $w_{1e_k}^{\varepsilon}$ and $w_{2e_k}^{\varepsilon}$ be as in (7.15). In view of Remark 7.1, the weak formulation (7.16) still holds if we remove the integral over $\gamma \in \mathcal{O}$. Take $(\overline{\varphi u_{\varepsilon}^+}, \overline{\varphi u_{\varepsilon}^-})$ as a test function there; we get

$$\begin{split} \int_{\Omega} (Q\eta_{e_{k}}^{\varepsilon+}) \cdot \nabla(\overline{\varphi}u_{\varepsilon}^{+}) dx &+ \int_{\Omega} (Q\eta_{e_{k}}^{\varepsilon-}) \cdot \nabla(\overline{\varphi}u_{\varepsilon}^{-}) dx \\ &+ \frac{1}{\varepsilon\beta} \int_{\Gamma_{\varepsilon}} (w_{1e_{k}}^{\varepsilon} - w_{2e_{k}}^{\varepsilon}) \overline{\varphi}(u_{\varepsilon}^{+} - u_{\varepsilon}^{-}) ds = 0. \end{split}$$

Similarly, in the weak formulation (3.5), take $(\varphi \overline{w_{1e_{\iota}}^{\varepsilon}}, \varphi \overline{w_{2e_{\iota}}^{\varepsilon}})$ as the test function; we get

$$\begin{split} \int_{\Omega} (Q\xi_{\varepsilon}^{+}) \cdot \nabla(\overline{\varphi} w_{1e_{k}}^{\varepsilon}) dx &+ \int_{\Omega} (Q\xi_{\varepsilon}^{-}) \cdot \nabla(\overline{\varphi} w_{2e_{k}}^{\varepsilon}) dx \\ &+ \frac{1}{\varepsilon\beta} \int_{\Gamma_{\varepsilon}} (u_{\varepsilon}^{+} - u_{\varepsilon}^{-}) \overline{\varphi} (w_{1e_{k}}^{\varepsilon} - w_{2e_{k}}^{\varepsilon}) ds = 0. \end{split}$$

Note that the integrating domains in the first formula can be taken as above because φ is compactly supported in E_{ε} , which implies that $\varepsilon \Phi(\Gamma) \cap \operatorname{supp} \varphi = \Gamma_{\varepsilon} \cap \operatorname{supp} \varphi$. Subtracting the two formulas above and noticing in particular that the interface terms cancel out, we get

$$\begin{split} \int_{\Omega} (Q\eta_{e_{k}}^{\varepsilon+}) \cdot \overline{\nabla\varphi} u_{\varepsilon}^{\mathrm{ext}} dx + \int_{\Omega} (Q\eta_{e_{k}}^{\varepsilon-}) \cdot \overline{\nabla\varphi} u_{\varepsilon}^{\mathrm{ext}} dx + \int_{\Omega} (Q\eta_{e_{k}}^{\varepsilon-}) \cdot \overline{\nabla\varphi} (Qu_{\varepsilon}^{-} - u_{\varepsilon}^{\mathrm{ext}} \chi_{\varepsilon}^{-}) dx \\ - \int_{\Omega} (Q\xi_{\varepsilon}^{+}) \cdot \overline{\nabla\varphi} w_{1e_{k}}^{\varepsilon} dx + \int_{\Omega} (Q\xi_{\varepsilon}^{-}) \cdot \overline{\nabla\varphi} w_{2e_{k}}^{\varepsilon} dx = 0. \end{split}$$

By the convergence results (7.19), (7.17), (7.18), (7.20) and (7.21), we observe that each integrand above is a product of a strong converging term with a weak converging term. Therefore, we can pass the above to the limit $\varepsilon \rightarrow 0$ and get

$$\int_{\Omega} (\eta_{1e_k} + \eta_{2e_k}) u_0 \cdot \overline{\nabla \varphi} dx = \int_{\Omega} (\xi_1 + \xi_2) x_k \cdot \overline{\nabla \varphi} dx, \qquad (7.25)$$

where η_{1e_k} (resp. η_{2e_k}) is defined as the right-hand side of (7.19) with the "+" (resp. "-") sign. The integral on the right can be written as

$$\int_{\Omega} (\xi_1 + \xi_2) \cdot [\overline{\nabla(\varphi x_k)} - e_k \overline{\varphi}] dx = -\int_{\Omega} (\xi_1 + \xi_2) \cdot e_k \overline{\varphi} dx,$$

where we have used (7.22). For the integral involving $\eta_{1e_k} + \eta_{2e_k}$, we check that

$$e_i \cdot (\eta_{1e_k} + \eta_{2e_k}) = k_0 \varrho \mathbb{E} \int_{\Phi(Y)} \left(\chi_{\Phi(Y^+)} e_i \cdot \nabla w_{e_k}^+(x, \cdot) + \chi_{\Phi(Y^-)} e_i \cdot \nabla w_{e_k}^-(x, \cdot) + \delta_{ij} \right) dx.$$

This shows that

$$(\eta_{1e_k}+\eta_{2e_k})u_0\cdot\overline{\nabla\varphi}=\sum_{j=1}^2e_j\cdot(\eta_{1e_k}+\eta_{2e_k})u_0\frac{\partial\overline{\varphi}}{\partial x_j}=K_{kj}^*u_0\frac{\partial\overline{\varphi}}{\partial x_j},$$

where we have used the definition of the matrix (K_{ij}^*) in (2.26). Now (7.25) becomes

$$\int_{\Omega} \sum_{j=1}^{2} K_{kj}^{*} u_0 \frac{\partial \overline{\varphi}}{\partial x_j} dx = -\int_{\Omega} (\xi_1 + \xi_2) \cdot e_k \, \overline{\varphi} dx.$$

Since $\varphi \in C_0^{\infty}(\Omega)$) is arbitrary we conclude that

$$(\xi_1 + \xi_2) \cdot e_k = \sum_{j=1}^2 K_{kj}^* \frac{\partial u_0}{\partial x_j}, \quad \text{for all } k.$$

Substitute this relation in (7.22) and (7.24); one obtains

$$\int_{\Omega} K^* \nabla u_0 \cdot \nabla \overline{\varphi} dx = \int_{\partial \Omega} g \overline{\varphi}, \quad \text{for all } \varphi \in H^1(\Omega)).$$
(7.26)

