# A mathematical model for Alzheimer's disease: An approach via stochastic homogenization of the 

 Smoluchowski equationBruno Franchi, ${ }^{*}$ Martin Heida $^{\dagger}$ and Silvia Lorenzani ${ }^{\ddagger}$


#### Abstract

In this note, we apply the theory of stochastic homogenization to find the asymptotic behavior of the solution of a set of Smoluchowski's coagulation-diffusion equations with non-homogeneous Neumann boundary conditions. This system is meant to model the aggregation and diffusion of $\beta$-amyloid peptide $(\mathrm{A} \beta)$ in the cerebral tissue, a process associated with the development of Alzheimer's disease. In contrast to the approach used in our previous works, in the present paper we account for the non-periodicity of the cellular structure of the brain by assuming a stochastic model for the spatial distribution of neurons. Further, we consider non-periodic random diffusion coefficients for the amyloid aggregates and a random production of $\mathrm{A} \beta$ in the monomeric form at the level of neuronal membranes.


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## 1 Introduction

The primary feature of several neurological diseases, such as Prion diseases, Alzheimer's disease, Parkinson's disease, Creutzfeldt-Jacob disease is the pathological presence of misfolded protein aggregates (that is, proteins that fail to configure properly, becoming structurally abnormal) [7], 21]. In this paper, we focus our interest in Alzheimer's disease (AD). Indeed, AD has a huge social and economic impact. Until 2040 its worldwide global prevalence (estimated as high as 44 millions in 2015) is expected to double every 20 years. In particular, existing clinical data support the idea that amyloid- $\beta$ peptide (A $\beta$ ) has a critical role as initiator of a complex network of pathologic changes in the brain, ultimately leading to Alzheimer's disease ('amyloid hypothesis', see e.g. [12], [19], [23]). Although there is no doubt that the presence of fibrillar $\mathrm{A} \beta$ deposition (senile plaques) is the hallmark of the clinical syndrome of AD , the bulk of human biomarker data reveals the existence of a discrepancy between the appearance of amyloid deposits and clinical dementia, with $\mathrm{A} \beta$ plaques anatomically disconnected from areas of severe neuronal loss. One of the most reliable explanations, which also supports the amyloid hypothesis, is that, in addition to fibrillar plaques, oligomeric forms of $\mathrm{A} \beta$ can play a dominant role in triggering a wide variety of pathogenic effects. Mice which accumulate $\mathrm{A} \beta$ oligomers, but not fibrillar plaques, develop synaptic damage, inflammation and cognitive impairment [29], [31. Despite the biological relevance of the negative effects produced, the exact mechanisms of misfolded protein aggregation and propagation, as well as their toxicity, are still not well understood. Furthermore, the complexity of the underlying processes makes it difficult to extrapolate the effects of protein misfolding from the microscopic (e.g. molecular) to the macroscopic (e.g. organs) scale, preventing the development of effective therapeutic interventions. In order to complement the medical and biological research, the last few decades have seen the emergence of several mathematical models that can help to provide a better insight into the laws governing the processes of protein aggregation and the effects of toxicity spreading. The mathematical approaches considered so far can be predominantly divided into two different classes: on the one hand, there are the models designed to describe processes at the molecular (microscopic) scale (aggregation
kinetics, short-range spatial spreading, etc...) [7, [1, [9, [16] while, on the other side, there are models that account for large-scale events characterizing the progression of neurodegenerative misfolded protein-related diseases [7], 3, 4], 5].

### 1.1 A mathematical model for the aggregation of $\beta$-amyloid based on Smoluchowski's equations.

In 2013, Achdou et al. proposed in [1] a mathematical model for the aggregation and diffusion of $\beta$-amyloid peptide ( $\mathrm{A} \beta$ ) in the brain affected by Alzheimer's disease ( AD ) at a microscopic scale (the size of a single neuron). In particular, these authors considered a portion of cerebral tissue, represented by a bounded smooth region $\boldsymbol{Q} \subset \mathbb{R}^{3}$, and described the neurons as a family of regular regions $G_{j}$ such that:
(i) $\bar{G}_{j} \subset \boldsymbol{Q} \quad$ if $j=1,2, \ldots \bar{M}$;
(ii) $\bar{G}_{i} \cap \bar{G}_{j}=\varnothing \quad$ if $i \neq j$.

Then, the following system of Smoluchowski equations has been introduced:

$$
\begin{cases}\frac{\partial u_{1}}{\partial t}(t, x)-d_{1} \triangle_{x} u_{1}(t, x)+u_{1}(t, x) & \sum_{j=1}^{M} a_{1, j} u_{j}(t, x)=0 \\ \frac{\partial u_{1}}{\partial \nu} \equiv \nabla_{x} u_{1} \cdot n=0 & \text { on } \partial \boldsymbol{Q} \\ \frac{\partial u_{1}}{\partial \nu} \equiv \nabla_{x} u_{1} \cdot n=\eta_{j} & \text { on } \Gamma_{j}, \quad j=1, \ldots, \bar{M} \\ & \\ u_{1}(0, x)=U_{1} \geq 0 & \end{cases}
$$

if $1<m<M$

$$
\begin{cases}\frac{\partial u_{m}}{\partial t}(t, x)-d_{m} \triangle_{x} u_{m}(t, x)+u_{m}(t, x) & \sum_{j=1}^{M} a_{m, j} u_{j}(t, x)=  \tag{2}\\ & \frac{1}{2} \sum_{j=1}^{m-1} a_{j, m-j} u_{j} u_{m-j} \\ \frac{\partial u_{m}}{\partial \nu} \equiv \nabla_{x} u_{m} \cdot n=0 & \text { on } \partial \boldsymbol{Q} \\ & \\ \frac{\partial u_{m}}{\partial \nu} \equiv \nabla_{x} u_{m} \cdot n=0 & \text { on } \Gamma_{j}, j=1, \ldots, \bar{M} \\ u_{m}(0, x)=0 & \end{cases}
$$

and

$$
\begin{cases}\frac{\partial u_{M}}{\partial t}(t, x)-d_{M} \triangle_{x} u_{M}(t, x)=\frac{1}{2} \sum_{\substack{k<M(i f) \\ j<M) \\ j<M(\text { if } k=M)}}^{j+k \geq M}  \tag{3}\\ \frac{\partial u_{M}}{\partial \nu} \equiv \nabla_{x} u_{M} \cdot n=0 & a_{j, k} u_{j} u_{k} \\ \frac{\partial u_{M}}{\partial \nu} \equiv \nabla_{x} u_{M} \cdot n=0 & \text { on } \partial \boldsymbol{Q} \\ & \\ u_{M}(0, x)=0 & \text { on } \Gamma_{j}, j=1, \ldots, \bar{M}\end{cases}
$$

where $u_{j}(t, x)(1 \leq j<M-1)$ is the molar concentration at the point $x$ and at the time $t$ of an $\mathrm{A} \beta$ assembly of $j$ monomers, while $u_{M}$ takes into account aggregations of more than $M-1$ monomers. The production of $\mathrm{A} \beta$ in monomeric form from the neuron membranes has been modeled by coupling the evolution equation for $u_{1}$ with a nonhomogeneous Neumann condition on the boundaries of the sets $G_{j}$, indicated by $\Gamma_{j}$. In particular, $0 \leq \eta_{j} \leq 1$ is a smooth function for $j=1, \ldots, \bar{M}$. These monomers, by binary coalescence, give rise to larger assemblies, which can diffuse in the cerebral tissue with a
diffusion coefficient $d_{j}$ that depends on their size. Let us remark that, Achdou's model has been formulated to be valid on small spatial domains, therefore isotropic diffusion has been assumed. The coagulation rates $a_{i, j}$ are symmetric $a_{i, j}=a_{j, i}>0, i, j=1, \ldots, M$, but $a_{M M}=0$, since it is assumed that long fibrils, characterized by a very slow diffusion, do not coagulate with each other.

### 1.2 Stochastic homogenization.

Since the development of modern imaging techniques (useful to evaluate the progression of Alzheimer's disease) requires the need to test the predictions of mathematical modeling at the macroscale, in the present paper, we have applied the homogenization method to the model presented by Achdou et al. [1], in order to describe the effects of the production and agglomeration of the $\mathrm{A} \beta$ at the macroscopic level. The homogenization theory, introduced by the mathematicians in the seventies to perform a sort of averaging procedure on the solutions of partial differential equations with rapidly varying coefficients or describing media with microstructures, has been already successfully applied in 9, [10] to derive a limiting model from that proposed by Achdou et al. [1], in the context of a periodically perforated domain. In particular, in [9], [10] we have constructed our set $\boldsymbol{Q}^{\varepsilon}$, starting from a fixed bounded domain $\boldsymbol{Q}$ (which represents a portion of cerebral tissue) and removing from it many small holes of characteristic size $\varepsilon$ (the neurons) distributed periodically. Then, we have rewritten the model problem (11)-(3) as a family of equations in $\boldsymbol{Q}^{\varepsilon}$ and we have performed the limit $\varepsilon \rightarrow 0$ in the framework of the two-scale convergence, first introduced by Nguetseng [24] and Allaire [2]. The peculiarity of the two-scale convergence method, used in [9], 10] to study the limiting behavior of the Smoluchowskitype equations, is that in a single process, one can find the homogenized equations and prove the convergence of a sequence of solutions to the problem at hand. Since the picture presented in our previous works [9], [10] is a too crude oversimplification of the biomedical reality, in the present paper we have chosen to resort to a stochastic parametrization of the model equations: that is, we account for the non-periodic cellular structure of the brain. In particular, the distribution of neurons is modeled in the following way: it exists
a family of predominantly genetic causes, not wholly deterministic, which influences the position of neurons and the microscopic structure of the parenchyma in a portion of the brain tissue $\boldsymbol{Q}$. Also, we consider non-periodic random diffusion coefficients and a random production of $\mathrm{A} \beta$ in the monomeric form at the level of neuronal membranes. This together defines a probability space $(\Omega, \mathcal{F}, \mathbb{P})$.

Denoting by $\omega \in \Omega$ the random variable in our model, the set of random holes in $\mathbb{R}^{m}$ (representing the neurons) is labeled by $G(\omega)$. The production of $\beta$-amyloid at the boundary $\Gamma(\omega)$ of $G(\omega)$ is described by a random scalar function $\eta(x, \omega)$ and the diffusivity, in the brain parenchyma, of clusters of different sizes $s$ is modeled by random matrices $D_{s}(x, \omega)$ on $\Omega$. For technical reasons, we assume that the randomness of the medium is stationary, that is, the probability distribution of the random variables observed in a set $A \subset \mathbb{R}^{m}$ is shift invariant ( all variables share the same distribution in $A$ and $A+x$, $\left.x \in \mathbb{R}^{m}\right)$. As shown by Papanicolaou and Varadhan [26] (who introduced this concept), the assumption of stationarity provides a family of mappings $\left(\tau_{x}\right)_{x \in \mathbb{R}^{m}}: \Omega \rightarrow \Omega$ such that $\eta(x, \omega)=\eta\left(\tau_{x} \omega\right)$ and $D_{s}(x, \omega)=D_{s}\left(\tau_{x} \omega\right)$. The periodic homogenization can be recovered in this frame considering $\Omega=[0,1)^{m}$ with $\tau_{x} \omega=x+\omega \bmod [0,1)^{m}$, where one canonically chooses $\omega=0$ (see also [13]).

The above mentioned findings can be interpreted in the sense that the stationarity of the coefficients and the resulting dynamical system $\tau_{x}$ transfer some structural properties from $\mathbb{R}^{m}$ to $\Omega$ such that we could formally identify $\Omega \approx \mathbb{R}^{m}$. Accordingly, a stationary random set in $\mathbb{R}^{m}$ corresponds to a subset of $\Omega$ and a random Hausdorff measure on $\mathbb{R}^{m}$ corresponds to a measure on $\Omega$. In order to prevent confusion, let us note that, all the similarities we mention here are of algebraic and measure-theoretic nature and not in the sense of a vector space isomorphism. With the above short overview, we just want to point out that many useful tools in periodic homogenization find their counterpart in the stochastic setting. The stochastic homogenization theorems can be formulated in a very similar way to their periodic version, if we rely on the above connections and similarities, though the mathematics behind differs sometimes significantly. In this framework, we have studied the limiting behavior of the system of nonlinear Smoluchowski-type equations
describing our model by using a sort of stochastic version of the two-scale convergence method.

The rest of the paper is organized as follows. In Section 2, we give a brief survey of the probabilistic background behind the theory of stochastic homogenization and in Section 3 we present all the main definitions and theorems related to the stochastic twoscale convergence method. In Section 4 we give a detailed description of our model and derive all the a priori estimates needed to apply the two-scale homogenization technique. Then, Section 5 is devoted to the presentation and the proof of our main results on the stochastic homogenization of the nonlinear Smoluchowski coagulation-diffusion equations in a randomly perforated domain. Finally, Appendix A is introduced to summarize some basic concepts on the realization of random sets.

## 2 Random media

The method of stochastic two-scale convergence introduced by Zhikov and Piatnitsky [32] is based on a setting that was originally introduced by Papanicolaou and Varadhan [26]. The connection between the abstract setting on random singular measures in 32] and the theory of random sets was worked out in [13]. Hence we will first introduce the setting of [26] and explain the ideas pointed out in [13] before we move on to the definition of two-scale convergence.

### 2.1 Stationary ergodic dynamical systems.

This section has the intention to provide a probabilistic background for the theory of stochastic homogenization, and more particularly for stochastic two-scale convergence. We follow the formulation given by Papanicolaou and Varadhan [26], enriched by the ideas presented in [18], [32] and [13, 14].

The whole theory is based on the concept of dynamical systems.

Definition 2.1 (Dynamical system). Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space. An $m$-dimensional dynamical system is defined as a family of measurable bijective mappings $\tau_{x}: \Omega \rightarrow \Omega$,
$x \in \mathbb{R}^{m}$, satisfying the following conditions:
(i) the group property: $\tau_{0}=\mathbb{1}\left(\mathbb{1}\right.$ is the identity mapping), $\tau_{x+y}=\tau_{x} \circ \tau_{y} \quad \forall x, y \in \mathbb{R}^{m} ;$
(ii) the mappings $\tau_{x}: \Omega \rightarrow \Omega$ preserve the measure $\mathbb{P}$ on $\Omega$, i.e., for every $x \in \mathbb{R}^{m}$, and every $\mathbb{P}$-measurable set $F \in \mathcal{F}$, we have $\mathbb{P}\left(\tau_{x} F\right)=\mathbb{P}(F)$;
(iii) the map $\mathcal{T}: \Omega \times \mathbb{R}^{m} \rightarrow \Omega:(\omega, x) \mapsto \tau_{x} \omega$ is measurable (for the standard $\sigma$-algebra on the product space, where on $\mathbb{R}^{m}$ we take the Borel $\sigma$-algebra).

Note that (i) and (iii) imply that, for every $x \in \mathbb{R}^{m}$ and measurable $F \subset \Omega$, the set $\tau_{x} F$ is measurable: since $\tau_{-x}\left(\tau_{x} F\right)=F$ we find that $\tau_{x} F$ is the projection of $\mathcal{T}^{-1}(F) \cap\{-x\} \times \Omega$ onto $\Omega$. We define the notion of ergodicity for the dynamical system.

