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Nils Berglund, Christian Kuehn

## To cite this version:

Nils Berglund, Christian Kuehn. Regularity structures and renormalisation of FitzHugh-Nagumo SPDEs in three space dimensions. 2015. hal-01141380v1

HAL Id: hal-01141380
https://hal.science/hal-01141380v1
Preprint submitted on 12 Apr 2015 (v1), last revised 16 Jun 2015 (v2)

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# Regularity structures and renormalisation of FitzHugh-Nagumo SPDEs in three space dimensions 

Nils Berglund and Christian Kuehn


#### Abstract

We prove local existence of solutions for a class of suitably renormalised coupled SPDE-ODE systems driven by space-time white noise, where the space dimension is equal to 2 or 3 . This class includes in particular the FitzHugh-Nagumo system describing the evolution of action potentials of a large population of neurons, as well as models with multidimensional gating variables. The proof relies on the theory of regularity structures recently developed by M. Hairer, which is extended to include situations with semigroups that are not regularising in space. We also provide explicit expressions for the renormalisation constants, for general cubic nonlinearities.


April 12, 2015.
Preprint version: Comments and suggestions welcome.
2010 Mathematical Subject Classification. 60H15, 35K57 (primary), 81S20, 82C28 (secondary).
Keywords and phrases. Stochastic partial differential equations, parabolic equations, reactiondiffusion equations, FitzHugh-Nagumo equation, regularity structures, renormalisation.

## 1 Introduction

Many spatially extended systems can be described by a reaction-diffusion equation coupled to an ordinary differential equation, of the form

$$
\begin{align*}
\partial_{t} u & =\mathfrak{D} \Delta_{x} u+F(u, v), \\
\partial_{t} v & =G(u, v), \tag{1.1}
\end{align*}
$$

where $u=u(t, x)$ and $v=v(t, x)$ are functions of time $t \geqslant 0$ and space $x \in \mathbb{R}^{d}$ and take values in $\mathbb{R}^{m}$ and $\mathbb{R}^{n}$ respectively, $\Delta_{x}$ denotes the usual Laplacian, $\mathfrak{D} \in \mathbb{R}^{m \times m}$ is a positive definite matrix of diffusion coefficients, and $F: \mathbb{R}^{m} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ and $G: \mathbb{R}^{m} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ are given sufficiently smooth functions. A well-studied situation of this kind arises in neuroscience, where $x \in \mathbb{R}$ measures the position along a neuron's axon, $u(t, x) \in \mathbb{R}$ is the membrane potential at time $t$ and position $x$, and the so-called gating variables $v \in \mathbb{R}^{n}$ describe the number of open ion channels of $n$ different types. Classical examples are the Hodgkin-Huxley model [24] and its simplifications such as the Morris-Lecar model [32] and the FitzHugh-Nagumo model $[15,33]$. However, there are many more interesting situations outside neuroscience that can be described by coupled systems of the form (1.1), for instance in chemical kinetics [40], in population dynamics [8], and in pattern formation [2]. Furthermore, many other classical reaction-diffusion equations from mathematical biology, such as the Keller-Segel [25] or the Gray-Scott model [16] are of the form (1.1) if one regards the model parameters $v$ not as constants via $\partial_{t} v=0$ but as (slowly-)changing parameters via $\partial_{t} v=\tilde{\varepsilon} \widetilde{G}(u, v)$ for some small parameter $0<\tilde{\varepsilon} \ll 1[28]$ and some smooth vector field
$\widetilde{G}$. In fact, the number of applications of systems of the form (1.1) goes far beyond the areas we mentioned here.

However, in some cases, a deterministic model of the form (1.1) is not sufficient to describe the dynamics, and one has to add a stochastic term to capture the effect of thermal fluctuations and other unresolved dynamical processes affecting the system. In the case of neuron dynamics, see for instance the book and recent survey by Tuckwell [38, 37] and the recent reviews by Bressloff $[6,7]$. In particular, stochastic versions of the FitzHughNagumo PDE have attracted considerable recent attention [36, 39, 34] as well as related models of stochastic neural fields [29,27]. In the derivation of these models and in related numerical studies, one frequently uses correlated [31] as well as space-time white noise [1] stochastic forcing.

Before studying properties of solutions to a stochastic partial differential equation (SPDE) with space-time white noise, one has to ensure that such an equation is welldefined, i.e., one has to attach a meaning of "solution" to a given SPDE. Depending on the space dimension $d$, this problem may in fact be extremely hard. Consider the case of a one-component system ( $m=1, n=0, \mathfrak{D} \equiv 1$ ) of the form

$$
\begin{equation*}
\partial_{t} u=\Delta_{x} u+F(u)+\xi, \tag{1.2}
\end{equation*}
$$

where $\xi$ denotes space-time white noise (precise mathematical definitions will be given below). If $d=1$, it is possible to define a notion of mild solution via the Duhamel principle in a quite general setting (see for instance [11]). In higher space dimension, the problem is much more difficult, owing to the fact that space-time white noise is extremely singular.

Consider for instance the case of the Allen-Cahn equation given by (1.2) with $F(u)=$ $a u-u^{3}$ (if $a=0$ this is also known in quantum field theory as the dynamical $\Phi_{d}^{4}$ model, cf [13]). In the case $d=1$, the proof of existence of a unique solution goes back to Faris and Jona-Lasinio [12], and many quantitative properties of this solution are known (see for instance $[9,4,3])$. The two-dimensional case was solved by Da Prato and Debussche using a particular class of Besov spaces [10]. The case $d=3$, however, was only solved very recently by Hairer, using his theory of regularity structures [19]. One of the difficulties is that a renormalisation procedure has to be used to properly define solutions. This is achieved by replacing space-time white noise $\xi$ by a mollified version $\xi^{\varepsilon}$, solving the resulting regularised equation, and passing to the limit of vanishing regularisation (for an alternative approach, see [30]).

The theory of regularity structures provides a framework to study SPDEs with very singular noise, including but not limited to parabolic equations of the form (1.2). The basic idea is to construct an abstract space in which one can define algebraic operations on distributions (in the sense of generalised functions), such as multiplication, composition with a smooth function, and convolution with a kernel. The fixed-point equation obtained by applying Duhamel's principle to the SPDE with mollified noise is then lifted to the abstract space, solved in that space, and finally projected down to a distribution in "physical" space. In addition, the theory allows to incorporate the renormalisation procedure that is needed in most cases; see for instance [18] for an introduction. Furthermore, there are already several recent applications of regularity structures, e.g., to large deviation theory of the Allen-Cahn equation [23], to the dynamical sine-Gordon model [22], to Wong-Zakai approximation of nonlinear parabolic SPDEs [21], to the parabolic Anderson model on bounded [19] as well as unbounded [20] domains, and to the KPZ equation [19, 17].

The purpose of the present work is to extend the theory of regularity structures to multicomponent systems of the form

$$
\begin{align*}
\partial_{t} u & =\Delta_{x} u+F(u, v)+\xi, \\
\partial_{t} v & =G(u, v) . \tag{1.3}
\end{align*}
$$

The main result is the proof of existence, in space dimensions $d=2$ and 3 , of local solutions to the system (1.3) when $F$ is a cubic polynomial and $G$ is linear. This includes the case of the standard FitzHugh-Nagumo model, but also other equations such as the Koper model [26] with diffusion in the fast component, which features vectorial variables $v$. We will mostly focus on the case $d=3$, first because this is the physically relevant case in many applications, but also because it is the technically more challenging case (the Allen-Cahn equation is known not to be renormalisable for $d \geqslant 4$, so the same will hold a fortiori for multicomponent systems with cubic nonlinearities).

The main technical difficulty that has to be overcome to obtain these results comes from the fact that the semigroup associated with the second equation in (1.3) is not at all regularising in space. Though one expects that this loss of regularisation is somehow compensated by the fact that no singular noise term acts on the equation for $v$, the general theory in [19] cannot be applied directly, because it uses in an essential way the assumption that the heat kernel is smooth everywhere except at the origin. Therefore, we have to extend the regularity structure constructed for the Allen-Cahn equation with new abstract symbols representing integration with respect to a singular kernel, and to prove suitable bounds involving these symbols. Furthermore, one has to analyse which role new symbols play in the renormalisation. We find that certain terms involving the singular kernel do not have to be renormalised, while on the other hand, general cubic nonlinearities yield renormalisation terms which do not appear in the Allen-Cahn case.

The remainder of this work is organised as follows. Section 2 contains all main local existence results. Section 3 gives a summary of the most important aspects of the theory of regularity structures contained in [19], illustrated in the case of the Allen-Cahn equation. In Section 4, we present our results allowing to represent the operation of integration with respect to a singular kernel. Section 5 contains the fixed-point argument proving local existence and uniqueness of solutions for the SPDE with mollified noise, and Section 6 deals with the renormalisation procedure which is necessary to pass to the limit of vanishing mollification. Finally, in Section 7 we complete the proofs of the main results.
Notations: We write $|x|$ to denote either the absolute value of $x \in \mathbb{R}$ or the $\ell^{1}$-norm of $x \in \mathbb{R}^{n}$, while $\|x\|$ denotes the Euclidean norm of $x \in \mathbb{R}^{n}$. If $a, b \in \mathbb{R}$ we use $a \wedge b:=$ $\min \{a, b\}$ and $a \vee b:=\max \{a, b\}$. If $f, g$ are two real-valued functions depending on small parameters $\varepsilon, \delta, \ldots$ (which will be clear from the context), we write $f \lesssim g$ to indicate that there exists a constant $C>0$ such that $f \leqslant C g$ holds uniformly in the small parameters. We use the notation $f \asymp g$ to indicate that we have both $f \lesssim g$ and $g \lesssim f$.
Acknowledgements: N.B. would like to thank the Institute for Analysis and Scientific Computing at the Technical University in Vienna for kind hospitality and financial support. CK would like to thank the Austrian Academy of Sciences (ÖAW) for support via an APART fellowship. CK also acknowledges support of the European Commission (EC/REA) via a Marie-Curie International Re-integration Grant.

## 2 Results

### 2.1 The standard FitzHugh-Nagumo equation

We start by considering the particular case of the standard FitzHugh-Nagumo SPDE, given by

$$
\begin{align*}
\partial_{t} u & =\Delta_{x} u+u-u^{3}+\xi, \\
\partial_{t} v & =a_{1} u+a_{2} v . \tag{2.1}
\end{align*}
$$

Here $\Delta_{x}=\sum_{j=1}^{d} \partial_{x_{j}}^{2}$ is the usual Laplacian, $u$ and $v$ are functions of time $t \geqslant 0$ and space $x \in \mathbb{T}^{d}$ (the torus in dimension $d=2$ or $d=3$ ), $\xi$ stands for space-time white noise on $\mathbb{R} \times \mathbb{T}^{d}$, and $a_{1}, a_{2}$ a real parameters. As such, the system (2.1) is not well-posed, and a renormalisation procedure is required to define a notion of solution. To do this, we first choose a rescaled mollifier

$$
\begin{equation*}
\varrho_{\varepsilon}(t, x)=\frac{1}{\varepsilon^{d+2}} \varrho\left(\frac{t}{\varepsilon^{2}}, \frac{x}{\varepsilon}\right), \tag{2.2}
\end{equation*}
$$

where $\varrho: \mathbb{R}^{d+1} \rightarrow \mathbb{R}$ is a smooth compactly supported function of integral 1 . We set $\xi^{\varepsilon}=\varrho_{\varepsilon} * \xi$, where the star stands for space-time convolution, and consider the sequence of equations

$$
\begin{align*}
\partial_{t} u^{\varepsilon} & =\Delta_{x} u^{\varepsilon}+[1+C(\varepsilon)] u^{\varepsilon}-\left(u^{\varepsilon}\right)^{3}+\xi^{\varepsilon}, \\
\partial_{t} v^{\varepsilon} & =a_{1} u^{\varepsilon}+a_{2} v^{\varepsilon}, \tag{2.3}
\end{align*}
$$

where $C(\varepsilon) \in \mathbb{R}$. Then our first result is the following, which is close in spirit to [19, Thm. 1.15].

Theorem 2.1 (Standard FitzHugh-Nagumo SPDE). Assume $u_{0}$ belongs to the Hölder space $\mathcal{C}^{\eta}$ for some $\eta>-\frac{2}{3}$, and $v_{0}$ belongs to the Hölder space $\mathcal{C}^{\gamma}$ for some $\gamma>0$. Then there exists a choice of the constant $C(\varepsilon)$ such that the system (2.3) with initial condition $\left(u_{0}, v_{0}\right)$ admits a sequence of local solutions $\left(u^{\varepsilon}, v^{\varepsilon}\right)$, converging in probability to a limit $(u, v)$ as $\varepsilon \rightarrow 0$. This limit is independent of the choice of mollifier $\varrho$.

The proof of this theorem and the two next ones is given in Section 7. For the precise definition of Hölder spaces $\mathcal{C}^{\eta}$ with negative index, we refer to Definition 3.3 below. By local solution we mean that for any cut-off $L>0$, the solution is defined up to the random time when the $\mathcal{C}^{\eta}$-norm of $\left(u^{\varepsilon}, v^{\varepsilon}\right)$ first reaches $L$. In other words, we cannot exclude the possibility of finite-time blow-up.

The constant $C(\varepsilon)$ admits an explicit expression in terms of the heat kernel ${ }^{1}$

$$
\begin{equation*}
G^{(d)}(t, x)=\frac{1}{|4 \pi t|^{d / 2}} \mathrm{e}^{-|x|^{2} / 4 t} 1_{\{t>0\}} . \tag{2.4}
\end{equation*}
$$

In dimension $d=3, C(\varepsilon)$ is the same as for the dynamical $\Phi_{3}^{4}$ model considered in [19]. Namely, setting $G_{\varepsilon}^{(d)}=G^{(d)} * \varrho_{\varepsilon}$, one has

$$
\begin{equation*}
C(\varepsilon)=3 \int_{\mathbb{R}^{4}} G_{\varepsilon}^{(3)}(z)^{2} \mathrm{~d} z-18 \int_{\mathbb{R}^{4}} G^{(3)}(z) \mathcal{Q}_{0}^{\varepsilon}(z)^{2} \mathrm{~d} z \tag{2.5}
\end{equation*}
$$

[^0]where $\mathcal{Q}_{0}^{\varepsilon}(z)=\int G_{\varepsilon}^{(3)}\left(z_{1}\right) G_{\varepsilon}^{(3)}\left(z_{1}-z\right) \mathrm{d} z_{1}$. In particular, it has been known since the works of Feldman and Osterwalder $[13,14]$ that
\[

$$
\begin{equation*}
C(\varepsilon)=\frac{C_{1}}{\varepsilon}+C_{2} \log \left(\varepsilon^{-1}\right)+C_{3} \tag{2.6}
\end{equation*}
$$

\]

for some constants $C_{1}, C_{2}, C_{3} \in \mathbb{R}$. The constants $C_{1}$ and $C_{3}$ depend on the choice of mollifier $\varrho$, while $C_{2}$ is independent of $\varrho(c f[18$, Rem. 6.2]).

In dimension $d=2$, one can show (cf [23, Rem. 2.14]) that the renormalisation constants are given by

$$
\begin{equation*}
C(\varepsilon)=3 \int_{\mathbb{R}^{3}} G_{\varepsilon}^{(2)}(z)^{2} \mathrm{~d} z=\frac{3}{4 \pi} \log \left(\varepsilon^{-1}\right)+C_{3}, \tag{2.7}
\end{equation*}
$$

for some constant $C_{3} \in \mathbb{R}$, depending again on the choice of $\varrho$.
The proof of Theorem 2.1 will also provide some information on the structure of the solutions. Indeed, we always have

$$
\begin{equation*}
u^{\varepsilon}(t, x)=\chi_{\varepsilon}(t, x)+\varphi^{\varepsilon}(t, x), \tag{2.8}
\end{equation*}
$$

where $\chi_{\varepsilon}=G^{(d)} * \xi^{\varepsilon}$ is the stochastic convolution of heat kernel and noise, and $\varphi^{\varepsilon}$ converges to a function (as opposed to a distribution). In other words, the only singular term in the limit $\varepsilon \rightarrow 0$ is given by $\chi_{\varepsilon}$, which is independent of the nonlinear term in the equation. The function $v^{\varepsilon}(t, x)$ has the same structure, as it can be represented in terms of $u^{\varepsilon}$ by solving a linear inhomogeneous equation.

### 2.2 General cubic nonlinearities

We can now extend the above results to a more general class of systems, of the form

$$
\begin{align*}
\partial_{t} u & =\Delta_{x} u+F(u, v)+\xi, \\
\partial_{t} v & =a_{1} u+a_{2} v, \tag{2.9}
\end{align*}
$$

where $F$ is a general cubic polynomial of the form

$$
\begin{equation*}
F(u, v)=\alpha_{1} u+\alpha_{2} v+\beta_{1} u^{2}+\beta_{2} u v+\beta_{3} v^{2}+\gamma_{1} u^{3}+\gamma_{2} u^{2} v+\gamma_{3} u v^{2}+\gamma_{4} v^{3} . \tag{2.10}
\end{equation*}
$$

In that case, the renormalised equations take the form

$$
\begin{align*}
\partial_{t} u^{\varepsilon} & =\Delta_{x} u^{\varepsilon}+\left[F\left(u^{\varepsilon}, v^{\varepsilon}\right)+C_{0}(\varepsilon)+C_{1}(\varepsilon) u^{\varepsilon}+C_{2}(\varepsilon) v^{\varepsilon}\right]+\xi^{\varepsilon}, \\
\partial_{t} v^{\varepsilon} & =a_{1} u^{\varepsilon}+a_{2} v^{\varepsilon}, \tag{2.11}
\end{align*}
$$

where $C_{0}(\varepsilon), C_{1}(\varepsilon), C_{2}(\varepsilon) \in \mathbb{R}$, and we have the following result:
Theorem 2.2 (General cubic nonlinearities). Assume $u_{0}$ and $v_{0}$ satisfy the same assumptions as in Theorem 2.1. Then there exists a choice of constants $C_{0}(\varepsilon), C_{1}(\varepsilon), C_{2}(\varepsilon)$ such that the system (2.11) with initial condition ( $u_{0}, v_{0}$ ) admits a sequence of local solutions $\left(u^{\varepsilon}, v^{\varepsilon}\right)$, converging in probability to a limit $(u, v)$ as $\varepsilon \rightarrow 0$. This limit is independent of the choice of mollifier $\varrho$.

The renormalisation constants can again be computed explicitly to leading order. They are given by

$$
\begin{align*}
& C_{0}(\varepsilon)=\frac{1}{3} \beta_{1} C(\varepsilon), \\
& C_{1}(\varepsilon)=\gamma_{1} C(\varepsilon),  \tag{2.12}\\
& C_{2}(\varepsilon)=\frac{1}{3} \gamma_{2} C(\varepsilon),
\end{align*}
$$

where $C(\varepsilon)$ is the constant defined in (2.5) for $d=3$ and in (2.7) for $d=2$. Note in particular that these constants depend only on the coefficients of $u^{2}, u^{3}$ and $u^{2} v$ in $F$.

### 2.3 Vectorial gating variables

Another generalisation of interest is to systems with vectorial variables $v$, of the form

$$
\begin{align*}
& \partial_{t} u=\Delta_{x} u+F(u, v)+\xi, \\
& \partial_{t} v=u A_{1}+A_{2} v, \tag{2.13}
\end{align*}
$$

where $v(t, x)$ takes values in $\mathbb{R}^{n}$. Here $A_{1} \in \mathbb{R}^{n}$ is a constant vector, and $A_{2} \in \mathbb{R}^{n \times n}$ is a square matrix, while $F$ is again a cubic polynomial in $u$ and the $v_{i}$. This allows for instance to consider the Koper model [26] with spatial diffusion and space-time white noise in the fast variable; for the stochastic Koper model without the Laplacian representing spatial diffusion see [5]. The reaction terms of the Koper model can be written in the form (2.13) with

$$
F(u, v)=3 u+v_{1}-u^{3}, \quad A_{1}=\binom{k}{0}, \quad A_{2}=\left(\begin{array}{cc}
-2 & 1  \tag{2.14}\\
\epsilon_{2} & \epsilon_{2}
\end{array}\right)
$$

where $k, \epsilon_{2} \in \mathbb{R}$ are model parameters.
The natural candidate for the renormalised system associated to (2.13) is given by

$$
\begin{align*}
& \partial_{t} u^{\varepsilon}=\Delta_{x} u^{\varepsilon}+\left[F\left(u^{\varepsilon}, v^{\varepsilon}\right)+C_{0}(\varepsilon)+C_{1}(\varepsilon) u^{\varepsilon}+\sum_{i=1}^{n} C_{2, i}(\varepsilon) v_{i}^{\varepsilon}\right]+\xi^{\varepsilon} \\
& \partial_{t} v^{\varepsilon}=u^{\varepsilon} A_{1}+A_{2} v^{\varepsilon} \tag{2.15}
\end{align*}
$$

for constants $C_{0}(\varepsilon), C_{1}(\varepsilon), C_{2, i}(\varepsilon) \in \mathbb{R}$. Indeed, we have the following result.
Theorem 2.3 (Vectorial variables $v$ ). Assume $u_{0}$ and the components of $v_{0}$ satisfy the same assumptions as $u_{0}$ and $v_{0}$ in Theorem 2.1. Then there exists a choice of constants $C_{0}(\varepsilon), C_{1}(\varepsilon), C_{2, i}(\varepsilon)$ such that the system (2.15) with initial condition $\left(u_{0}, v_{0}\right)$ admits a sequence of local solutions $\left(u^{\varepsilon}, v^{\varepsilon}\right)$, converging in probability to a limit $(u, v)$ as $\varepsilon \rightarrow 0$. This limit is independent of the choice of mollifier $\varrho$.