Finally, we recall that for all $\gamma \in \mathcal{O}_1$,

$$\int_{\Omega} Qu_{\varepsilon}^{+}(x,\gamma)dx = 0, \text{ and } Qu_{\varepsilon}^{-}(\cdot,\gamma) \rightharpoonup \theta u_{0}(\cdot,\gamma) \text{ weakly in } L^{2}(\Omega)$$

indicate that

$$\int_{\Omega} u_0(x,\gamma) dx = \lim_{\varepsilon \to 0} \int_{\Omega} (Qu_{\varepsilon}^+(x,\gamma) + Qu_{\varepsilon}^-(x,\gamma)) dx$$
$$= \lim_{\varepsilon \to 0} \int_{\Omega} Qu_{\varepsilon}^-(x,\gamma) dx = \theta \int_{\Omega} u_0(x,\gamma) dx.$$

Since $\theta < 1$, we obtain

$$\int_{\Omega} u_0(x,\gamma) dx = 0. \tag{7.27}$$

In summary, the weak limit $u_0(x, \gamma)$ provides a solution to the problem (7.26)(7.27). Thanks to this normalization condition and the fact that K^* is uniformly elliptic, the proof of which is not difficult and is omitted, the solution to this problem is unique.

We check that the unique deterministic solution to the homogenized equation (2.25) solves the problem (7.26)(7.27). By uniqueness of the latter problem, we conclude that $u_0(x, \gamma)$ obtained in step one for all $\gamma \in O_1$ must agree with the deterministic solution to (2.25). We denote this solution as $u_0(x)$. Consequently, all converging subsequences of $(u_{\varepsilon}^+, u_{\varepsilon}^-)$ converge to $u_0(x)$ and hence the whole sequence converges to this limit. This completes the proof.

7.3 Effective admittivity of a dilute suspension

In this subsection, we consider the case when the cells are dilute. We aim to derive a formal first order asymptotic expansion of the effective admittivity in terms of the volume fraction of the dilute cells.

In the formula of the homogenized coefficient (2.26), the integral term has the form

$$J_{ij} = \mathbb{E} \int_{\Phi(Y^+)} e_j \cdot \nabla w_{e_i}^+(y, \cdot) \, dy + \mathbb{E} \int_{\Phi(Y^-)} e_j \cdot \nabla w_{e_i}^-(y, \cdot) \, dy.$$

Thanks to the ergodic theorem, J_{ij} also takes the form

$$J_{ij} = \lim_{N \to \infty} \frac{1}{N^2} \sum_{n \in \mathcal{I}(N)} \left(\int_{\Phi(Y_n^+)} e_j \cdot \nabla w_{e_i}^+(y, \cdot) \, dy + \int_{\Phi(Y_n^-)} e_j \cdot \nabla w_{e_i}^-(y, \cdot) \, dy \right).$$

Here, $\mathcal{I}(N)$ is the indices for the cubes $\{Y_n\}$ inside the big cube *NY*. Now using integration by parts, we simplify the above expression to

$$J_{ij} = \lim_{N \to \infty} \frac{1}{N^2} \left(\int_{\partial \Phi(NY)} n_j w_{e_i}^+(y, \cdot) \, ds(y) - \sum_{n \in \mathcal{I}(N)} \int_{\Gamma_n} (w_{e_i}^+ - w_{e_i}^-)(y, \cdot) n_j \, ds(y) \right).$$

Here, *n* denotes the outer normal vector along the boundary of $\Phi(NY)$ and $\Phi(Y_n^-)$, $n \in \mathcal{I}(N)$; $n_j = n \cdot e_j$ denotes its *j*-th component. Note that the boundary terms at $\{\partial \Phi(Y_n)\}_{n \in \mathcal{I}(N)} \cap \Phi(NY)$ are canceled because two adjacent cubes share the same outer normal vector at their common boundary except for reversed signs.

Finally, we have seen that $w_{e_i}^+$ has sub-linear growth. Since the surface $\Phi(NY)$ has volume of order O(N), the sub-linear growth indicates that the boundary integral at $\partial \Phi(NY)$ is of order $o(N^2)$. Consequently, when divided by N^2 this term goes to zero. By applying the ergodic theorem again, we obtain that

$$J_{ij} = \lim_{N \to \infty} \frac{1}{N^2} \sum_{n \in \mathcal{I}(N)} \int_{\Gamma_n} (w_{e_i}^- - w_{e_i}^+)(y, \cdot) n_j \, ds(y) = \mathbb{E} \int_{\partial \Phi(Y^-)} (w_{e_i}^- - w_{e_i}^+)(y, \cdot) n_j \, ds(y).$$
(7.28)

In the following, we investigate this integral further by deriving a formal representation for the jump $w_{e_i}^+ - w_{e_i}^-$ in the case when the inclusions are dilute, *i.e.*, small and far away from each other.

To model the dilute suspension, we assume that the reference cell Y^- is of the form ρB , where *B* is a domain of unit length scale and unit volume, and $\rho := \sqrt{|Y^-|} \ll 1$ denotes the

small length scale of the dilute inclusions. Due to the assumptions (2.17) and (2.18), the length scale of the cell $\Phi(Y^-)$ is still of order ρ . Further, by our construction the distance of the cell $\Phi(Y^-)$ from the "boundary" $\partial \Phi(Y)$ is of order one, which is much larger than the size of the inclusion.

Since the distances between the inclusions are much larger than their sizes, we may use the single inclusion approximation. That is, $w_{e_i}^{\pm}$ can be approximated by the solutions to the following interface problem:

$$\begin{cases} \nabla \cdot k_0 \nabla w_{e_i}^{\pm} = 0 \text{ in } \Phi(Y^-) \text{ and } \mathbb{R}^2 \setminus \Phi(Y^-), \\ \frac{\partial w_{e_i}^+}{\partial n} = \frac{\partial w_{e_i}^-}{\partial n}, \text{ and } w_{e_i}^+ - w_{e_i}^- = \rho \beta k_0 (\frac{\partial w_{e_i}^-}{\partial n} + n \cdot e_i) \text{ on } \Phi(\Gamma), \\ w_{e_i}^+ \to 0 \text{ at } \infty. \end{cases}$$