Definition 2.2 (Ergodicity). A dynamical system is called ergodic if one of the following equivalent conditions is fulfilled
(i) given a measurable and invariant function $f$ in $\Omega$, that is

$$
\forall x \in \mathbb{R}^{m} \quad f(\omega)=f\left(\tau_{x} \omega\right)
$$

almost everywhere in $\Omega$, then

$$
f(\omega)=\text { const. } \quad \text { for } \quad \mathbb{P}-\text { a.e. } \omega \in \Omega \text {; }
$$

(ii) if $F \in \mathcal{F}$ is such that $\tau_{x} F=F \quad \forall x \in \mathbb{R}^{m}$, then $\mathbb{P}(F)=0$ or $\mathbb{P}(F)=1$.

Definition 2.3 (Stationarity). Given a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, a real valued process is a measurable function $f: \mathbb{R}^{m} \times \Omega \rightarrow \mathbb{R}$. We will say $f$ is stationary if the distribution of the random variable $f(y, \cdot): \Omega \rightarrow \mathbb{R}$ is independent of $y$, i.e., for all $a \in \mathbb{R}, \mathbb{P}(\{\omega: f(y, \omega)>a\})$ is independent of $y$. This is qualified by assuming the existence of a dynamical system $\tau_{y}: \Omega \rightarrow \Omega\left(y \in \mathbb{R}^{m}\right)$ and saying that $f: \mathbb{R}^{m} \times \Omega \rightarrow \mathbb{R}$ is stationary if

$$
f\left(y+y^{\prime}, \omega\right)=f\left(y, \tau_{y^{\prime}} \omega\right) \text { for all } y, y^{\prime} \in \mathbb{R}^{m} \text { and } \omega \in \Omega
$$

Finally, we say that a random variable $f: \mathbb{R}^{m} \times \Omega \rightarrow \mathbb{R}$ is stationary ergodic if it is stationary and the underlying dynamical system is ergodic. Naturally, if $f$ is taking values in a finite dimensional space, we will say it is stationary if all of its components in a given basis are stationary with respect to the same dynamical system. This property is also called jointly stationary.

Remark 2.4. [26] A function $f$ is stationary ergodic if and only if there is some measurable function $\tilde{f}: \Omega \rightarrow \mathbb{R}$ such that

$$
f(x, \omega)=\tilde{f}\left(\tau_{x} \omega\right)
$$

For a fixed $\omega \in \Omega$ the function $x \mapsto \tilde{f}\left(\tau_{x} \omega\right)$ of argument $x \in \mathbb{R}^{m}$ is said to be a realization of function $\tilde{f}$.

Let $L^{p}(\Omega)(1 \leq p<\infty)$ denote the space formed by (the equivalence classes of) measurable functions that are $\mathbb{P}$-integrable with exponent $p$ and $L^{\infty}(\Omega)$ be the space of measurable essentially bounded functions. If $f \in L^{p}(\Omega)$, then $\mathbb{P}$-almost all realizations $f\left(\tau_{x} \omega\right)$ belong to $L_{\mathrm{loc}}^{p}\left(\mathbb{R}^{m}\right)$ [18].

We define the following $m$-parameter group of operators in the space $L^{2}(\Omega)$ :

$$
U(x): L^{2}(\Omega) \rightarrow L^{2}(\Omega), \quad f \mapsto[U(x) f](\omega):=f\left(\tau_{x} \omega\right) .
$$

It is known [18] that the operator $U(x)$ is unitary for each $x \in \mathbb{R}^{m}$ and the group $U(x)$ is strongly continuous, i.e.

$$
\forall f \in L^{2}(\Omega): \quad \lim _{x \rightarrow 0}\|U(x) f-f\|_{L^{2}(\Omega)}=0
$$

For $x=\left\{0,0, \ldots, x_{i}, 0, \ldots, 0\right\}$ we obtain a one-parameter group whose infinitesimal generator will be denoted by $\mathrm{D}_{i}$ with domain $\mathcal{D}_{i}(\Omega)$. The unitarity of the group $U(x)$ implies that the operators $\mathrm{D}_{i}$ are skew-symmetric:

$$
\begin{equation*}
\forall f, g \in \mathcal{D}_{i}(\Omega): \quad \int_{\Omega}\left(\mathrm{D}_{i} f\right) g d \mathbb{P}=-\int_{\Omega} f\left(\mathrm{D}_{i} g\right) d \mathbb{P} \tag{4}
\end{equation*}
$$

and by definition of the generators we have

$$
\begin{equation*}
\mathrm{D}_{i} f=\lim _{x_{i} \neq 0, x_{i} \rightarrow 0} \frac{f\left(\tau_{x_{i}} \omega\right)-f(\omega)}{x_{i}} \tag{5}
\end{equation*}
$$

in the sense of convergence in $L^{2}(\Omega)$. As Papanicolaou and Varadhan [26] have shown, almost every realization possesses a weak derivative and it holds

$$
\frac{\partial}{\partial x_{i}} f\left(\tau_{x} \omega\right)=\left(\mathrm{D}_{i} f\right)\left(\tau_{x} \omega\right) \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{m}\right)
$$

Also we have that $\left(i \mathrm{D}_{1}, \ldots, i \mathrm{D}_{m}\right)$ are commuting, self-adjoint, closed, and densely defined linear operators on $L^{2}(\Omega)$ [17], and we may define

$$
\mathrm{D}_{\omega} f:=\left(\mathrm{D}_{1} f, \ldots, \mathrm{D}_{m} f\right)^{\top} .
$$

We introduce the space $W^{1,2}(\Omega)$ with norm $\|\cdot\|_{1,2}$ through

$$
\begin{aligned}
W^{1,2}(\Omega) & :=\mathcal{D}_{1}(\Omega) \bigcap \ldots \bigcap \mathcal{D}_{m}(\Omega) \\
\|f\|_{1,2} & :=\|f\|_{L^{2}(\Omega)}+\sum_{i=1}^{m}\left\|\mathrm{D}_{i} f\right\|_{L^{2}(\Omega)} .
\end{aligned}
$$

Further let $L_{\text {loc }}^{2}\left(\mathbb{R}^{m} ; \mathbb{R}^{m}\right)$ be the set of measurable functions $f: \mathbb{R}^{m} \rightarrow \mathbb{R}^{m}$ such that $\left.f\right|_{\mathbf{U}} \in L^{2}\left(\mathbf{U} ; \mathbb{R}^{m}\right)$ for every bounded domain $\mathbf{U}$ and we define

$$
\begin{aligned}
L_{\mathrm{pot}, \mathrm{loc}}^{2}\left(\mathbb{R}^{m}\right) & :=\left\{f \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{m} ; \mathbb{R}^{m}\right) \mid \forall \mathbf{U} \text { bounded domain, } \exists \varphi \in H^{1}(\mathbf{U}): f=\nabla \varphi\right\}, \\
L_{\mathrm{sol}, \mathrm{loc}}^{2}\left(\mathbb{R}^{m}\right) & :=\left\{f \in L_{\mathrm{loc}}^{2}\left(\mathbb{R}^{m} ; \mathbb{R}^{m}\right) \mid \int_{\mathbb{R}^{m}} f \cdot \nabla \varphi=0 \forall \varphi \in C_{c}^{1}\left(\mathbb{R}^{m}\right)\right\} .
\end{aligned}
$$

Recalling the notion of a realization $f_{\omega}(x):=f\left(\tau_{x} \omega\right)$ for $f \in L^{2}(\Omega)$, we can then define corresponding spaces on $\Omega$ through

$$
\begin{align*}
L_{\mathrm{pot}}^{2}(\Omega) & :=\left\{f \in L^{2}\left(\Omega ; \mathbb{R}^{m}\right): f_{\omega} \in L_{\mathrm{pot}, \mathrm{loc}}^{2}\left(\mathbb{R}^{m}\right) \text { for } \mathbb{P}-\text { a.e. } \omega \in \Omega\right\}, \\
L_{\mathrm{sol}}^{2}(\Omega) & :=\left\{f \in L^{2}\left(\Omega ; \mathbb{R}^{m}\right): f_{\omega} \in L_{\mathrm{sol}, \mathrm{loc}}^{2}\left(\mathbb{R}^{m}\right) \text { for } \mathbb{P} \text { - a.e. } \omega \in \Omega\right\},  \tag{6}\\
\mathcal{V}_{\mathrm{pot}}^{2}(\Omega) & :=\left\{f \in L_{\mathrm{pot}}^{2}(\Omega): \int_{\Omega} f d \mathbb{P}=0\right\} .
\end{align*}
$$

It has been shown in Chapter 7 of [18] that all of these spaces are closed and that $L^{2}\left(\Omega ; \mathbb{R}^{m}\right)=L_{\text {sol }}^{2}(\Omega) \oplus \mathcal{V}_{\text {pot }}^{2}(\Omega)$. This has been proved using the continuous smoothing operator

$$
\begin{equation*}
\mathscr{I}_{\delta}: L^{2}(\Omega) \rightarrow W^{1,2}(\Omega), \quad \mathscr{I}_{\delta} f(\omega):=\int_{\mathbb{R}^{m}} \eta_{\delta}(x) f\left(\tau_{x} \omega\right) \mathrm{d} x, \tag{7}
\end{equation*}
$$

where $\eta_{\delta}$ is a Dirac-sequence of smooth functions. It can be shown that, for every $f \in$ $L^{2}(\Omega)$, it holds $\mathscr{I}_{\delta} f \rightarrow f$ as $\delta \rightarrow 0$ and the continuity of $\mathscr{I}_{\delta}$ implies $\mathrm{D}_{i} \mathscr{I}_{\delta} f=\mathscr{I}_{\delta} \mathrm{D}_{i} f$ for all $f \in W^{1,2}(\Omega)$. Thus, if we consider

$$
\tilde{\mathcal{V}}:=\operatorname{closure}_{L^{2}(\Omega)}\left\{\mathrm{D}_{\omega} f: f \in W^{1,2}(\Omega)\right\}
$$

we first obtain $\tilde{\mathcal{V}} \subseteq \mathcal{V}_{\text {pot }}^{2}(\Omega)$ and for $g \in \tilde{\mathcal{V}}^{\perp}$, we have for every $\delta>0$

$$
\forall f \in W^{1,2}(\Omega): \quad 0=\left\langle g, \mathrm{D}_{\omega} \mathscr{I}_{\delta} f\right\rangle=\left\langle\mathscr{I}_{\delta} g, \mathrm{D}_{\omega} f\right\rangle=-\sum_{i}\left\langle\mathrm{D}_{i} \mathscr{\mathscr { I }}_{\delta} g, f\right\rangle
$$

and hence $\sum_{i} \mathrm{D}_{i} \mathscr{I}_{\delta} g=0$. In particular, $\mathscr{I}_{\delta} g \in L_{\text {sol }}^{2}(\Omega)$ and since $L_{\text {sol }}^{2}(\Omega)$ is closed we find $\tilde{\mathcal{V}}^{\perp} \subseteq L_{\text {sol }}^{2}(\Omega)$. This implies $\tilde{\mathcal{V}} \supseteq \mathcal{V}_{\text {pot }}^{2}(\Omega)$ and hence

$$
\begin{equation*}
\tilde{\mathcal{V}}=\operatorname{closure}_{L^{2}(\Omega)}\left\{\mathrm{D}_{\omega} f: f \in W^{1,2}(\Omega)\right\}=\mathcal{V}_{\mathrm{pot}}^{2}(\Omega) . \tag{8}
\end{equation*}
$$

In what follows, we will often impose the following assumption:

Assumption 2.5. Assume that $\Omega$ is a separable metric space and $(\Omega, \mathcal{F}, \mathbb{P})$ is a probability space with countably generated $\sigma$-algebra and let $\tau_{x}, x \in \mathbb{R}^{m}$, be a dynamical system in the sense of Definition 2.1 that is ergodic in the sense of Definition 2.2.

It was discussed in [13] that the latter assumption is not a restriction to our choice of parameters.

By $\mathcal{M}\left(\mathbb{R}^{m}\right)$ we denote the space of finitely bounded Borel measures on $\mathbb{R}^{m}$ equipped with the Vague topology, which makes $\mathcal{M}\left(\mathbb{R}^{m}\right)$ a separable metric space [8]. The $\sigma$-field defined by this topology is denoted by $\mathcal{B}(\mathcal{M})$ since it is a Borel $\sigma$-field on $\mathcal{M}$. A random measure is a measurable mapping

$$
\mu_{\bullet}: \Omega \rightarrow \mathcal{M}\left(\mathbb{R}^{m}\right), \quad \omega \mapsto \mu_{\omega}
$$

which is equivalent to the measurability of all mappings $\omega \mapsto \mu_{\omega}(A)$, where $A \subset \mathbb{R}^{m}$ are arbitrary bounded Borel sets. A random measure is stationary if the distribution of $\mu_{\omega}(A)$ is invariant under translations of $A$. In particular, random measures satisfy $\mu_{\tau_{x} \omega}(A)=\mu_{\omega}(A+x)$. For stationary random measures we find the following important property.

Theorem 2.6 (8] Existence of Palm measure and Campbell's Formula). Let $\mathcal{L}$ be the Lebesgue-measure on $\mathbb{R}^{m}$ with $d x:=d \mathcal{L}(x)$ and $(\Omega, \mathcal{F}, \mathbb{P})$ and $\tau$ as in Assumption 2.5. Then there exists a unique measure $\mu_{\mathcal{P}}$ on $\Omega$ such that

$$
\int_{\Omega} \int_{\mathbb{R}^{m}} f\left(x, \tau_{x} \omega\right) \mathrm{d} \mu_{\omega}(x) \mathrm{d} \mathbb{P}(\omega)=\int_{\mathbb{R}^{m}} \int_{\Omega} f(x, \omega) \mathrm{d} \mu_{\mathcal{P}}(\omega) \mathrm{d} x
$$

for all $\mathcal{B}\left(\mathbb{R}^{m}\right) \times \mathcal{B}(\Omega)$-measurable non negative functions and all $\mu_{\mathcal{P}} \times \mathcal{L}$ - integrable functions. Furthermore

$$
\begin{align*}
\mu_{\mathcal{P}}(A) & =\int_{\Omega} \int_{\mathbb{R}^{m}} g(s) \chi_{A}\left(\tau_{s} \omega\right) \mathrm{d} \mu_{\omega}(s) \mathrm{dP}(\omega),  \tag{9}\\
\int_{\Omega} f(\omega) \mathrm{d} \mu_{\mathcal{P}} & =\int_{\Omega} \int_{\mathbb{R}^{m}} g(s) f\left(\tau_{s} \omega\right) \mathrm{d} \mu_{\omega}(s) \mathrm{dP}(\omega) \tag{10}
\end{align*}
$$

for an arbitrary $g \in L^{1}\left(\mathbb{R}^{m}, \mathcal{L}\right)$ with $\int_{\mathbb{R}^{m}} g(x) \mathrm{d} x=1$ and $\mu_{\mathcal{P}}$ is $\sigma$-finite.
The measure $\mu_{\mathcal{P}}$ from Theorem [2.6 is called Palm measure. By (9) $\mu_{\mathcal{P}}$ can be interpreted as the push-forward measure of $g(x) \mathrm{d} \mu_{\omega}(x) \mathrm{dP}(\omega)$ under $(x, \omega) \mapsto \tau_{x} \omega$. Stationarity implies that this push-forward is independent of the choice of $g$. We say that the random measure $\mu_{\omega}$ has finite intensity if

$$
\begin{equation*}
+\infty>\int_{\Omega} \int_{\mathbb{R}^{m}} \chi_{\Omega \times[0,1]^{m}}\left(\tau_{x} \omega, x\right) \mathrm{d} \mu_{\omega}(x) \mathrm{dP}(\omega)=\mu_{\mathcal{P}}(\Omega) . \tag{11}
\end{equation*}
$$

Definition 2.7. Given a stationary random measure $\mu_{\omega}$, we introduce the scaled measure $\mu_{\omega}^{\varepsilon}$ through

$$
\begin{equation*}
\mu_{\omega}^{\varepsilon}(A):=\varepsilon^{m} \mu_{\omega}\left(\varepsilon^{-1} A\right) \tag{12}
\end{equation*}
$$

One important property of random measures is the following generalization of the Birkhoff ergodic theorem.