The renormalisation constants can again be expressed in terms of the coefficients of the initial equation (2.13). Writing

$$
\begin{equation*}
F(u, v)=\alpha_{1} u+\beta_{1} u^{2}+\gamma_{1} u^{3}+\sum_{i=1}^{n}\left[\alpha_{2, i} v_{i}+\beta_{2, i} u v_{i}+\gamma_{2, i} u^{2} v_{i}\right]+R(u, v) \tag{2.16}
\end{equation*}
$$

with $|R(u, v)| \leqslant C\|v\|^{2}$, we have

$$
\begin{align*}
C_{0}(\varepsilon) & =\frac{1}{3} \beta_{1} C(\varepsilon) \\
C_{1}(\varepsilon) & =\gamma_{1} C(\varepsilon),  \tag{2.17}\\
C_{2, i}(\varepsilon) & =\frac{1}{3} \gamma_{2, i} C(\varepsilon), \quad i=1, \ldots, n
\end{align*}
$$

where $C(\varepsilon)$ is again the constant defined in (2.5) for $d=3$ and in (2.7) for $d=2$. Note that in the particular case of the Koper model, $C_{1}(\varepsilon)$ is equal to the constant $C(\varepsilon)$ obtained for the FitzHugh-Nagumo equation, while all other constants vanish. This is not surprising, since both models have the same nonlinearity.

Remark 2.4. The notion of local subcriticality given in [19, Assumption 8.3] suggests that the class of renormalisable SPDE-ODE models of the form (1.3) is larger than the one considered here. More precisely, in dimension $d=3$, one would expect models with quartic $F$ and linear $G$ to be renormalisable, as well as models with $F$ quartic in $u$ and quadratic in $v$ and $G$ quadratic in $u$ and linear in $v$. The reason we did not include them in our analysis is a technical one: quartic $F$ would produce solutions whose singular component is not homogeneous in space and time, as in (2.8). The way in which we lift the singular kernel to the regularity structure is not able to deal with such situations (see Section 4.2 below). We plan to further investigate this issue in future work. $\diamond$

## 3 Regularity structure for the Allen-Cahn equation

This section serves the double purpose of giving a very brief account of the theory of regularity structures contained in [19], and of describing a regularity structure for the Allen-Cahn equation

$$
\begin{equation*}
\partial_{t} u=\Delta_{x} u+u-u^{3}+\xi \tag{3.1}
\end{equation*}
$$

where $\Delta_{x}=\sum_{j=1}^{d} \partial_{x_{j}}^{2}$ is the usual Laplacian, $\xi=\xi(t, x)$ denotes space-time white noise on $\mathbb{R} \times \mathbb{T}^{d}$ where $\mathbb{T}^{d}$ denotes the $d$-dimensional torus with $d=2$ or $d=3$, and we seek a solution $u:[0, T] \times \mathbb{T}^{d} \rightarrow \mathbb{R}$ for a given initial condition $u_{0}=u(0, x)$ in a suitable function space. The regularity structure for (3.1) will serve as our starting point to build a larger structure allowing to represent the FitzHugh-Nagumo equation.

By Duhamel's principle, one possible solution concept for (3.1) is to consider the following integral equation

$$
\begin{equation*}
u_{t}=\int_{0}^{t} S(t-s)\left[u_{s}-u_{s}^{3}+\xi_{s}\right] \mathrm{d} s+S(t) u_{0} \tag{3.2}
\end{equation*}
$$

where $S(t)=\mathrm{e}^{t \Delta_{x}}$ denotes the semigroup of the heat equation compatible with the boundary conditions and $u_{s}:=u(s, \cdot), \xi_{s}:=(s, \cdot)$. The purpose of a regularity structure is to provide an abstract space in which one can construct a fixed point of (3.2) when $\xi$ is replaced by a mollified version $\xi^{\varepsilon}$. Then the idea is to take the limit $\varepsilon \rightarrow 0$, and to project the solution to a distribution on $[0, T] \times \mathbb{T}^{d}$, where "distribution" is understood in the sense of a "generalized function" representing a sample path.

### 3.1 Regularity structures

Definition 3.1 ([19, Def. 2.1]). A regularity structure is a triple $(A, T, \mathcal{G})$ consisting of (R1) an index set $A \subset \mathbb{R}$, containing 0 , which is bounded from below and locally finite; (R2) a model space $T$, which is a graded vector space $T=\bigoplus_{\alpha \in A} T_{\alpha}$, where each $T_{\alpha}$ is a Banach space; the space $T_{0}$ is isomorphic to $\mathbb{R}$ and its unit is denoted 1;
(R3) a structure group $\mathcal{G}$ of linear operators acting on $T$, such that

$$
\begin{equation*}
\Gamma \tau-\tau \in \bigoplus_{\beta<\alpha} T_{\beta}=: T_{\alpha}^{-} \tag{3.3}
\end{equation*}
$$

holds for every $\Gamma \in \mathcal{G}$, every $\alpha \in A$ and every $\tau \in T_{\alpha}$; furthermore, $\Gamma \mathbf{1}=\mathbf{1}$ for every $\Gamma \in \mathcal{G}$.

Example 3.2. A simple but important example of regularity structure is the polynomial regularity structure for $d$ variables. In that case, $A=\mathbb{N}_{0}$ is the set of non-negative integers. For $\ell \in \mathbb{N}_{0}, T_{\ell}$ is the space of homogeneous polynomials in $d$ variables of degree $\ell$. It is spanned by the monomials $X^{k}=X_{1}^{k_{1}} \ldots X_{d}^{k_{d}}$ for which $|k|=k_{1}+\cdots+k_{d}=\ell$. Finally, the structure group $\mathcal{G}$ is defined by

$$
\begin{equation*}
\Gamma_{h}\left(X^{k}\right)=(X-h)^{k}, \quad h \in \mathbb{R}^{d} . \tag{3.4}
\end{equation*}
$$

This group is isomorphic to $\mathbb{R}^{d}$, and one easily sees that $\Gamma_{h}$ satisfies the requirement (3.3). The interpretation of $\mathcal{G}$ is that it allows to convert a Taylor expansion around a point $x \in \mathbb{R}^{d}$ into the expansion around another point $x+h$.

We will henceforth denote by $\bar{T}$ the polynomial regularity structure with $d+1$ variables $X_{0}, \ldots, X_{d}$, where $X_{0}$ represents the time variable. Since the linear part of the Allen-Cahn equation is given by a parabolic operator, it turns out to be useful to make the time variable "count double". This is done by defining the parabolic scaling

$$
\begin{equation*}
\mathfrak{s}=(2,1, \ldots, 1) \in \mathbb{N}^{d+1}, \tag{3.5}
\end{equation*}
$$

acting on $\mathbb{R}^{d+1}$ and declaring that a monomial $X^{k}$ has homogeneity $\left|X^{k}\right|=|k|_{\mathfrak{s}}$ where the scaled degree is defined by $|k|_{\mathfrak{s}}=2 k_{0}+\sum_{i=1}^{d} k_{i}$.

The regularity structure of the Allen-Cahn equation (3.3) in $\mathbb{R}^{d}$ is built by enlarging $\bar{T}$, i.e., by adding new symbols other than the polynomial symbols $X^{k}$ to $\bar{T}$ (see [23]). The noise is represented by a symbol $\Xi$. In order to account for the fact that space-time white noise has Hölder regularity $\alpha$ for any $\alpha<-(d+2) / 2$, we set

$$
\begin{equation*}
\alpha_{0}=-\frac{d+2}{2}-\kappa \tag{3.6}
\end{equation*}
$$

where $\kappa>0$ will be chosen sufficiently small in the sequel, and declare that $\Xi$ has homogeneity $|\Xi|=\alpha_{0}$.

The set of symbols is equipped by a product, which by definition is commutative and associative with unit 1. The product of two elements $\tau, \sigma \in T$ has homogeneity $|\tau \sigma|:=|\tau|+|\sigma|$. Furthermore, integration against the heat kernel is represented by a map $\mathcal{I}: T \rightarrow T$, which by definition satisfies $|\mathcal{I}(\tau)|:=|\tau|+2$, in order to account for the regularizing effect of the heat kernel.

Let $\mathcal{F}$ be the set containing all possible products of symbols $\mathbf{1}, \Xi, X_{i}$ and their images by $\mathcal{I}$, and denote by $\mathcal{H}=\operatorname{span}(\mathcal{F})$ the vector space spanned by all these symbols. This is

| $\tau$ | Symb | $\|\tau\|$ | $d=3$ | $d=2$ | $\Delta(\tau)$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $\Xi$ | $\Xi$ | $\alpha_{0}$ | $-\frac{5}{2}-\kappa$ | $-2-\kappa$ | $\Xi \otimes 1$ |
| $\mathcal{I}(\Xi)^{3}$ | V | $3 \alpha_{0}+6$ | $-\frac{3}{2}-3 \kappa$ | $0-3 \kappa$ | $* \otimes 1$ |
| $\mathcal{I}(\Xi)^{2}$ | $v$ | $2 \alpha_{0}+4$ | $-1-2 \kappa$ | $0-2 \kappa$ | $v \otimes 1$ |
| $\mathcal{I}\left(\mathcal{I}(\Xi)^{3}\right) \mathcal{I}(\Xi)^{2}$ | * | $5 \alpha_{0}+12$ | $-\frac{1}{2}-5 \kappa$ | $2-5 \kappa$ | $\stackrel{\sim}{\bullet} \cdot \otimes 1+\mathscr{V} \otimes \mathcal{J}(\stackrel{\bullet}{*})$ |
| $\mathcal{I}(\Xi)$ | $\dagger$ | $\alpha_{0}+2$ | $-\frac{1}{2}-\kappa$ | $0-\kappa$ | $\bullet \otimes 1$ |
| $\mathcal{I}\left(\mathcal{I}(\Xi)^{3}\right) \mathcal{I}(\Xi)$ | $\stackrel{\circ}{\circ}$ | $4 \alpha_{0}+10$ | $0-4 \kappa$ | $2-4 \kappa$ | $\mathscr{\bullet} \otimes 1+\boldsymbol{i} \otimes \mathcal{J}(\stackrel{\bullet}{*})$ |
| $\mathcal{I}\left(\mathcal{I}(\Xi)^{2}\right) \mathcal{I}(\Xi)^{2}$ | 2 | $4 \alpha_{0}+10$ | $0-4 \kappa$ | $2-4 \kappa$ | $\mathfrak{Q} \otimes 1+\boldsymbol{V} \otimes \mathcal{J}(\mathcal{V})$ |
| $\mathcal{I}(\Xi)^{2} X_{i}$ | $\cdots X_{i}$ | $2 \alpha_{0}+5$ | $0-2 \kappa$ | $1-2 \kappa$ | $\mathcal{V} \otimes X_{i}+\mathcal{V} X_{i} \otimes \mathbf{1}$ |
| 1 | 1 | 0 | 0 | 0 | $1 \otimes 1$ |
| $\mathcal{I}\left(\mathcal{I}(\Xi)^{3}\right)$ | $\stackrel{*}{*}$ | $3 \alpha_{0}+8$ | $\frac{1}{2}-3 \kappa$ | $2-3 \kappa$ | $\stackrel{\oplus}{\bullet} \otimes 1+1 \otimes \mathcal{J}(\stackrel{\bullet}{*})$ |
| $\mathcal{I}\left(\mathcal{I}(\Xi)^{2}\right) \mathcal{I}(\Xi)$ | $\%$ | $3 \alpha_{0}+8$ | $\frac{1}{2}-3 \kappa$ | $2-3 \kappa$ | $\mathfrak{V} \otimes 1+\boldsymbol{i} \otimes \mathcal{J}(\mathscr{V})$ |
| $\mathcal{I}(\mathcal{I}(\Xi)) \mathcal{I}(\Xi)^{2}$ | i. | $3 \alpha_{0}+8$ | $\frac{1}{2}-3 \kappa$ | $2-3 \kappa$ | $\begin{aligned} & \dot{\mathcal{V}} \otimes \mathbf{1}+\boldsymbol{V} \otimes \mathcal{J}(\boldsymbol{\imath})+ \\ & \boldsymbol{V} X_{i} \otimes \mathcal{J}_{i}(\boldsymbol{\varphi})+\boldsymbol{V} \otimes X_{i} \mathcal{J}_{i}(\boldsymbol{\imath}) \end{aligned}$ |
| $\mathcal{I}\left(\mathcal{I}(\Xi)^{2}\right)$ | 9 | $2 \alpha_{0}+6$ | $1-2 \kappa$ | $2-2 \kappa$ | $\widehat{Y} \otimes 1+1 \otimes \mathcal{J}(\mathcal{V})$ |
| $\mathcal{I}(\mathcal{I}(\Xi)) \mathcal{I}(\Xi)$ | $i$. | $2 \alpha_{0}+6$ | $1-2 \kappa$ | $2-2 \kappa$ | $\begin{aligned} & \mathfrak{\bullet} \otimes \mathbf{1}+\boldsymbol{\imath} \otimes \mathcal{J}(\mathfrak{\imath})+ \\ & \boldsymbol{\imath} X_{i} \otimes \mathcal{J}_{i}(\boldsymbol{\mathfrak { i }})+\boldsymbol{i} \otimes X_{i} \mathcal{J}_{i}(\boldsymbol{\imath}) \end{aligned}$ |
| $X_{i}$ | ${ }^{\text {a }}$ |  |  |  | $X_{i} \otimes \mathbf{1}+\mathbf{1} \otimes X_{i}$ |
| $\mathcal{I}(\mathcal{I}(\Xi))$ | 1 | $\alpha_{0}+4$ | $\frac{3}{2}-\kappa$ | $2-\kappa$ | $\begin{aligned} & \mathfrak{l} \otimes \mathbf{1}+\mathbf{1} \otimes \mathcal{J}(\boldsymbol{\bullet})+ \\ & X_{i} \otimes \mathcal{J}_{i}(\boldsymbol{\imath})+\mathbf{1} \otimes X_{i} \mathcal{J}_{i}(\boldsymbol{\uparrow}) \end{aligned}$ |

Table 1. Elements of $\mathcal{F}_{F}$ of lowest homogeneity for the Allen-Cahn equation. They have been ordered in increasing order of homogeneity for the case $d=3$. The expressions for $\Delta(\tau)$ are shown for the case $d=3$. We have written $\mathcal{J}$ instead of $\mathcal{J}_{0}$ and $\mathcal{J}_{i}$ instead of $\mathcal{J}_{e}$, where the $e_{i}$ are canonical basis vectors of $\mathbb{Z}_{+}^{d+1}$. Summation over the index $i$ is understood.
an infinite-dimensional space but in practice only a finite-dimensional subspace of $\mathcal{H}$ will be needed. Let $\mathcal{U}$ be the smallest set containing $\mathbf{1}, X_{i}$ and $\mathcal{I}(\Xi)$, and such that

$$
\begin{equation*}
\tau_{1}, \tau_{2}, \tau_{3} \in \mathcal{U} \quad \Rightarrow \quad \mathcal{I}\left(\tau_{1} \tau_{2} \tau_{3}\right) \in \mathcal{U} \tag{3.7}
\end{equation*}
$$

We then set

$$
\begin{equation*}
\mathcal{F}_{F}=\{\Xi\} \cup\left\{\tau_{1} \tau_{2} \tau_{3}: \tau_{i} \in \mathcal{U}\right\}, \tag{3.8}
\end{equation*}
$$

and define the model space as being the vector space $T=\operatorname{span}\left(\mathcal{F}_{F}\right)$ spanned by $\mathcal{F}_{F}$. The index set is then defined as

$$
\begin{equation*}
A=\left\{|\tau|: \tau \in \mathcal{F}_{F}\right\} . \tag{3.9}
\end{equation*}
$$

The model space $T$ admits a natural grading

$$
\begin{equation*}
T=\bigoplus_{\gamma \in A} T_{\gamma} \tag{3.10}
\end{equation*}
$$

obtained by letting $T_{\gamma}$ be the vector space spanned by elements of homogeneity $\gamma$. We equip each $T_{\gamma}$ with a norm $\|\cdot\|_{\gamma}$; the choice of norm is irrelevant since [19, Lemma 8.10] shows the (nontrivial) fact that all $T_{\gamma}$ are finite-dimensional.

Table 1 shows the elements in $\mathcal{F}_{F}$ of lowest homogeneity, using a graphical representation introduced by Hairer. Each symbol $\Xi$ is denoted by a dot, and integration with respect to $\mathcal{I}$ is denoted by a vertical line pointing downwards. For instance, $\mathcal{I}(\Xi)=\uparrow$. Multiplication of symbols is represented by joining them at the base, so for instance $\mathcal{I}(\Xi)^{2}=\boldsymbol{V}$.

Finally, we have to extend the definition of the structure group $\mathcal{G}$, which is the most involved part of the construction. The first step is to introduce a set $\mathcal{F}_{+} \subset \mathcal{F}$, which contains all formal expressions of the type

$$
\begin{equation*}
X^{k} \prod_{j} \mathcal{J}_{k_{j}} \tau_{j} \tag{3.11}
\end{equation*}
$$

where $\tau_{j} \in \mathcal{F}$ and the multiindices $k_{j}$ are such that $\left|\tau_{j}\right|+2-\left|k_{j}\right|>0$. By definition, an expression of the form (3.11) has homogeneity $|k|+\sum_{j}\left(\left|\tau_{j}\right|+2-\left|k_{j}\right|\right)$, which is always strictly positive, except for the element 1 , which has homogeneity 0 . The set $\mathcal{F}_{F}^{+} \subset \mathcal{F}_{+}$is defined similarly, but with $\tau_{j} \in \mathcal{F}_{F}$. We set $\mathcal{H}_{+}=\operatorname{span}\left(\mathcal{F}_{+}\right)$and $T^{+}=\operatorname{span}\left(\mathcal{F}_{F}^{+}\right)$.

The last ingredient is given by a bilinear map $\Delta: \mathcal{H} \rightarrow \mathcal{H} \otimes \mathcal{H}_{+}$defined by setting

$$
\begin{equation*}
\Delta \mathbf{1}=\mathbf{1} \otimes \mathbf{1}, \quad \Delta X_{i}=X_{i} \otimes \mathbf{1}+\mathbf{1} \otimes X_{i}, \quad \Delta \Xi=\Xi \otimes \mathbf{1} \tag{3.12}
\end{equation*}
$$

and extended inductively to all of $\mathcal{H}$ by the rules

$$
\begin{align*}
& \Delta(\tau \sigma)=(\Delta \tau)(\Delta \sigma) \\
& \Delta(\mathcal{I} \tau)=(\mathcal{I} \otimes \operatorname{Id}) \Delta(\tau)+\sum_{\ell, m} \frac{X^{\ell}}{\ell!} \otimes \frac{X^{m}}{m!} \mathcal{J}_{\ell+m} \tau \tag{3.13}
\end{align*}
$$

where Id denotes the identity map and we impose that

$$
\begin{equation*}
\mathcal{J}_{k} \tau=0 \quad \text { if }|k|_{\mathfrak{s}} \geq|\tau|_{\mathfrak{s}}+2, \tag{3.14}
\end{equation*}
$$

while $\mathcal{J}_{k} \tau$ is a new formal symbol otherwise. Table 1 also shows $\Delta(\tau)$ for the first few elements of $\mathcal{F}_{F}$.

Denote by $\mathcal{H}_{+}^{\star}$ the dual of $\mathcal{H}_{+}$, that is, the set of linear maps $h: \mathcal{H}_{+} \rightarrow \mathbb{R}$, that we will write as $\tau \mapsto\langle h, \tau\rangle$. Let $\mathcal{G}$ be the set of grouplike elements $g \in \mathcal{H}_{+}^{\star}$, that is, those satisfying $\langle g, \tau \sigma\rangle=\langle g, \tau\rangle\langle g, \sigma\rangle$ for all $\tau, \sigma \in \mathcal{H}_{+}$. Then $\mathcal{G}_{F}$ is obtained by identifying elements of $\mathcal{G}$ acting the same way on $T^{+}$. The action of $\mathcal{G}_{F}$ on $T$ is defined by

$$
\begin{equation*}
g \mapsto \Gamma_{g} \tau=(\operatorname{Id} \otimes g) \Delta \tau \tag{3.15}
\end{equation*}
$$

The fact that $\left(A, T, \mathcal{G}_{F}\right)$ constructed in this way is indeed a regularity structure is proved in [19, Thm. 8.24].