Here, $\Phi(\Gamma)$ denotes the boundary of the inclusion. Note that the extra ρ in the jump condition is due to the fact that the length scale of the inclusion $\Phi(\Upsilon^-)$ is of order ρ . Using double layer potentials, we represent $w_{e_i}^+$ and $w_{e_i}^-$ as $\mathcal{D}_{\Phi(\Gamma)}[\phi_i]$ restricted to $\Phi(\Upsilon^-)$ and $\mathbb{R}^2 \setminus \Phi(\Upsilon^-)$ respectively. Due to the trace formula of $\mathcal{D}_{\Phi(\Gamma)}$ and the jump conditions above, the function ϕ_i is determined by

$$-\phi_i = \rho \beta k_0 \left(\frac{\partial \mathcal{D}_{\Phi(\Gamma)}[\phi_i]}{\partial n} + n_i\right). \tag{7.29}$$

Let us define the operator $\mathcal{L}_{\Phi(\Gamma)}$ by $\frac{\partial \mathcal{D}_{\Phi(\Gamma)}}{\partial n}$, then we have that

$$w_{e_i}^+ - w_{e_i}^- = -\phi_i = \rho \beta k_0 (I + \rho \beta k_0 \mathcal{L}_{\Phi(\Gamma)})^{-1} [n_i], \text{ on } \Phi(\Gamma).$$

As a consequence, we have also that

$$J_{ij} \simeq -\rho \beta k_0 \mathbb{E} \int_{\Phi(\Gamma)} (I + \rho \beta k_0 \mathcal{L}_{\Phi(\Gamma)})^{-1} [n_i] n_j ds.$$

Let us define ψ_i to be $-(I + \rho\beta k_0 \mathcal{L}_{\Phi(\Gamma)})^{-1}[n_i]$, that is $\psi_i + \rho\beta k_0 n \cdot \nabla \mathcal{D}_{\Phi(\Gamma)}[\psi_i](x) = -n_i$. Define the scaled function $\tilde{\psi}_i(\tilde{x}) = \psi_i(\rho\tilde{x})$ on the scaled curve $\rho^{-1}\Phi(\Gamma)$. Using the homogeneity of the gradient of the Newtonian potential, we verify that

$$\mathcal{D}_{\Phi(\Gamma)}[\psi_i](x) = \mathcal{D}_{\rho^{-1}\Phi(\Gamma)}[\tilde{\psi}_i](\tilde{x}), \quad \text{and} \quad \rho n \cdot \nabla \mathcal{D}_{\Phi(\Gamma)}[\psi_i](x) = n \cdot \nabla \mathcal{D}_{\rho^{-1}\Phi(\Gamma)}[\tilde{\psi}_i](\tilde{x}),$$

where $\tilde{x} = \rho^{-1}x$. This shows that $\tilde{\psi}_i = -(I + \beta k_0 \mathcal{L}_{\rho^{-1}\Phi(\Gamma)})^{-1}[n_i]$. Using the change of variable $y \to \rho \tilde{y}$ in the previous integral representation of J_{ij} , we rewrite it as

$$J_{ij} \simeq \rho \beta k_0 \mathbb{E} \int_{\rho^{-1} \Phi(\Gamma)} \psi_i(\rho \tilde{y}) n_j ds(\rho \tilde{y}) = \rho^2 \beta k_0 \mathbb{E} \int_{\rho^{-1} \Phi(\Gamma)} \tilde{\psi}_i n_j ds(\tilde{y}).$$

Finally, the approximation (2.27) of the effective permittivity for the dilute suspension holds, where $f = \rho \rho^2$ is the volume fraction where ρ accounts for the averaged change of volume due to the random diffeomorphism; the polarization matrix *M* is defined by (2.28) and is associated to the deformed inclusion scaled to the unit length scale. Note that the imaging approach developed in subsection 6.2 can be applied here as well.

8 Numerical simulations

We present in this section some numerical simulations to illustrate the fact that the Debye relaxation times are characteristics of the microstructure of the tissue.

We take realistic values for our parameters, which are the same as those used in Subsection 2.2 and let the frequency $\omega \in [10^4, 10^9]$ Hz.

We first want to retrieve the invariant properties of the Debye relaxation times. We consider (Figure 8.1) an elliptic cell (in green) that we translate (to obtain the red one), rotate (to obtain the purple one) and scale (to obtain the dark blue one). We compute the membrane polarization tensor, its imaginary part, and the associated eigenvalues which are plotted as a function of the frequency (Figure 8.2). The frequency is here represented on a logarithmic scale. One can see that for the two eigenvalues the maximum of the curves occurs at the same frequency, and hence that the Debye relaxation times are identical for the four elliptic cells. Note that the red and green curves are even superposed; this comes from the fact that M is invariant by translation.

Next, we are interested in the effect of the shape of the cell on the Debye relaxation times. We consider for this purpose, (Figure 8.3) a circular cell (in green), an elliptic cell (in red) and a very elongated elliptic cell (in blue). We compute similarly as in the preceding case, the polarization tensors associated to the three cells, take their imaginary part and plot the two eigenvalues of these imaginary parts with respect to the frequency. As shown in Figure 8.4, the maxima occur at different frequencies for the first and second eigenvalues. Hence, we can distinguish with the Debye relaxation times between these three shapes.

Finally, we study groups of one (in green), two (in blue) and three cells (in red) in the unit period (Figure 8.5) and the corresponding polarization tensors for the homogenized media. The associated relaxation times are different in the three configurations (Figure 8.6) and hence can be used to differentiate tissues with different cell density or organization.

These simulations prove that the Debye relaxation times are characteristics of the shape and organization of the cells. For a given tissue, the idea is to obtain by spectroscopy the frequency dependence spectrum of its effective admittivity. One then has access to the membrane polarization tensor and the spectra of the eigenvalues of its imaginary part. One compares the associated Debye relaxation times to the known ones of healthy and cancerous tissues at different levels. Then one would be able to know using statical tools with which probability the imaged tissue is cancerous and at which level.

9 Concluding remarks

In this paper, we have explained how the dependence of the effective electrical admittivity measures the complexity of the cellular organization of the tissue. We have derived new formulas for the effective admittivity of suspensions of cells and characterized their dependence with respect to the frequency in terms of membrane polarization tensors. We have applied the formulas in the dilute case to image suspensions of cells from electrical boundary measurements. We have presented numerical results to illustrate the use of the Debye relaxation time in classifying microstructures. We also developed a selective spectroscopic imaging approach. We have shown that specifying the pulse shape in terms of the relaxation times of the dilute suspensions gives rise to selective imaging.