Lemma 2.8. ([14], Lemma 2.14) Let Assumption 2.5 hold for $(\Omega, \mathcal{F}, \mathbb{P}, \tau)$. Let $\boldsymbol{Q} \subset \mathbb{R}^{m}$ be a bounded domain, $\phi \in C(\overline{\boldsymbol{Q}})$ and $f \in L^{1}\left(\Omega ; \mu_{\mathcal{P}}\right)$. Then, for almost every $\omega \in \Omega$

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \int_{Q} \phi(x) f\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} \mu_{\omega}^{\varepsilon}(x)=\int_{Q} \int_{\Omega} \phi(x) f(\tilde{\omega}) \mathrm{d} \mu_{\mathcal{P}}(\tilde{\omega}) \mathrm{d} x \tag{13}
\end{equation*}
$$

A further useful result towards this direction is the following.
Lemma 2.9. ([14], Lemma 2.15) Let Assumption 2.5 hold for $(\Omega, \mathcal{F}, \mathbb{P}, \tau)$. Let $\boldsymbol{Q} \subset \mathbb{R}^{m}$ be a bounded domain and let $f \in L^{\infty}\left(\boldsymbol{Q} \times \Omega ; \mathcal{L} \otimes \mu_{\mathcal{P}}\right)$. Then, f has a $\mathcal{B}(\boldsymbol{Q}) \otimes \mathcal{F}$-measurable representative which is an ergodic function in the sense that for almost every $\omega \in \Omega$

$$
\begin{align*}
\lim _{\varepsilon \rightarrow 0} \int_{Q} f\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} \mu_{\omega}^{\varepsilon}(x) & =\int_{Q} \int_{\Omega} f(x, \tilde{\omega}) \mathrm{d} \mu_{\mathcal{P}}(\tilde{\omega}) \mathrm{d} x \\
\lim _{\varepsilon \rightarrow 0} \int_{Q}\left|f\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right)\right|^{p} \mathrm{~d} \mu_{\omega}^{\varepsilon}(x) & =\int_{Q} \int_{\Omega}|f(x, \tilde{\omega})|^{p} \mathrm{~d} \mu_{\mathcal{P}}(\tilde{\omega}) \mathrm{d} x \tag{14}
\end{align*}
$$

for every $1 \leq p<\infty$.

Based on the previous lemma, we can get the following result:
Lemma 2.10. Let Assumption 2.5 hold for $(\Omega, \mathcal{F}, \mathbb{P}, \tau)$. Let $\boldsymbol{Q} \subset \mathbb{R}^{m}$ be a bounded domain and let $f \in L^{\infty}\left(\boldsymbol{Q} \times \Omega ; \mathcal{L} \otimes \mu_{\mathcal{P}}\right)$. Then, $f$ has a $\mathcal{B}(\boldsymbol{Q}) \otimes \mathcal{F}$-measurable representative which is an ergodic function in the sense that for almost every $\omega \in \Omega$ and for all $\varphi \in C(\overline{\boldsymbol{Q}})$ it holds

$$
\begin{align*}
\lim _{\varepsilon \rightarrow 0} \int_{\boldsymbol{Q}} f\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right) \varphi(x) \mathrm{d} \mu_{\omega}^{\varepsilon}(x) & =\int_{\boldsymbol{Q}} \int_{\Omega} f(x, \tilde{\omega}) \varphi(x) \mathrm{d} \mu_{\mathcal{P}}(\tilde{\omega}) \mathrm{d} x \\
\lim _{\varepsilon \rightarrow 0} \int_{\boldsymbol{Q}}\left|f\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right)\right|^{p} \varphi(x) \mathrm{d} \mu_{\omega}^{\varepsilon}(x) & =\int_{\boldsymbol{Q}} \int_{\Omega}|f(x, \tilde{\omega})|^{p} \varphi(x) \mathrm{d} \mu_{\mathcal{P}}(\tilde{\omega}) \mathrm{d} x \tag{15}
\end{align*}
$$

for every $1 \leq p<\infty$.

Proof. This follows from the fact that $C(\overline{\boldsymbol{Q}})$ is separable and Lemma 2.9 yields (15) for a countable subset of $C(\overline{\boldsymbol{Q}})$ and a set of full measure $\tilde{\Omega} \subset \Omega$. By an approximation $\left\|\varphi-\varphi_{\delta}\right\|_{\infty}<\delta$ and Lemma 2.9 we obtain the claim.

### 2.2 Random measures and random sets.

In this paper, we consider random sets of the following form. For every $\omega \in \Omega$ the set $G(\omega)$ is an open subset of $\mathbb{R}^{m}$. The boundary $\Gamma(\omega)=\partial G(\omega)$ is a $(m-1)$-dimensional piece-wise Lipschitz manifold. Furthermore, we assume that the measures

$$
\mu_{\omega}(A):=\int_{A \cap G^{\mathrm{C}}(\omega)} \mathrm{d} x, \quad \mu_{\Gamma(\omega)}(A):=\mathcal{H}^{m-1}(A \cap \Gamma(\omega))
$$

are stationary. Hence, by Theorem [2.6 there exist corresponding Palm measures $\mu_{\mathcal{P}}$ for $\mu_{\omega}$ and $\mu_{\Gamma, \mathcal{P}}$ for $\mu_{\Gamma(\omega)}$ and by Lemma 2.14 of [13] there exists a measurable set $\Gamma \subset \Omega$ with $\chi_{\Gamma(\omega)}(x)=\chi_{\Gamma}\left(\tau_{x} \omega\right)$ for $\mathcal{L}+\mu_{\Gamma(\omega)}$-almost every $x$ for $\mathbb{P}$-almost every $\omega$ and $\mathbb{P}(\Gamma)=0$, $\mu_{\Gamma, \mathcal{P}}(\Omega \backslash \Gamma)=0$. Also it was observed there that, if for every $\omega$ we have $\mu_{\omega}=\mathcal{L}$, then also $\mu_{\mathcal{P}}=\mathbb{P}$. From the corresponding proofs in [13], as well as the fact that $\mu_{\omega}$ has a RadonNikodym derivative with respect to $\mathcal{L}$, we find $G \subset \Omega$ such that $\mu_{\mathcal{P}}(A)=\mathbb{P}\left(A \cap G^{\complement}\right)$, $\chi_{G^{\mathrm{C}}(\omega)}(x)=\chi_{G^{\mathrm{C}}}\left(\tau_{x} \omega\right)$ and

$$
\begin{equation*}
\chi_{G^{\mathrm{c}}} \mathrm{~d} \mu_{\mathcal{P}}=d \mu_{\mathcal{P}}=\chi_{G^{\mathrm{c}}} \mathrm{dP} . \tag{16}
\end{equation*}
$$

Remark 2.11. If $A$ is a bounded Borel set, then

$$
\begin{equation*}
\mu_{\Gamma(\omega)}^{\varepsilon}(A):=\varepsilon^{m} \mu_{\Gamma(\omega)}\left(\varepsilon^{-1} A\right)=\varepsilon \mathcal{H}^{m-1}\left(A \cap \Gamma^{\varepsilon}(\omega)\right) . \tag{17}
\end{equation*}
$$

It was shown in [13] that for random measures such as $\mu_{\omega}$ or $\mu_{\Gamma(\omega)}$ the underlying probability space can be assumed to be separable and metric, since the boundedly finite Borel measures equipped with the Vague topology form a separable metric space [8]. It was also pointed out in [13] that $\tau:(x, \omega) \mapsto \tau_{x} \omega$ is continuous.

Remark 2.12. If $\Omega$ is separable and metric, this implies that $L^{2}(\Omega ; \mathbb{P})$ and $L^{2}\left(\Omega, \mu_{\Gamma, P}\right)$ are separable and that the bounded continuous functions $C_{b}(\Omega)$ are dense in both spaces. Therefore, there exists a countable set $\Psi:=\left(\psi_{i}\right)_{i \in \mathbb{N}}$ such that $\psi_{i} \in C_{b}(\Omega)$ for every $i$ and such that $\Psi$ lies dense in $L^{2}(\Omega ; \mathbb{P})$ and $L^{2}\left(\Omega, \mu_{\Gamma, \mathcal{P}}\right)$. Furthermore, recalling (7) and approximating $\psi_{i}$ with the sequence $\mathscr{I}_{\frac{1}{n}} \psi, n \in \mathbb{N}$, we can assume that $\psi_{i} \in W^{1,2}(\Omega) \cap$ $C_{b}(\Omega)$. The space $\mathcal{V}_{\mathrm{pot}}^{2}(\Omega)$ is a subspace of a separable space and hence has to be separable, too. In particular $\nabla \psi_{i}$ can be assumed to be dense in $\mathcal{V}_{\text {pot }}^{2}(\Omega)$. We then define

$$
\Psi=\left(\psi_{i}\right)_{i \in \mathbb{N}} \bigcup_{j=1}^{m}\left(\mathrm{D}_{j} \psi_{i}\right)_{i \in \mathbb{N}}
$$

Since $\Omega$ is assumed to be separable metric, we can also make the following definition.

Definition 2.13. The space of bounded continuously differentiable functions on $\Omega$ is

$$
\begin{aligned}
C_{b}^{1}(\Omega) & :=\left\{f \in C_{b}(\Omega): \operatorname{D} f \in C_{b}(\Omega)\right\} \\
\|f\|_{C_{b}^{1}(\Omega)} & :=\|f\|_{\infty}+\|\mathrm{D} f\|_{\infty} .
\end{aligned}
$$

Let us remark that, since $(x, \omega) \mapsto \tau_{x} \omega$ is continuous, $f \in C_{b}^{1}(\Omega)$ implies $f\left(\tau_{x} \omega\right) \in$ $C_{b}^{1}\left(\mathbb{R}^{m}\right) \forall \omega \in \Omega$. Concerning the random geometries considered in this work, we make the assumptions listed below.

Definition 2.14 (See [11]). An open set $G \subset \mathbb{R}^{m}$ is said to be minimally smooth with constants $(\delta, N, M)$ if we may cover $\Gamma=\partial G$ by a countable sequence of open sets $\left(U_{i}\right)_{i \in \mathbb{N}}$ such that

1) Each $x \in \mathbb{R}^{m}$ is contained in at most $N$ of the open sets $U_{i}$.
2) For any $x \in \Gamma$, the ball $B_{\delta}(x)$ is contained in at least one $U_{i}$.
3) For any $i$, the portion of the boundary $\Gamma$ inside $U_{i}$ agrees (in some Cartesian system of coordinates) with the graph of a Lipschitz function whose Lipschitz semi-norm is at most $M$.

In particular a set $G \subset \mathbb{R}^{m}$ is minimally smooth if and only if $\mathbb{R}^{m} \backslash G$ is minimally smooth.

Let $\boldsymbol{Q}$ be a bounded domain in $\mathbb{R}^{m}$. For given constants $(\delta, N, M)$, we consider $G(\omega)$ a random open set which is a.s. minimally smooth with constants ( $\delta, N, M$ ) (uniformly minimally smooth). We furthermore assume that $G(\omega):=\bigcup_{i \in \mathbb{N}} G_{i}(\omega)$ is a countable union of disjoint open balls $G_{i}(\omega)$ with a maximal diameter $d_{0}$.

We then consider $G^{\varepsilon}(\omega):=\varepsilon G(\omega)$ and

$$
\begin{equation*}
\boldsymbol{Q}^{\varepsilon}(\omega):=\boldsymbol{Q} \backslash\left(\bigcup_{i \in I_{\varepsilon}(\omega)} \varepsilon G_{i}(\omega)\right), \quad \Gamma_{\boldsymbol{Q}}^{\varepsilon}(\omega):=\bigcup_{i \in I_{\varepsilon}(\omega)} \partial\left(\varepsilon G_{i}(\omega)\right) \tag{18}
\end{equation*}
$$

where

$$
I_{\varepsilon}(\omega):=\left\{i: \varepsilon G_{i}(\omega) \subset \boldsymbol{Q} \text { and } \varepsilon d_{0}<\min \left\{d(x, y): x \in \partial\left(\varepsilon G_{i}(\omega)\right), y \in \partial \boldsymbol{Q}\right\}\right\}
$$

Remark 2.15. Note that we constructed the micro structures $\boldsymbol{Q} \backslash \boldsymbol{Q}^{\varepsilon}(\omega)$ such that they do not intersect with the boundary of $\boldsymbol{Q}$ and such that every hole in $\boldsymbol{Q}^{\varepsilon}(\omega)$ has a minimal distance $\varepsilon d_{0}$ to $\partial \boldsymbol{Q}$. This is because we require in our proofs that $\varepsilon^{-1} \boldsymbol{Q}^{\varepsilon}(\omega)$ is a $(\delta, N, M)$ - minimal set (or $\boldsymbol{Q}^{\varepsilon}(\omega)$ is a $\left(\delta \varepsilon, N, \varepsilon^{-1} M\right)$ minimal set, respectively). In particular, without the minimal distance between two disjoint parts of the boundary, the resulting set $\boldsymbol{Q}^{\varepsilon}(\omega)$ would violate condition 3) from Definition 2.14, i.e. $\partial \boldsymbol{Q}^{\varepsilon}(\omega)$ would not be a $\varepsilon^{-1} M$-Lipschitz graph inside balls of diameter $\frac{\varepsilon}{2} d_{0}$.

Assumption 2.16. There are constants $d_{0}, \delta, N, M$ (independent of $\omega$ ) such that $\mathbb{P}$-a.s. the set $G(\omega)$ consists of a countable union of bounded sets $G_{k}(\omega)(k \in \mathbb{N})$ such that the sets $\mathbb{R}^{m} \backslash G_{k}(\omega)$ are all connected, while

$$
d\left(G_{k}(\omega), G_{j}(\omega)\right) \geq d_{0} \quad \text { whenever } \quad k \neq j
$$

and each set $G_{k}(\omega)$ is minimally smooth with constants $(\delta, N, M)$ and has a diameter smaller than $d_{0}$. The Lipschitz constant is uniformly over all $G_{k}$.

Remark 2.17. In particular, this guarantees that $\mathbb{R}^{m} \backslash G(\omega)$ is connected and has a Lipschitz boundary $\partial G$, which represents the union of the boundaries of the holes. Furthermore, the distance condition ensures that the boundary of $G(\omega)$ is locally representable as a graph.