### 3.2 Models and reconstruction theorem

A regularity structure can be used to represent an SPDE as a fixed-point problem in an abstract space. One also has to provide a connection between the abstract space and the "physical" space the solution lives in. First, we need to introduce a version of the classical Hölder spaces $\mathcal{C}^{\alpha}$ on $\mathbb{R}^{d+1}$. In order to take the parabolic scaling into account, the Euclidean norm is replaced by

$$
\begin{equation*}
\|(t, x)\|_{\mathfrak{s}}=|t|^{1 / 2}+\sum_{i=1}^{d}\left|x_{i}\right| . \tag{3.16}
\end{equation*}
$$

Denote by $\langle\mathbf{1}, \cdot\rangle$ the element of the dual of $T$ defined by $\langle\mathbf{1}, \mathbf{1}\rangle=1$ and $\langle\mathbf{1}, \tau\rangle=0$ for all $\tau \in \bigoplus_{\gamma \neq 0} T_{\gamma}$. If $\alpha>0$, the space $\mathcal{C}_{\mathfrak{s}}^{\alpha}$ is defined as the set of all functions $\varphi: \mathbb{R}^{d+1} \rightarrow \mathbb{R}$ such that there exists a function $\hat{\varphi}: \mathbb{R}^{d+1} \rightarrow T_{\alpha}^{-}$such that $\langle\mathbf{1}, \hat{\varphi}(z)\rangle=\varphi(z)$ for every $z=(t, x) \in \mathbb{R}^{d+1}$ and such that for every compact set $\mathfrak{K} \subset \mathbb{R}^{d+1}$ the estimate

$$
\begin{equation*}
\left|\hat{\varphi}(z+h)-\Gamma_{h} \hat{\varphi}(z)\right|_{\beta} \lesssim|h|_{\mathfrak{s}}^{\alpha-\beta} \tag{3.17}
\end{equation*}
$$

holds uniformly over $\beta<\alpha,\|h\|_{\mathfrak{s}} \leq 1$ and $z \in \mathfrak{K}$ (cf [19, Def. 2.14]). Lemma 2.12 in [19] shows that in the case of the Euclidean scaling $\mathfrak{s}=(1, \ldots, 1)$, this definition coincides with the usual definition of Hölder spaces.

If $\alpha<0, \mathcal{C}_{5}^{\alpha}$ can be defined in a natural way as a subset of Schwarz distributions $\mathcal{S}^{\prime}\left(\mathbb{R}^{d+1}\right)$ as follows. Given $r \in \mathbb{N}$, denote by $\mathcal{B}_{\mathfrak{s}, 0}^{r}$ the space of functions of class $\mathcal{C}^{r}$ which are supported in the set $\left\{z \in \mathbb{R}^{d+1}:\|z\|_{\mathfrak{s}} \leqslant 1\right\}$.

Definition 3.3 ([19, Def. 3.7]). Assume $\alpha<0$ and let $r=-\lfloor\alpha\rfloor$. A Schwarz distribution $v \in \mathcal{S}^{\prime}\left(\mathbb{R}^{d+1}\right)$ belongs to $\mathcal{C}_{\mathfrak{s}}^{\alpha}$ if it belongs to the dual of $\mathcal{C}_{0}^{r}$ and for every compact set $\mathfrak{K} \subset \mathbb{R}^{d+1}$, there exists a constant $C$ such that

$$
\begin{equation*}
\left\langle v, \mathcal{S}_{\mathfrak{s}, z}^{\delta} \eta\right\rangle \leqslant C \delta^{\alpha} \tag{3.18}
\end{equation*}
$$

holds for all $\eta \in \mathcal{B}_{\mathfrak{s}, 0}^{r}$ with $\|\eta\|_{\mathcal{C}^{r}} \leqslant 1$, all $\delta \in(0,1]$ and all $z \in \mathfrak{K}$. Here

$$
\begin{equation*}
\left(\mathcal{S}_{\mathfrak{s},(t, x)}^{\delta} \eta\right)(\bar{t}, \bar{x}):=\delta^{-(d+2)} \eta\left(\delta^{-2}(\bar{t}-t), \delta^{-1}(\bar{x}-x)\right) . \tag{3.19}
\end{equation*}
$$

Next we introduce the notion of a model associated with a regularity structure $(A, T, \mathcal{G})$.
Definition 3.4 ([19, Def. 2.17]). A model for a regularity structure $(A, T, \mathcal{G})$ with scaling $\mathfrak{s}$ is a pair $(\Pi, \Gamma)$, defined by a collection $\left\{\Pi_{z}: T \rightarrow \mathcal{S}^{\prime}\left(\mathbb{R}^{d+1}\right)\right\}_{z \in \mathbb{R}^{d+1}}$ of continuous linear maps and a map $\Gamma: \mathbb{R}^{d+1} \times \mathbb{R}^{d+1} \rightarrow \mathcal{G}$ with the following properties.
(M1) $\Gamma_{z z}=\operatorname{Id}$ is the identity of $\mathcal{G}$ and $\Gamma_{z z^{\prime}} \Gamma_{z^{\prime} z^{\prime \prime}}=\Gamma_{z z^{\prime \prime}}$ for all $z, z^{\prime}, z^{\prime \prime} \in \mathbb{R}^{d+1}$.
(M2) $\Pi_{z^{\prime}}=\Pi_{z} \Gamma_{z z^{\prime}}$ for all $z, z^{\prime} \in \mathbb{R}^{d+1}$.
(M3) For any $\gamma \in \mathbb{R}$ and any compact set $\mathfrak{K} \subset \mathbb{R}^{d+1}$, one has

$$
\begin{equation*}
\|\Pi\|_{\gamma ; \mathfrak{\Re}}:=\sup _{z \in \mathfrak{K}} \sup _{\alpha<\gamma} \sup _{\tau \in T_{\alpha}} \sup _{\eta \in \mathcal{B}_{\mathfrak{s}, 0}^{r}} \sup _{0<\delta \leqslant 1} \frac{\left|\left\langle\Pi_{z} \tau, \mathcal{S}_{\mathfrak{s}, z}^{\delta} \eta\right\rangle\right|}{\delta^{\alpha}\|\tau\|}<\infty . \tag{3.20}
\end{equation*}
$$

(M4) For any $\gamma \in \mathbb{R}$ and any compact set $\mathfrak{K} \subset \mathbb{R}^{d+1}$, one has

$$
\begin{equation*}
\|\Gamma\|_{\gamma ; \mathfrak{K}}:=\sup _{z, z^{\prime} \in \mathfrak{K}} \sup _{\alpha<\gamma} \sup _{\beta<\alpha} \sup _{\tau \in T_{\alpha}} \frac{\left\|\Gamma_{z z^{\prime}} \tau\right\|_{\beta}}{\|\tau\|\left\|z-z^{\prime}\right\|_{\mathfrak{s}}^{\alpha-\beta}}<\infty . \tag{3.21}
\end{equation*}
$$

In the case of the Allen-Cahn equation with mollified noise $\xi^{\varepsilon}:=\xi * \varrho_{\varepsilon}$, where $\varrho(t, x)$ is a mollifier and $\varrho_{\varepsilon}(t, x)=\varepsilon^{-(d+2)} \varrho\left(t / \varepsilon^{2}, x / \varepsilon\right)$, a canonical way of building a model proceeds as follows. First, we can just define

$$
\begin{align*}
\left(\Pi_{z} \Xi\right)(\bar{z}) & =\xi^{\varepsilon}(\bar{z}) \\
\left(\Pi_{z} X^{k}\right)(\bar{z}) & =(\bar{z}-z)^{k} \\
\left(\Pi_{z} \tau \bar{\tau}\right)(\bar{z}) & =\left(\Pi_{z} \tau\right)(\bar{z})\left(\Pi_{z} \bar{\tau}\right)(\bar{z}) \quad \forall \tau, \bar{\tau} \in T \tag{3.22}
\end{align*}
$$

To also include the integration map in the fixed point problem let $G=G(t, x)$ be the heat kernel defined as the fundamental solution associated to $\partial_{t} u=\Delta_{x} u$. We split the heat kernel $G$ as

$$
\begin{equation*}
G=R+K, \tag{3.23}
\end{equation*}
$$

where $K:\left(\mathbb{R}^{d+1} \backslash\{0\}\right) \rightarrow \mathbb{R}$ is a singular part satisfying specific algebraic properties, while $R: \mathbb{R}^{d+1} \rightarrow \mathbb{R}$ is a smooth part. The properties of $K$ are:

- $K$ is supported in $\left\{|x|^{2}+|t| \leqslant 1\right\}$ where $|x|=\sum_{j=1}^{d}\left|x_{j}\right|$;
- $K(t, x)=0$ for $t \leqslant 0$ and $K(t,-x)=K(t, x)$ for all $(t, x)$;
- we have

$$
\begin{equation*}
K(t, x)=\frac{1}{|4 \pi t|^{d / 2}} \mathrm{e}^{-|x|^{2} / 4 t} \quad \text { for }|x|^{2}+|t| \leqslant \frac{1}{2} \tag{3.24}
\end{equation*}
$$

and $K(t, x)$ is smooth for $|x|^{2}+|t|>\frac{1}{2}$;

- furthermore, one has vanishing moments

$$
\begin{equation*}
\int_{\mathbb{R}^{d+1}} K(t, x) P(t, x) \mathrm{d} x \mathrm{~d} t=0 \tag{3.25}
\end{equation*}
$$

for all polynomials $P$ of parabolic degree less or equal some fixed $\zeta \geqslant 2$.
Lemma 5.5 in [19] shows that such a splitting indeed exists and also the vanishing moments condition holds for $K$; both parts $K, R$ then satisfy a number of derivative bounds. Furthermore, it follows from Lemma 5.5 in [19] that there exists a decomposition

$$
\begin{equation*}
K\left(z-z^{\prime}\right)=\sum_{n \geqslant 0} K_{n}\left(z-z^{\prime}\right) \tag{3.26}
\end{equation*}
$$

where $n \in \mathbb{N}_{0}$ and suitable derivative bounds and vanishing moment conditions hold for each of the kernels $K_{n}$. With this construction, the abstract integration map is represented by the Taylor-series-like expression

$$
\begin{equation*}
\left(\Pi_{z} \mathcal{I} \tau\right)(\bar{z})=\int K\left(\bar{z}-\bar{z}^{\prime}\right)\left(\Pi_{z} \tau\right)\left(\mathrm{d} \bar{z}^{\prime}\right)+\sum_{\ell} \frac{(\bar{z}-z)^{\ell}}{\ell!}\left\langle f_{z}, \mathcal{J}_{\ell} \tau\right\rangle \tag{3.27}
\end{equation*}
$$

where the linear forms $f_{z} \in \mathcal{G} \subset \mathcal{H}_{+}^{\star}$ are constructed as follows:

$$
\begin{align*}
\left\langle f_{z}, \mathbf{1}\right\rangle & =1 \\
\left\langle f_{z}, X_{i}\right\rangle & =-z_{i}, \\
\left\langle f_{z}, \tau \bar{\tau}\right\rangle & =\left\langle f_{z}, \tau\right\rangle\left\langle f_{z}, \bar{\tau}\right\rangle \quad \forall \tau, \bar{\tau} \in T \\
\left\langle f_{z}, \mathcal{J}_{\ell} \tau\right\rangle & =-\int D^{\ell} K(z-\bar{z})\left(\Pi_{z} \tau\right)(\mathrm{d} \bar{z}), \tag{3.28}
\end{align*}
$$

where in the more general kernel case $K=K(z, \bar{z})$ one has to replace $D^{\ell}$ by $D_{1}^{\ell}$ and the subscript indicates derivative with respect to the first argument. The second term in the definition (3.27) of $\Pi_{z} \mathcal{I} \tau$ may seem somewhat strange, and will be motivated in the next section. Note that it should really be interpreted by applying both sides to test functions, i.e., for all smooth compactly supported functions $\psi$ and all $\tau \in T_{\alpha}$ one requires

$$
\begin{equation*}
\left(\Pi_{z} \mathcal{I} \tau\right)(\psi)=\sum_{n \geq 0} \int_{\mathbb{R}^{d+1}} \psi\left(z^{\prime}\right)\left(\Pi_{z} \tau\right)\left(K_{n ; z z^{\prime}}^{\alpha}\right) \mathrm{d} z^{\prime} \tag{3.29}
\end{equation*}
$$

where $K_{n ; z z^{\prime}}^{\alpha}$ is defined via

$$
\begin{equation*}
K_{n ; z z^{\prime}}^{\alpha}\left(z^{\prime \prime}\right)=K_{n}\left(z^{\prime}-z^{\prime \prime}\right)-\sum_{|k| s<\alpha+2} \frac{\left(z^{\prime}-z\right)^{k}}{k!} D^{k} K_{n}\left(z-z^{\prime \prime}\right) \tag{3.30}
\end{equation*}
$$

The group elements $\Gamma_{z \bar{z}}$ are then simply defined by

$$
\begin{equation*}
\Gamma_{z \bar{z}}=\left(F_{z}\right)^{-1} F_{\bar{z}} \quad \text { where } F_{z}=\Gamma_{f_{z}} \tag{3.31}
\end{equation*}
$$

(recall that $F_{z}$ is invertible because $f_{z}$ is grouplike). In this way, the algebraic property (M1) of Definition 3.4 is automatically satisfied. Regarding the algebraic property (M2), we can note the following relations. Using Sweedler's notation [35], we write $\Delta \tau=\tau^{(1)} \otimes$ $\tau^{(2)}$, although $\Delta \tau$ is usually a sum of such terms. Let further $\gamma_{z \bar{z}}$ denote the element of the structure group such that $\Gamma_{z \bar{z}}=\Gamma_{\gamma_{z \bar{z}}}$. Then we have by (3.15)

$$
\begin{equation*}
\Pi_{z} \Gamma_{z \bar{z}} \tau=\Pi_{z}\left(\operatorname{Id} \otimes \gamma_{z \bar{z}}\right) \Delta \tau=\Pi_{z} \tau^{(1)}\left\langle\gamma_{z \bar{z}}, \tau^{(2)}\right\rangle \tag{3.32}
\end{equation*}
$$

Property (M2) thus amounts to the relation

$$
\begin{equation*}
\Pi_{\bar{z}} \tau=\Pi_{z} \tau^{(1)}\left\langle\gamma_{z \bar{z}}, \tau^{(2)}\right\rangle \tag{3.33}
\end{equation*}
$$

which provides some intuition for the meaning of $\Delta$. It is easy to check that this relation holds for elements $\tau$ of the polynomial regularity structure. In the general case, the fact that $(\Pi, \Gamma)$ is indeed a model is a nontrivial fact, proved in [19, Prop. 8.27].

We can now introduce the spaces $\mathcal{D}^{\gamma}$, which play an analogous role as the $\mathcal{C}_{\mathfrak{s}}^{\alpha}$ on the level of the regularity structure.
Definition 3.5 ([19, Def. 3.1]). Let $\gamma \in \mathbb{R}$. The space $\mathcal{D}^{\gamma}$ consists of all function $f$ : $\mathbb{R}^{d+1} \rightarrow T_{\gamma}^{-}$such that for every compact set $\mathfrak{K} \subset \mathbb{R}^{d+1}$ one has

$$
\begin{equation*}
\|f\|_{\gamma ; \mathfrak{R}}:=\sup _{z \in \mathfrak{K}} \sup _{\beta<\gamma}\|f(z)\|_{\beta}+\sup _{\substack{z, \bar{z} \in \mathfrak{K} \\\|z-\bar{z}\|_{\mathfrak{s}} \leqslant 1}} \sup _{\beta<\gamma} \frac{\left\|f(z)-\Gamma_{z \bar{z}} f(\bar{z})\right\|_{\beta}}{\|z-\bar{z}\|_{\mathfrak{s}}^{\gamma-\beta}}<\infty . \tag{3.34}
\end{equation*}
$$

In the particular case of the polynomial regularity structure, it is again quite straightforward to check that the requirement (3.34) is equivalent to $f$ being the Taylor expansion of an element of the Hölder space $\mathcal{C}_{\mathfrak{s}}^{\gamma}$. The spaces $\mathcal{D}^{\gamma}$ depend on the model via $\Gamma$, but not on $\Pi$, cf [19, Remark 3.4]. In order to compare elements of $\mathcal{D}^{\gamma}$ for different models, it is useful to introduce

$$
\begin{equation*}
\|f ; \bar{f}\|_{\gamma ; \mathfrak{K}}=\|f-\bar{f}\|_{\gamma ; \mathfrak{K}}+\sup _{\substack{z, \bar{z} \in \mathfrak{K} \\\|z-\bar{z}\|_{\mathfrak{s}} \leqslant 1}} \sup _{\beta<\gamma} \frac{\left\|f(z)-\bar{f}(z)-\Gamma_{z \bar{z}} f(\bar{z})+\bar{\Gamma}_{z \bar{z}} \bar{f}(\bar{z})\right\|_{\beta}}{\|z-\bar{z}\|_{\mathfrak{s}}^{\gamma-\beta}} \tag{3.35}
\end{equation*}
$$

(which is in general not a function of $f-\bar{f}$ ), where

$$
\begin{equation*}
\|f-\bar{f}\|_{\gamma ; \mathfrak{K}}=\sup _{z \in \mathfrak{K}} \sup _{\beta<\gamma}\|f(z)-\bar{f}(z)\|_{\beta} . \tag{3.36}
\end{equation*}
$$

The central result allowing to link elements in $\mathcal{D}^{\gamma}$ and in $\mathcal{C}_{\mathfrak{s}}^{\alpha}$ is the reconstruction theorem [19, Thm 3.10]. It states that if $\alpha_{*}=\min A$ and given $r>\left|\alpha_{*}\right|$ there exists, for any $\gamma \in \mathbb{R}$, a unique continuous linear map $\mathcal{R}: \mathcal{D}^{\gamma} \rightarrow \mathcal{C}_{\mathfrak{s}}^{\alpha_{*}}$ such that

$$
\begin{equation*}
\left|\left\langle\mathcal{R} f-\Pi_{z} f(z), \mathcal{S}_{\mathfrak{s}, z}^{\delta} \eta\right\rangle\right| \leqslant C \delta^{\gamma}\|\Pi\|_{\gamma ; \overline{\mathfrak{R}}}\|f\|_{\gamma ; \overline{\mathfrak{\beta}}} \tag{3.37}
\end{equation*}
$$

holds uniformly over all test functions $\eta \in \mathcal{B}_{\mathfrak{s}, 0}^{r}$, all $\delta \in(0,1]$, all $f \in \mathcal{D}^{\gamma}$ and all $z \in \mathfrak{K}$. The constant $C>0$ depends only on $\gamma$ and on the regularity structure, and the set $\overline{\mathfrak{\kappa}}$ is the 1 fattening of $\mathfrak{K}$ (i.e., the points at distance at most 1 from $\mathfrak{K}$ ). The reconstruction theorem also provides bounds on the dependence of $\mathcal{R}$ on the model ( $\Pi, \Gamma)$. Heuristically, (3.37) states that $\mathcal{R} f$ locally looks like $\Pi_{z} f(z)$ near any point $z \in \mathbb{R}^{d+1}$, up to terms of order $\gamma$.

### 3.3 Lifting the convolution maps

We can now lift the operation of convolution with the heat kernel to the space $\mathcal{D}^{\gamma}$. The requirement (3.25) is for compatibility with the condition that $\mathcal{I}$ should define an abstract integration map of order 2 in the sense of [19, Def. 5.7], namely

- I $: T_{\alpha} \rightarrow T_{\alpha+2}$ for every $\alpha \in A$;
- $\mathcal{I} \tau=0$ for every $\tau$ in $\bar{T}$;
- $\mathcal{I} \Gamma \tau-\Gamma \mathcal{I} \tau \in \bar{T}$ for every $\tau \in T$ and every $\Gamma \in \mathcal{G}$.

The central result is given by the so-called multilevel Schauder estimates [19, Thm 5.12], which state in particular that for all $\gamma \in \mathbb{R}$ such that $\gamma+2 \notin \mathbb{N}$, there exists a map $\mathcal{K}_{\gamma}: \mathcal{D}^{\gamma} \rightarrow \mathcal{D}^{\gamma+2}$ such that

$$
\begin{equation*}
\mathcal{R} \mathcal{K}_{\gamma} f=K * \mathcal{R} f \tag{3.38}
\end{equation*}
$$

holds for all $f \in \mathcal{D}^{\gamma}$. In other words, the following diagram commutes:


The map $\mathcal{K}_{\gamma}$ has the following expression. For any $f \in \mathcal{D}^{\gamma}$,

$$
\begin{equation*}
\left(\mathcal{K}_{\gamma} f\right)(z)=\mathcal{I} f(z)+\mathcal{J}(z) f(z)+\left(\mathcal{N}_{\gamma} f\right)(z), \tag{3.39}
\end{equation*}
$$

where for each $\tau \in T_{\alpha}$,

$$
\begin{align*}
\mathcal{J}(z) \tau & =\sum_{|k|_{s}<\alpha+2} \frac{X^{k}}{k!} \int_{\mathbb{R}^{4}} D^{k} K(z-\bar{z})\left(\Pi_{z} \tau\right)(\mathrm{d} \bar{z}),  \tag{3.40}\\
\left(\mathcal{N}_{\gamma} f\right)(z) & =\sum_{|k|_{s}<\gamma+2} \frac{X^{k}}{k!} \int_{\mathbb{R}^{4}} D^{k} K(z-\bar{z})\left(\mathcal{R} f-\Pi_{z} f(z)\right)(\mathrm{d} \bar{z}) . \tag{3.41}
\end{align*}
$$

Note that the last two operators have values in $\bar{T}$, the polynomial part of the regularity structure, and that the only nonlocal operator is $\mathcal{N}_{\gamma}$.