Figure 8.1: An ellipse translated, rotated and scaled.



Figure 8.2: Frequency dependence of the eigenvalues of $\Im M$ for the 4 ellipses in Figure 8.1.



Figure 8.3: A circle, an ellipse and a very elongated ellipse.



Figure 8.4: Frequency dependence of the eigenvalues of $\Im M$ for the 3 different cell shapes in Figure 8.3.



Figure 8.5: Groups of one, two and three cells.



Figure 8.6: Frequency dependence of the eigenvalues of *SM* in the 3 different cases.

An important problem is to investigate the frequency dependence of the effective electrical properties of dense periodic arrangements of cells such as skin cells. Our classification approach proposed in this paper is expected to be applicable for nondilute suspensions but at given volume fraction. Another challenging problem is to extend our results to elasticity models of the cell. In [10, 8], formulas for the effective shear modulus and effective viscosity of dilute suspensions of elastic inclusions were derived. On the other hand, it was observed experimentally that the dependence of the viscosity of a biological tissue with respect to the frequency characterizes the microstructure [16, 24]. A mathematical justification and modeling for this important finding are under investigation and would be the subject of a forthcoming paper.

Finally, it would be very interesting to develop a physics-based learning approach, based on Debye relaxation times, for classifying tissue organizations at the cell scale from macroscopic spectroscopic admittivity measurements. One can learn from training examples such as biopsies the underlying microstructures and then, classify unseen ones from spectroscopic measurements of their admittivities. For doing so, it is important to construct a distance between spectroscopic measurements which allows to statistically classify or separate different microstructures into different groups.

A Extension lemmas

Due to the problem settings of this paper, we need to study convergence properties of functions that are defined on the multiple connected sets \mathbb{R}_2^+ , $\Phi(\mathbb{R}_2^+)$ and $\varepsilon \Phi(\mathbb{R}_2^-)$. Extension operators becomes useful to treat such functions.

Consider two open sets $U, V \subset \mathbb{R}^2$ with the relation $U \subset V$, and two Sobolev spaces $W^{1,p}(U)$ and $W^{1,p}(V)$, $p \in [1, \infty]$. What we call an *extension operator* is a bounded linear map $P: W^{1,p}(U) \to W^{1,p}(V)$, such that Pu = u a.e. on U for all $u \in W^{1,p}(U)$. In this section, we introduce several extension operators of this kind that are needed in the paper. They extend functions that are defined on Y^- , \mathbb{R}_2^+ , $\Phi(\mathbb{R}_2^+)$ and $\varepsilon \Phi(\mathbb{R}_2^+)$ (hence Ω_{ε}^+) respectively.

Throughout this section, the short hand notion $\mathcal{M}_A(f)$ for a measurable set $A \subset \mathbb{R}^2$ with positive volume and a function $f \in L^1(A)$ denotes the mean value of f in A, that is

$$\mathcal{M}_A(f) = \frac{1}{|A|} \int_A f(x) dx. \tag{A.1}$$

We start with an extension operator inside the unit cube *Y*. Since Y^- has smooth boundary, there exists an extension operator $S : W^{1,p}(Y^+) \to W^{1,p}(Y)$ such that for all $f \in W^{1,p}(Y^+)$ and $p \in [1, \infty)$,

$$\|Sf\|_{L^{p}(Y)} \leq C \|f\|_{L^{p}(Y^{+})}, \quad \|Sf\|_{W^{1,p}(Y)} \leq C \|f\|_{W^{1,p}(Y^{+})},$$
(A.2)

where *C* only depends on *p* and Y^- . Such an *S* is given in [27, section 5.4], where the second estimate above is given; the first estimate easily follows from their construction as well. Cioranescu and Saint Paulin [22] constructed another extension operator which refines the second estimate above. For the reader's convenience, we state and prove their result in the following. Similar results can be found in [34] as well.

Theorem A.1. Let Y, Y^+ and Y^- be as defined in section 2; in particular, ∂Y^- is smooth. Then there exists an extension operator $P : W^{1,p}(Y^+) \to W^{1,p}(Y)$ satisfying that for any $f \in W^{1,p}(Y^+)$ and

 $p \in [1, \infty)$,

$$\|\nabla Pf\|_{L^{p}(Y)} \leq C \|\nabla f\|_{L^{p}(Y^{+})}, \quad \|Pf\|_{L^{p}(Y)} \leq C \|f\|_{L^{p}(Y^{+})},$$
(A.3)

where C only depends on the dimension and the set Y^- .

Proof. Recall the mean operator \mathcal{M} in (A.1) and the extension operator S in (A.2). Given f, we define Pf by

$$Pf = \mathcal{M}_{Y^{+}}(f) + S(f - \mathcal{M}_{Y^{+}}(f)).$$
(A.4)

Then by setting $\psi = f - \mathcal{M}_{Y^+}(f)$, we have that

$$\|\nabla Pf\|_{L^{p}(Y)} = \|\nabla S\psi\|_{L^{p}(Y)} \le C\|\psi\|_{W^{1,p}(Y^{+})} \le C\|\nabla\psi\|_{L^{p}(Y^{+})} = C\|\nabla f\|_{L^{p}(Y^{+})}.$$

In the second inequality above, we used the Poincaré–Wirtinger inequality for ψ and the fact that ψ is mean-zero on Y^+ . The L^2 bound of Pf follows from the observation

$$\|\mathcal{M}_{Y^+}(f)\|_{L^p(Y)} \le \left(\frac{|Y|}{|Y^+|}\right)^{\frac{1}{p}} \|f\|_{L^p(Y^+)}$$

and the L^p estimate of Sf in (A.2). This completes the proof.

Apply the extension operator on each translated cubes in \mathbb{R}_2^+ , we get the following.

Corollary A.1. Recall the definition of Y_n , Y_n^+ and Y_n^- in section 2. Abuse notations and define

$$(Pu)|_{Y_n} = P(u|_{Y_n^+}), \quad n \in \mathbb{Z}^2, u \in W^{1,p}_{\text{loc}}(\mathbb{R}_2^+).$$
 (A.5)

Then P is an extension operator from $W_{loc}^{1,p}(\mathbb{R}_2^+)$ to $W_{loc}^{1,p}(\mathbb{R}^2)$. Further, with the same positive constant C in (A.3) and for any $n \in \mathbb{Z}^2$, we have

$$\|\nabla Pu\|_{L^{p}(Y_{n})} \leq C \|\nabla u\|_{L^{p}(Y_{n}^{-})}, \quad \|Pu\|_{L^{p}(Y_{n})} \leq C \|u\|_{L^{p}(Y_{n}^{-})}.$$
(A.6)

Given a diffeomorphism, the extension operator *P* can be transformed as follows. In the same manner, under the map of scaling, the extension operator is naturally defined.