Lemma 2.18. Suppose that Assumption [2.16 is satisfied. Then, there exists a family of linear continuous extension operators

$$
\mathcal{E}_{\varepsilon}: W^{1, p}\left(\boldsymbol{Q}^{\varepsilon}\right) \rightarrow W^{1, p}(\boldsymbol{Q})
$$

and a constant $C>0$ independent of $\varepsilon$ such that

$$
\mathcal{E}_{\varepsilon} \phi=\phi \text { in } \boldsymbol{Q}^{\varepsilon}(\omega)
$$

and

$$
\begin{gather*}
\int_{Q}\left|\mathcal{E}_{\varepsilon} \phi\right|^{p} d x \leq C \int_{Q^{\varepsilon}}|\phi|^{p} d x  \tag{19}\\
\int_{Q}\left|\nabla\left(\mathcal{E}_{\varepsilon} \phi\right)\right|^{p} d x \leq C \int_{Q^{\varepsilon}}|\nabla \phi|^{p} d x \tag{20}
\end{gather*}
$$

$\mathbb{P}$-a.s. for any $\phi \in W^{1, p}\left(\boldsymbol{Q}^{\varepsilon}\right)$ and for any $p \in(1,+\infty)$.

Proof. Following the line of the proof reported in 11 (Proposition 3.3, p. 230), for any $k \in \mathbb{N}, \omega \in \Omega$, we denote by $\hat{G}_{k}(\omega)$ a $d_{0} / 4$-neighborhood of $G_{k}(\omega)$ (the sets $G_{k}(\omega)$ are defined in Assumption [2.16). Since, under our assumptions, the set $\hat{G}_{k}(\omega) \backslash G_{k}(\omega)$ has Lipschitz boundary, then, according to Theorem 5, p. 181 in [28], there exists an extension operator $E_{k}$

$$
\begin{equation*}
E_{k}: W^{1, p}\left(\hat{G}_{k}(\omega) \backslash G_{k}(\omega)\right) \rightarrow W^{1, p}\left(\hat{G}_{k}(\omega)\right) \tag{21}
\end{equation*}
$$

such that: $E_{k} \phi=\phi$ a.e. in $\hat{G}_{k}(\omega) \backslash G_{k}(\omega)$ and, for some constant $C$ independent of $k$, we have

$$
\begin{equation*}
\left\|E_{k} \phi\right\|_{L^{p}\left(\hat{G}_{k}(\omega)\right)} \leq C\|\phi\|_{L^{p}\left(\hat{G}_{k}(\omega) \backslash G_{k}(\omega)\right)} \tag{22}
\end{equation*}
$$

$$
\begin{equation*}
\left\|E_{k} \phi\right\|_{W^{1, p}\left(\hat{G}_{k}(\omega)\right)} \leq C\|\phi\|_{W^{1, p}\left(\hat{G}_{k}(\omega) \backslash G_{k}(\omega)\right)} . \tag{23}
\end{equation*}
$$

Let us define new extensions

$$
\begin{equation*}
\hat{E}_{k}: W^{1, p}\left(\hat{G}_{k}(\omega) \backslash G_{k}(\omega)\right) \rightarrow W^{1, p}\left(G_{k}(\omega)\right) \tag{24}
\end{equation*}
$$

by

$$
\begin{equation*}
\hat{E}_{k} \phi:=E_{k}\left(\phi-(\phi)_{k}\right)+(\phi)_{k} \tag{25}
\end{equation*}
$$

where

$$
\begin{equation*}
(\phi)_{k}:=\int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)} \phi d y \tag{26}
\end{equation*}
$$

Putting them all together, we define an extension

$$
\begin{equation*}
\mathcal{E}: W^{1, p}\left(G^{\complement}(\omega)\right) \rightarrow W^{1, p}(\boldsymbol{Q}) \tag{27}
\end{equation*}
$$

given by

$$
\mathcal{E} \phi(y):=\left\{\begin{align*}
\phi(y) & \text { whenever } y \in G^{\complement}(\omega)  \tag{28}\\
\hat{E}_{k} \phi(y) & \text { whenever } y \in \hat{G}_{k}(\omega) .
\end{align*}\right.
$$

Now, in $\hat{G}_{k}(\omega) \backslash G_{k}(\omega)$ we have

$$
\begin{equation*}
\hat{E}_{k} \phi=\left(\phi-(\phi)_{k}\right)+(\phi)_{k}=\phi \tag{29}
\end{equation*}
$$

Moreover, by (22) and Hölder's inequality, we have

$$
\begin{align*}
& \int_{G_{k}(\omega)}\left|\hat{E}_{k} \phi\right|^{p} d y=\int_{G_{k}(\omega)}\left|E_{k}\left(\phi-(\phi)_{k}\right)+(\phi)_{k}\right|^{p} d y \\
& \leq C \int_{G_{k}(\omega)}\left|E_{k}\left(\phi-(\phi)_{k}\right)\right|^{p} d y+C \int_{G_{k}(\omega)}\left|(\phi)_{k}\right|^{p} d y \\
& \leq C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}\left|\phi-(\phi)_{k}\right|^{p} d y+C \int_{G_{k}(\omega)}\left|(\phi)_{k}\right|^{p} d y  \tag{30}\\
& \leq C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\phi|^{p} d y+C^{\prime} \int_{G_{k}(\omega)}\left|(\phi)_{k}\right|^{p} d y \\
& \leq C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\phi|^{p} d y
\end{align*}
$$

where, for simplicity, the letter $C$ denotes a positive constant (independent of $k$ ) that can change from line to line. Due to Assumption [2.16, the following Poincaré inequality holds:

$$
\begin{equation*}
\int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}\left|\phi-(\phi)_{k}\right|^{p} d y \leq C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\nabla \phi|^{p} d y \tag{31}
\end{equation*}
$$

Therefore, by using (231), (25) and (31), we get

$$
\begin{align*}
& \int_{G_{k}(\omega)}\left|\nabla\left(\hat{E}_{k} \phi\right)\right|^{p} d y=\int_{G_{k}(\omega)}\left|\nabla\left(E_{k}\left(\phi-(\phi)_{k}\right)\right)\right|^{p} d y \\
& \leq C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}\left|\left(\phi-(\phi)_{k}\right)\right|^{p} d y+C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\nabla \phi|^{p} d y  \tag{32}\\
& \leq C \int_{\hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\nabla \phi|^{p} d y
\end{align*}
$$

Since this holds for every $k$ with the same $C$ we have proved that

$$
\begin{gather*}
\int_{\cup_{k} \hat{G}_{k}(\omega)}|\mathcal{E} \phi|^{p} d y \leq C \int_{\cup_{k} \hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\phi|^{p} d y  \tag{33}\\
\int_{\cup_{k} \hat{G}_{k}(\omega)}|\nabla(\mathcal{E} \phi)|^{p} d y \leq C \int_{\cup_{k} \hat{G}_{k}(\omega) \backslash G_{k}(\omega)}|\nabla \phi|^{p} d y \tag{34}
\end{gather*}
$$

that is,

$$
\begin{gather*}
\int_{Q}|\mathcal{E} \phi|^{p} d y \leq C \int_{G^{\mathrm{C}}(\omega)}|\phi|^{p} d y  \tag{35}\\
\int_{Q}|\nabla(\mathcal{E} \phi)|^{p} d y \leq C \int_{G^{\mathrm{C}}(\omega)}|\nabla \phi|^{p} d y . \tag{36}
\end{gather*}
$$

By performing the change of variable $y=x / \varepsilon$, with $x \in \boldsymbol{Q}^{\varepsilon}(\omega)$, it is easy to obtain the corresponding re-scaled estimates (19) and (20), where $\mathcal{E}_{\varepsilon}$ is the re-scaled extension operator.

As a matter of fact, we can describe a portion of the cerebral cortex as a bounded open set $\boldsymbol{Q} \subset \mathbb{R}^{3}$, whereas the neurons are represented by a family of holes distributed
randomly in $\boldsymbol{Q}$ and having a characteristic size $\varepsilon$. A detailed construction of random domains that satisfy the assumptions listed in this section is reported in Appendix A.

## 3 Two-scale convergence

We will use a slightly modified version of stochastic two-scale convergence compared to the one presented in [14]. Let $\Psi:=\left(\psi_{i}\right)_{i \in \mathbb{N}}$ be the countable dense family of $C_{b}(\Omega)$ functions according to Remark 2.12,

Lemma 3.1. Let $\left(f_{i}\right)_{i \in \mathbb{N}}$ be a countable family in $L^{\infty}(\boldsymbol{Q} \times \Omega ; \mathcal{L} \times \mathbb{P})$ and $\left(g_{i}\right)_{i \in \mathbb{N}}$ be a countable family in $L^{\infty}\left(\boldsymbol{Q} \times \Gamma ; \mathcal{L} \times \mu_{\Gamma, \mathcal{P}}\right)$. Then there exists a set of full measure $\Omega_{\Psi} \subset \Omega$ such that for almost every $\omega \in \Omega_{\Psi}$, every $i \in \mathbb{N}$, every $\psi \in \Psi$ and every $\varphi \in C_{b}(\overline{\boldsymbol{Q}})$ the following holds:

$$
\begin{gather*}
\lim _{\varepsilon \rightarrow 0} \int_{\boldsymbol{Q}} \varphi^{2}(x) \psi^{2}\left(\tau_{\frac{x}{\varepsilon}} \omega\right) f_{i}^{2}\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} x=\int_{\boldsymbol{Q}} \int_{\Omega} \varphi^{2}(x) \psi^{2}(\tilde{\omega}) f_{i}^{2}(x, \tilde{\omega}) \mathrm{dP}(\tilde{\omega}) \mathrm{d} x  \tag{37}\\
\lim _{\varepsilon \rightarrow 0} \int_{\boldsymbol{Q}} g_{i}\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right) \varphi(x) \psi\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} \mu_{\Gamma(\omega)}^{\varepsilon}(x)=\int_{\boldsymbol{Q}} \int_{\Omega} g_{i}(x, \tilde{\omega}) \varphi(x) \psi(\tilde{\omega}) \mathrm{d} \mu_{\Gamma, \mathcal{P}}(\tilde{\omega}) \mathrm{d} x . \tag{38}
\end{gather*}
$$

Remark 3.2. The first equality (37) is needed for the proof of existence of the two-scale limits. Therefore we put the square here. The second limit (38) is needed directly in the proof of the main homogenization theorem. Therefore we study the convergence of $g_{i}$ tested with $\varphi \psi$.

Proof. (Proof of Lemma (3.1) For fixed $i$ the limits (37) and (38) hold for a.e. $\omega \in \Omega$ due to Lemma 2.10. Since the family $\left(f_{i}\right)_{i \in \mathbb{N}}$ is countable, we conclude.

Definition 3.3. Let $\Psi$ be the set of Remark 2.12 and let $\omega \in \Omega_{\Psi}$. Let $u^{\varepsilon} \in L^{2}(\boldsymbol{Q})$ for all $\varepsilon>0$. We say that $\left(u^{\varepsilon}\right)$ converges (weakly) in two scales to $u \in L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega)\right)$ and write $u^{\varepsilon} \stackrel{2 s}{\longrightarrow} u$ if $\sup _{\varepsilon>0}\left\|u^{\varepsilon}\right\|_{L^{2}(\boldsymbol{Q})}<\infty$ and if for every $\psi \in \Psi, \varphi \in C(\overline{\boldsymbol{Q}})$ there holds with $\phi_{\omega, \varepsilon}(x):=\varphi(x) \psi\left(\tau_{\frac{x}{\varepsilon}} \omega\right)$ that

$$
\lim _{\varepsilon \rightarrow 0} \int_{\boldsymbol{Q}} u^{\varepsilon}(x) \phi_{\omega, \varepsilon}(x) \mathrm{d} x=\int_{\boldsymbol{Q}} \int_{\Omega} u(x, \tilde{\omega}) \varphi(x) \psi(\tilde{\omega}) \mathrm{dP}(\tilde{\omega}) \mathrm{d} x
$$

Furthermore, we say that $u^{\varepsilon}$ converges strongly in two scales to $u$, written $u^{\varepsilon} \xrightarrow{2 s} u$, if for every weakly two-scale converging sequence $v^{\varepsilon} \in L^{2}(\boldsymbol{Q})$ with $v^{\varepsilon} \stackrel{2 s}{ } v \in L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega)\right)$ as $\varepsilon \rightarrow 0$ there holds

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \int_{Q} u^{\varepsilon} v^{\varepsilon} \mathrm{d} x=\int_{Q} \int_{\Omega} u v \mathrm{~d} \mathbb{P}(\tilde{\omega}) \mathrm{d} x \tag{39}
\end{equation*}
$$

Remark 3.4. Let us remark that the notion of two-scale convergence strongly depends on the choice of $\omega$. Also, let us note that $\phi_{\omega, \varepsilon} \xrightarrow{2 s} \varphi \psi$ strongly in two scales by definition.

Lemma 3.5. ([14], Lemma 4.4-1.) Let $u^{\varepsilon} \in L^{2}(\boldsymbol{Q})$ be a sequence of functions such that $\left\|u^{\varepsilon}\right\|_{L^{2}(\boldsymbol{Q})} \leq C$ for some $C>0$ independent of $\varepsilon$. Then there exists a subsequence $\left(u^{\varepsilon^{\prime}}\right)_{\varepsilon^{\prime} \rightarrow 0}$ and $u \in L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega)\right)$ such that $u^{\varepsilon^{\prime}} \xrightarrow{2 s} u$ and

$$
\begin{equation*}
\|u\|_{L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega)\right)} \leq \liminf _{\varepsilon^{\prime} \rightarrow 0}\left\|u^{\varepsilon^{\prime}}\right\|_{L^{2}(\boldsymbol{Q})} . \tag{40}
\end{equation*}
$$

Furthermore, let $\left(f_{i}\right)_{i \in \mathbb{N}}$ be a family of functions such as in Lemma 3.1. Then for every $i \in \mathbb{N}, \varphi \in C(\bar{\Omega})$ and $\psi \in \Psi$ it holds

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \int_{\boldsymbol{Q}} u^{\varepsilon}(x) \phi_{\omega, \varepsilon}(x) f_{i}\left(x, \tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} x=\int_{\boldsymbol{Q}} \int_{\Omega} u(x, \tilde{\omega}) \varphi(x) \psi(\tilde{\omega}) f_{i}(x, \tilde{\omega}) \mathrm{d} \mathbb{P}(\tilde{\omega}) \mathrm{d} x \tag{41}
\end{equation*}
$$

Proof. Let $\left(\varphi_{j}\right)_{j \in \mathbb{N}}$ be a countable dense subset of $C(\overline{\boldsymbol{Q}})$ and write $\Psi=\left(\psi_{k}\right)_{k \in \mathbb{N}}$. Then the span of $\varphi_{j} \psi_{k} f_{i}$ is dense in $L^{2}(\boldsymbol{Q} \times \Omega)$ (assuming w.l.o.g. that $1 \in\left(f_{i}\right)_{i \in \mathbb{N}}$ ). Thus (41) follows from [14], Lemma 4.4-1, using (37), for all $\varphi_{j} \psi_{k} f_{i}$. The statement follows eventually from a density argument to conclude for general $\varphi \in C(\overline{\boldsymbol{Q}})$.