The rôle of the operators $\mathcal{J}$ and $\mathcal{N}_{\gamma}$ is to ensure that $\mathcal{K}_{\gamma} f$ has the properties required to belong to $\mathcal{D}^{\gamma+2}$, which would not be the case if one simply sets $\mathcal{K}_{\gamma} f=\mathcal{I} f$. The maps $\mathcal{J}$ are related to the coefficients $\mathcal{J}_{k_{j}}$ appearing in (3.11); in fact, the $\mathcal{J}_{k_{j}} \tau_{j}$ play the rôle of placeholders for the $\mathcal{J}(z) \tau$.

As for the smooth part $R$ of the heat kernel, it can be lifted as in [19, (7.7)]. Namely, with a smooth kernel $R$ we associate maps $R_{\gamma}: \mathcal{C}_{\mathfrak{s}}^{\alpha_{*}} \rightarrow \mathcal{D}^{\gamma}$ given by

$$
\begin{equation*}
\left(R_{\gamma} v\right)(z)=\sum_{|k|_{s}<\gamma} \frac{X^{k}}{k!} \int_{\mathbb{R}^{d+1}} D^{k} R(z-\bar{z}) v(\bar{z}) \mathrm{d} \bar{z}:=\sum_{|k|_{s}<\gamma} \frac{X^{k}}{k!}\left\langle v, D^{k} R(z-\cdot)\right\rangle \tag{3.42}
\end{equation*}
$$

It follows from [19, Prop 3.28] that for $v \in \mathcal{C}_{\mathfrak{s}}^{\alpha_{*}}$ one has

$$
\left(\mathcal{R} R_{\gamma} v\right)(z)=\langle v, R(z-\cdot)\rangle,
$$

and thus

$$
\begin{equation*}
\mathcal{R} R_{\gamma} \mathcal{R} f=R * \mathcal{R} f \tag{3.43}
\end{equation*}
$$

In other words, the following diagram commutes:


Assume that for an appropriate choice of $\gamma$, one can find an element $U \in \mathcal{D}^{\gamma}$ satisfying the fixed-point equation

$$
\begin{equation*}
U=\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+}\left(\Xi+U-U^{3}\right)+G u_{0}, \tag{3.44}
\end{equation*}
$$

for some $\bar{\gamma} \geqslant \gamma-2$, where $\boldsymbol{R}^{+}(t, x)=1_{\{t>0\}}$ and $G u_{0}$ denotes a suitable lift of the convolution in space of heat kernel and initial condition. Applying the reconstruction operator $\mathcal{R}$ to both sides of this equation, one can show that $u=\mathcal{R} U$ satisfies

$$
\begin{equation*}
u=(K+R) *\left(\boldsymbol{R}^{+}\left[u-u^{3}+\xi^{\varepsilon}\right]\right)+G u_{0}, \tag{3.45}
\end{equation*}
$$

which is equivalent to (3.2); in (3.45) we also use $G u_{0}$ as a notation for the usual convolution in space of heat kernel and initial condition. This is basically the strategy implemented in Theorem 7.8 and Section 9.4 of [19], except that one has to deal with two additional technical difficulties. The first one is that due to the singular behaviour of the heat kernel as time goes to 0 , the definition of the space $\mathcal{D}^{\gamma}$ has to be modified. We will apply this modification to our case in Section 4.3. The second difficulty is that although a fixed point exists for every mollification parameter $\varepsilon>0$, it cannot converge as $\varepsilon \rightarrow 0$. This is why a renormalisation procedure is needed, which we will adapt to our case in Section 6.

## 4 Extension of the regularity structure

Our aim is now to extend the regularity structure built for the Allen-Cahn equation, in order to allow to represent a family of coupled SPDEs-ODEs of the form

$$
\begin{align*}
& \partial_{t} u=\Delta_{x} u+F(u, v)+\xi^{\varepsilon}, \\
& \partial_{t} v=u A_{1}+A_{2} v, \tag{4.1}
\end{align*}
$$

where $F(u, v)$ is a cubic polynomial, and $A_{1} \in \mathbb{R}^{n}, A_{2} \in \mathbb{R}^{n \times n}$ are either scalars as for the classical FitzHugh-Nagumo case with $n=1$, or a vector and a matrix if $v$ has multiple components. Duhamel's formula allows us to represent (mild) solutions of (4.1) as

$$
\begin{align*}
& u_{t}=\int_{0}^{t} S(t-s)\left[\xi_{s}^{\varepsilon}+F\left(u_{s}, v_{s}\right)\right] \mathrm{d} s+S(t) u_{0} \\
& v_{t}=\int_{0}^{t} u_{s} Q(t-s) \mathrm{d} s+\mathrm{e}^{t A_{2}} v_{0} \tag{4.2}
\end{align*}
$$

where $Q(t):=\mathrm{e}^{t A_{2}} A_{1}$. We thus have to lift to the regularity structure the operation of time-integration with respect to $Q$, which has no smoothing effect in space.

In the case where $v$ has values in $\mathbb{R}^{n}$ with $n>1$, the kernel $Q(t)$ is in fact a vector of dimension $n$. In what follows, we will mainly deal with the case of scalar $Q$, since the generalisation to the vectorial case is rather straightforward.

### 4.1 Extension theorem

Let us fix a finite time horizon $T$. Then we can always assume that the kernel $Q$ satisfies the following properties:

- $Q(t)$ is supported on $[0,2 T]$ and smooth for $t>0$;
- $Q(t)=\mathrm{e}^{t A_{2}} A_{1}$ for all $t \in[0, T]$;
- one has

$$
\begin{equation*}
\int_{0}^{\infty} t^{n} Q(t) \mathrm{d} t=0 \tag{4.3}
\end{equation*}
$$

for $n=0,1,2$.
The reason why this is allowed is that we will be interested in showing existence of solutions on a sufficiently small interval $[0, T]$, so that the behaviour of $Q$ outside this interval is not going to matter. In particular, it will not be necessary to decompose $Q$ as a sum of $Q_{n}$ concentrated in sets of radius $2^{-n}$, as in the case of the heat kernel in [19, Section 5], cf (3.26).

In order to be able to represent the time-integration map, we will have to extend our regularity structure and the associated model. We do this by adding to $T$ new symbols denoted $\mathcal{E}(\tau), \tau \in T \backslash \bar{T}$, which we represent by an open blue dot. Thus for instance we write $\mathcal{I}(\mathcal{E}(\Xi))=\dot{i}$. In practice, we will only apply $\mathcal{E}$ to elements $\tau$ of homogeneity $|\tau|<2$. We thus set $V=T_{2}^{-}=\{\tau \in T:|\tau|<2\}$ ( $V$ is called a sector of $T$ ). If $\tau \in V \cap \bar{T}$, where $\bar{T}$ is the polynomial part of the regularity structure, we simply set $\mathcal{E}(\tau)=0$ (the value of $\mathcal{E}(\tau)$ on other elements of $\bar{T}$ does not matter). We also postulate that for $\tau \in T \backslash \bar{T}, \mathcal{E}(\tau)$ has the same homogeneity as $\tau$, and we choose the norm on the vector space generated by the new symbols in such a way that

$$
\begin{equation*}
\|\mathcal{E}(\tau)\|_{\alpha}=\|\tau\|_{\alpha} \quad \forall \alpha \in A \tag{4.4}
\end{equation*}
$$

The operator $\mathcal{E}$ defined in this way is an abstract integration map of order 0 on $V$ in the sense of $[19$, Def. 5.7], i.e.,

- $\mathcal{E}: V \cap T_{\alpha} \rightarrow T_{\alpha}$ for every $\alpha \in A ;$
- $\mathcal{E} \tau=0$ for every $\tau \in V \cap \bar{T}$;
- $\mathcal{E} \Gamma \tau-\Gamma \mathcal{E} \tau \in \bar{T}$ for every $\tau \in V$ and every $\Gamma \in \mathcal{G}$.

The first two properties are obvious by definition of $\mathcal{E} \tau$. Regarding the third property, we first extend the structure group by setting

$$
\begin{equation*}
\Delta(\mathcal{E} \tau)=(\mathcal{E} \otimes \mathrm{Id}) \Delta \tau \tag{4.5}
\end{equation*}
$$

In this way, we have in fact $\mathcal{E} \Gamma \tau-\Gamma \mathcal{E} \tau=0$ for all $\tau$ because

$$
\begin{equation*}
\Gamma_{g} \mathcal{E} \tau=(\operatorname{Id} \otimes g) \Delta(\mathcal{E} \tau)=(\mathcal{E} \otimes g) \Delta \tau=\mathcal{E} \tau^{(1)}\left\langle g, \tau^{(2)}\right\rangle=\mathcal{E} \Gamma_{g} \tau \tag{4.6}
\end{equation*}
$$

Since $\mathcal{E}$ is not regularity-increasing, there can be in principle infinitely many symbols of given homogeneity. The trick, however, will be to build the fixed-point map in such a way that only finitely many new symbols are needed.

The concrete extension procedure is an inductive one. Given a subset $W \subset T$, on which the model $Z=(\Pi, \Gamma)$ is defined, we want to extend the model to a larger set $\widehat{W}=W \cup\{\mathcal{E}(\tau): \tau \in W\}$. We can then apply the usual extension theorem [19, Thm 5.14] to extend the model to $\widehat{W} \cup\{\mathcal{I}(\tau): \tau \in \widehat{W}\}$, and so on as often as needed.

The inductive step from $W$ to $\widehat{W}$ goes as follows. Assume that Definition 3.4 already holds on a regularity structure $(A, W, \mathcal{G})$. To extend the model, we set, for any test function $\psi$ and any $\tau \in W$

$$
\begin{equation*}
\left\langle\Pi_{t, x} \mathcal{E} \tau, \psi\right\rangle=\left\langle\Pi_{t, x} \tau, \hat{\psi}\right\rangle \tag{4.7}
\end{equation*}
$$

where for all $(\bar{t}, y) \in \mathbb{R}^{d+1}$

$$
\begin{equation*}
\hat{\psi}(\bar{t}, y):=\int_{\mathbb{R}} Q(s-\bar{t}) \psi(s, y) \mathrm{d} s=\int_{\bar{t}}^{\bar{t}+2 T} Q(s-\bar{t}) \psi(s, y) \mathrm{d} s . \tag{4.8}
\end{equation*}
$$

Note that in case the model is a continuous function, this is equivalent to

$$
\begin{equation*}
\left(\Pi_{z} \mathcal{E} \tau\right)(\bar{t}, \bar{x})=\int_{\bar{t}-2 T}^{\bar{t}} Q(\bar{t}-s)\left(\Pi_{z} \tau\right)(s, \bar{x}) \mathrm{d} s \tag{4.9}
\end{equation*}
$$

Relation (4.5) allows us to extend the definition of the structure group to $\widehat{W}$, setting

$$
\begin{equation*}
F_{z}(\mathcal{E} \tau)=\Gamma_{f_{z}}(\mathcal{E} \tau)=\left(\operatorname{Id} \otimes f_{z}\right) \Delta(\mathcal{E} \tau)=\left(\mathcal{E} \otimes f_{z}\right) \Delta(\tau)=\mathcal{E} F_{z}(\tau) \tag{4.10}
\end{equation*}
$$

and thus using (4.10) and (4.6) yields

$$
\begin{equation*}
\Gamma_{z z^{\prime}}(\mathcal{E} \tau)=\left(F_{z}\right)^{-1} F_{z^{\prime}}(\mathcal{E} \tau) \tag{4.11}
\end{equation*}
$$

for all $z, z^{\prime} \in \mathbb{R}^{d+1}$.
Remark 4.1. In cases where $v(t, x)$ takes values in $\mathbb{R}^{n}$ and $Q(t) \in \mathbb{R}^{n}$, we can in fact introduce $n$ commuting symbols $\mathcal{E}_{1}, \ldots, \mathcal{E}_{n}$, and define the extended model by the relations

$$
\begin{equation*}
\left\langle\Pi_{t, x} \mathcal{E}_{i} \tau, \psi\right\rangle=\left\langle\Pi_{t, x} \tau, \hat{\psi}_{i}\right\rangle \tag{4.12}
\end{equation*}
$$

for $i=1, \ldots, n$, where

$$
\begin{equation*}
\hat{\psi}_{i}(\bar{t}, y)=\int_{\bar{t}}^{\bar{t}+2 T} Q_{i}(s-\bar{t}) \psi(s, y) \mathrm{d} s \tag{4.13}
\end{equation*}
$$

The results that follow remain true when $\mathcal{E}$ and $Q$ are replaced by $\mathcal{E}_{i}$ and $Q_{i} . \diamond$
Remark 4.2. In the case of the heat kernel, when defining $\Pi_{t, x}(\mathcal{I} \tau)$, one has to subtract a suitable Taylor approximation from the kernel, cf (3.27). The reason for this is that the model for $\mathcal{I} \tau$ should "vanish at the right order", in cases where $|\tau|<0<|\mathcal{I} \tau|$. In the present case, this would amount to replacing $Q(s-\bar{t})$ in the integral (4.8) by

$$
\tilde{Q}_{t \bar{t}}^{\alpha}(s)=Q(s-\bar{t})-\sum_{k_{0}<\alpha / 2} \frac{(s-t)^{k_{0}}}{k_{0}!} Q^{\left(k_{0}\right)}(t-\bar{t})
$$

where $Q^{\left(k_{0}\right)}$ is the $k_{0}$-th derivative of $Q$. One would then also have to introduce additional symbols $\mathcal{J}^{v}$ to represent elements of the Taylor series, and to modify accordingly the definition of $\Delta(\mathcal{E} \tau)$. However, since the kernel is singular in space, there is no obvious way of defining these symbols. We will thus avoid the introduction of such modifications by setting up the fixed-point map in such a way that the only symbols $\tau$ for which $\Pi_{t, x}(\mathcal{E} \tau)$ occurs have negative homogeneity. $\diamond$

The following result is an adaptation of the extension theorem [19, Thm. 5.14] to our degenerate situation.

Proposition 4.3 (Extension theorem for $\mathcal{E}$ ). Let $Z=(\Pi, \Gamma)$ be a model for the regularity structure $(A, W, \mathcal{G})$, where $W \subset T$, and let $\widehat{W}=W \cup\{\mathcal{E} \tau: \tau \in W\}$. Then $\widehat{Z}=(\widehat{\Pi}, \widehat{\Gamma})$ obtained by extending $Z$ in the above way is a model for $(A, \widehat{W}, \mathcal{G})$. In particular, we have

$$
\begin{equation*}
\|\widehat{\Pi}\|_{\gamma ; \mathfrak{R}} \leqslant\|\Pi\|_{\gamma ; \overline{\mathfrak{R}}}\left[I_{1}+I_{2}\|\Gamma\|_{\gamma ; \overline{\mathfrak{R}}}\right], \quad\|\widehat{\Gamma}\|_{\gamma, \mathfrak{R}}=\|\Gamma\|_{\gamma ; \overline{\mathfrak{R}}} \tag{4.14}
\end{equation*}
$$

where $\overline{\mathfrak{K}}$ is the 1-fattening of $\mathfrak{K}$, and

$$
\begin{equation*}
I_{1}=\int_{0}^{2 T}|Q(t)| \mathrm{d} t, \quad I_{2}=\sup _{\substack{\alpha<\gamma \\ \alpha \in A \\ \beta<\alpha}} \sum_{\substack{\beta<\alpha}} \int_{0}^{2 T} t^{(\alpha-\beta) / 2}|Q(t)| \mathrm{d} t \tag{4.15}
\end{equation*}
$$

Furthermore, if $\bar{Z}=(\bar{\Pi}, \bar{\Gamma})$ is a second model and $(\hat{\bar{\Pi}}, \widehat{\bar{\Gamma}})$ is its extension, then

$$
\begin{align*}
& \|\widehat{\Pi}-\widehat{\bar{\Pi}}\|_{\gamma ; \mathfrak{\kappa}} \leqslant\|\Pi-\bar{\Pi}\|_{\gamma ; \overline{\mathfrak{\Omega}}}\left[I_{1}+I_{2}\|\Gamma\|_{\gamma ; \bar{\Omega}}\right]+I_{2}\|\Pi\|_{\gamma ; \overline{\mathfrak{\Omega}}}\|\Gamma-\bar{\Gamma}\|_{\gamma ; \bar{\Omega}},  \tag{4.16}\\
& \|\widehat{\Gamma}-\widehat{\bar{\Gamma}}\|_{\gamma, \mathfrak{K}}=\|\Gamma-\bar{\Gamma}\|_{\gamma ; \overline{\mathfrak{\Omega}}} . \tag{4.17}
\end{align*}
$$

Proof: We first show that the definition of $\Pi_{t, x} \mathcal{E} \tau$ is compatible with the requirement that $\mathcal{E} \tau=0$ for $\tau \in V \cap \bar{T}$. Indeed, since for $k=\left(k_{0}, \underline{k}\right) \in \mathbb{N}_{0}^{d+1}$ one has $\Pi_{t, x} X^{k}(s, y)=$ $(s-t)^{k_{0}}(y-x)^{\underline{k}}$, we obtain by (4.7) that

$$
\left(\Pi_{t, x} \mathcal{E} X^{k}\right)(s, y)=\int_{0}^{2 T} Q(\bar{t})(s-\bar{t}-t)^{k_{0}} \mathrm{~d} \bar{t}(y-x)^{\underline{k}}
$$

which vanishes for all $k_{0} \leqslant 2$ as a consequence of the assumption (4.3) on $Q$.
Next we prove that $\widehat{\Pi}$ satisfies property (M3) of Definition 3.4. Writing $\eta_{t, x}^{\delta}=\mathcal{S}_{\mathfrak{s},(t, x)}^{\delta} \eta$, we have for any $\tau \in W_{\alpha}$ with $\alpha<\gamma$

$$
\left\langle\Pi_{t, x} \mathcal{E} \tau, \eta_{t, x}^{\delta}\right\rangle=\left\langle\Pi_{t, x} \tau, \hat{\eta}_{t, x}^{\delta}\right\rangle
$$

where

$$
\begin{aligned}
\hat{\eta}_{t, x}^{\delta}(\bar{t}, y) & =\int_{\bar{t}}^{\bar{t}+2 T} Q(s-\bar{t}) \eta_{t, x}^{\delta}(s, y) \mathrm{d} s \\
& =\int_{0}^{T} Q(\bar{s}) \eta_{t, x}^{\delta}(\bar{t}+\bar{s}, y) \mathrm{d} \bar{s} \\
& =\int_{0}^{T} Q(\bar{s}) \eta_{t-\bar{s}, x}^{\delta}(\bar{t}, y) \mathrm{d} \bar{s}
\end{aligned}
$$

This implies

$$
\begin{aligned}
\left\langle\Pi_{t, x} \mathcal{E} \tau, \eta_{t, x}^{\delta}\right\rangle & =\int_{0}^{T} Q(\bar{s})\left\langle\Pi_{t, x} \tau, \eta_{t-\bar{s}, x}^{\delta}\right\rangle \mathrm{d} \bar{s} \\
& =\int_{0}^{T} Q(\bar{s})\left\langle\Pi_{t-\bar{s}, x} \Gamma_{t-\bar{s}, x ; t, x} \tau, \eta_{t-\bar{s}, x}^{\delta}\right\rangle \mathrm{d} \bar{s}
\end{aligned}
$$

By property (M3) for the original model,

$$
\left|\left\langle\Pi_{t-\bar{s}, x} \Gamma_{t-\bar{s}, x ; t, x} \tau, \eta_{t-\bar{s}, x}^{\delta}\right\rangle\right| \leqslant\|\Pi\|_{\gamma ; \boldsymbol{\Omega}}\left\|\Gamma_{t-\bar{s}, x ; t, x} \tau\right\| \delta^{\alpha} .
$$

By property (M4) and the fact that $\Gamma \tau-\tau \in T_{\alpha}^{-}$we further have

$$
\begin{aligned}
\left\|\Gamma_{t-\bar{s}, x ; t, x} \tau\right\| & \leqslant\|\tau\|_{\alpha}+\sum_{\beta<\alpha}\left\|\Gamma_{t-\bar{s}, x ; ;, x} \tau\right\|_{\beta} \\
& \leqslant\|\tau\|_{\alpha}+\sum_{\beta<\alpha}\|\Gamma\|_{\gamma ; \mathfrak{R}}\|\tau\|_{\alpha}\|(t-\bar{s}, x)-(t, x)\|_{\mathfrak{s}}^{\alpha-\beta} \\
& =\|\tau\|_{\alpha}\left[1+\|\Gamma\|_{\gamma ; \mathfrak{K}} \sum_{\beta<\alpha} \bar{s}^{(\alpha-\beta) / 2}\right] .
\end{aligned}
$$