Corollary A.2. Let $\Phi(\cdot, \gamma)$ be a random diffeomorphism satisfying (2.17) and (2.18). Denote the inverse function Φ^{-1} by Ψ . Define P_{γ} as

$$P_{\gamma}u = [P(u \circ \Phi)] \circ \Psi, \quad u \in W^{1,p}_{\text{loc}}(\Phi(\mathbb{R}^+_2)).$$
(A.7)

Then P_{γ} is an extension operator from $W^{1,p}_{loc}(\Phi(\mathbb{R}^+_2))$ to $W^{1,p}_{loc}(\Phi(\mathbb{R}^2))$ which satisfies that

$$\|\nabla P_{\gamma}u\|_{L^{p}(\Phi(Y_{n}))} \leq C \|\nabla u\|_{L^{p}(\Phi(Y_{n}^{-}))}, \quad \|P_{\gamma}u\|_{L^{p}(\Phi(Y_{n}))} \leq C \|u\|_{L^{p}(\Phi(Y_{n}^{-}))}, \quad (A.8)$$

where the constant C depends further on the constants in (2.17) and (2.18).

Corollary A.3. Let $\Phi(\cdot, \gamma)$ and Ψ be as above. For each $\varepsilon > 0$, define P_{γ}^{ε} as follows: for any $u \in W_{\text{loc}}^{1,p}(\varepsilon \Phi(\mathbb{R}^+_2))$, $P_{\gamma}^{\varepsilon}u$ is defined on each deformed and scaled cube $\varepsilon \Phi(Y_n)$ by

$$P^{\varepsilon}_{\gamma}u(x) = \varepsilon P\tilde{u}(\Psi(\frac{x}{\varepsilon})), \tag{A.9}$$

where $\tilde{u} = \varepsilon^{-1} u \circ \varepsilon \Phi$ and P is as in (A.6). Then P_{γ}^{ε} is an extension operator from $W_{loc}^{1,p}(\varepsilon \Phi(\mathbb{R}^{2}))$ to $W_{loc}^{1,p}(\varepsilon \Phi(\mathbb{R}^{2}))$ which satisfies that for any $n \in \mathbb{Z}^{2}$,

$$\|\nabla P_{\gamma}^{\varepsilon}u\|_{L^{p}(\varepsilon\Phi(Y_{n}))} \leq C \|\nabla u\|_{L^{p}(\varepsilon\Phi(Y_{n}^{-}))}, \quad \|P_{\gamma}^{\varepsilon}u\|_{L^{p}(\varepsilon\Phi(Y_{n}))} \leq C \|u\|_{L^{p}(\varepsilon\Phi(Y_{n}^{-}))}, \quad (A.10)$$

where the constant C depends on the same parameters as stated below (A.8).

Proof. We focus on proving (A.10). Under the change of variable $x = \varepsilon \Phi(y)$, we have

$$\nabla_{x} P^{\varepsilon}_{\gamma} u(x) = \nabla \Psi(\frac{x}{\varepsilon}) \nabla_{y} P \tilde{u}(\Phi^{-1}(\frac{x}{\varepsilon})) = \nabla \Psi(\Phi(y)) \nabla_{y} P \tilde{u}(y),$$

On each deformed and scaled cube $\varepsilon \Phi(Y_n)$, we calculate

$$\begin{split} \|\nabla P_{\gamma}^{\varepsilon}u\|_{L^{p}(\varepsilon\Phi(Y_{n}))}^{p} &= \int_{\varepsilon\Phi(Y_{n})} |\nabla_{x}P_{\gamma}^{\varepsilon}u(x)|^{p}dx = \int_{Y_{n}} |\nabla\Psi(\Phi(y))\nabla_{y}P\tilde{u}(y)|^{p}\varepsilon^{2} \det(\nabla\Phi(y))dy \\ &\leq \varepsilon^{2}\int_{Y_{n}} |\nabla\Psi(\Phi(y))|^{p}|\nabla_{y}P\tilde{u}(y)|^{p} \det(\nabla\Phi(y))dy \leq C\varepsilon^{2}\int_{Y_{n}} |\nabla_{y}P\tilde{u}(y)|^{p}dy. \end{split}$$

Here, we have used the Cauchy–Schwarz inequality and the bounds (2.17)-(2.18) on the Jacobian matrix and its determinant. Upon applying (A.3), we get

$$\|\nabla P^{\varepsilon}_{\gamma} u\|_{L^{p}(\varepsilon \Phi(Y_{n}))}^{p} \leq C \varepsilon^{2} \|\nabla_{y} \tilde{u}\|_{L^{p}(Y_{n}^{+})}^{p}.$$

Since $\tilde{u}(y) = \frac{1}{\varepsilon}u(\varepsilon\Phi(y))$, we have $\nabla_y \tilde{u}(y) = \nabla_y \Phi(y) \nabla_x u(\varepsilon\Phi(y))$. Change variables in the last integral and repeat the analysis above to get

$$\|\nabla_y \tilde{u}\|_{L^p(Y_n^+)}^p \le C\varepsilon^{-d} \|\nabla_x u\|_{L^p(\varepsilon\Phi(Y_n^+))}^p$$

Combining the above estimates, one finds some *C* independent of ε or γ such that (A.10) holds. Moreover, the constant *C* is uniform for all $\varepsilon \Phi(Y_n)$. The L^2 estimate for $P_{\gamma}^{\varepsilon}u$ is simpler and ignored. This completes the proof.

Finally, we define the extension operator from $W^{1,p}(\Omega_{\varepsilon}^+)$ to $W^{1,p}(\Omega)$. This is essentially the same operator in Corollary A.3. Indeed, recall that Ω is decomposed to the cushion K_{ε} and the cell containers E_{ε} ; see (2.20). We only need to apply P_{γ}^{ε} in E_{ε} .