Remark 3.6. As already observed in [14], Lemma 3.5 implies that for every $f \in L^{\infty}(\Omega)$, the class of test-functions $\Psi$ can be enriched by a countable subset $f \Psi \subset L^{2}(\Omega)$ changing $\Omega_{\Psi}$ only by a set of measure 0 .

We note that the definition of two-scale convergence in [14] is formulated in a different way. However, due to Lemma 4.6 of [14], we can recover our Definition 3.3. In particular, the original version of Lemma 3.5 yields two-scale convergence in the sense of 14 ] [Definition 4.2], and by Lemma 4.6 of [14] one infers Lemma 3.5. Finally, if $\Omega$ is compact, we recover the statements of [32] by separability of $C_{b}(\Omega)=C(\Omega)$.

Lemma 3.7. There exists $\tilde{\Omega} \subset \Omega_{\Psi}$ of full measure such that for all $\omega \in \tilde{\Omega}$ the following holds: If $u^{\varepsilon} \in H^{1}\left(\boldsymbol{Q} ; \mathbb{R}^{m}\right)$ for all $\varepsilon$, with $\left\|\nabla u^{\varepsilon}\right\|_{L^{2}(\boldsymbol{Q})}<C$ for $C$ independent from $\varepsilon>$ 0 , then there exists a subsequence denoted by $u^{\varepsilon}$, functions $u \in H^{1}\left(\boldsymbol{Q} ; \mathbb{R}^{m}\right)$ and $v \in$ $L^{2}\left(\boldsymbol{Q} ; L_{\text {pot }}^{2}(\Omega)\right)$ such that $u^{\varepsilon} \rightharpoonup u$ weakly in $H^{1}(\boldsymbol{Q})$ and

$$
\nabla u^{\varepsilon} \stackrel{2 s}{\longrightarrow} \nabla u+v \quad \text { as } \varepsilon \rightarrow 0
$$

The original version of the above Lemma in [14] was formulated in $H_{0}^{1}(\boldsymbol{Q})$. However, the proof applies for all sequences in $H^{1}(\boldsymbol{Q})$.

We are also interested in the convergence behavior of functions $u^{\varepsilon}:[0, T] \rightarrow L^{2}(\boldsymbol{Q})$. In particular, we provide the following definition:

Definition 3.8. Let $\Psi$ be the set of Remark [2.12, $\Lambda=\left(\varphi_{i}\right)_{i \in \mathbb{N}}$ be a countable dense subset of $C(\overline{\boldsymbol{Q}}), \omega \in \Omega_{\Psi}$ and $u^{\varepsilon} \in L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$ for all $\varepsilon>0$. We say that $\left(u^{\varepsilon}\right)$ converges (weakly) in two scales to $u \in L^{2}\left(0, T ; L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega, \mathbb{P})\right)\right.$ ), and write $u^{\varepsilon} \stackrel{2 s}{ } u$, if for all continuous and piece-wise affine functions $\phi:[0, T] \rightarrow \operatorname{span} \Psi \times \Lambda$ there holds, with

$$
\begin{aligned}
\phi_{\omega, \varepsilon}(t, x): & =\phi\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right), \\
& \lim _{\varepsilon \rightarrow 0} \int_{0}^{T} \int_{Q} u^{\varepsilon} \phi_{\omega, \varepsilon} d x d t=\int_{0}^{T} \int_{\boldsymbol{Q}} \int_{\Omega} u(t, x, \tilde{\omega}) \phi(t, x, \tilde{\omega}) d \mathbb{P}(\tilde{\omega}) d x d t .
\end{aligned}
$$

Note that the test functions now have values in the vector space span $\Psi$ since they are affine. Similar to the stationary case, we obtain the following lemma.

Lemma 3.9. ([14], Lemma 4.16) Let $T>0$. Then, every sequence $\left(u^{\varepsilon}\right)_{\varepsilon>0}$ with $u^{\varepsilon} \in$ $L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$ satisfying $\left\|u^{\varepsilon}\right\|_{L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)} \leq C$ for some $C>0$ independent from $\varepsilon$ has a weakly two-scale convergent subsequence with limit function $u \in L^{2}\left(0, T ; L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega, \mathbb{P})\right)\right)$. Furthermore, if $\left\|\partial_{t} u^{\varepsilon}\right\|_{L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)} \leq C$ uniformly for $1<p \leq \infty$, then also $\partial_{t} u \in$ $L^{2}\left(0, T ; L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega, \mathbb{P})\right)\right)$ and $\partial_{t} u^{\varepsilon} \xrightarrow{2 s} \partial_{t} u$ in the sense of Definition 3.8 as well as $u^{\varepsilon}(t) \xrightarrow{2 s}$ $u(t)$ for all $t \in[0, T]$.

As a special case of the last result, we have
Lemma 3.10. ([14], Lemma 4.17) Let $\Psi$ and $\Omega_{\Psi}$ be given by Remark 2.12 and $\omega \in \Omega_{\Psi}$. Let $u^{\varepsilon} \in C^{\text {Lip }}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$ for all $\varepsilon>0$ such that $\left\|u^{\varepsilon}\right\|_{C^{L i p}\left(0, T ; L^{2}(\boldsymbol{Q})\right)} \leq C$ for some $C$ independent from $\varepsilon>0$. Then, there exists $u \in C^{L i p}\left(0, T ; L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega, \mathbb{P})\right)\right)$ and a subsequence $u^{\varepsilon^{\prime}}$ of $u^{\varepsilon}$ such that $u^{\varepsilon^{\prime}}(t) \xrightarrow{2 s} u(t)$ for all $t \in[0, T]$.

### 3.1 Domains with holes.

Since $G(\omega)$ is a random set, there exists, by the considerations in Section 2.2, a set $G \subset \Omega$ such that $\chi_{G(\omega)}(x)=\chi_{G}\left(\tau_{x} \omega\right)$. Based on $G$, respectively its complement $G^{\complement}$, we obtain the following generalized concept of two-scale convergence.

Lemma 3.11. Let $u^{\varepsilon} \in L^{2}(\boldsymbol{Q})$ be a sequence of functions such that $\sup _{\varepsilon>0}\left\|u^{\varepsilon}\right\|_{L^{2}(\boldsymbol{Q})}<$ $\infty$. If $\left(u^{\varepsilon^{\prime}}\right)_{\varepsilon^{\prime} \rightarrow 0}$ is a subsequence such that $u^{\varepsilon^{\prime}} \xrightarrow{2 s} u$ for some $u \in L^{2}\left(\boldsymbol{Q} ; L^{2}(\Omega)\right)$, then $u^{\varepsilon} \chi_{\boldsymbol{Q}^{\varepsilon}} \stackrel{2 s}{ } u \chi_{G^{\mathrm{C}}}$.

Proof. Let $\left(u^{\varepsilon^{\prime}}\right)_{\varepsilon^{\prime} \rightarrow 0}$ be a subsequence such that $u^{\varepsilon^{\prime}} \stackrel{\text { 2s }}{ } u$. Then the definition of two-scale convergence in $L^{2}(\boldsymbol{Q})$ together with Remark 3.6 implies that, for every $\varphi \in C(\overline{\boldsymbol{Q}})$ and $\psi \in \Psi$, it holds

$$
\lim _{\varepsilon \rightarrow 0} \int_{Q} u^{\varepsilon}(x) \chi_{G^{\mathrm{c}}}\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \varphi(x) \psi\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} x=\int_{Q} \int_{\Omega} u(x, \tilde{\omega}) \chi_{G^{\mathrm{c}}}(\tilde{\omega}) \varphi(x) \psi(\tilde{\omega}) \mathrm{dP}(\tilde{\omega}) \mathrm{d} x
$$

Furthermore, for $\delta>0$, let us consider the ball $B_{\delta}(x)$ of radius $\delta$ and center $x$. For $\varepsilon>0$ small enough and with $\phi_{\omega, \varepsilon}(x):=\varphi(x) \psi\left(\tau_{\frac{x}{\varepsilon}} \omega\right)$ it holds that

$$
\begin{aligned}
& \left|\int_{\boldsymbol{Q}} u^{\varepsilon}(x)\left(\chi_{G^{\mathrm{C}}}\left(\tau_{\frac{x}{\varepsilon}} \omega\right)-\chi_{\boldsymbol{Q}^{\varepsilon}}(x)\right) \phi_{\omega, \varepsilon}(x) \mathrm{d} x\right| \\
& \leq\left(\int_{B_{\delta}(\partial \boldsymbol{Q})}\left|u^{\varepsilon}\right|^{2} d x\right)\left(\int_{B_{\delta}(\partial \boldsymbol{Q})} \varphi^{2}(x) \psi^{2}\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} x\right) \\
& \leq\left(\int_{B_{\delta}(\partial \boldsymbol{Q})} \varphi^{2}(x) \psi^{2}\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \mathrm{d} x\right) \sup _{\varepsilon>0}\left\|u^{\varepsilon}\right\|_{L^{2}(\boldsymbol{Q})} \\
& \rightarrow\left(\int_{B_{\delta}(\partial \boldsymbol{Q})} \int_{\Omega} \varphi^{2}(x) \psi^{2}(\omega) \mathrm{dP}(\omega) \mathrm{d} x\right) \sup _{\varepsilon>0}\left\|u^{\varepsilon}\right\|_{L^{2}(\boldsymbol{Q})}
\end{aligned}
$$

Since $\delta>0$ is arbitrary, the statement follows.

Lemma 3.12. Let $u^{\varepsilon} \in H^{1}\left(\boldsymbol{Q}^{\varepsilon}(\omega)\right)$ be a sequence of functions such that $\sup _{\varepsilon>0}\left\|u^{\varepsilon}\right\|_{H^{1}\left(\boldsymbol{Q}^{\varepsilon}(\omega)\right)}<$ $\infty$. Then there exist functions $u \in H^{1}(\boldsymbol{Q})$ and $v \in L^{2}\left(\boldsymbol{Q} ; L_{\mathrm{pot}}^{2}(\Omega)\right)$ such that $\mathcal{E}_{\varepsilon} u^{\varepsilon} \rightharpoonup u$ weakly in $H^{1}(\boldsymbol{Q})$ as well as $u^{\varepsilon} \stackrel{2 s}{ } \chi_{G^{\mathrm{c}}} u$ and $\nabla u^{\varepsilon} \stackrel{2 s}{ } \chi_{G^{\mathrm{c}}} \nabla u+\chi_{G^{\mathrm{c}}} v$.

Proof. Lemma 2.18 implies that $\sup _{\varepsilon>0}\left\|\mathcal{E}_{\varepsilon} u^{\varepsilon}\right\|_{H^{1}(\boldsymbol{Q})}<\infty$. Hence, due to Lemma 3.7 there exists $u \in H^{1}(\boldsymbol{Q})$ and $v \in L^{2}\left(\boldsymbol{Q} ; L_{\text {pot }}^{2}(\Omega)\right)$ such that $\mathcal{E}_{\varepsilon} u^{\varepsilon} \rightharpoonup u$ weakly in $H^{1}(\boldsymbol{Q})$ and $\nabla\left(\mathcal{E}_{\varepsilon} u^{\varepsilon}\right) \xrightarrow{2 s} \nabla u+v$. Lemma 3.11 now implies $u^{\varepsilon} \stackrel{2 s}{\longrightarrow} \chi_{G^{\mathrm{c}}} u$ and $\nabla u^{\varepsilon} \stackrel{2 s}{\longrightarrow} \chi_{G^{\mathrm{c}}} \nabla u+\chi_{G^{\mathrm{c}}} v$.

Lemma 3.13. Let $u^{\varepsilon} \in L^{2}\left(0, T ; H^{1}\left(\boldsymbol{Q}^{\varepsilon}(\omega)\right)\right)$ be a sequence of functions such that

$$
\sup _{\varepsilon>0}\left\|u^{\varepsilon}\right\|_{L^{2}\left(0, T ; H^{1}\left(\boldsymbol{Q}^{\varepsilon}(\omega)\right)\right)}+\left\|\partial_{t} u^{\varepsilon}\right\|_{L^{2}\left(0, T ; L^{2}\left(\boldsymbol{Q}^{\varepsilon}(\omega)\right)\right)}<\infty
$$

Then there exist functions $u \in L^{2}\left(0, T ; H^{1}(\boldsymbol{Q})\right)$ with $\partial_{t} u \in L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$ and $v \in$ $L^{2}\left(0, T ; L^{2}\left(\boldsymbol{Q} ; L_{\mathrm{pot}}^{2}(\Omega)\right)\right)$ such that $\mathcal{E}_{\varepsilon} u^{\varepsilon} \rightharpoonup u$ weakly in $L^{2}\left(0, T ; H^{1}(\boldsymbol{Q})\right)$ and $\mathcal{E}_{\varepsilon} u^{\varepsilon} \rightarrow u$ strongly in $L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$ as well as

$$
u^{\varepsilon} \stackrel{2 s}{-} \chi_{G^{\mathrm{c}}} u, \quad \partial_{t} u^{\varepsilon} \stackrel{2 s}{\sim} \chi_{G^{\mathrm{c}}} \partial_{t} u, \quad \text { and } \quad \nabla u^{\varepsilon} \stackrel{2 s}{\sim} \chi_{G^{\mathrm{c}}} \nabla u+\chi_{G^{\mathrm{c}}} v .
$$

Proof. We only have to prove $\mathcal{E}_{\varepsilon} u^{\varepsilon} \rightarrow u$ strongly in $L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$ since the remaining part of the statement has either been demonstrated above or can be obtained by generalizing previous considerations.

We first observe that, for every times $t_{1}, t_{2} \in[0, T]$, it holds by Lemma 2.18 that

$$
\begin{aligned}
\left\|\int_{t_{1}}^{t_{2}} \mathcal{E}_{\varepsilon} u^{\varepsilon}(t) \mathrm{d} t\right\|_{H^{1}(\boldsymbol{Q})} & \leq\left\|\mathcal{E}_{\varepsilon} \int_{t_{1}}^{t_{2}} u^{\varepsilon}(t) \mathrm{d} t\right\|_{H^{1}(\boldsymbol{Q})} \\
& \leq C\left\|\int_{t_{1}}^{t_{2}} u^{\varepsilon}(t) \mathrm{d} t\right\|_{H^{1}\left(\boldsymbol{Q}^{\varepsilon}\right)} \\
& \leq C T^{\frac{1}{2}}\left\|u^{\varepsilon}\right\|_{L^{2}\left(0, T ; H^{1}\left(\boldsymbol{Q}^{\varepsilon}\right)\right)}
\end{aligned}
$$

and hence $\left\{\int_{t_{1}}^{t_{2}} \mathcal{E}_{\varepsilon} u^{\varepsilon}(t) \mathrm{d} t\right\}_{\varepsilon>0}$ is precompact in $L^{2}(\boldsymbol{Q})$. Next, one can write by using again Lemma 2.18:

$$
\begin{aligned}
\int_{0}^{T-h}\left\|\mathcal{E}_{\varepsilon}\left(u^{\varepsilon}(t)-u^{\varepsilon}(t+h)\right)\right\|_{L^{2}(\boldsymbol{Q})}^{2} d t & \leq C \int_{0}^{T-h}\left\|u^{\varepsilon}(t)-u^{\varepsilon}(t+h)\right\|_{L^{2}\left(\boldsymbol{Q}^{\varepsilon}\right)}^{2} d t \\
& \leq C \int_{0}^{T-h}\left\|\int_{t}^{t+h} \partial_{t} u^{\varepsilon}(s) \mathrm{d} s\right\|_{L^{2}\left(\boldsymbol{Q}^{\varepsilon}\right)}^{2} d t \\
& \leq C \int_{0}^{T-h} h\left\|\partial_{t} u^{\varepsilon}\right\|_{L^{2}\left(t, t+h ; L^{2}\left(\boldsymbol{Q}^{\varepsilon}\right)\right)}^{2} d t \\
& \leq C h\left\|\partial_{t} u^{\varepsilon}\right\|_{L^{2}\left(0, T ; L^{2}\left(\boldsymbol{Q}^{\varepsilon}\right)\right)}^{2}
\end{aligned}
$$

where the constant $C$ changes in the last step. Since it holds $\mathcal{E}_{\varepsilon} u^{\varepsilon} \rightharpoonup u$ in $L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$, we conclude from Simon's compactness theorem (see Theorem 1 of [27]).