This proves property (M3) for the new model, as well as the first bound in (4.14).
Property (M1) is automatically satisfied by the new model thanks to (4.11). Regarding property (M2), it is known [19, Sec. 8.3] that there exists a linear map $\Pi: T \rightarrow \mathcal{S}^{\prime}\left(\mathbb{R}^{d+1}\right)$ such that $\Pi \tau=\Pi_{z} F_{z}^{-1} \tau$ for $\tau \in T$ is independent of $z$. Property (M2) holds if we can find for any $\tau \in W$ a distribution $\boldsymbol{\Pi}(\mathcal{E} \tau)$, independent of $(t, x)$, such that

$$
\Pi_{t, x}(\mathcal{E} \tau)=\boldsymbol{\Pi} F_{t, x}(\mathcal{E} \tau) \quad \forall z=(t, x) \in \mathbb{R}^{d+1}
$$

To achieve this, we simply define $\boldsymbol{\Pi}$ on the extended structure by

$$
\langle\boldsymbol{\Pi} \mathcal{E} \tau, \psi\rangle=\langle\boldsymbol{\Pi} \tau, \hat{\psi}\rangle
$$

where $\hat{\psi}$ is defined by (4.8), and check that it has the required property. Indeed, by (4.10), we have

$$
\begin{aligned}
\left\langle\boldsymbol{\Pi} F_{t, x}(\mathcal{E} \tau), \psi\right\rangle & =\left\langle\boldsymbol{\Pi} \mathcal{E}\left(F_{t, x} \tau\right), \psi\right\rangle \\
& =\left\langle\boldsymbol{\Pi} F_{t, x} \tau, \hat{\psi}\right\rangle \\
& =\left\langle\Pi_{t, x} \tau, \hat{\psi}\right\rangle \\
& =\left\langle\Pi_{t, x} \mathcal{E} \tau, \psi\right\rangle,
\end{aligned}
$$

which is precisely what we needed to show.
Regarding Property (M4), we note that (4.6) and (4.4) imply that for all $\tau \in W_{\alpha}$,

$$
\left\|\Gamma_{z \bar{z}} \mathcal{E} \tau\right\|_{\beta}=\left\|\mathcal{E} \Gamma_{z \bar{z}} \tau\right\|_{\beta}=\left\|\Gamma_{z \bar{z}} \tau\right\|_{\beta} \leqslant\|\Gamma\|_{\gamma ; \mathfrak{\beta}}\|\tau\|^{\alpha}\|z-\bar{z}\|_{\mathfrak{s}}^{\alpha-\beta},
$$

which also proves the second relation in (4.14).
It remains to prove the bounds involving two different models. For this we note that

$$
\left\langle\Pi_{z} \mathcal{E} \tau-\bar{\Pi}_{z} \mathcal{E} \tau, \eta_{z}^{\delta}\right\rangle=\int_{0}^{2 T} Q(s)\left\langle\left[\Pi_{z^{\prime}(s)} \Gamma_{z^{\prime}(s) z}-\bar{\Pi}_{z^{\prime}(s)} \bar{\Gamma}_{z^{\prime}(s) z}\right] \tau, \eta_{z^{\prime}(s)}^{\delta}\right\rangle \mathrm{d} s
$$

where we have used the notation $z^{\prime}(s)=(t-s, x)$. We add and subtract the term $\Pi_{z^{\prime}(s)} \bar{\Gamma}_{z^{\prime}(s) z}$ inside the bracket and note that for $\tau \in T_{\alpha}$,

$$
\begin{aligned}
& \left.\left|\left\langle\Pi_{z^{\prime}(s)}\right| \Gamma_{z^{\prime}(s) z}-\bar{\Gamma}_{z^{\prime}(s) z}\right] \tau, \eta_{z^{\prime}(s)}^{\delta}\right\rangle \mid \leqslant\|\Pi\|_{\gamma ; \bar{\Omega}} \delta^{\alpha}\|\tau\|_{\alpha} \sum_{\beta<\alpha} s^{(\alpha-\beta) / 2}\|\Gamma-\bar{\Gamma}\|_{\gamma ; \bar{\kappa}} \\
& \mid\left\langle\left[\Pi_{z^{\prime}(s)}-\bar{\Pi}_{z^{\prime}(s)} \bar{\Gamma}_{z^{\prime}(s) z} \tau, \eta_{z^{\prime}(s)}^{\delta}\right\rangle\right| \leqslant\|\Pi-\bar{\Pi}\|_{\gamma ; \overline{\bar{\kappa}}} \delta^{\alpha}\|\tau\|_{\alpha}\left[1+\sum_{\beta<\alpha} s^{(\alpha-\beta) / 2}\|\bar{\Gamma}\|_{\gamma ; \bar{\beta}}\right],
\end{aligned}
$$

from which (4.16) follows easily. Finally, (4.17) is a direct consequence of the fact that $\mathcal{E}$ and $\Gamma_{z z^{\prime}}$ commute.

### 4.2 Multilevel Schauder estimates

We now would like to construct an operator $\mathcal{K}_{\gamma}^{v}: \mathcal{D}^{\gamma} \rightarrow \mathcal{D}^{\gamma}$ which lifts the integration against $Q$ to the regularity structure. Formally, this means that we should have

$$
\begin{equation*}
\mathcal{R} \mathcal{K}_{\gamma}^{v} f=K_{v} * \mathcal{R} f, \quad K_{v}(t, x, \bar{t}, \bar{x})=Q(t-\bar{t}) \delta(x-\bar{x}) \tag{4.18}
\end{equation*}
$$

for all $f \in \mathcal{D}^{\gamma}$. In fact, this should really be interpreted as

$$
\begin{equation*}
\left(\mathcal{R} \mathcal{K}_{\gamma}^{v} f\right)(t, x)=\int_{t-2 T}^{t} Q(t-\bar{t})(\mathcal{R} f)(\bar{t}, x) \mathrm{d} \bar{t} \tag{4.19}
\end{equation*}
$$

or in terms of test functions

$$
\begin{equation*}
\left\langle\mathcal{R} \mathcal{K}_{\gamma}^{v} f, \psi\right\rangle=\langle\mathcal{R} f, \hat{\psi}\rangle \tag{4.20}
\end{equation*}
$$

where $\hat{\psi}$ is defined in (4.8).
The problem with such a plan is that the kernel $K_{v}$ being singular in space, we cannot apply an expansion as in (3.41). However, for the class of equations we are interested in, we do not actually need to define $\mathcal{K}_{\gamma}^{v}$ on all of $\mathcal{D}^{\gamma}$. It will be quite sufficient to define it on a subset of $\mathcal{D}^{\gamma}$, which is given by the functions whose components with negative homogeneity do not depend on $(t, x)$. There we can simply set $\mathcal{K}_{\gamma}^{v}=\mathcal{E}$. We thus make the following definition.

Definition 4.4. The space $\mathcal{D}_{0}^{\gamma}(\Gamma)$ is the space of functions $f \in \mathcal{D}_{\gamma}(\Gamma)$ of the form

$$
\begin{equation*}
f(z)=\sum_{\tau \in T:|\tau|<0} c_{\tau} \tau+\sum_{\tau \in T:|\tau| \geqslant 0} f_{\tau}(z) \tau=: f_{-}+f_{+}(z), \tag{4.21}
\end{equation*}
$$

where the $c_{\tau}$ do not depend on $z$.
Remark 4.5. Not all constant functions belong to $\mathcal{D}^{\gamma}$, as they still have to satisfy the analytical bound (3.34), which includes the requirement $\left\|f(z)-\Gamma_{z \bar{z}} f(\bar{z})\right\|_{\beta}=\mathcal{O}\left(\|z-\bar{z}\|^{\gamma-\beta}\right)$ as $z \rightarrow \bar{z}$ for all $\beta<\gamma$. In fact, since $\Gamma_{z \bar{z}} \tau=\left(\operatorname{Id} \otimes \gamma_{z \bar{z}}\right) \Delta \tau$, we see that a sufficient condition for having $\Gamma_{z \bar{z}} \tau=\tau$ is

$$
\Delta \tau=\tau \otimes \mathbf{1}
$$

Thus the sum defining $f_{-}$should involve only terms $c_{\tau} \tau$ such that $\Delta \tau=\tau \otimes \mathbb{1} . \diamond$
The following result shows that the multilevel Schauder estimates contained in [19, Thm. 5.12] also hold in the case of $\mathcal{E}$, in a similar form.

Proposition 4.6 (Mutlilevel Schauder estimates on $\mathcal{E}$ ). The operator $\mathcal{E}$ maps $\mathcal{D}_{0}^{\gamma}$ into $\mathcal{D}_{0}^{\gamma}$, and satisfies

$$
\begin{equation*}
\|\mathcal{E} f\|_{\gamma ; \mathfrak{R}}=\|f\|_{\gamma ; \mathfrak{R}} . \tag{4.22}
\end{equation*}
$$

The identity

$$
\begin{equation*}
\langle\mathcal{R E} f, \psi\rangle=\langle\mathcal{R} f, \hat{\psi}\rangle \tag{4.23}
\end{equation*}
$$

holds for all $f \in \mathcal{D}_{0}^{\gamma}$ and all test functions $\psi$, where $\hat{\psi}$ is defined in (4.8). Furthermore, if $(\bar{\Pi}, \bar{\Gamma})$ is a second model compatible with $\mathcal{E},{ }^{2}$ then

$$
\begin{equation*}
\|\mathcal{E} f ; \mathcal{E} \bar{f}\|_{\gamma ; \mathfrak{R}}=\|f ; \bar{f}\|_{\gamma ; \mathfrak{R}} \tag{4.24}
\end{equation*}
$$

for all $f \in \mathcal{D}_{0}^{\gamma}(\Gamma)$ and $\bar{f} \in \mathcal{D}_{0}^{\gamma}(\bar{\Gamma})$.

[^1]Proof: We start by showing that relation (4.23) holds, assuming $f, \mathcal{E} f \in \mathcal{D}_{0}^{\gamma}$. It follows from the reconstruction theorem and its corollary [19, Prop. 3.28] that ${ }^{3}$

$$
\mathcal{R} f_{+}(z)=\left\langle\mathbf{1}, f_{+}(z)\right\rangle=\Pi_{z} f_{+}(z)
$$

for all $z \in \mathbb{R}^{d+1}$. Since $\Pi_{z} f_{-}=\sum_{|\tau|<0} c_{\tau} \Pi_{z}(\tau)$, by the uniqueness part of the reconstruction theorem, we also have

$$
\mathcal{R} f_{-}(z)=\sum_{|\tau|<0} c_{\tau} \Pi_{z}(\tau)=\Pi_{z} f_{-} .
$$

Thus in fact $\mathcal{R} f(z)=\Pi_{z} f(z)$ holds for all $z \in \mathbb{R}^{d+1}$. It follows that for any test function $\eta_{z}^{\delta}$ localised around $z$,

$$
\begin{aligned}
\left\langle\mathcal{R E} f, \eta_{z}^{\delta}\right\rangle & =\left\langle\mathcal{R E} f-\Pi_{z} \mathcal{E} f(z), \eta_{z}^{\delta}\right\rangle+\left\langle\Pi_{z} \mathcal{E} f, \eta_{z}^{\delta}\right\rangle \\
& =\left\langle\Pi_{z} \mathcal{E} f, \eta_{z}^{\delta}\right\rangle \\
& =\left\langle\Pi_{z} f, \hat{\eta}_{z}^{\delta}\right\rangle \\
& =\left\langle\mathcal{R} f, \hat{\eta}_{z}^{\delta}\right\rangle
\end{aligned}
$$

We have used the fact that $\mathcal{E} f \in \mathcal{D}_{0}^{\gamma}$, (4.7), and finally the fact that $f \in \mathcal{D}_{0}^{\gamma}$. This proves (4.23).

Next we check that $\mathcal{E}$ indeed maps $\mathcal{D}_{0}^{\gamma}$ into itself. Since $\mathcal{E}$ preserves the homogeneity, it is in fact sufficient to show that $\mathcal{E}$ maps $\mathcal{D}_{0}^{\gamma}$ into $\mathcal{D}^{\gamma}$. Since $\mathcal{E}$ is norm-preserving, we obviously have for any compact $\mathfrak{K}$,

$$
\sup _{z \in \mathfrak{K}} \sup _{\beta<\gamma}\|\mathcal{E} f(z)\|_{\beta}=\sup _{z \in \mathfrak{K}} \sup _{\beta<\gamma}\|f(z)\|_{\beta} .
$$

Furthermore, since $\mathcal{E}$ and $\Gamma_{z \bar{z}}$ commute,

$$
\begin{aligned}
\left\|\mathcal{E} f(z)-\Gamma_{z \bar{z}} \mathcal{E} f(\bar{z})\right\|_{\beta} & =\left\|\mathcal{E}\left(f(z)-\Gamma_{z \bar{z}} f(\bar{z})\right)\right\|_{\beta} \\
& =\left\|f(z)-\Gamma_{z \bar{z}} f(\bar{z})\right\|_{\beta} \\
& \leqslant\|f\|_{\gamma ; \mathfrak{R}}\|z-\bar{z}\|_{\mathfrak{s}}^{\gamma-\beta} .
\end{aligned}
$$

This shows that indeed $\mathcal{E} f \in \mathcal{D}^{\gamma}$, and that the bound (4.4) holds.
It remains to prove the relations (4.24) involving the quantity $\|\cdot ; \cdot\|_{\gamma, \mathfrak{\Re}}$ defined in (3.35). Since $\mathcal{E}$ is norm-preserving in each $T_{\beta}$, we see that $\|\cdot ; \cdot\|_{\gamma, \mathfrak{K}}$ is actually invariant under the $\operatorname{map}(f, \bar{f}) \mapsto(\mathcal{E} f, \mathcal{E} \bar{f})$ by using that $\mathcal{E}$ commutes with $\Gamma_{z \bar{z}}$ for all $z, \bar{z} \in \mathbb{R}^{d+1}$.

### 4.3 Extension to the space $\mathcal{D}^{\gamma, \eta}$

As already mentioned at the end of Section 3.3, the definition of the spaces $\mathcal{D}^{\gamma}$ has to be modified in order to be able to deal with the time-zero singularity of the heat kernel. In particular, we have to check that the estimates on $\mathcal{E}$ established in Proposition 4.6 still hold in these modified spaces.

We recall a few notations from [19, Section 6]. Let $P=\left\{(t, x) \in \mathbb{R}^{d+1}: t=0\right\}$ denote the time-zero hyperplane, and introduce the scaled norms

$$
\begin{equation*}
\|z\|_{P}=1 \wedge|t|^{1 / 2}, \quad\|z, \bar{z}\|_{P}=\|z\|_{P} \wedge\|\bar{z}\|_{P} \tag{4.25}
\end{equation*}
$$

[^2]For a compact $\mathfrak{K} \subset \mathbb{R}^{d+1}$, let

$$
\begin{equation*}
\mathfrak{K}_{P}=\left\{(z, \bar{z}) \in(\mathfrak{K} \backslash P)^{2}: z \neq \bar{z},\|z-\bar{z}\|_{\mathfrak{s}} \leqslant\|z, \bar{z}\|_{P}\right\} . \tag{4.26}
\end{equation*}
$$

Definition 4.7 ([19, Def. 6.2]). Given a model $(\Pi, \Gamma)$ and constants $\gamma>0$ and $\eta \in \mathbb{R}$, set

$$
\begin{equation*}
\|f\|_{\gamma, \eta ; \mathfrak{R}}:=\sup _{z \in \mathfrak{K} \backslash P} \sup _{\beta<\gamma} \frac{\|f(z)\|_{\beta}}{\|z\|_{P}^{(n-\beta) \wedge 0}}, \quad \square f \rrbracket_{\gamma, \eta ; \mathfrak{R}}:=\sup _{z \in \mathfrak{K} \backslash P} \sup _{\beta<\gamma} \frac{\|f(z)\|_{\beta}}{\|z\|_{P}^{\eta-\beta}} . \tag{4.27}
\end{equation*}
$$

The space $\mathcal{D}^{\gamma, \eta}$ consists of all functions $f: \mathbb{R}^{d+1} \backslash P \rightarrow T_{\gamma}^{-}$such that

$$
\begin{equation*}
\|f\|_{\gamma, \eta ; \mathfrak{\Re}}:=\|f\|_{\gamma, \eta ; \mathfrak{R}}+\sup _{(z, \bar{z}) \in \mathfrak{R}_{P}} \sup _{\beta<\gamma} \frac{\left\|f(z)-\Gamma_{z \bar{z}} f(\bar{z})\right\|_{\beta}}{\|z-\bar{z}\|_{\mathfrak{s}}^{\gamma-\beta}\|z, \bar{z}\|_{P}^{\eta-\gamma}}<\infty . \tag{4.28}
\end{equation*}
$$

Given a second model $(\bar{\Pi}, \bar{\Gamma})$ and $\bar{f} \in \mathcal{D}^{\gamma, \eta}(\bar{\Gamma})$, we also set

$$
\begin{equation*}
\|f ; \bar{f}\|_{\gamma, \eta ; \mathfrak{\Re}}:=\|f-\bar{f}\|_{\gamma, \eta ; \mathfrak{\beta}}+\sup _{(z, \bar{z}) \in \mathfrak{R}_{P}} \sup _{\beta<\gamma} \frac{\left\|f(z)-\bar{f}(z)-\Gamma_{z \bar{z}} f(\bar{z})+\bar{\Gamma}_{z \bar{z}} \bar{f}(\bar{z})\right\|_{\beta}}{\|z-\bar{z}\|_{\mathfrak{s}}^{\gamma-\beta}\|z, \bar{z}\|_{P}^{\eta-\gamma}} . \tag{4.29}
\end{equation*}
$$

Finally, similarly to Definition 4.4, we define $\mathcal{D}_{0}^{\gamma, \eta}$ to be the set of $f \in \mathcal{D}^{\gamma, \eta}$ whose components with negative homogeneity do not depend on $z$.

Most of the results in [19, Section 6] are directly applicable to the present setting. The only result that has to be adapted is Proposition 6.16. Fortunately, this turns out to be very easy.

Proposition 4.8. Let $f \in \mathcal{D}_{0}^{\gamma, \eta}$ with $\eta<\gamma$, and assume that $\eta \wedge \alpha>-2$, where $\alpha$ is the lowest homogeneity of components of $f$. Then $\mathcal{E} f \in \mathcal{D}_{0}^{\gamma, \eta}$, and one has

$$
\begin{align*}
\|\mathcal{E} f\|_{\gamma, \eta ; \mathfrak{R}} & =\|f\|_{\gamma, \eta ; \mathfrak{R}}, \\
\|\mathcal{E} f ; \mathcal{E} \bar{f}\|_{\gamma, \eta ; \mathfrak{R}} & =\|f ; \bar{f}\|_{\gamma, \eta ; \mathfrak{R}} \tag{4.30}
\end{align*}
$$

for any $\bar{f} \in \mathcal{D}_{0}^{\gamma, \eta}(\bar{\Gamma})$, where $(\bar{\Pi}, \bar{\Gamma})$ is any other model for the regularity structure.
Proof: The proof is straightforward, using the facts that $\mathcal{E}$ commutes with $\Gamma_{z \bar{z}}$ and preserves the norm in each $T_{\beta}$. In particular, we have $\|\mathcal{E} f\|_{\gamma, \eta ; \mathfrak{R}}=\|f\|_{\gamma, \eta ; \mathfrak{K}}$.

## 5 The fixed-point equation

In this section, we construct a fixed-point equation lifting (4.2) to an appropriate space $\mathcal{D}^{\gamma, \eta}$, and prove that it admits a unique fixed point, depending Lipschitz-continuously on the initial data and the noise. The procedure is very similar to the one for the AllenCahn equation or the $\Phi_{3}^{4}$ model carried out in [19, Section 9], with only a few adjustments required by the additional variable $v$.