Theorem A.2. Let the domains $\Omega_{\varepsilon}^{\pm}$, K_{ε} and E_{ε} be as defined in section 2. Let $\Phi(\cdot, \gamma)$ be a random diffeomorphism satisfying (2.17)-(2.18). Define the linear operator P_{γ}^{ε} as follows: for $u \in W^{1,p}(\Omega_{\varepsilon}^{+})$, let $P_{\gamma}^{\varepsilon}u$ be given by (A.9) in E_{ε} , and let $P_{\gamma}^{\varepsilon}u = u$ in K_{ε} . Then P_{γ}^{ε} is an extension operator from $W^{1,p}(\Omega_{\varepsilon}^{+})$ to $W^{1,p}(\Omega)$ and it satisfies

$$\|\nabla P^{\varepsilon}_{\gamma} u\|_{L^{p}(\Omega)} \leq C \|\nabla u\|_{L^{p}(\Omega^{+}_{\varepsilon})}, \quad \|P^{\varepsilon}_{\gamma} u\|_{L^{p}(\Omega)} \leq C \|u\|_{L^{p}(\Omega^{+}_{\varepsilon})}, \tag{A.11}$$

where the constants C's do not depend on ε or γ .

Proof. Since P_{γ}^{ε} leaves *u* unchanged in K_{ε} and it satisfies the estimates (A.10) uniformly in the cubes $E_{\varepsilon} = \bigcup_{n \in \mathcal{I}_{\varepsilon}} \varepsilon \Phi(Y_n)$, we have the following:

$$\begin{aligned} \|\nabla P^{\varepsilon}_{\gamma} u\|^{p}_{L^{p}(\Omega)} &= \|\nabla u\|^{p}_{L^{p}(K_{\varepsilon})} + \sum_{n \in \mathcal{I}_{\varepsilon}} \|\nabla P^{\varepsilon}_{\gamma} u\|^{p}_{L^{p}(\varepsilon \Phi(Y_{n}))} \\ &\leq \|\nabla u\|^{p}_{L^{p}(K_{\varepsilon})} + C \sum_{n \in \mathcal{I}_{\varepsilon}} \|\nabla u\|^{p}_{L^{p}(\varepsilon \Phi(Y_{n}^{+}))} \leq C \|\nabla u\|^{p}_{L^{p}(\Omega^{+}_{\varepsilon})}.\end{aligned}$$

This completes the proof of the first estimate in (A.11). The second estimate follows in the same manner, completing the proof. \Box

B Poincaré–Wirtinger inequality

Our next goal is to derive a Poincaré–Wirtinger inequality for functions in $H^1(\Omega_{\varepsilon}^+)$ with a constant independent of ε and γ . The following fact of the fluctuation of a function is useful.

Lemma B.1. Let $X \subset \mathbb{R}^2$ be an open bounded domain with positive volume and $f \in L^1(X)$. Assume that $X_1 \subset X$ is a subset with positive volume, then we have

$$\|f - \mathcal{M}_{X_1}(f)\|_{L^2(X_1)} \le \|f - \mathcal{M}_X(f)\|_{L^2(X)}.$$
(B.1)

Proof. To simplify notations, let f_1 be the restriction of f on X_1 , $m_1 = \mathcal{M}_{X_1}(f_1)$ and $\theta_1 = |X_1|/|X|$. Similarly, let f_2 be the restriction of f on $X_2 = X \setminus X_1$, $m_2 = \mathcal{M}_{X_2}(f_2)$. Let $m = \mathcal{M}_X(f)$. Then we have that

$$f - m = \begin{cases} f_1 - m_1 + (1 - \theta)(m_1 - m_2), & x \in X_1, \\ f_2 - m_2 + \theta(m_2 - m_1), & x \in X_2. \end{cases}$$

Then basic computation plus the observation that $f_i - m_i$ integrates to zero on X_i for i = 1, 2 yield the following:

$$\|f - m\|_{L^{2}(X)}^{2} = \|f_{1} - m_{1}\|_{L^{2}(X_{1})}^{2} + \|f_{2} - m_{2}\|_{L^{2}(X_{2})}^{2} + (1 - \theta)\theta|X|(m_{2} - m_{1})^{2}.$$

Since the items on the right-hand side are all non-negative, we obtain (B.1).

Corollary B.1. Assume the same conditions as in Theorem A.2. Then for any $u \in H^1_{\mathbb{C}}(\Omega^+_{\varepsilon})$, we have that

$$\|u\|_{L^2(\Omega_{\varepsilon}^+)} \le C \|\nabla u\|_{L^2(\Omega_{\varepsilon}^+)'}$$
(B.2)

where the constant C does not depend on ε or γ .

Proof. Thanks to Theorem A.2, we extend *u* to $P^{\varepsilon}_{\gamma}u$ which is in $H^1(\Omega)$. Use (B.1) and the fact that $\mathcal{M}_{\Omega^+_{\varepsilon}}(u) = 0$ to get

$$\|u\|_{L^2(\Omega^+_{\varepsilon})} \leq \|P^{\varepsilon}_{\gamma}u - \mathcal{M}_{\Omega}(P^{\varepsilon}_{\gamma}u)\|_{L^2(\Omega)}.$$

Now apply the standard Poincaré–Wirtinger inequality for functions in $H^1(\Omega)$, and then use (A.11). We get

$$P_{\gamma}^{\varepsilon}u - \mathcal{M}_{\Omega}(P_{\gamma}^{\varepsilon}u)\|_{L^{2}(\Omega)} \leq C \|\nabla P_{\gamma}^{\varepsilon}u\|_{L^{2}(\Omega)} \leq C \|\nabla u\|_{L^{2}(\Omega_{\varepsilon}^{+})}$$

The constant *C* depends on Ω and the parameters stated in Theorem A.2 but not on ε or γ . The proof is now complete.

Another corollary of the extension lemma is that we have the following uniform estimate when taking the trace of $u \in W_{\varepsilon}$ on the fixed boundary $\partial \Omega$.

Corollary B.2. Assume the same conditions as in Theorem A.2. Then there exists a constant C depending on Ω and the parameters as stated in Theorem A.2 but independent of ε and γ such that

$$\|u\|_{H^{\frac{1}{2}}(\partial\Omega)} \le C \|\nabla u\|_{L^{2}(\Omega^{+}_{\varepsilon})},$$
(B.3)

for any $u \in H^1(\Omega_{\varepsilon}^+)$.