## 4 Setting of the problem and estimates

Throughout this paper, $\varepsilon$ will denote the general term of a sequence of positive reals which converges to zero. We consider in the following a system of anisotropic diffusioncoagulation Smoluchowski-type equations which describes the dynamics of cluster growth. In particular, we introduce the vector-valued random function $u^{\varepsilon}:[0, T] \times \boldsymbol{Q}^{\varepsilon} \rightarrow \mathbb{R}^{M}$, $u^{\varepsilon}=\left(u_{1}^{\varepsilon}, \ldots, u_{M}^{\varepsilon}\right)$ (with $M \in \mathbb{N}$ being fixed) where the variable $u_{s}^{\varepsilon} \geq 0(1 \leq s<$ $M)$ represents the concentration of $s$-clusters, that is, clusters consisting of $s$ identical elementary particles (monomers), while $u_{M}^{\varepsilon} \geq 0$ takes into account aggregations of more than $M-1$ monomers. We assume that the only reaction allowing clusters to coalesce to form larger clusters is a binary coagulation mechanism, while the movement of clusters results only from a diffusion process described by a stationary ergodic random matrix

$$
\left(d_{i, j}^{s}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right)\right)_{i, j=1, \ldots, m}=: D_{s}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \quad 1 \leq s \leq M,
$$

where $(t, x) \in[0, T] \times \boldsymbol{Q}$. Here $D_{s}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right)$ is the realization (see Remark 2.4) of a random matrix. Indeed, aging (as well as the AD itself) yields atrophy of the cerebral parenchyma, inducing changes in the diffusion rate of the amyloid agglomerates. In addition, this rate may vary for different regions of the brain. Finally, we have to take into account that $\mathrm{A} \beta$ aggregates do not diffuse freely in an uniform fluid: the cerebral tissue consists of large non-neuronal support cells (the macroglia) and the $\mathrm{A} \beta$ polymers move within the cerebrospinal fluid along the interstices between these cells that, in turn, are stochastically distributed.

With these notations, our system reads:

$$
\begin{cases}\frac{\partial u_{1}^{\varepsilon}}{\partial t}-\operatorname{div}\left(D_{1}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{1}^{\varepsilon}\right)+u_{1}^{\varepsilon} \sum_{j=1}^{M} a_{1, j} u_{j}^{\varepsilon}=0 & \text { in }[0, T] \times \boldsymbol{Q}^{\varepsilon}  \tag{42}\\ {\left[D_{1}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{1}^{\varepsilon}\right] \cdot n=0} & \text { on }[0, T] \times \partial \boldsymbol{Q} \\ {\left[D_{1}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{1}^{\varepsilon}\right] \cdot \nu_{\Gamma_{Q}^{\varepsilon}}=\varepsilon \eta\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right)} & \text { on }[0, T] \times \Gamma_{\boldsymbol{Q}}^{\varepsilon} \\ u_{1}^{\varepsilon}(0, x)=U_{1} & \text { in } \boldsymbol{Q}^{\varepsilon} \\ \text { if } 1<s<M & \end{cases}
$$


and eventually

$$
\begin{cases}\frac{\partial u_{M}^{\varepsilon}}{\partial t}-\operatorname{div}\left(D_{M}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{M}^{\varepsilon}\right)=g^{\varepsilon} & \text { in }[0, T] \times \boldsymbol{Q}^{\varepsilon}  \tag{44}\\ {\left[D_{M}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{M}^{\varepsilon}\right] \cdot n=0} & \text { on }[0, T] \times \partial \boldsymbol{Q} \\ {\left[D_{M}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{M}^{\varepsilon}\right] \cdot \nu_{\Gamma_{Q}}=0} & \text { on }[0, T] \times \Gamma_{\boldsymbol{Q}}^{\varepsilon} \\ u_{M}^{\varepsilon}(0, x)=0 & \text { in } \boldsymbol{Q}^{\varepsilon}\end{cases}
$$

where the gain terms $f^{\varepsilon}$ and $g^{\varepsilon}$ in (43) and (44) are given by

$$
\begin{gather*}
f^{\varepsilon}=\frac{1}{2} \sum_{j=1}^{s-1} a_{j, s-j} u_{j}^{\varepsilon} u_{s-j}^{\varepsilon}  \tag{45}\\
g^{\varepsilon}=\frac{1}{2} \sum_{\substack{j+k \geq M \\
k<M(\text { if } j=M) \\
j<M(\text { if } k=M)}} a_{j, k} u_{j}^{\varepsilon} u_{k}^{\varepsilon} . \tag{46}
\end{gather*}
$$

The kinetic coefficients $a_{i, j}$ represent a reaction in which an $(i+j)$-cluster is formed from an $i$-cluster and a $j$-cluster. Therefore, they can be interpreted as "coagulation rates" and are symmetric $a_{i, j}=a_{j, i}>0(i, j=1, \ldots, M)$, but $a_{M, M}=0$. Let us remark that the meaning of $u_{M}^{\varepsilon}$ differs from that of $u_{s}^{\varepsilon}(s<M)$, since it describes the sum of the densities of all the 'large' assemblies. It is assumed that large assemblies exhibit all the same coagulation properties and do not coagulate with each other.

The production of $\beta$-amyloid peptide by the malfunctioning neurons is described imposing a non-homogeneous Neumann condition on the boundary of the holes, randomly selected within our domain. To this end, we consider on $\Gamma_{Q}^{\varepsilon}$ in Eq. (42) a stationary ergodic random function $\eta=\eta\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right)$. Here $\eta\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right)$ is the realization (see Remark 2.4) of a random function:

$$
\begin{equation*}
\eta:[0, T] \times \overline{\boldsymbol{Q}} \times \Omega \rightarrow[0,1] \tag{47}
\end{equation*}
$$

where the value ' 0 ' is assigned to 'healthy' neurons while all the other values in $] 0,1$ ] indicate different degrees of malfunctioning. Moreover, we assume that $\eta$ is an increasing function of time, since once the neuron has become 'ill', it can no longer regain its original state of health.

Further hypotheses are listed below:
(H.1) the diffusion coefficients satisfy $d_{i, j}^{s} \in C^{1}\left([0, T] \times \boldsymbol{Q} ; C_{b}^{1}(\Omega)\right)$ for $i, j=1, \ldots, m$, $s=1, \ldots, M$. We put

$$
\Lambda^{\star}:=\sup _{i, j, s}\left\|d_{i, j}^{s}\right\|_{C^{1}\left([0, T] \times \boldsymbol{Q} ; C_{b}^{1}(\Omega)\right)}
$$

In particular, the map $(t, x, \omega) \rightarrow D_{s}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right)$ is continuously differentiable;
(H.2) $d_{i, j}^{s}=d_{j, i}^{s}$, for $i, j=1, \ldots, m, s=1, \ldots, M$;
(H.3) there exists $0<\lambda \leq \Lambda$ such that

$$
\lambda|\xi|^{2} \leq \sum_{i, j=1}^{m} d_{i, j}^{s}\left(t, x, \tau_{\frac{x}{e}} \omega\right) \xi_{i} \xi_{j} \leq \Lambda|\xi|^{2}
$$

for all $s=1, \ldots, M, \xi \in \mathbb{R}^{m},(t, x) \in[0, T] \times \overline{\boldsymbol{Q}}$ and for $\mathbb{P}$-a.e. $\omega \in \Omega$.

Moreover, the function $\eta$, appearing in (42), is a given bounded function satisfying the following conditions:
(Н.4) $\eta \in C^{1}\left([0, T] \times \boldsymbol{Q} ; C_{b}^{1}(\Omega)\right)$;
(H.5) $\eta(0, \cdot, \cdot)=0$ and $U_{1}$ is a positive constant such that

$$
\begin{equation*}
U_{1} \leq\|\eta\|_{L^{\infty}([0, T] \times \overline{\boldsymbol{Q}} \times \Omega)} \tag{48}
\end{equation*}
$$

We can repeat now almost verbatim the arguments of [10], Theorems 2.1, 2.2, 2.3 and 2.4 to obtain the following "deterministic" (i.e. for fixed $\omega \in \Omega$ ) existence and regularity result.

Theorem 4.1. Suppose Assumption 2.16 (where additionally $G(\omega)$ has a smooth boundary) and (H.1) - (H.5) hold. Then for $\mathbb{P}$-a.e. $\omega \in \Omega$ and for any $\varepsilon>0$ the system (42) (44) admits a unique maximal classical solution

$$
u_{\omega}^{\varepsilon}=\left(u_{\omega, 1}^{\varepsilon}, \ldots, u_{\omega, M}^{\varepsilon}\right)
$$

such that
(i) there exists $\alpha \in(0,1), \alpha$ depending only on $N, \lambda, \Lambda^{\star}, \varepsilon$ and $\omega$, such that $u^{\varepsilon} \in$ $C^{1+\alpha / 2,2+\alpha}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}, \mathbb{R}^{M}\right)$ for $\mathbb{P}$-a.e. $\omega \in \Omega$ and

$$
\begin{equation*}
\left\|u_{\omega}^{\varepsilon}\right\|_{C^{1+\alpha / 2,2+\alpha}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}, \mathbb{R}^{M}\right)} \leq C_{0}=C_{0}\left(U_{1},\|\eta\|_{L^{\infty}([0, T] \times \overline{\boldsymbol{Q}} \times \Omega)}, K, \varepsilon, \omega, \alpha\right) \tag{49}
\end{equation*}
$$

(ii) $u_{\omega, j}^{\varepsilon}(t, x)>0$ for $(t, x) \in[0, T] \times \boldsymbol{Q}^{\varepsilon}$, $\mathbb{P}$-a.e. $\omega \in \Omega$ and $j=1, \ldots, M$.

In the sequel we shall rely on the fact that statements that hold $\mathbb{P}$-a.e. can be seen as deterministic assertions, since they hold whenever $\boldsymbol{Q}^{\varepsilon}$ is a set enjoying the regularity properties described in Remark [2.15, Assumption 2.16 and Remark 2.17,

Arguing as in [10], the first and crucial step will consist of proving that the $u_{\omega, j}^{\varepsilon}$ are equibounded in $L^{\infty}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}\right)$ for $\mathbb{P}$-a.e. $\omega \in \Omega$ and $j=1, \ldots, M$.

In particular, an uniform bound for $u_{\omega}^{\varepsilon}$ in $L^{\infty}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}\right)$ is provided by the following statement:

Theorem 4.2. Let $u_{\omega}^{\varepsilon}=\left(u_{\omega, 1}^{\varepsilon}, \ldots, u_{\omega, M}^{\varepsilon}\right)$ be as in Theorem 4.1. Then

$$
\begin{equation*}
\left\|u_{\omega, 1}^{\varepsilon}\right\|_{L^{\infty}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}\right)} \leq\left|U_{1}\right|+c\|\eta\|_{L^{\infty}([0, T] \times \overline{\boldsymbol{Q}} \times \Omega)} \tag{50}
\end{equation*}
$$

for $\mathbb{P}$-a.e. $\omega \in \Omega$, where $c$ is independent of $\varepsilon>0$.
In addition, there exists $K>0$ such that

$$
\begin{equation*}
\left\|u_{\omega, j}^{\varepsilon}\right\|_{L^{\infty}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}\right)} \leq K \tag{51}
\end{equation*}
$$

for $\mathbb{P}$-a.e. $\omega \in \Omega$, uniformly with respect to $\varepsilon>0$.

Proof. Thanks to extension Lemma 2.18, the function $u_{\omega}^{\varepsilon}$ can be continued on all $[0, T] \times \boldsymbol{Q}$. Therefore we can repeat step by step the arguments of [10], Theorems 2.2 and 2.3, that in turn rely on [20] (see also [25] and [30]).

Therefore

Theorem 4.3 ([10], Theorems 3.1. and 3.2). The sequence $\left(\nabla_{x} u_{\omega, j}^{\varepsilon}\right)_{\varepsilon>0}(1 \leq j \leq M)$ is bounded in $L^{2}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}\right)$ for $\mathbb{P}$-a.e. $\omega \in \Omega$, uniformly in $\varepsilon$.

In addition, the sequence $\left(\partial_{t} u_{\omega, j}^{\varepsilon}\right)_{\varepsilon>0}(1 \leq j \leq M)$ is bounded in $L^{2}\left([0, T] \times \boldsymbol{Q}^{\varepsilon}\right)$ for $\mathbb{P}$-a.e. $\omega \in \Omega$, uniformly in $\varepsilon$.

## 5 Homogenization

Our main statement shows that it is possible to homogenize the set of Eqs. (42)-(44) as $\varepsilon \rightarrow 0$.