### 5.1 The set-up

From now on, we assume that the space dimension is $d=3$, since this is the most difficult case which is still renormalisable. Before lifting the fixed-point equation (4.2) to the space $\mathcal{D}^{\gamma, \eta}$, we will reformulate it somewhat. Although the resulting equation is seemingly more
complicated, it will turn out to be easier to analyse, because $\Pi_{z}(\mathcal{E} \tau)$ will only have to be evaluated for $\tau$ of negative homogeneity (cf Remark 4.2). We introduce the notation

$$
\begin{equation*}
S^{Q}(t, s)=\int_{s}^{t} Q(t-r) S(r-s) \mathrm{d} r=\int_{s}^{t} S(t-r) Q(r-s) \mathrm{d} r . \tag{5.1}
\end{equation*}
$$

Inserting the first equation in (4.2) into the second one and interchanging the integrals, we obtain

$$
\begin{align*}
v_{t} & =\int_{0}^{t} Q(t-s)\left[\int_{0}^{s} S(s-r)\left[\xi_{r}^{\varepsilon}+F\left(u_{r}, v_{r}\right)\right] \mathrm{d} r+S(s) u_{0}\right] \mathrm{d} s+\mathrm{e}^{t A_{2}} v_{0} \\
& =\int_{0}^{t} S^{Q}(t, s)\left[\xi_{s}^{\varepsilon}+F\left(u_{s}, v_{s}\right)\right] \mathrm{d} s+S^{Q}(t, 0) u_{0}+\mathrm{e}^{t A_{2}} v_{0} \\
& =\int_{0}^{t} S(t-s)\left[\int_{0}^{s} Q(s-r)\left[\xi_{r}^{\varepsilon}+F\left(u_{r}, v_{r}\right)\right] \mathrm{d} r\right] \mathrm{d} s+S^{Q}(t, 0) u_{0}+\mathrm{e}^{t A_{2}} v_{0} . \tag{5.2}
\end{align*}
$$

Consider now the fixed-point equation on $\mathcal{D}^{\gamma, \eta}$ given by

$$
\begin{align*}
U & =\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+}[\Xi+F(U, V)]+G u_{0}, \\
V & =\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \mathcal{E} \boldsymbol{R}^{+}[\Xi+F(U, V)]+G^{Q} u_{0}+\widehat{Q} v_{0} \tag{5.3}
\end{align*}
$$

where $\boldsymbol{R}^{+}(t, x)=1_{\{t>0\}}$ and $\bar{\gamma} \geqslant \gamma-2$ as before, and by definition, $\mathcal{R} \boldsymbol{R}^{+} \Xi$ is given by the distribution $\xi^{\varepsilon} 1_{\{t>0\}}$. Furthermore,

$$
\begin{align*}
\left(G u_{0}\right)(t, x) & =\left(S(t) u_{0}\right)(x)=\int_{\mathbb{R}^{3}} G(t, x-y) u_{0}(y) \mathrm{d} y \\
\left(G^{Q} u_{0}\right)(t, x) & =\left(S^{Q}(t, 0) u_{0}\right)(x)=\int_{0}^{t} Q(t-s) \int_{\mathbb{R}^{3}} G(s, x-y) u_{0}(y) \mathrm{d} y \mathrm{~d} s \\
\left(\widehat{Q} v_{0}\right)(t, x) & =\mathrm{e}^{t A_{2}} v_{0}(x) 1_{\{t>0\}} \tag{5.4}
\end{align*}
$$

In fact, these distributions have to be lifted to $\mathcal{D}^{\gamma, \eta}$, which is always possible by construction, cf [19, Def. 2.14]. In order to understand the structure of the solution of the fixed-point equation (5.3), let us rewrite it in the form

$$
\begin{align*}
U & =\mathcal{I} \boldsymbol{R}^{+}[\Xi+F(U, V)]+\widehat{R}_{u}(U, V)+G u_{0}, \\
V & =\mathcal{I E} \boldsymbol{R}^{+}[\Xi+F(U, V)]+\widehat{R}_{v}(U, V)+G^{Q} u_{0}+\widehat{Q} v_{0} \tag{5.5}
\end{align*}
$$

where $\widehat{R}_{u}(U, V)$ contains the polynomial part of $\mathcal{K}_{\bar{\gamma}}[\Xi+F(U, V)]$ as well as the term $R_{\gamma} \mathcal{R}[\Xi+F(U, V)]$ (which always belongs to $\bar{T}$ ), and $\widehat{R}_{v}(U, V)$ is defined analogously. Assume $F$ is a cubic polynomial of the form

$$
\begin{equation*}
F(u, v)=\alpha_{1} u+\alpha_{2} v+\beta_{1} u^{2}+\beta_{2} u v+\beta_{3} v^{2}+\gamma_{1} u^{3}+\gamma_{2} u^{2} v+\gamma_{3} u v^{2}+\gamma_{4} v^{3} . \tag{5.6}
\end{equation*}
$$

Further assume for simplicity that the initial conditions $u_{0}$ and $v_{0}$ belong to some $\mathcal{C}^{\alpha}$ with $\alpha \geqslant 0$. Iterating the map (5.5), starting for instance with identically zero functions, it is not difficult to see that any fixed point of (5.3) must have the form
where the coefficients satisfy $a_{i}=\gamma_{i}$ for $i=1,2,3,4$ and

$$
\begin{align*}
& b_{1}=\beta_{1}+3 \varphi \gamma_{1}+\psi \gamma_{2}, \\
& b_{2}=\beta_{2}+2 \varphi \gamma_{2}+2 \psi \gamma_{3},  \tag{5.8}\\
& b_{3}=\beta_{3}+\varphi \gamma_{3}+3 \psi \gamma_{4} .
\end{align*}
$$

In (5.7), $\varrho_{U}, \varrho_{V}$ denote terms of homogeneity at least $\frac{3}{2}-\kappa$, and $\varphi(z), \psi(z), \nabla \varphi(z), \nabla \psi(z)$ are independent functions - the notation is not supposed to suggest that $\varphi$ and $\psi$ are differentiable. The scalar product notation $\langle\nabla \varphi, X\rangle$ is a shorthand for $\sum_{i=0}^{3}(\nabla \varphi)_{i} X_{i}$.

Proposition 5.1 (Equivalence of fixed-point equations). Assume that $\varepsilon>0$ and that (5.3) admits a fixed point $(U, V)$ where $U, V \in \mathcal{D}^{\gamma, \eta}$ for some $\gamma, \eta \in \mathbb{R}$. Then $(u, v)=(\mathcal{R} U, \mathcal{R} V)$ satisfies the fixed-point equations (4.2) and (5.2).

Proof: The first step is to note that (3.27) implies

$$
\begin{aligned}
& \left(\Pi_{\bar{z}}^{\varepsilon} \uparrow\right)(z)=\int K\left(z-z_{1}\right) \xi^{\varepsilon}\left(z_{1}\right) \mathrm{d} z_{1}=: \chi_{\varepsilon}(z) \\
& \left(\Pi_{\bar{z}}^{\varepsilon} \boldsymbol{\vartheta}\right)(z)=\int K^{Q}\left(z-z_{1}\right) \xi^{\varepsilon}\left(z_{1}\right) \mathrm{d} z_{1}=: \chi_{\varepsilon}^{Q}(z)
\end{aligned}
$$

independently of $\bar{z}$, where we have used the notation

$$
K^{Q}(t, x)=\int_{0}^{2 T} K(t-s, x) Q(s) \mathrm{d} s
$$

Applying the reconstruction operator $\mathcal{R}$ to (5.7), since all functions are smooth for $\varepsilon>0$, we obtain

$$
\begin{aligned}
& u(z)=\chi_{\varepsilon}(z)+\varphi(z), \\
& v(z)=\chi_{\varepsilon}^{Q}(z)+\psi(z) .
\end{aligned}
$$

Here we have used the fact that if $f_{+} \in \mathcal{D}^{\gamma, \eta}$ has only components of strictly positive homogeneity, one has $\mathcal{R} f_{+}(z)=\Pi_{z}^{\varepsilon} f_{+}(z)=0$. On the other hand, using (3.38) and (3.43), we obtain for all $f \in \mathcal{D}^{\gamma, \eta}$

$$
\left(\mathcal{R}\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) f\right)(z)=(G * \mathcal{R} f)(z) .
$$

In a similar way, using (4.9) we see that

$$
\begin{aligned}
\left(\mathcal{R}\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \mathcal{E} f\right)(t, x) & =(G * \mathcal{R E} f)(t, x) \\
& =\iint G(t-s, x-y) \int Q(s-\bar{s}) \mathcal{R} f(\bar{s}, y) \mathrm{d} \bar{s} \mathrm{~d} s \mathrm{~d} y \\
& =\iint\left[\int G\left(s_{1}, x-y\right) Q\left(t-\bar{s}-s_{1}\right) \mathrm{d} s_{1}\right] \mathcal{R} f(\bar{s}, y) \mathrm{d} \bar{s} \mathrm{~d} y \\
& =\left(G^{Q} * \mathcal{R} f\right)(t, x) .
\end{aligned}
$$

Applying $\mathcal{R}$ to the equations (5.3), we thus have

$$
\begin{aligned}
u & =G * \mathcal{R} \boldsymbol{R}^{+}[\Xi+F(U, V)]+G u_{0}, \\
v & =G^{Q} * \mathcal{R} \boldsymbol{R}^{+}[\Xi+F(U, V)]+G^{Q} u_{0}+\widehat{Q} v_{0}
\end{aligned}
$$

It is not difficult to check that $\mathcal{R}$ and $\boldsymbol{R}^{+}$commute. Thus the result follows if we are able to show that

$$
\mathcal{R} F(U, V)=F(\mathcal{R} U, \mathcal{R} V)
$$

Computing the part of nonpositive homogeneity of the difference $\mathcal{R} F(U, V)=F(\mathcal{R} U, \mathcal{R} V)$, we see that it is a linear combination of terms of the form

$$
\left(\Pi_{z}^{\varepsilon} \mathscr{\bullet}\right)(z), \quad\left(\Pi_{z}^{\varepsilon} \mathscr{\bullet} \cdot\right)(z), \quad\left(\Pi_{z}^{\varepsilon} \dot{\mathscr{\bullet}}\right)(z) \quad \text { and } \quad\left(\Pi_{z}^{\varepsilon} \bullet\right)(z)^{2}\left\langle\nabla \varphi,\left(\Pi_{z}^{\varepsilon} X\right)(z)\right\rangle,
$$

and all possible terms obtained from these by substitutions of the form $\boldsymbol{\dagger} \mapsto \dot{\boldsymbol{i}}$. It follows directly from (3.22) that $\left(\Pi_{z}^{\varepsilon} X\right)(z)=0$. It thus remains to show that all terms of the three other types vanish as well. Applying (3.27), (3.28), and recalling (3.14) we obtain

$$
\left(\Pi_{\bar{z}}^{\varepsilon} \mathscr{Y}\right)(z)=\int\left[K\left(z-z_{1}\right)-K\left(\bar{z}-z_{1}\right)\right]\left(\Pi_{\bar{z}}^{\varepsilon} \mathscr{V}\right)\left(z_{1}\right) \mathrm{d} z_{1}
$$

which vanishes in $z=\bar{z}$. It follows that $\left(\Pi_{z}^{\varepsilon} \mathscr{\bullet}\right)(z)=\left(\Pi_{z}^{\varepsilon} \Upsilon^{\bullet}\right)(z)\left(\Pi_{z}^{\varepsilon} \mathscr{V}\right)(z)=0$. Similarly, noting $\mathcal{E} \mathscr{V}=\mathfrak{\zeta}$, we have

$$
\left(\Pi_{\bar{z}}^{\varepsilon} \hat{\zeta}\right)(z)=\int\left[K\left(z-z_{1}\right)-K\left(\bar{z}-z_{1}\right)\right]\left(\Pi_{\bar{z}}^{\varepsilon} \hat{\gamma}\right)\left(z_{1}\right) \mathrm{d} z_{1}
$$

which also vanishes in $z=\bar{z}$. It follows that $\left(\Pi_{z}^{\varepsilon} \mathfrak{\varepsilon} \cdot\right)(z)=\left(\Pi_{z}^{\varepsilon} \not{ }^{\bullet}\right)(z)\left(\Pi_{z}^{\varepsilon} \mathcal{V}\right)(z)=0$. All other terms can be treated similarly, and the result follows.

The expansion (5.7) shows in particular that both $U$ and $V$ have regularity $-\frac{1}{2}-\kappa$, and that their components of negative homogeneity do not depend on $(t, x)$. We will thus look for a fixed point of (5.7) in a space $\mathcal{D}_{0}^{\gamma, \eta}$ for appropriate values of $\gamma$ and $\eta$.

### 5.2 Local existence and uniqueness of the fixed point

In order to describe the effect of $F$, it will be useful to introduce the following slight generalisation of Definition 4.4. Given a model $(\Pi, \Gamma)$, the space $\mathcal{D}_{\beta, \alpha}^{\gamma, \eta}(\Gamma)$ is the space of functions $f \in \mathcal{D}^{\gamma, \eta}(\Gamma)$ of the form

$$
\begin{equation*}
f(z)=\sum_{\tau \in T: \alpha \leqslant|\tau|<\beta} c_{\tau} \tau+\sum_{\tau \in T:|\tau| \geqslant \beta} f_{\tau}(z) \tau, \tag{5.9}
\end{equation*}
$$

where the $c_{\tau}$ do not depend on $z$.
The following result establishes the strong local Lipschitz continuity of $F$ in the sense of [19, Section 7.3].
Proposition 5.2 (Strong local Lipschitz continuity). Fix $\alpha<0<\gamma$ and set $\bar{\gamma}=\gamma+2 \alpha$ and $\bar{\eta}=3 \eta \wedge(\eta+2 \alpha)$. Let $\mathfrak{K} \subset \mathbb{R}^{d+1}$ be compact, and assume $U, V \in \mathcal{D}_{0, \alpha}^{\gamma, \eta}(\Gamma)$ are such that

$$
\begin{equation*}
\|U\|_{\gamma, \eta ; \mathfrak{\Re}}+\|V\|_{\gamma, \eta ; \mathfrak{R}} \leqslant R . \tag{5.10}
\end{equation*}
$$

Then $F(U, V) \in \mathcal{D}_{2 \alpha, 3 \alpha}^{\bar{\gamma}, \bar{\eta}}(\Gamma)$, and there exists a constant $C(R)$ such that

$$
\begin{equation*}
\|F(U, V)\|_{\bar{\gamma}, \bar{\eta} ; \mathfrak{R}} \leqslant C(R) . \tag{5.11}
\end{equation*}
$$

Furthermore, if $(\bar{\Pi}, \bar{\Gamma})$ is a possibly different model and $\bar{U}, \bar{V} \in \mathcal{D}_{0, \alpha}^{\gamma, \eta}(\bar{\Gamma})$ satisfy the bound $\|\bar{U}\|_{\gamma, \eta ; \mathfrak{\Re}}+\|\bar{V}\|_{\gamma, \eta ; \mathfrak{\Re}} \leqslant R$, then

$$
\begin{equation*}
\|F(U, V) ; F(\bar{U}, \bar{V})\|_{\bar{\gamma}, \bar{\eta} ; \mathfrak{\Re}} \leqslant C(R)\left[\|U ; \bar{U}\|_{\gamma, \eta ; \mathfrak{\kappa}}+\|V ; \bar{V}\|_{\gamma, \eta ; \mathfrak{\Re}}+\|\Gamma-\bar{\Gamma}\|_{2 \gamma+\alpha ; \mathfrak{\kappa}}\right] . \tag{5.12}
\end{equation*}
$$

Proof: Writing $U(z)=U_{-}+U_{+}(z)$ where $U_{-}$contains all terms of negative homogeneity and applying [19, Prop. 6.12], we see that

$$
\begin{array}{ll}
U(z)^{2}=U_{-}^{2}+2 U_{-} U_{+}(z)+U_{+}(z)^{2} & \in \mathcal{D}_{\alpha, 2 \alpha}^{\gamma+\alpha, 2 \eta \wedge(\eta+\alpha)} \\
U(z)^{3}=U_{-}^{3}+3 U_{-}^{2} U_{+}(z)+3 U_{-} U(z)^{2}+U_{+}(z)^{3} \in \mathcal{D}_{2 \alpha, 3 \alpha}^{\gamma+2 \alpha, 3 \eta \wedge(\eta+2 \alpha)}
\end{array}
$$

where the multiplication of elements in $\mathcal{D}^{\gamma,, \eta}$ is well-defined due to the results in [19, Sec. 4]. Similar relations hold for powers of $V$ and cross-terms, proving (5.11). The proof of (5.12) is similar, using the bound on $\|f ; g\|$ given in [19, Prop. 6.12].

Remark 5.3. In cases where the two models $(\Pi, \Gamma)$ and $(\bar{\Pi}, \bar{\Gamma})$ are identical, (5.12) automatically provides the bound

$$
\begin{equation*}
\|F(U, V)-F(\bar{U}, \bar{V})\|_{\bar{\gamma}, \bar{\eta} ; \mathfrak{R}} \leqslant C(R)\left[\|U-\bar{U}\|_{\gamma, \eta ; \mathfrak{R}}+\|V-\bar{V}\|_{\gamma, \eta ; \mathfrak{R}}\right] . \tag{5.13}
\end{equation*}
$$

In order to proceed, it is convenient to rewrite the fixed-point equation (5.3) in a slightly different form. First we set

$$
\begin{align*}
& W_{1}:=\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+} \Xi, \\
& W_{2}:=\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \mathcal{E} \boldsymbol{R}^{+} \Xi . \tag{5.14}
\end{align*}
$$

For the time being, let us assume that $W_{1}, W_{2} \in \mathcal{D}_{0, \alpha_{0}+2}^{\gamma, \eta}$ for some $\gamma \in \mathbb{R}$ and all $\eta<\alpha_{0}+2$ (recall that $\alpha_{0}=-\frac{5}{2}-\kappa$ in dimension $d=3$ ); we come back to this point in Remark 5.5 below. Then we define the maps

$$
\begin{align*}
& \mathcal{M}_{1}:(U, V) \mapsto\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+} F(U, V)+W_{1}+G u_{0}, \\
& \mathcal{M}_{2}:(U, V) \mapsto\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \mathcal{E} \boldsymbol{R}^{+} F(U, V)+W_{2}+G^{Q} u_{0}+\widehat{Q} v_{0} . \tag{5.15}
\end{align*}
$$

If $(U, V)$ is a fixed point of $\mathcal{M}=\left(\mathcal{M}_{1}, \mathcal{M}_{2}\right)$ then it is a fixed point of (5.3).
Following [19, Section 7.1], we write $O=[-1,2] \times \mathbb{R}^{3}$ and $O_{T}=(-\infty, T] \times \mathbb{R}^{3}$, and introduce the shorthand $\|\cdot\|_{\gamma, \eta ; T}$ for $\|\cdot\|_{\gamma, \eta ; O_{T}}$, and define $\|\cdot ; \cdot\|_{\gamma, \eta ; T}$ similarly. The following result is an adaptation of [19, Thm. 7.8].

Proposition 5.4 (Existence and uniqueness of the fixed point). Assume $-\frac{2}{3}<\eta<-\frac{1}{2}$ and $\alpha>-1$ are such that $\eta+2 \alpha>-2$, and assume $\gamma>-2 \alpha$. For any $u_{0}, v_{0}, W_{1}, W_{2}$ such that $W_{1}+G u_{0}, W_{2}+G^{Q} u_{0}+\widehat{Q} v_{0} \in \mathcal{D}_{0, \alpha}^{\gamma, \eta}$, there exists a time $T>0$ such that $\mathcal{M}$ admits a unique fixed point $\left(U^{*}, V^{*}\right) \in \mathcal{D}^{\gamma, \eta} \times \mathcal{D}^{\gamma, \eta}$ on $(0, T)$. Furthermore, the solution map $\mathcal{S}_{T}:\left(u_{0}, v_{0}, W_{1}, W_{2}, Z\right) \mapsto\left(U^{*}, V^{*}\right)$ is jointly Lipschitz continuous in the sense that if $Z$ and $\bar{Z}$ are two models and $\left(u_{0}, v_{0}, W_{1}, W_{2}\right)$ and $\left(\bar{u}_{0}, \bar{v}_{0}, \bar{W}_{1}, \bar{W}_{2}\right)$ are two sets of initial conditions and forcing terms such that

$$
\begin{equation*}
\left\|G u_{0} ; G \bar{u}_{0}\right\|_{\gamma, \eta ; T}+\left\|\widehat{Q} v_{0} ; \widehat{Q} \bar{v}_{0}\right\|_{\gamma, \eta ; T}+\sum_{i=1,2}\left\|W_{i} ; \bar{W}_{i}\right\|_{\gamma, \eta ; T}+\|Z ; \bar{Z}\|_{2 \gamma+\alpha ; O} \leqslant \delta \tag{5.16}
\end{equation*}
$$

then one has

$$
\begin{equation*}
\left\|U^{*} ; \bar{U}^{*}\right\|_{\gamma, \eta ; T}+\left\|V^{*} ; \bar{V}^{*}\right\|_{\gamma, \eta ; T} \leqslant C_{0} \delta \tag{5.17}
\end{equation*}
$$

for some constant $C_{0}>0$.