Proof. Thanks to Theorem A.2 we extend *u* to $P_{\gamma}^{\varepsilon}u$ which is in $H^1(\Omega)$. The trace inequality on Ω shows

$$\|P_{\gamma}^{\varepsilon}u\|_{H^{\frac{1}{2}}(\partial\Omega)} \le C(\Omega)\|P_{\gamma}^{\varepsilon}u\|_{H^{1}(\Omega)}.$$
(B.4)

The desired estimate then follows from (A.11) and (B.2).

C Equivalence of the two norms on W_{ε}

In this section, we prove Proposition 3.2 which establishes the equivalence between the two norms on W_{ε} . We essentially follow [45] where the periodic case was considered. The random deformation setting requires certain modification. The details of such modifications are provided here for the reader's convenience.

The first inequality of the proposition is proved by the following lemma together with the Poincaré–Wirtinger inequality (B.2):

Lemma C.1. There exists a constant C independent of ε or γ , such that

$$\|v^{\pm}\|_{L^{2}(\Gamma_{\varepsilon})}^{2} \leq C(\varepsilon^{-1}\|v^{\pm}\|_{L^{2}(\Omega_{\varepsilon}^{\pm})}^{2} + \varepsilon\|\nabla v^{\pm}\|_{L^{2}(\Omega_{\varepsilon}^{\pm})}^{2})$$
(C.1)

for any $v^+ \in H^1(\Omega^+_{\varepsilon})$ and $v^- \in H^1(\Omega^-_{\varepsilon})$.

Proof. According to the set-up, the interface Γ_{ε} consists of $\varepsilon \Phi(\Gamma_i)$ where $i = 1, \dots, N(\varepsilon)$ are the labels for the deformed cubes $\{\varepsilon \Phi(Y_i)\}$ inside Ω and Γ_i are the corresponding unit scale interfaces.

Let us consider the case of $v^+ \in H^1(\Omega_{\varepsilon}^+)$; the other case is proved in the same manner. Denote by v_i the restriction of v^+ on the deformed cube $\varepsilon \Phi(Y_i)$. We lift this function to $\tilde{v}_i(y) = v_i(\varepsilon \Phi(y))$ which is now defined on Y_i^+ . For this function, we have the trace inequality

$$\|\tilde{v}_i\|_{L^2(\Gamma_i)}^2 \le C(\|\tilde{v}_i\|_{L^2(Y_i^+)}^2 + \|\nabla\tilde{v}_i\|_{L^2(Y_i^+)}^2).$$
(C.2)

Note that this constant depends on the reference shape Y^- but is uniform in *i*.

On the other hand, because for any $\gamma \in O$, the diffeomorphism Φ satisfies (2.17) and (2.18), the Lebesgue measures ds(x) on the curve $\varepsilon \Phi(\Gamma_i)$ and ds(y) on Γ_i , which are related by the change of variable $x = \varepsilon \Phi(y)$, satisfy

$$C_1 ds(x) \le \varepsilon ds(y) \le C_2 ds(x)$$

for some constant $C_{1,2}$ which depend only on the constants in the assumptions and Y^- but uniform in ε and γ ; see e.g. [35, Proposition A.1] for a precise relation between ds(x) and ds(y).

Consequently, we have

$$\|v^+\|_{L^2(\Gamma_{\varepsilon})}^2 = \sum_{i=1}^{N(\varepsilon)} \int_{\varepsilon \Phi(\Gamma_i)} |v_i(x)|^2 ds(x) \le C \varepsilon \sum_{i=1}^{N(\varepsilon)} \int_{\Gamma_i} |\tilde{v}_i(y)|^2 ds(y).$$

Apply (C.2) and change the variable back; use again $dx \sim \varepsilon^2 dy$ and $\nabla_y \tilde{v}_i = \varepsilon \nabla_x v_i$ to get

$$\begin{split} \|v^+\|_{L^2(\Gamma_{\varepsilon})}^2 &\leq C\varepsilon \sum_{i=1}^{N(\varepsilon)} \int_{Y_i^+} |\tilde{v}_i(y)|^2 + |\nabla_y \tilde{v}(y)|^2 dy \\ &\leq C\varepsilon^{-1} \sum_{i=1}^{N(\varepsilon)} \int_{\varepsilon \Phi(Y_i^+)} |v_i(x)|^2 + \varepsilon^2 |\nabla v(x)|^2 dx \end{split}$$

This completes the proof of (C.1).

The other inequality in (3.10) is implied by the following lemma:

Lemma C.2. There exists a constant C > 0 independent of ε or γ such that

$$\|v\|_{L^{2}(\Omega_{\varepsilon}^{-})} \leq C\left(\sqrt{\varepsilon}\|v\|_{L^{2}(\Gamma_{\varepsilon})} + \varepsilon\|\nabla v\|_{L^{2}(\Omega_{\varepsilon}^{-})}\right)$$
(C.3)

for all $v \in H^1(\Omega_{\varepsilon}^-)$.

Proof. We first observe that on the reference cube Y with reference cell Y^- , we have that

$$\|v\|_{L^{2}(Y^{-})}^{2} \leq C\left(\|v\|_{L^{2}(\Gamma_{0})}^{2} + \|\nabla v\|_{L^{2}(Y^{-})}^{2}\right),$$
(C.4)

for any $v \in H^1(Y^-)$ where *C* only depends on Y^- and the dimension. Indeed, suppose otherwise, we could find a sequence $\{v_n\} \subset H^1(Y^-)$ such that $\|v_n\|_{L^2(Y^-)} \equiv 1$ but

$$\|v_n\|_{L^2(\Gamma_0)} + \|\nabla v_n\|_{L^2(Y^-)} \longrightarrow 0, \quad \text{as } n \to \infty$$

Then since $||v_n||_{H^1}$ is uniformly bounded, there exists a subsequence, still denoted as $\{v_n\}$, and a function $v \in H^1(Y^-)$ such that

$$v_n \rightarrow v$$
 weakly in $H^1(\Upsilon^-)$, $\nabla v_n \rightarrow \nabla v$ weakly in $L^2(\Upsilon^-)$.