Theorem 5.1. Let $u_{s}^{\varepsilon}(t, x)(1 \leq s \leq M)$ be a family of nonnegative classical solutions to the system (42)-(44). Denote by a tilde the extension by zero outside $\boldsymbol{Q}^{\varepsilon}(\omega)$ and let $\chi_{G^{\mathrm{C}}}$ represent the characteristic function of the random set $G^{\complement}(\omega)$. Then, the sequences $\left(\widetilde{u_{s}^{\varepsilon}}\right)_{\varepsilon>0}$, $\left(\widetilde{\nabla_{x} u_{s}^{\varepsilon}}\right)_{\varepsilon>0}$ and $\left(\widetilde{\partial_{t} u_{s}^{\varepsilon}}\right)_{\varepsilon>0}(1 \leq s \leq M)$ stochastically two-scale converge to: $\left[\chi_{G^{\complement}} u_{s}(t, x)\right]$, $\left[\chi_{G^{\mathrm{C}}}\left(\nabla_{x} u_{s}(t, x)+v_{s}(t, x, \omega)\right)\right],\left[\chi_{G^{\mathrm{C}}} \partial_{t} u_{s}(t, x)\right](1 \leq s \leq M)$, respectively. The limiting functions $\left[(t, x) \mapsto u_{s}(t, x),(t, x, \omega) \mapsto v_{s}(t, x, \omega)\right](1 \leq s \leq M)$ are the unique solutions lying in $L^{2}\left(0, T ; H^{1}(\boldsymbol{Q})\right) \times L^{2}\left([0, T] \times \boldsymbol{Q} ; L_{\mathrm{pot}}^{2}(\Omega)\right)$ of the following two-scale homogenized systems:

$$
\text { If } s=1 \text { : }
$$

$$
\begin{cases}\theta \frac{\partial u_{1}}{\partial t}(t, x)-\operatorname{div}_{x}\left[D_{1}^{\star}(t, x) \nabla_{x} u_{1}(t, x)\right]  \tag{52}\\ +\theta u_{1}(t, x) \sum_{j=1}^{M} a_{1, j} u_{j}(t, x)=\int_{\Omega} \chi_{\Gamma_{G^{\mathrm{C}}}} \eta(t, x, \omega) d \mu_{\Gamma, \mathcal{P}}(\omega) & \text { in }[0, T] \times \boldsymbol{Q} \\ {\left[D_{1}^{\star}(t, x) \nabla_{x} u_{1}(t, x)\right] \cdot n=0} & \text { on }[0, T] \times \partial \boldsymbol{Q} \\ u_{1}(0, x)=U_{1} & \text { in } \boldsymbol{Q}\end{cases}
$$

If $1<s<M$ :

$$
\begin{cases}\theta \frac{\partial u_{s}}{\partial t}(t, x)-\operatorname{div_{x}}\left[D_{s}^{\star}(t, x) \nabla_{x} u_{s}(t, x)\right] &  \tag{53}\\ +\theta u_{s}(t, x) \sum_{j=1}^{M} a_{s, j} u_{j}(t, x) & \text { in }[0, T] \times \boldsymbol{Q} \\ =\frac{\theta}{2} \sum_{j=1}^{s-1} a_{j, s-j} u_{j}(t, x) u_{s-j}(t, x) & \text { on }[0, T] \times \partial \boldsymbol{Q} \\ & \text { in } \boldsymbol{Q}\end{cases}
$$

If $s=M$ :

$$
\begin{cases}\theta \frac{\partial u_{M}}{\partial t}(t, x)-\operatorname{div}_{x}\left[D_{M}^{\star}(t, x) \nabla_{x} u_{M}(t, x)\right] &  \tag{54}\\ =\frac{\theta}{2} \sum_{\substack{k<M(i f j=M) \\ j<M(i f k=M)}}^{j+k \geq M} a_{j, k} u_{j}(t, x) u_{k}(t, x) & \text { in }[0, T] \times \boldsymbol{Q} \\ {\left[D_{M}^{\star}(t, x) \nabla_{x} u_{M}(t, x)\right] \cdot n=0} & \text { on }[0, T] \times \partial \boldsymbol{Q} \\ u_{M}(0, x)=0 & \text { in } \boldsymbol{Q}\end{cases}
$$

where

$$
\theta=\int_{\Omega} \chi_{G^{\complement}} d \mu_{\mathcal{P}}(\omega)=\mathbb{P}\left(G^{\complement}\right)
$$

represents the fraction of volume occupied by $G^{\complement}$ and, for every $1 \leq s \leq M, D_{s}^{\star}(t, x)$ is a deterministic matrix, called "effective diffusivity", defined by

$$
\left(D_{s}^{\star}\right)_{i j}(t, x)=\int_{\Omega} \chi_{G^{\mathrm{c}}} D_{s}(t, x, \omega)\left(w_{i}(t, x, \omega)+\hat{e}_{i}\right) \cdot\left(w_{j}(t, x, \omega)+\hat{e}_{j}\right) d \mathbb{P}(\omega)
$$

with $\hat{e}_{i}$ being the $i$-th canonical unit vector in $\mathbb{R}^{m}$, and $\left(w_{i}\right)_{1 \leq i \leq m} \in L^{2}\left([0, T] \times \boldsymbol{Q} ; L_{\mathrm{pot}}^{2}\left(G^{\complement}\right)\right)$ the family of solutions of the following microscopic problem

$$
\left\{\begin{array}{lc}
-\operatorname{div}_{\omega}\left[D_{s}(t, x, \omega)\left(w_{i}(t, x, \omega)+\hat{e}_{i}\right)\right]=0 & \text { in } G^{\complement}  \tag{55}\\
D_{s}(t, x, \omega)\left[w_{i}(t, x, \omega)+\hat{e}_{i}\right] \cdot \nu_{\Gamma_{G^{\mathrm{C}}}}=0 & \text { on } \Gamma_{G^{\mathrm{C}}}
\end{array}\right.
$$

Finally,

$$
v_{s}(t, x, \omega)=\sum_{i=1}^{m} w_{i}(t, x, \omega) \frac{\partial u_{s}}{\partial x_{i}}(t, x) \quad(1 \leq s \leq M)
$$

Proof. In view of Theorems 4.2 and 4.3, the sequences $\widetilde{\left(u_{s}^{\varepsilon}\right)_{\varepsilon>0}}, \widetilde{\left(\nabla_{x} u_{s}^{\varepsilon}\right)_{\varepsilon>0}}$ and $\widetilde{\left(\frac{\partial u_{s}^{\varepsilon}}{\partial t}\right)_{\varepsilon>0}}$ $(1 \leq s \leq M)$ are bounded in $L^{2}([0, T] \times \boldsymbol{Q})$. Using Lemma 3.13, they two-scale converge, up to a subsequence, respectively, to: $\left[\chi_{G^{\mathrm{c}}} u_{s}(t, x)\right],\left[\chi_{G^{\mathrm{C}}}\left(\nabla_{x} u_{s}(t, x)+v_{s}(t, x, \omega)\right)\right]$, $\left[\chi_{G^{\mathrm{r}}} \partial_{t} u_{s}(t, x)\right]$, where $u_{s} \in L^{2}\left(0, T ; H^{1}(\boldsymbol{Q})\right)$ and $v_{s} \in L^{2}\left([0, T] \times \boldsymbol{Q} ; L_{\text {pot }}^{2}(\Omega)\right)$. As test functions for homogenization, let us take

$$
\begin{equation*}
\phi^{\varepsilon}(t, x, \omega):=\phi_{0}(t, x)+\varepsilon \phi(t, x) \psi\left(\tau_{\frac{x}{\varepsilon}} \omega\right) \tag{56}
\end{equation*}
$$

where $\phi_{0}, \phi \in C^{1}([0, T] \times \overline{\boldsymbol{Q}})$ and $\psi \in \Psi$, with $\Psi$ being the set of Remark 2.12,
In the case when $s=1$, let us multiply the first equation of (42) by the test function $\phi^{\varepsilon}$. Integrating, the divergence theorem yields

$$
\begin{align*}
& \int_{0}^{T} \int_{\boldsymbol{Q}^{\varepsilon}(\omega)} \frac{\partial u_{1}^{\varepsilon}}{\partial t} \phi^{\varepsilon}(t, x, \omega) d x d t+\int_{0}^{T} \int_{\boldsymbol{Q}^{\varepsilon}(\omega)}\left\langle D_{1}\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \nabla_{x} u_{1}^{\varepsilon}, \nabla \phi^{\varepsilon}\right\rangle d x d t \\
& +\int_{0}^{T} \int_{\boldsymbol{Q}^{\varepsilon}(\omega)} u_{1}^{\varepsilon} \sum_{j=1}^{M} a_{1, j} u_{j}^{\varepsilon} \phi^{\varepsilon}(t, x, \omega) d x d t=\varepsilon \int_{0}^{T} \int_{\Gamma_{\boldsymbol{Q}}^{\varepsilon}(\omega)} \eta\left(t, x, \tau_{\frac{x}{\varepsilon}} \omega\right) \phi^{\varepsilon}(t, x, \omega) d \mathcal{H}^{m-1} d t \tag{57}
\end{align*}
$$

Passing to the two-scale limit, as $\varepsilon \rightarrow 0$, we get, taking into account (17):

$$
\begin{gather*}
\int_{0}^{T} \int_{Q} \int_{\Omega} \chi_{G^{\mathrm{c}}} \frac{\partial u_{1}}{\partial t}(t, x) \phi_{0}(t, x) d \mathbb{P}(\omega) d x d t \\
+\int_{0}^{T} \int_{Q} \int_{\Omega} \chi_{G^{\mathrm{c}}} D_{1}(t, x, \omega)\left[\nabla_{x} u_{1}(t, x)+v_{1}(t, x, \omega)\right] \\
\cdot\left[\nabla_{x} \phi_{0}(t, x)+\phi(t, x) \nabla_{\omega} \psi(\omega)\right] d \mathbb{P}(\omega) d x d t \\
+\int_{0}^{T} \int_{Q} \int_{\Omega} \chi_{G^{\mathrm{c}}} u_{1}(t, x) \sum_{j=1}^{M} a_{1, j} u_{j}(t, x) \phi_{0}(t, x) d \mathbb{P}(\omega) d x d t \\
=\int_{0}^{T} \int_{Q} \int_{\Omega} \chi_{\Gamma_{G^{\mathrm{c}}}} \eta(t, x, \omega) \phi_{0}(t, x) d \mu_{\Gamma, \mathcal{P}}(\omega) d x d t . \tag{58}
\end{gather*}
$$

The term on the right-hand side follows from Eq. (38). The last term on the left-hand side of (58) has been obtained by observing that $\mathcal{E}_{\varepsilon} u_{j}^{\varepsilon} \rightarrow u_{j}$ strongly in $L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right.$ ) (see Lemma 3.13) and that the two-scale convergence of $u_{1}^{\varepsilon} \stackrel{2 s}{ } \chi_{G^{\mathrm{c}}} u_{1}$ implies weak convergence of $u_{1}^{\varepsilon} \phi^{\varepsilon}(\cdot, \cdot, \omega) \rightharpoonup u_{1} \phi_{0} \int_{\Omega} \chi_{G^{\mathrm{c}}} d \mathbb{P}(\omega)$ in $L^{2}\left(0, T ; L^{2}(\boldsymbol{Q})\right)$.

An integration by parts shows that (58) can be put in the strong form associated with the following homogenized system:

$$
\begin{array}{ll}
\quad-\operatorname{div}_{\omega}\left[D_{1}(t, x, \omega)\left(\nabla_{x} u_{1}(t, x)+v_{1}(t, x, \omega)\right)\right]=0 & \text { in }[0, T] \times \boldsymbol{Q} \times G^{\complement} \\
{\left[D_{1}(t, x, \omega)\left(\nabla_{x} u_{1}(t, x)+v_{1}(t, x, \omega)\right)\right] \cdot \nu_{\Gamma_{G^{\mathrm{C}}}}=0} & \text { on }[0, T] \times \boldsymbol{Q} \times \Gamma_{G^{\mathrm{C}}} \\
\quad+\theta u_{1}(t, x) \sum_{j=1}^{M} a_{1, j} u_{j}(t, x)-\int_{\Omega} \chi_{\Gamma_{G^{\mathrm{C}}}} \eta(t, x, \omega) d \mu_{\Gamma, \mathcal{P}}(\omega)=0 & \text { in }[0, T] \times \boldsymbol{Q} \\
\theta \frac{\partial u_{1}}{\partial t}(t, x)-\operatorname{div}_{x}\left[\int_{\Omega} \chi_{G^{\mathrm{C}}} D_{1}(t, x, \omega)\left(\nabla_{x} u_{1}(t, x)+v_{1}(t, x, \omega)\right) d \mathbb{P}(\omega)\right] \\
{\left[\int_{\Omega} \chi_{G^{\mathrm{C}}} D_{1}(t, x, \omega)\left(\nabla_{x} u_{1}(t, x)+v_{1}(t, x, \omega)\right) d \mathbb{P}(\omega)\right] \cdot n=0} & \text { on }[0, T] \times \partial \boldsymbol{Q} \tag{62}
\end{array}
$$

where

$$
\begin{equation*}
\theta=\int_{\Omega} \chi_{G^{\complement}} d \mathbb{P}(\omega)=\mathbb{P}\left(G^{\complement}\right) \tag{63}
\end{equation*}
$$

represents the fraction of volume occupied by $G^{\complement}$. To conclude, by continuity, we have that

$$
u_{1}(0, x)=U_{1} \quad \text { in } \boldsymbol{Q}
$$

The function $v_{1}(t, x, \omega)$, satisfying (59) and (60), can be expressed as follows

$$
\begin{equation*}
v_{1}(t, x, \omega):=\sum_{i=1}^{m} w_{i}(t, x, \omega) \frac{\partial u_{1}}{\partial x_{i}}(t, x) \tag{64}
\end{equation*}
$$

where $\left(w_{i}\right)_{1 \leq i \leq m} \in L^{2}\left([0, T] \times \boldsymbol{Q} ; L_{\text {pot }}^{2}\left(G^{\complement}\right)\right)$ is the family of solutions of the microscopic problem

$$
\begin{cases}-\operatorname{div}_{\omega}\left[D_{1}(t, x, \omega)\left(w_{i}(t, x, \omega)+\hat{e}_{i}\right)\right]=0 & \text { in } G^{\complement}  \tag{65}\\ D_{1}(t, x, \omega)\left[w_{i}(t, x, \omega)+\hat{e}_{i}\right] \cdot \nu_{\Gamma_{G^{\mathrm{C}}}}=0 & \text { on } \Gamma_{G^{\mathrm{C}}}\end{cases}
$$

and $\hat{e}_{i}$ is the $i$-th unit vector of the canonical basis of $\mathbb{R}^{m}$. The system (65) represents the stochastic version of the "cell problem" defined in periodic homogenization. By using the relation (64) in Eqs.(61) and (62), we get

$$
\begin{align*}
& \theta \frac{\partial u_{1}}{\partial t}(t, x)-\operatorname{div}_{x}\left[D_{1}^{\star}(t, x) \nabla_{x} u_{1}(t, x)\right]+\theta u_{1}(t, x) \sum_{j=1}^{M} a_{1, j} u_{j}(t, x)  \tag{66}\\
& -\int_{\Omega} \chi_{\Gamma_{G^{\mathrm{C}}}} \eta(t, x, \omega) d \mu_{\Gamma, \mathcal{P}}(\omega)=0 \quad \text { in }[0, T] \times \boldsymbol{Q} \\
& {\left[D_{1}^{\star} \nabla_{x} u_{1}(t, x)\right] \cdot n=0 \quad \text { on }[0, T] \times \partial \boldsymbol{Q}} \tag{67}
\end{align*}
$$

where the entries of the matrix $D_{1}^{\star}$ (called "effective diffusivity") are given by

$$
\begin{equation*}
\left(D_{1}^{\star}\right)_{i j}(t, x)=\int_{\Omega} \chi_{G^{\mathrm{c}}} D_{1}(t, x, \omega)\left[w_{i}(t, x, \omega)+\hat{e}_{i}\right] \cdot\left[w_{j}(t, x, \omega)+\hat{e}_{j}\right] d \mathbb{P}(\omega) \tag{68}
\end{equation*}
$$

The proof for the case $1<s \leq M$ is achieved by applying exactly the same arguments.

## A Appendix A

We review some basic results on the realization of random domains based on continuum percolation theory [22].

## A. 1 Stationary ergodic point processes.

Since in percolation theory, random modeling is based on the occurrences of stationary point processes, in this section, we state their definition and some basic properties [8].

Definition A.1. Denote the $\sigma$-algebra of Borel sets in $\mathbb{R}^{m}$ by $\mathcal{B}^{m}$.
(i) A Borel measure $\mu$ on $\mathbb{R}^{m}$ is boundedly finite if $\mu(A)<\infty$ for every bounded Borel set $A$.
(ii) Let $N$ be the space of all boundedly finite integer-valued measures on $\mathcal{B}^{m}$, called counting measures for short.