Proof: Given $R>0$ let $\mathcal{B}(R)=\left\{f \in \mathcal{D}_{0, \alpha}^{\gamma, \eta}:\|f\|_{\gamma, \eta ; T} \leqslant R\right\}$ and pick $U, V \in \mathcal{B}(R)$. Applying Proposition 5.2, we see that $F(U, V) \in \mathcal{D}_{2 \alpha, 3 \alpha}^{\bar{\gamma}, \bar{\eta}}$ and satisfies

$$
\|F(U, V)\|_{\bar{\gamma}, \bar{\eta} ; T} \leqslant C(2 R) .
$$

The assumptions on $\gamma, \eta$ and $\alpha$ guarantee that $\bar{\gamma}>0$ and $\bar{\eta}>-2$. We can thus apply Theorem 7.1 and Lemma 7.3 in [19] to obtain the existence of constants $C_{1}, \kappa>0$ such that

$$
\begin{aligned}
& \left\|\mathcal{M}_{1}(U)\right\|_{\gamma, \eta ; T} \leqslant C_{1} T^{\kappa / 2} C(2 R)+\left\|W_{1}\right\|_{\gamma, \eta ; T}+\left\|G u_{0}\right\|_{\gamma, \eta ; T}, \\
& \left\|\mathcal{M}_{2}(U)\right\|_{\gamma, \eta ; T} \leqslant C_{1} T^{\kappa / 2} C(2 R)+\left\|W_{2}\right\|_{\gamma, \eta ; T}+\left\|G^{Q} u_{0}\right\|_{\gamma, \eta ; T}+\left\|\widehat{Q} v_{0}\right\|_{\gamma, \eta ; T},
\end{aligned}
$$

for every $T \in(0,1]$. Choosing

$$
R=1+\left[\left\|W_{1}\right\|_{\gamma, \eta ; T}+\left\|G u_{0}\right\|_{\gamma, \eta ; T}\right] \vee\left[\left\|W_{2}\right\|_{\gamma, \eta ; T}+\left\|G^{Q} u_{0}\right\|_{\gamma, \eta ; T}+\left\|\widehat{Q} v_{0}\right\|_{\gamma, \eta ; T}\right]
$$

and $T$ sufficiently small, we see that $\mathcal{M}$ maps the set $\mathcal{B}(R) \times \mathcal{B}(R)$ into itself.
Next we show that $\mathcal{M}$ is a contraction in $\mathcal{B}(R) \times \mathcal{B}(R)$. For $U, \bar{U}, V, \bar{V} \in \mathcal{B}(R)$, we have

$$
\mathcal{M}_{1}(U, V)-\mathcal{M}_{1}(\bar{U}, \bar{V})=\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+}[F(U, V)-F(\bar{U}, \bar{V})] .
$$

Using again Theorem 7.1 and Lemma 7.3 in [19] and applying (5.13) to bound the term $[F(U, V)-F(\bar{U}, \bar{V})]$, we get

$$
\left\|\mathcal{M}_{1}(U, V)-\mathcal{M}_{1}(\bar{U}, \bar{V})\right\|_{\gamma, \eta ; T} \leqslant C_{1} T^{\kappa / 2} C(2 R)\left[\|U-\bar{U}\|_{\gamma, \eta ; T}+\|V-\bar{V}\|_{\gamma, \eta ; T}\right]
$$

as well as an analogous bound for $\mathcal{M}_{2}$. Taking again $T$ small enough, we obtain that $\mathcal{M}$ is indeed a contraction, and the existence of a unique fixed point in a ball in $\mathcal{D}^{\gamma, \eta} \times \mathcal{D}^{\gamma, \eta}$ follows from Banach's fixed-point theorem.

Finally, in order to prove the bound (5.17), we set $U_{1}=\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+} F(U, V)$ and $\bar{U}_{1}=\left(\mathcal{K}_{\bar{\gamma}}+R_{\gamma} \mathcal{R}\right) \boldsymbol{R}^{+} F(\bar{U}, \bar{V})$. Using the inequality

$$
\|f+g ; \bar{f}+\bar{g}\|_{\gamma, \eta ; T} \leqslant\|f ; \bar{f}\|_{\gamma, \eta ; T}+\|g ; \bar{g}\|_{\gamma, \eta ; T}
$$

we obtain

$$
\left\|\mathcal{M}_{1}(U, V) ; \mathcal{M}_{1}(\bar{U}, \bar{V})\right\|_{\gamma, \eta ; T} \leqslant\left\|U_{1} ; \bar{U}_{1}\right\|_{\gamma, \eta ; T}+\left\|W_{1} ; \bar{W}_{1}\right\|_{\gamma, \eta ; T}+\left\|G u_{0} ; G \bar{u}_{0}\right\|_{\gamma, \eta ; T}
$$

where the first term on the right-hand side can be estimated with the help of (5.12). As before, an analogous estimate holds for $\mathcal{M}_{2}$. The required bound thus follows from [19, Thm. 7.1 and Lemma 7.3].

## Remark 5.5.

1. The conditions on $\eta$ and $\alpha$ in fact impose $\alpha>-\frac{3}{4}$. A possibility is to further impose that $\alpha \geqslant \eta$, in which case the only remaining conditions are those on $\eta$ and $\gamma$.
2. Lemma 7.5 in [19] shows that if $u_{0} \in \mathcal{C}^{\eta}$ (with the Euclidean scaling and we recall we take $\left.-\frac{2}{3}<\eta<-\frac{1}{2}\right)$ then $G u_{0} \in \mathcal{D}^{\gamma, \eta}$ for any $\gamma>(\eta \vee 0)$.
3. As discussed in Section 9.4 of [19] (see Equation (9.18) and below), the assumption $W_{1} \in \mathcal{D}_{0, \alpha_{0}+2}^{\gamma, \eta}$ is indeed satisfied for $\gamma, \eta$ as in Proposition 5.4, provided we define $\mathcal{R} \boldsymbol{R}^{+} \Xi$ as the distribution $\xi^{\varepsilon} 1_{\{t \geqslant 0\}} . \diamond$

Let $Z^{\varepsilon}=\Psi\left(\xi^{\varepsilon}\right)$ denote the canonical model for mollified noise $\xi^{\varepsilon}$ constructed in Sections 3.2 and 4.1. In exactly the same way as in [19, Prop. 9.8 and 9.10], we can now extend the solution map $\mathcal{S}_{T}$ to a maximal solution map $\mathcal{S}^{L}$, defined until the first time $\|(\mathcal{R} U)(t, \cdot)\|_{\eta}+\|(\mathcal{R} V)(t, \cdot)\|_{\eta}$ reaches a (large) cut-off value $L$. This map satisfies

$$
\begin{equation*}
\mathcal{R S}^{L}\left(u_{0}, v_{0}, Z^{\varepsilon}\right)=\overline{\mathcal{S}}^{L}\left(u_{0}, v_{0}, \xi^{\varepsilon}\right), \tag{5.18}
\end{equation*}
$$

where $\overline{\mathcal{S}}^{L}$ is the solution map of the original SPDE-ODE system (4.1) with mollified noise $\xi^{\varepsilon}$. In other words, the following diagram commutes:


## 6 Renormalisation of the equation

### 6.1 The renormalisation group

Consider a general model $Z=(\Pi, \Gamma)$ for our regularity structure, which is admissible in the sense that is satisfies (3.27), the last identity in (3.28) and (4.7). The model may also be specified by the pair ( $\boldsymbol{\Pi}, f)$ such that for all $z, \bar{z} \in \mathbb{R}^{d+1}$, one has

$$
\begin{equation*}
\Pi_{z}=\Pi \Gamma_{f_{z}} \quad \text { and } \quad \Gamma_{z \bar{z}}=\left(\Gamma_{f_{z}}\right)^{-1} \Gamma_{f_{\bar{z}}}=: \Gamma_{\gamma_{z \bar{z}}} . \tag{6.1}
\end{equation*}
$$

To build a renormalisation transformation, one first needs to define sets $\mathcal{F}_{\star} \subset \mathcal{F}_{0} \subset \mathcal{F}_{F}$ such that

- $\mathcal{F}_{0}$ contains all $\tau \in \mathcal{F}_{F}$ of strictly negative homogeneity;
- for all $\tau \in \mathcal{F}_{0}$, one has $\Delta \tau \in \operatorname{span}\left(\mathcal{F}_{0}\right) \otimes \operatorname{span}\left(\mathcal{F}_{0}^{+}\right)$, where $\mathcal{F}_{0}^{+} \subset \mathcal{F}_{F}^{+}$contains all symbols of the form (3.11) with $\tau_{j} \in \mathcal{F}_{\star}$ and $\left|\mathcal{J}_{k_{j}} \tau_{j}\right|>0$ (including the empty product).
For the Allen-Cahn model, a possible choice is

$$
\begin{align*}
& \mathcal{F}_{\star}=\{\stackrel{\bullet}{\bullet}, \mathscr{V}, \bullet\}, \tag{6.2}
\end{align*}
$$

For FitzHugh-Nagumo type equations, we choose the same sets, enriched by all symbols obtained by substitutions $\uparrow \rightarrow \boldsymbol{i}$.

Let $\mathcal{H}_{0}=\operatorname{span}\left(\mathcal{F}_{0}\right)$ and $\mathcal{H}_{0}^{+}=\operatorname{span}\left(\mathcal{F}_{0}^{+}\right)$. By definition, a renormalisation transformation is a linear map $M: \mathcal{H}_{0} \rightarrow \mathcal{H}_{0}$ such that $M X^{k}=X^{k}$ and $M \mathcal{I} \tau=\mathcal{I} M \tau$ for all $\tau \in \mathcal{F}_{0}$ such that $\mathcal{I} \tau \in \mathcal{F}_{0}$ (we impose a similar condition involving $\mathcal{E}$ ).

With $M$, we would like to associate another admissible model ( $\boldsymbol{\Pi}^{M}, f^{M}$ ) by setting

$$
\begin{equation*}
\boldsymbol{\Pi}^{M} \tau=\boldsymbol{\Pi} M \tau \quad \forall \tau \in \mathcal{F}_{0} \tag{6.3}
\end{equation*}
$$

As shown in [19, Section 8.3], this is indeed possible provided one is able to find linear maps $\Delta^{M}: \mathcal{H}_{0} \rightarrow \mathcal{H}_{0} \otimes \mathcal{H}_{0}^{+}$, and $\hat{M}: \mathcal{H}_{0}^{+} \rightarrow \mathcal{H}_{0}^{+}$such that

$$
\begin{align*}
\Pi_{z}^{M} \tau & =\left(\Pi_{z} \otimes f_{z}\right) \Delta^{M} \tau, \\
\left\langle f_{z}^{M}, \bar{\tau}\right\rangle & =\left\langle f_{z}, \hat{M} \bar{\tau}\right\rangle \tag{6.4}
\end{align*}
$$

holds for all $\tau \in \mathcal{F}_{0}$ and all $\bar{\tau} \in \mathcal{F}_{0}^{+}$, where $\Pi_{z}^{M}$ and $f_{z}^{M}$ satisfy relations analogous to (6.1). The map $\hat{M}$ should also be a multiplicative morphism leaving the $X^{k}$ invariant, that is,

$$
\begin{equation*}
\hat{M}\left(\tau_{1} \tau_{2}\right)=\left(\hat{M} \tau_{1}\right)\left(\hat{M} \tau_{2}\right), \quad \hat{M}\left(X^{k}\right)=X^{k} \tag{6.5}
\end{equation*}
$$

In addition, the elements $\Gamma_{z \bar{z}}^{M}$ can be built using a map $\hat{\Delta}^{M}: \mathcal{H}_{0}^{+} \rightarrow \mathcal{H}_{0}^{+} \otimes \mathcal{H}_{0}^{+}$such that

$$
\begin{equation*}
\gamma_{z \bar{z}}^{M}=\left(\gamma_{z \bar{z}} \otimes f_{\bar{z}}\right) \hat{\Delta}^{M} \tag{6.6}
\end{equation*}
$$

It is shown in [19, Prop. 8.36] that given a linear map $M$, there is a unique choice of maps $\Delta^{M}, \hat{\Delta}^{M}$ and $\hat{M}$ such that (6.4) and (6.6) hold, making ( $\boldsymbol{\Pi}^{M}, f^{M}$ ) an admissible model.

Definition 6.1 ([19, Def. 8.41]). The renormalisation group $\mathfrak{R}$ is the set of linear maps $M: \mathcal{H}_{0} \rightarrow \mathcal{H}_{0}$ satisfying $M \mathcal{I}=\mathcal{I} M, M \mathcal{E}=\mathcal{E} M, M\left(X^{k}\right)=X^{k}$, and such that for all $\tau \in \mathcal{F}_{0}$ and all $\bar{\tau} \in \mathcal{F}_{0}^{+}$one has

$$
\begin{align*}
& \Delta^{M} \tau-\tau \otimes \mathbf{1} \in \bigoplus_{\beta>\alpha}\left(T_{\beta} \cap \mathcal{F}_{0}\right) \otimes T^{+}, \\
& \hat{\Delta}^{M} \bar{\tau}-\bar{\tau} \otimes \mathbf{1} \in \bigoplus_{\beta>\alpha}\left(T_{\beta} \cap \mathcal{F}_{0}^{+}\right) \otimes T^{+} \tag{6.7}
\end{align*}
$$

for all $\tau \in T_{\alpha}$ and all $\bar{\tau} \in T_{\alpha} \cap \mathcal{F}_{0}^{+}$.
It is shown in [19, Lemma 8.43] that $\mathfrak{R}$ is indeed a group. The main result [19, Thm. 8.44] states that given $M \in \mathfrak{R}$, and an admissible model $(\boldsymbol{\Pi}, f)$, the model $\left(\boldsymbol{\Pi}^{M}, f^{M}\right)$ constructed as above is indeed an admissible model. In particular, it satisfies the required analytic bounds and behaves well under the extension theorem [19, Thm. 5.14].

Let us now introduce a particular subgroup of $\mathfrak{R}$ suited to our problem. In the case of the Allen-Cahn equation (cf [19, Sections 9.2 and 10.5]), a two-parameter group of transformations $M=\exp \left\{-C_{1} L_{1}-C_{2} L_{2}\right\}$ is sufficient, where the generators $L_{1}$ and $L_{2}$ are defined by the substitution rules

$$
\begin{equation*}
L_{1}: \mathcal{V} \mapsto \mathbf{1}, \quad L_{2}: \mathbb{D} \mapsto \mathbf{1} \tag{6.8}
\end{equation*}
$$

Each substitution rule is applied as often as possible, so for instance $L_{1}(\mathscr{\bullet})=3 \boldsymbol{\uparrow}$, because there are 3 ways of picking $\mathcal{V}$ in $\stackrel{*}{*}$.

In our case, we may a priori have to apply similar substitutions for all elements obtained from $\mathcal{V}$ and $\mathbb{Z}$ by replacing one or several $\mathcal{I}$ by $\mathcal{I E}$. We thus introduce the sets of symbols

$$
\begin{align*}
& \mathcal{F}_{1}=\{\mathfrak{V}, \mathfrak{V}, \dot{\mathfrak{P}}\}, \\
& \mathcal{F}_{2}=\left\{\mathcal{I}\left(\tau_{1}\right) \tau_{2}, \mathcal{E} \mathcal{I}\left(\tau_{1}\right) \tau_{2}: \tau_{1}, \tau_{2} \in \mathcal{F}_{1}\right\} . \tag{6.9}
\end{align*}
$$

Note that $\mathcal{F}_{2}$ has 18 elements. This yields a 21 -parameter subgroup of $\mathfrak{R}$ given by

$$
\begin{equation*}
M=\exp \left\{-\sum_{\tau \in \mathcal{F}_{1}} C_{1}(\tau) L_{\tau}-\sum_{\tau \in \mathcal{F}_{2}} C_{2}(\tau) L_{\tau}\right\}, \tag{6.10}
\end{equation*}
$$

where $L_{\tau}$ denotes the substitution $\tau \mapsto \mathbf{1}$. In practice, however, we will see that only 2 of these parameters are really needed to renormalise the model.

We introduce the shorthands $C_{1}=C_{1}(\mathscr{V}), C_{1}^{\prime}=C_{1}\left(\mathfrak{v}^{\circ}\right)$ and $C_{1}^{\prime \prime}=C_{1}\left(\mathfrak{q}^{\circ}\right)$. The action of $M$ on a few elements that will occur in the computations below is given in the following list:

$$
\begin{align*}
& M(\boldsymbol{\imath})=\boldsymbol{\imath}, \\
& M(\mathcal{V})=\boldsymbol{V}-C_{1} \mathbf{1}, \\
& M(\ddot{*})=\ddot{\vee}-3 C_{1} \uparrow, \\
& M(\dot{\mathscr{V}})=\mathscr{\mathscr { V }}-C_{1} \mathscr{Y}-C_{2}(\dot{\mathscr{V}}) \mathbf{1} . \tag{6.11}
\end{align*}
$$

The action of $M$ on elements such as $\mathfrak{\beta}$ is obtained by obvious substitutions, e.g. we have $M(\ddot{\theta})=\ddot{\psi}-C_{1}^{\prime \prime} \bullet-2 C_{1}^{\prime} \boldsymbol{ף}$. As in [19, Section 9.2], one can determine the map $\Delta^{M}$ and thus, via (6.4), the renormalised model $\Pi^{M}$. Again, we just list a few expressions that will play a rôle in the computations below:

$$
\begin{align*}
& \Pi_{z}^{M}(\bullet)=\Pi_{z}(\bullet), \\
& \Pi_{z}^{M}(\mathscr{V})=\Pi_{z}(\mathfrak{\imath})^{2}-C_{1}, \\
& \Pi_{z}^{M}(\mathscr{\bullet})=\Pi_{z}(\boldsymbol{\varphi})^{3}-3 C_{1} \Pi_{z}(\boldsymbol{\uparrow}), \\
& \Pi_{z}^{M}(\mathscr{Y})=\Pi_{z}(\boldsymbol{Y}), \\
& \Pi_{z}^{M}(\mathcal{M})=\Pi_{z}(\mathscr{Y}) \Pi_{z}^{M}(\mathscr{V})-C_{2}(\mathbb{V} \cdot \mathbf{0}) . \tag{6.12}
\end{align*}
$$

Similar relations with obvious substitutions hold for the elements obtained by having some $\mathcal{I}$ act on $\mathcal{E}$.

### 6.2 Convergence of the renormalised models

From now on, we assume that the driving noise $\xi$ is Gaussian white noise on $D=\mathbb{R} \times \mathbb{T}^{d}$, where $\mathbb{T}^{d}=\mathbb{R}^{d} / \mathbb{Z}^{d}$ is the $d$-dimensional torus. This means that we are given a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, the Hilbert space $H=L^{2}(D)$, and a collection of centred Gaussian random variables $\left\{W_{h}\right\}_{h \in H}$ such that

$$
\begin{equation*}
\mathbb{E}\left\{W_{h} W_{\bar{h}}\right\}=\langle h, \bar{h}\rangle \tag{6.13}
\end{equation*}
$$

for all $h, \bar{h} \in H$. Then $\xi$ is the distribution defined by $\langle\xi, \varphi\rangle=W_{\pi \varphi}$ for every test function $\varphi$, where $\pi$ is the canonical projection from $L^{2}\left(\mathbb{R}^{d+1}\right)$ to $L^{2}(D)$.