Consequently, $\|\nabla v\|_{L^2} \leq \liminf \|\nabla v_n\|_{L^2} = 0$, which implies that v = C for some constant. Moreover, since the embedding $H^1(Y^-) \hookrightarrow L^2(\Gamma_0)$ is compact, the convergence $v_n \to v$ holds strongly in $L^2(\Gamma_0)$ and $\|v\|_{L^2(\Gamma)} \leq \lim \|v_n\|_{L^2(\Gamma_0)} = 0$. Consequently $v \equiv 0$. On the other hand, $v_n \to v$ holds strongly in $L^2(Y^-)$ and hence $\|v\|_{L^2(Y^-)} = \lim \|v_n\|_{L^2(Y^-)} = 1$. This contradicts with the fact that $v \equiv 0$.

To prove (C.3), we lift functions in $\varepsilon \Phi(Y_i^-)$ to functions in Y_i^- as in the proof of the previous lemma, and use the scaling relations of the measures: $dx \sim \varepsilon^2 dy$ and $ds(x) \sim \varepsilon ds(y)$. We calculate

$$\|v\|_{L^2(\Omega_{\varepsilon}^-)}^2 = \sum_{i=1}^{N(\varepsilon)} \int_{\varepsilon \Phi(Y_i^-)} |v|^2 dx \le C\varepsilon^2 \int_{Y^-} |\tilde{v}|^2 dy \le C\varepsilon^2 \sum_{i=1}^{N(\varepsilon)} \int_{\Gamma_i} |\tilde{v}|^2 ds + \int_{Y_i^-} |\nabla \tilde{v}|^2 dy$$

where in the last inequality we used (C.4). Change the variables back to get

$$\|v\|_{L^2(\Omega_{\varepsilon}^-)}^2 \leq C\varepsilon^2 \sum_{i=1}^{N(\varepsilon)} \int_{\varepsilon \Phi(\Gamma_i)} \varepsilon^{-d+1} |v|^2 ds + \int_{\varepsilon \Phi(Y_i^-)} \varepsilon^{-d+2} |\nabla v|^2 dy.$$

Note that we used again $\nabla_y \tilde{v} = \varepsilon \nabla_x v$. The above inequality is precisely (C.3).

Proof of Proposition 3.2. To prove the first inequality, we apply Lemma C.1 to get

$$\begin{split} \varepsilon \| u^+ - u^- \|_{L^2(\Gamma_{\varepsilon})}^2 &\leq 2(\varepsilon \| u^+ \|_{L^2(\Gamma_{\varepsilon})}^2 + \varepsilon \| u^- \|_{L^2(\Gamma_{\varepsilon})}^2) \\ &\leq C(\| u^+ \|_{L^2(\Omega_{\varepsilon}^+)}^2 + \| u^- \|_{L^2(\Omega_{\varepsilon}^-)}^2 + \varepsilon^2 \| \nabla u^+ \|_{L^2(\Omega_{\varepsilon}^+)}^2 + \varepsilon^2 \| \nabla u^+ \|_{L^2(\Omega_{\varepsilon}^+)}^2). \end{split}$$

Only the first term in (B.2) does not show in $\|\cdot\|_{H^1_{\mathbb{C}}\times H^1}$, but it is controlled by $\|\nabla u^+\|_{L^2(\Omega^+_{\varepsilon})}$ uniformly in ε and γ thanks to (B.2).

For the second inequality, we only need to control $||u^-||_{L^2(\Omega_{\varepsilon}^-)}$. We apply Lemma C.2 and the triangle inequality:

$$\|u^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2} \leq C\left(\varepsilon \|u^{+}\|_{L^{2}(\Gamma_{\varepsilon})}^{2} + \varepsilon \|u^{+} - u^{-}\|_{L^{2}(\Gamma_{\varepsilon})}^{2} + \varepsilon^{2} \|\nabla u^{-}\|_{L^{2}(\Omega_{\varepsilon}^{-})}^{2}\right).$$

Only the first term does not appear in $\|\cdot\|_{W_{\epsilon}}$, but using Lemma C.1 and (B.2) we can bound it by

$$\varepsilon \|u^+\|_{L^2(\Gamma_{\varepsilon})}^2 \le C(\|u^+\|_{L^2(\Omega_{\varepsilon}^+)}^2 + \varepsilon^2 \|\nabla u^+\|_{L^2(\Omega_{\varepsilon}^+)}^2) \le C \|\nabla u^+\|_{L^2(\Omega_{\varepsilon}^+)}^2.$$

This completes the proof.

D Technical lemma

Lemma D.1. Let φ_1 be a function in $\mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^-))$. There exists at least one function θ in $(\mathcal{D}(\Omega, H^1_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, H^1_{\sharp}(Y^-)))^2$ solution of the following problem:

$$\begin{cases}
-\nabla_{y} \cdot \theta^{+}(x, y) = 0 & \text{in } Y^{+}, \\
-\nabla_{y} \cdot \theta^{-}(x, y) = 0 & \text{in } Y^{-}, \\
\theta^{+}(x, y) \cdot n = \theta^{-}(x, y) \cdot n & \text{on } \Gamma, \\
\theta^{+}(x, y) \cdot n = \varphi_{1}^{+}(x, y) - \varphi_{1}^{-}(x, y) & \text{on } \Gamma, \\
y \longmapsto \theta(x, y)Y - \text{periodic.}
\end{cases}$$
(D.1)

Proof. We look for a solution under the form $\theta = \nabla_y \eta$. We hence introduce the following variational problem:

Find
$$\eta \in (H^1_{\sharp}(Y^+)/\mathbb{C}) \times (H^1_{\sharp}(Y^-)/\mathbb{C})$$
 such that

$$\int_{Y^+} \nabla \eta^+(y) \cdot \overline{\psi}^+(y) dy + \int_{Y^-} \nabla \eta^-(y) \cdot \overline{\psi}^-(y) dy$$

$$= \frac{1}{\beta k_0} \int_{\Gamma} (\varphi_1^+ - \varphi_1^-) (\overline{\psi}^+ - \overline{\psi}^-)(y) ds(y),$$
for all $\psi \in (H^1_{\sharp}(Y^+)/\mathbb{C}) \times (H^1_{\sharp}(Y^-)/\mathbb{C}),$

for a fixed $x \in \Omega$. Lax-Milgram theorem gives us existence and uniqueness of such an η . Since $\varphi_1 \in \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, C^{\infty}_{\sharp}(Y^-))$, there exists at least one function $\theta \in (\mathcal{D}(\Omega, H^1_{\sharp}(Y^+)) \times \mathcal{D}(\Omega, H^1_{\sharp}(Y^-))^2$ solution of (D.1). Note that we do not have uniqueness of such a solution. \Box

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