Proposition A.2. A boundedly finite measure $X$ on $\mathcal{B}^{m}$ is a counting measure (i.e., $X \in N$ ) if and only if

$$
\begin{equation*}
X=\sum_{i} k_{i} \delta_{x_{i}} \tag{69}
\end{equation*}
$$

where $k_{i}$ are positive integers and $\left\{x_{i}\right\}$ is a countable set with at most finitely many $x_{i}$ in any bounded Borel set. In Eq. (69) we use Dirac measures defined for every $x_{i} \in \mathbb{R}^{m}$ by

$$
\delta_{x_{i}}(A)= \begin{cases}1 & \text { if } x_{i} \in A  \tag{70}\\ 0 & \text { otherwise }\end{cases}
$$

We equip $N$ with the $\sigma$-algebra $\mathcal{N}$ generated by sets of the form

$$
\{X \in N: X(A)=k\}
$$

where $A \in \mathcal{B}^{m}$ and $k$ is an integer. We finally introduce $N^{*}$ the set of all counting measures such that for all $i \in \mathbb{N}$ it holds $k_{i}=1$ in (69).

Definition A.3. A point process $X$ on state space $\mathbb{R}^{m}$ is a measurable mapping from a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ into $(N, \mathcal{N})$. It is called simple if $X(\omega) \in N^{*}$ a.s.. The distribution of $X$ is the measure $\mu$ on $\mathcal{N}$ induced by $X$, i.e.

$$
\begin{equation*}
\mu(G)=\mathbb{P}\left(X^{-1}(G)\right), \quad \text { for all } G \in \mathcal{N} \tag{71}
\end{equation*}
$$

The notation of Definition A. 3 is intended to imply that with every sample point $\omega \in \Omega$, we associate a particular realization that is a boundedly finite integer-valued Borel measure on $\mathbb{R}^{m}$. We denote it by $X(\cdot, \omega)$ or just $X(\cdot)$ (when we have no need to draw attention to the underlying spaces). A realization of a point process $X$ has the value $X(A, \omega)$ (or just $X(A)$ ) on the Borel set $A \in \mathcal{B}^{m}$. For each fixed $A, X_{A} \equiv X(A, \cdot)$ is a function mapping $\Omega$ into $\mathbb{R}_{+}$, and thus it is a candidate for a nonnegative random variable, as it is shown in the following proposition.

Proposition A.4. Let $X$ be a mapping from a probability space into $N$ and $\mathcal{A}$ a semiring of bounded Borel sets generating $\mathcal{B}^{m}$. Then $X$ is a point process if and only if $X_{A}$ is a random variable for each $A \in \mathcal{A}$.

Taking for $\mathcal{A}$ the semiring of all bounded sets in $\mathcal{B}^{m}$ we obtain the following corollary.

Corollary A.5. $X: \Omega \mapsto N$ is a point process if and only if $X(A)$ is a random variable for each bounded $A \in \mathcal{B}^{m}$.

We now consider invariance properties with respect to translations (or shifts). Let $T_{t}$ be the translation in $\mathbb{R}^{m}$ over the vector $t: T_{t}(s)=t+s$, for all $s \in \mathbb{R}^{m}$. Then $T_{t}$ induces a transformation

$$
S_{t}: N \rightarrow N
$$

through the relation

$$
\left(S_{t} n\right)(A)=n\left(T_{t}^{-1}(A)\right)
$$

for all $A \in \mathcal{B}^{m}$. It is easy to verify that $\left(S_{t}\right)_{t \in \mathbb{R}^{m}}$ form a group.

Definition A.6. The point process $X$ is said to be stationary if

$$
\begin{equation*}
\forall G \in \mathcal{N} \quad \mu\left(S_{t}^{-1}(G)\right)=\mu(G) \tag{72}
\end{equation*}
$$

In other words, a process is stationary if for every $A \subset \mathbb{R}^{m}$, the distribution of $n(A)$ is invariant under shifts $t+A$. This can be interpreted that $n \in N$ has the same probability as all its shifts $S_{t} n$.

Since $\mathbb{P}$ induces a probability measure $\mu$ on $(N, \mathcal{N})$ via (71), it is convenient to replace the space $(\Omega, \mathcal{F}, \mathbb{P})$ by $(N, \mathcal{N}, \mu)$ and to relabel formally $(\Omega, \mathcal{F}, \mathbb{P}):=(N, \mathcal{N}, \mu)$ so that any element $\omega \in \Omega$ represents a counting measure in $\mathbb{R}^{m}$. Identifying $\tau_{x}:=S_{x}$, by (72) we now have a measure-preserving (m. p.) dynamical system $\left(\Omega, \mathcal{F}, \mu, \tau_{x}\right)$.

Definition A.7. A stationary point process $\mu$ is said to be ergodic if $\left\{\tau_{x}: x \in \mathbb{R}^{m}\right\}$ acts ergodically on $(\Omega, \mathcal{F}, \mu)$ in the sense of Definition 2.2.

## A. 2 Percolation theory and random modeling.

The continuum percolation theory provides a general setting for the realization of random domains. In this framework, two common models are the Boolean model and the random-connection model.

## A.2.1 The Boolean model.

The Boolean model is driven by some stationary point process $X$. Each point of $X$ is the centre of a closed ball (in the usual Euclidean metric) with a random radius in such a way that radii corresponding to different points are independent of each other and identically distributed. The radii are also independent of $X$. Additionally, we want the resulting random model to be stationary. In order to assign independent random values to the radii, we partition $\mathbb{R}^{m}$ into binary cubes

$$
K(n, z):=\prod_{i=1}^{m}\left[z_{i} 2^{-n},\left(z_{i}+1\right) 2^{-n}\right]
$$

for all $n \in \mathbb{N}$ and $z \in \mathbb{Z}^{m}$. We call this a binary cube of order $n$. Each point $x \in X$ is contained in a unique binary cube of order $n, K(n, z(n, x))$ and for each point $x \in X$ there is a unique smallest number $n_{0}=n_{0}(x)$ such that $K\left(n_{0}, z\left(n_{0}, x\right)\right)$ contains no other points of $X$ (recall that $X$ is locally finite). We assign to each point $x_{i} \in X$ a random value in $[0, \infty)$ in the following way: For a probability measure $\mathbb{P}_{0}$ on $[0, \infty)$ we define

$$
\Omega_{2}:=\prod_{n \in \mathbb{N}} \prod_{z \in \mathbb{Z}^{m}}[0, \infty)
$$

with the corresponding product $\sigma$-algebra and product measure $\mathbb{P}_{2}:=\mathbb{P}_{0}^{\mathbb{N} \times \mathbb{Z}^{m}}$. Denoting by $\omega_{2} \in \Omega_{2}$ the elements of $\Omega_{2}$ we assign to each cube $K(n, z)$ the value $\omega_{2}(n, z)$ and to every $x \in X$ the radius $r=\omega_{2}\left(n_{0}, z\left(n_{0}, x\right)\right)$.

We now set $\Omega=\Omega_{1} \times \Omega_{2}$ and equip $\Omega$ with product measure $\mathbb{P}=\mathbb{P}_{1} \times \mathbb{P}_{2}$ and the usual product $\sigma$-algebra. A Boolean model is a measurable mapping from $\Omega$ into $N \times \Omega_{2}$.

The product structure of $\Omega$ implies that the radii are independent of the point process, and the product structure of $\Omega_{2}$ implies that different points have balls with independent, identically distributed radii.

Let the unit vectors in $\mathbb{R}^{m}$ be denoted by $e_{1}, \ldots, e_{m}$. The translation $T_{e_{i}}: \mathbb{R}^{m} \rightarrow \mathbb{R}^{m}$ defined by: $x \rightarrow x+e_{i}$ induces a transformation $U_{e_{i}}$ on $\Omega_{2}$ through the equation

$$
\begin{equation*}
\left(U_{e_{i}} \omega_{2}\right)(n, z)=\omega_{2}\left(n, z-2^{n} e_{i}\right) \tag{73}
\end{equation*}
$$

As before, $S_{e_{i}}$ is defined on $\Omega_{1}$ via the equation

$$
\begin{equation*}
\left(S_{e_{i}} \omega_{1}\right)(A)=\omega_{1}\left(T_{e_{i}}^{-1} A\right) \tag{74}
\end{equation*}
$$

Hence, $T_{e_{i}}$ induces a transformation $\tilde{T}_{e_{i}}$ on $\Omega=\Omega_{1} \times \Omega_{2}$ defined by

$$
\begin{equation*}
\tilde{T}_{e_{i}}(\omega)=\left(S_{e_{i}} \omega_{1}, U_{e_{i}} \omega_{2}\right) \tag{75}
\end{equation*}
$$

The transformation $\tilde{T}_{e_{i}}$ corresponds to a translation by the vector $e_{i}$ of a configuration of balls in space. The Boolean model is now stationary in the sense that $\mathbb{P}$ is shift invariant w.r.t. $\left(\tilde{T}_{x}\right)_{x \in \mathbb{Z}^{m}}$. If we replace $\Omega_{2}$ by $\Omega_{2} \times[0,1)^{m}$ as in Sections 2.6 and 3.2 of [14] we can construct a family of mappings $\left(\tau_{x}\right)_{x \in \mathbb{R}^{m}}$ on $\Omega$ such that we have stationarity of $\mathbb{P}$ w.r.t. $\tau_{x}$.

## A.2.2 The random-connection model.

As in Boolean models, a stationary point process $X$ is the first characteristic of the random-connection model (RCM) and it assigns randomly points in the space. The second characteristic of the model is a so-called connection function, which is a non-increasing function from the positive reals into $[0,1]$. Given a connection function $g$, the rule is as follows: for any two points $x_{1}$ and $x_{2}$ of the point process $X$, we insert an edge between $x_{1}$ and $x_{2}$ with probability $g\left(\left|x_{1}-x_{2}\right|\right)$, independently of all other pairs of points of $X$, where $|\cdot|$ denotes the usual Euclidean distance. The formal mathematical construction of a random-connection model is quite similar to the one of a Boolean model. First we assume that the point process $X$ is defined on a probability space $\left(\Omega_{1}, \mathcal{F}_{1}, \mathbb{P}_{1}\right)$. Next we consider a second probability space $\Omega_{2}$ defined as

$$
\Omega_{2}=\prod_{\left\{K(n, z), K\left(m, z^{\prime}\right)\right\}}[0,1]
$$

where the product is over all unordered pairs of binary cubes. An element $\omega_{2} \in \Omega_{2}$ is written as $\omega_{2}\left(\left\{(n, z),\left(m, z^{\prime}\right)\right\}\right)$. We equip $\Omega_{2}$ with product measure $\mathbb{P}_{2}$. As before, we set $\Omega=\Omega_{1} \times \Omega_{2}$ and we equip $\Omega$ with product measure $\mathbb{P}=\mathbb{P}_{1} \times \mathbb{P}_{2}$. A random-connection model is a measurable mapping from $\Omega$ into $N \times \Omega_{2}$ defined by

$$
\left(\omega_{1}, \omega_{2}\right) \rightarrow\left(X\left(\omega_{1}\right), \omega_{2}\right) .
$$

The realisation corresponding to $\left(\omega_{1}, \omega_{2}\right)$ is obtained as follows: for any two points $x$ and $y$ of $X\left(\omega_{1}\right)$, consider the binary cubes $K\left(n_{0}(x), z\left(n_{0}(x), x\right)\right)$ and $K\left(n_{0}(y), z\left(n_{0}(y), y\right)\right)$. We connect $x$ and $y$ if and only if

$$
\omega_{2}\left(\left\{\left(n_{0}(x), z\left(n_{0}(x), x\right)\right),\left(n_{0}(y), z\left(n_{0}(y), y\right)\right)\right\}\right)<g(|x-y|) .
$$

The dynamical system can be constructed similar to the Boolean model.

## A.2.3 The Poisson process.

Usually, both the Boolean and the random-connection models are based on occurrences of the Poisson point process.

Definition A.8. The point process $X$ is said to be a Poisson process with density $\lambda>0$ if (i) and (ii) below are satisfied:
(i) For mutually disjoint Borel sets $A_{1}, \ldots, A_{k}$, the random variables $X\left(A_{1}\right), \ldots, X\left(A_{k}\right)$ are mutually independent.
(ii) For any bounded Borel set $A \in \mathcal{B}^{m}$ we have for every $k \geq 0$

$$
\begin{equation*}
\mathbb{P}(X(A)=k)=e^{-\lambda \mathcal{L}(A)} \frac{\lambda^{k} \mathcal{L}(A)^{k}}{k!} \tag{76}
\end{equation*}
$$

where $\mathcal{L}(\cdot)$ denotes Lebesgue measure in $\mathbb{R}^{m}$.
Eq. (76) represents the probability that the number of points inside a bounded Borel set $A$ equals $k$. Condition (ii) guarantees that a Poisson process is stationary. Furthermore, one can prove [22]:

Proposition A.9. A Poisson point process is ergodic.

The following result shows that ergodicity of a point process carries over to a Boolean model or to a random-connection model driven by that process [22].

Proposition A.10. Suppose $X$ is ergodic. Then, any Boolean model or random-connection model driven by $X$ is also ergodic.

## A. 3 Realization of random perforated structures.

In Section 2.2, we have stated the main assumptions that our perforated domain should satisfy. We have also stressed how a random spherical structure can provide a rather realistic description of neurons in the cerebral tissue. Unfortunately, a Boolean model driven by the Poisson point process allows, in general, the perforations (i.e. the balls) to be generated arbitrarily close to each other so as to form large connected clusters and small angles. In this case Assumption 2.16 no longer holds and our method fails.

One way to construct domains in which the balls are non-intersecting and have a minimal positive distance between them is to combine the Boolean and the randomconnection model as follows. This procedure is known as Matern process (see [8], Example 10.4(d)). Let us consider a random-connection model driven by a Poisson process and applied on a bounded region of $\mathbb{R}^{m}$. Two points are connected with probability 1 if they have distance less than some constant $d_{0}$. All connected points are then deleted from the process. In case of the Poisson process this means that a point is deleted with probability $1-\exp \left(-\lambda d_{0}\right)$, where $\lambda$ is the intensity of the point process. Every remaining point will be assigned as the center of a ball of random radius $\rho(\omega)<\frac{d_{0}}{2}$. For simplicity, in our analysis we will consider balls with the same constant radius $r_{0}<\frac{d_{0}}{2}$. According to this construction, we obtain a domain randomly perforated with balls of the same radius and with minimal distance between them, which satisfies all the assumptions stated in Section 2.2. In particular, let $G(\omega)$ be the union of such random spheres, then our randomly perforated domain can be defined as

$$
\begin{equation*}
Q(\omega)=\mathbb{R}^{m} \backslash G(\omega) \tag{77}
\end{equation*}
$$

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[^0]:    ${ }^{*}$ Dipartimento di Matematica, Università di Bologna, Piazza di Porta S. Donato 5, 40126 Bologna, Italy, (bruno.franchi@unibo.it).
    ${ }^{\dagger}$ Weierstrass Institute for Applied Analysis and Stochastics, Mohrenstraße 39, 10117 Berlin, Germany, (martin.heida@wias-berlin.de)
    ${ }^{\ddagger}$ Dipartimento di Matematica, Politecnico di Milano, Piazza Leonardo da Vinci 32, 20133 Milano, Italy, (silvia.lorenzani@polimi.it).