Let $Z^{\varepsilon}=\left(\Pi^{\varepsilon}, \Gamma^{\varepsilon}\right)$ be the canonical model built for mollified noise $\xi^{\varepsilon}$. Our aim is now to find a specific sequence of renormalisation maps $M_{\varepsilon}$ in $\mathfrak{R}$ such that the sequence of models $\widehat{Z}^{\varepsilon}=M_{\varepsilon} Z^{\varepsilon}$ defined by

$$
\begin{equation*}
\widehat{\Pi}_{z}^{(\varepsilon)}=\left(\Pi_{z}^{\varepsilon}\right)^{M_{\varepsilon}} \tag{6.14}
\end{equation*}
$$

converges in a suitable sense to a limiting model $\widehat{Z}$.
The relevant result to achieve this is Theorem 10.7 in [19]. It states that if for all localised scaled test functions $\eta_{z}^{\lambda}=\mathcal{S}_{\mathfrak{s}, z}^{\lambda} \eta$ and all $\tau \in \mathcal{F}_{-}=\left\{\tau \in \mathcal{F}_{F}:|\tau|_{\mathfrak{s}}\right\}<0$, one can find random variables $\widehat{\Pi}_{z} \tau$ such that

$$
\begin{align*}
\mathbb{E}\left|\left\langle\widehat{\Pi}_{z} \tau, \eta_{z}^{\lambda}\right\rangle\right|^{2} & \lesssim \lambda^{2|\tau|_{s}+\kappa} \\
\mathbb{E}\left|\left\langle\widehat{\Pi}_{z} \tau-\widehat{\Pi}_{z}^{(\varepsilon)} \tau, \eta_{z}^{\lambda}\right\rangle\right|^{2} & \lesssim \varepsilon^{2 \theta} \lambda^{2|\tau|_{s}+\kappa} \tag{6.15}
\end{align*}
$$

holds for some $\kappa, \theta>0$, then there exists a unique admissible model $\widehat{Z}$ such that for all compact $\mathfrak{K} \subset \mathbb{R}^{d+1}$ and all $p \geqslant 1$ one has

$$
\begin{equation*}
\mathbb{E}\|\widehat{Z}\|_{\gamma ; \mathfrak{K}}^{p} \lesssim 1, \quad \mathbb{E}\left\|\widehat{Z} ; \widehat{Z}^{\varepsilon}\right\|_{\gamma ; \mathfrak{K}}^{p} \lesssim \varepsilon^{\theta p} \tag{6.16}
\end{equation*}
$$

Here $\|Z\|_{\gamma ; \mathfrak{\Re}}=\|\Pi\|_{\gamma ; \mathfrak{K}}+\|\Gamma\|_{\gamma ; \mathfrak{K}}$, and $\|Z ; \bar{Z}\|_{\mathfrak{K}}$ is defined similarly in terms of the associated "distances". In (6.15), each random variable $\widehat{\Pi}_{z} \tau$ should in addition belong to the inhomogeneous Wiener chaos of order $\|\tau\|$, where $\|\tau\|$ is the number of occurrences of $\Xi$ in $\tau(\mathrm{cf}[19$, Section 10.1]). More precisely, one has for every test function $\varphi$

$$
\begin{equation*}
\left\langle\widehat{\Pi}_{0}^{(\varepsilon)} \tau, \varphi\right\rangle=\sum_{k \leqslant\|\tau\|} I_{k}\left(\int_{\mathbb{R}^{d+1}} \varphi(z)\left(\widehat{\mathcal{W}}^{(\varepsilon ; k)} \tau\right)(z) \mathrm{d} z\right) \tag{6.17}
\end{equation*}
$$

for some functions $\left(\widehat{\mathcal{W}}^{(\varepsilon ; k)}\right)(z) \in H^{\otimes k}$, and a similar expression for the bare model $\Pi^{\varepsilon}$, with functions $\mathcal{W}^{(\varepsilon ; k)}$. The map $I_{k}: H^{\otimes k} \rightarrow L^{2}$ is constructed from Wiener's isometry as explained in [19, Section 10.1]. As usual, we will express (6.17) by the slightly more informal notation

$$
\begin{equation*}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \tau\right)(z)=\sum_{k \leqslant\|\tau\|} I_{k}\left(\left(\widehat{\mathcal{W}}^{(\varepsilon ; k)} \tau\right)(z)\right) . \tag{6.18}
\end{equation*}
$$

Proposition 10.11 in [19] provides conditions on the $\widehat{\mathcal{W}}^{(\varepsilon ; k)}$ from which the required conditions (6.15) follow. Namely, for each $\tau \in \mathcal{F}_{-}$, one should find functions $\widehat{\mathcal{W}}^{(k)}$ such that

$$
\begin{align*}
\left|\left\langle\left(\widehat{\mathcal{W}}^{(k)} \tau\right)(z),\left(\widehat{\mathcal{W}}^{(k)} \tau\right)(\bar{z})\right\rangle\right| & \leqslant C \sum_{\zeta}\left(\|z\|_{\mathfrak{s}}+\|\bar{z}\|_{\mathfrak{s}}\right)^{\zeta}\|z-\bar{z}\|_{\mathfrak{s}}^{\bar{\kappa}+2|\tau|_{\mathfrak{s}}-\zeta}, \\
\left|\left\langle\left(\delta \widehat{\mathcal{W}}^{(\varepsilon ; k)} \tau\right)(z),\left(\delta \widehat{\mathcal{W}}^{(\varepsilon ; k)} \tau\right)(\bar{z})\right\rangle\right| & \leqslant C \varepsilon^{2 \theta} \sum_{\zeta}\left(\|z\|_{\mathfrak{s}}+\|\bar{z}\|_{\mathfrak{s}}\right)^{\zeta}\|z-\bar{z}\|_{\mathfrak{s}}^{\bar{\kappa}+2|\tau|_{\mathfrak{s}}-\zeta} \tag{6.19}
\end{align*}
$$

for some $\bar{\kappa}, \theta>0$, where $\delta \widehat{\mathcal{W}}^{(\varepsilon ; k)}=\widehat{\mathcal{W}}^{(\varepsilon ; k)}-\widehat{\mathcal{W}}^{(k)}$ and the sums run over finitely many exponents $\zeta \in\left[0,2|\tau|_{\mathfrak{s}}+\kappa+|\mathfrak{s}|\right)$, where $|\mathfrak{s}|=d+2=5$.

In order to obtain such bounds for our system, we can to a large extent adapt the proof of [19, Thm. 10.22]. As a matter of fact, it is immediate to see that this theorem can be extended to our case, because the involved kernels have singularities which are not worse than in the Allen-Cahn case. The aim of the discussion that follows is thus only to determine which parameters of the renormalisation group introduced above are really needed to ensure convergence.

The first step is to observe that if $\varrho_{\varepsilon}$ is the mollifyer, i.e., if $\xi^{\varepsilon}=\varrho_{\varepsilon} * \xi$, noting $K_{\varepsilon}=K * \varrho_{\varepsilon}$ we have

$$
\begin{equation*}
\left(\Pi_{\bar{z}}^{\varepsilon} \uparrow\right)(z)=\left(K_{\varepsilon} * \xi\right)(z)=\int K_{\varepsilon}\left(z-z_{1}\right) \xi\left(z_{1}\right) \mathrm{d} z_{1}=: \chi_{\varepsilon}(z), \tag{6.20}
\end{equation*}
$$

which is independent of $\bar{z}$ and belongs to the first Wiener chaos with

$$
\begin{equation*}
\left(\left(\widehat{\mathcal{W}}^{(\varepsilon ; 1)} \mathfrak{\imath}\right)(z)\right)\left(z_{1}\right)=K_{\varepsilon}\left(z-z_{1}\right) . \tag{6.21}
\end{equation*}
$$

As a consequence, we simply set

$$
\begin{equation*}
\left(\widehat{\Pi}_{\bar{z}}^{(\varepsilon)} \bullet\right)(z)=\chi_{\varepsilon}(z), \quad\left(\widehat{\Pi}_{\bar{z}} \bullet\right)(z)=\chi_{0}(z):=(K * \xi)(z) . \tag{6.22}
\end{equation*}
$$

Due to translation invariance, it will henceforth be sufficient to evaluate models in $\bar{z}=0$. As already seen in the proof of Proposition 5.1, we have

$$
\begin{equation*}
\left(\Pi_{0}^{\varepsilon} \xi\right)(z)=\left(K^{Q} * \xi^{\varepsilon}\right)(z)=\int K^{Q}\left(z-z_{1}\right) \xi^{\varepsilon}\left(z_{1}\right) \mathrm{d} z_{1}=: \chi_{\varepsilon}^{Q}(z) . \tag{6.23}
\end{equation*}
$$

This can be rewritten in the form

$$
\begin{equation*}
\chi_{\varepsilon}^{Q}(z)=\left(K_{\varepsilon}^{Q} * \xi\right)(z)=\int K_{\varepsilon}^{Q}\left(z-z_{1}\right) \xi\left(z_{1}\right) \mathrm{d} z_{1}, \tag{6.24}
\end{equation*}
$$

where we have set

$$
\begin{equation*}
K_{\varepsilon}^{Q}(t, x)=\int_{0}^{2 T} Q(s) K_{\varepsilon}(t-s, x) \mathrm{d} s . \tag{6.25}
\end{equation*}
$$

It follows that we can set

$$
\begin{equation*}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \dot{\xi}\right)(z)=\chi_{\varepsilon}^{Q}(z), \quad\left(\widehat{\Pi}_{0} \dot{\phi}\right)(z)=\chi_{0}^{Q}(z):=\left(K^{Q} * \xi\right)(z) . \tag{6.26}
\end{equation*}
$$

The second step is to consider products of terms ${ }^{\bullet}$ and $\$$. To this end, we observe that by definition of the canonical model (see (3.22)),

$$
\begin{equation*}
\left(\Pi_{0}^{\varepsilon} \mathscr{V}\right)(z)=\chi_{\varepsilon}(z)^{2}=\iint K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \xi\left(z_{1}\right) \xi\left(z_{2}\right) \mathrm{d} z_{1} \mathrm{~d} z_{2}, \tag{6.27}
\end{equation*}
$$

which belongs to the second inhomogeneous Wiener chaos. Renormalisation is needed because the right-hand side is known to diverge in the limit $\varepsilon \rightarrow 0$. However, if one introduces the Wick product

$$
\begin{equation*}
\xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right)=\xi\left(z_{1}\right) \xi\left(z_{2}\right)-\delta\left(z_{1}-z_{2}\right), \tag{6.28}
\end{equation*}
$$

then one obtains

$$
\begin{equation*}
\left(\Pi_{0}^{\varepsilon} \mathscr{V}\right)(z)=\iint K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right) \mathrm{d} z_{1} \mathrm{~d} z_{2}+\int K_{\varepsilon}\left(z-z_{1}\right)^{2} \mathrm{~d} z_{1} \tag{6.29}
\end{equation*}
$$

The first term on the right-hand side belongs to the second homogeneous Wiener chaos, with

$$
\begin{equation*}
\left(\left(\widehat{\mathcal{W}}^{(\varepsilon ; 2)} \boldsymbol{V}\right)(z)\right)\left(z_{1}, z_{2}\right)=K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right), \tag{6.30}
\end{equation*}
$$

and is known to converge as $\varepsilon \rightarrow 0$. We can thus define

$$
\begin{equation*}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \vee\right)(z)=\iint K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right) \mathrm{d} z_{1} \mathrm{~d} z_{2}=: \chi_{\varepsilon}^{\diamond 2}(z) \tag{6.31}
\end{equation*}
$$

which is indeed compatible with the second equation in (6.12), provided we set

$$
\begin{equation*}
C_{1}(\varepsilon)=\int K_{\varepsilon}\left(z-z_{1}\right)^{2} \mathrm{~d} z_{1} . \tag{6.32}
\end{equation*}
$$

The limiting model $\left(\widehat{\Pi}_{0} \vee\right)(z)$ is then naturally defined by the same expression as in (6.31), but with $K_{\varepsilon}$ replaced by $K$.

In order to deal with the term $\mathcal{\bullet}$, the relevant Wick product formula is

$$
\begin{align*}
\xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right) \diamond \xi\left(z_{3}\right)= & \xi\left(z_{1}\right) \xi\left(z_{2}\right) \xi\left(z_{3}\right) \\
& -\xi\left(z_{1}\right) \delta\left(z_{2}-z_{3}\right)-\xi\left(z_{2}\right) \delta\left(z_{3}-z_{1}\right)-\xi\left(z_{3}\right) \delta\left(z_{1}-z_{2}\right) . \tag{6.33}
\end{align*}
$$

In view of the third equation in (6.12), this is indeed compatible with the definition

$$
\begin{equation*}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \stackrel{\diamond}{)}(z)=\iiint K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) K_{\varepsilon}\left(z-z_{3}\right) \xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right) \diamond \xi\left(z_{3}\right) \mathrm{d} z_{1} \mathrm{~d} z_{2} \mathrm{~d} z_{3}\right. \tag{6.34}
\end{equation*}
$$

without the need to introduce a further renormalisation constant.
Before being able to deal with terms such as $\mathbf{\Sigma}$, we have to compute the model associated with $⿳$ Y . Applying (3.27) (in which the sum only contains the term $\ell=0$ ) and the last relation in (3.28), we obtain

$$
\begin{align*}
\left(\Pi_{0}^{\varepsilon} \circlearrowleft\right)(z) & =\int\left[K\left(z-z_{1}\right)-K\left(-z_{1}\right)\right]\left(\Pi_{0}^{\varepsilon} \vee\right)\left(z_{1}\right) \mathrm{d} z_{1} \\
& =\int\left[K\left(z-z_{1}\right)-K\left(-z_{1}\right)\right] \chi_{\varepsilon}^{\diamond 2}\left(z_{1}\right) \mathrm{d} z_{1} \\
& =:\left(K * \chi_{\varepsilon}^{\diamond 2}\right)(z)-\left(K * \chi_{\varepsilon}^{\diamond 2}\right)(0), \tag{6.35}
\end{align*}
$$

where we have used the fact that the integrals of $K(z-\cdot)$ and $K$ are equal to cancel out the term $C_{1}$. It thus follows from the last equation in (6.12) that

$$
\begin{equation*}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \mathfrak{\Omega}\right)(z)=\left[\left(K * \chi_{\varepsilon}^{\diamond 2}\right)(z)-\left(K * \chi_{\varepsilon}^{\diamond 2}\right)(0)\right] \chi_{\varepsilon}^{\diamond 2}(z)-C_{2}(\stackrel{\sim}{\bullet}) . \tag{6.36}
\end{equation*}
$$

Recall the notation $\because \mathscr{E}=\mathcal{V}$. Using (4.9) and the fact that $\int_{0}^{2 T} Q(s) \mathrm{d} s=0$, we obtain

$$
\begin{align*}
\left(\Pi_{0}^{\varepsilon} \ddot{\gamma}\right)(t, x) & =\int_{t-2 T}^{t} Q(t-s)\left(\Pi_{0}^{\varepsilon} \mathscr{V}\right)(s, x) \mathrm{d} s \\
& =\int_{t-2 T}^{t} Q(t-s)\left[\chi_{\varepsilon}^{\diamond 2}(s, x)+C_{1}(\varepsilon)\right] \mathrm{d} s \\
& =\int_{t-2 T}^{t} Q(t-s) \chi_{\varepsilon}^{\diamond 2}(s, x) \mathrm{d} s=: \widetilde{\chi}_{\varepsilon}^{\diamond 2}(t, x) . \tag{6.37}
\end{align*}
$$

Note that this term does not need to be renormalised. It follows that

$$
\begin{align*}
\left(\Pi_{0}^{\varepsilon} \widehat{\varrho}\right)(t, x) & =\int\left[K\left(z-z_{1}\right)-K\left(-z_{1}\right)\right] \widetilde{\chi}_{\varepsilon}^{\diamond 2}\left(z_{1}\right) \mathrm{d} z_{1} \\
& =\left(K * \widetilde{\chi}_{\varepsilon}^{\Omega}\right)(z)-\left(K * \widetilde{\chi}_{\varepsilon}^{\diamond 2}\right)(0) \\
& =\left(K^{Q} * \chi_{\varepsilon}^{\diamond 2}\right)(z)-\left(K^{Q} * \chi_{\varepsilon}^{\diamond 2}\right)(0), \tag{6.38}
\end{align*}
$$

and thus

$$
\begin{equation*}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \therefore \cdot\right)(z)=\left[\left(K^{Q} * \chi_{\varepsilon}^{\diamond 2}\right)(z)-\left(K^{Q} * \chi_{\varepsilon}^{\diamond 2}\right)(0)\right] \chi_{\varepsilon}^{\diamond 2}(z)-C_{2}(\therefore \stackrel{\Omega}{\bullet}) . \tag{6.39}
\end{equation*}
$$

In order to compute the Wiener chaos decomposition of expressions such as (6.36) and (6.39), we will need the identity

$$
\begin{align*}
\left(\xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right)\right)\left(\xi\left(z_{3}\right) \diamond \xi\left(z_{4}\right)\right)= & \xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right) \diamond \xi\left(z_{3}\right) \diamond \xi_{4} \\
& +\xi\left(z_{1}\right) \diamond \xi\left(z_{3}\right) \delta\left(z_{2}-z_{4}\right)+\xi\left(z_{1}\right) \diamond \xi\left(z_{4}\right) \delta\left(z_{2}-z_{3}\right) \\
& +\xi\left(z_{2}\right) \diamond \xi\left(z_{3}\right) \delta\left(z_{1}-z_{4}\right)+\xi\left(z_{2}\right) \diamond \xi\left(z_{4}\right) \delta\left(z_{1}-z_{3}\right) \\
& +\delta\left(z_{1}-z_{3}\right) \delta\left(z_{2}-z_{4}\right)+\delta\left(z_{1}-z_{4}\right) \delta\left(z_{2}-z_{3}\right), \tag{6.40}
\end{align*}
$$

which follows from [19, Lemma 10.3] and the Wick rule $I_{k}(f) \diamond I_{\ell}(g)=I_{k+\ell}(f \otimes g)$. It turns out that the important term regarding renormalisation is the contribution to the zeroth Wiener chaos, resulting from the last line in (6.40), which is given in the case $\tau=\boldsymbol{2} \cdot$ by

$$
\begin{align*}
\left(\widehat{\mathcal{W}}^{(\varepsilon ; 0)} \mathfrak{S}_{\bullet}\right)(z)= & 2 \iiint K(z-\bar{z}) K_{\varepsilon}\left(\bar{z}-z_{1}\right) K_{\varepsilon}\left(\bar{z}-z_{2}\right) K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \mathrm{d} \bar{z} \mathrm{~d} z_{1} \mathrm{~d} z_{2} \\
& +2 \iiint K(-\bar{z}) K_{\varepsilon}\left(\bar{z}-z_{1}\right) K_{\varepsilon}\left(\bar{z}-z_{2}\right) K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \mathrm{d} \bar{z} \mathrm{~d} z_{1} \mathrm{~d} z_{2} \\
& -C_{2}(\dot{\mathfrak{V}} \cdot) . \tag{6.41}
\end{align*}
$$

As shown in [19, Thm. 10.22], only the first integral on the right-hand side (which is independent of $z$ by translation invariance) diverges as $\varepsilon \rightarrow 0$, which imposes a choice for $C_{2}(\mathbb{S} \cdot)$. In order to represent this quantity and related ones appearing when dealing with the other terms, we introduce the notations

$$
\begin{align*}
& \mathcal{Q}_{0}^{\varepsilon}(z)=\int K_{\varepsilon}\left(z_{1}\right) K_{\varepsilon}\left(z_{1}-z\right) \mathrm{d} z_{1} \\
& \mathcal{Q}_{1}^{\varepsilon}(z)=\int K_{\varepsilon}\left(z_{1}\right) K_{\varepsilon}^{Q}\left(z_{1}-z\right) \mathrm{d} z_{1} \\
& \mathcal{Q}_{2}^{\varepsilon}(z)=\int K_{\varepsilon}^{Q}\left(z_{1}\right) K_{\varepsilon}^{Q}\left(z_{1}-z\right) \mathrm{d} z_{1} . \tag{6.42}
\end{align*}
$$

In particular, we have

$$
\begin{align*}
& C_{2}(\mathfrak{\mathcal { S }})=2 \int K(z) \mathcal{Q}_{0}^{\varepsilon}(z)^{2} \mathrm{~d} z \\
& C_{2}(\mathcal{Q} \bullet)=2 \int K^{Q}(z) \mathcal{Q}_{0}^{\varepsilon}(z)^{2} \mathrm{~d} z \tag{6.43}
\end{align*}
$$

Before proceeding to the convergence result, we collect a number of bounds on integrals involving the kernels $K_{\varepsilon}$ and $K_{\varepsilon}^{Q}$. In what follows we will sometimes write $K_{0}$ in place of $K$ as well as $K_{0}^{Q}$ in place of $K^{Q}$, while $|x|=\left|x_{1}\right|+\left|x_{2}\right|+\left|x_{3}\right|$ denotes the $\ell^{1}$-norm of $x \in \mathbb{R}^{3}$, so that $\|(t, x)\|_{\mathfrak{s}}=|t|^{1 / 2}+|x|$, cf (3.16). Note however that $|x|$ and the Euclidean norm $\|x\|$ are equivalent.
Lemma 6.2. The following bounds hold for all $z=(t, x) \in \mathbb{R}^{4}$ and all $\varepsilon \in[0,1]$ :

$$
\begin{equation*}
\left|K_{\varepsilon}(z)\right| \lesssim \frac{1}{\left(\|z\|_{\mathfrak{s}}+\varepsilon\right)^{3}}, \quad\left|K_{\varepsilon}^{Q}(z)\right| \lesssim \frac{1}{|x|+\varepsilon} \tag{6.44}
\end{equation*}
$$

and

$$
\begin{equation*}
\int K_{\varepsilon}(z)^{2} \mathrm{~d} z \lesssim \frac{1}{\varepsilon}, \quad \int K_{\varepsilon}(z) K_{\varepsilon}^{Q}(z) \mathrm{d} z \lesssim 1, \quad \int K_{\varepsilon}^{Q}(z)^{2} \mathrm{~d} z \lesssim 1 \tag{6.45}
\end{equation*}
$$

Furthermore, one has

$$
\begin{equation*}
\left|\mathcal{Q}_{0}^{\varepsilon}(z)\right| \lesssim \frac{1}{\|z\|_{\mathfrak{s}}+\varepsilon}, \quad\left|\mathcal{Q}_{1}^{\varepsilon}(z)\right| \lesssim 1, \quad\left|\mathcal{Q}_{2}^{\varepsilon}(z)\right| \lesssim 1 \tag{6.46}
\end{equation*}
$$

Finally, for $i, j \in\{0,1,2\}$ let

$$
\begin{equation*}
I_{i j}(\varepsilon)=\int K(z) \mathcal{Q}_{i}^{\varepsilon}(z) \mathcal{Q}_{j}^{\varepsilon}(z) \mathrm{d} z, \quad I_{i j}^{Q}(\varepsilon)=\int K^{Q}(z) \mathcal{Q}_{i}^{\varepsilon}(z) \mathcal{Q}_{j}^{\varepsilon}(z) \mathrm{d} z \tag{6.47}
\end{equation*}
$$

Then $I_{00}(\varepsilon) \asymp|\log \varepsilon|$, while all other $I_{i j}(\varepsilon)$ and all $I_{i j}^{Q}(\varepsilon)$ are bounded uniformly in $\varepsilon$.
The proof is given in Appendix A. We can now state the main result of this section, which is an adaptation of [19, Thm. 10.22].
Proposition 6.3 (Convergence of the renormalised models). There exists a random model $\widehat{Z}$, independent of the choice of mollifier $\varrho$, and a family of renormalisation transformations $M_{\varepsilon} \in \mathfrak{R}$ such that for any $\theta<-\frac{5}{2}-\alpha_{0}$, any compact set $\mathfrak{K}$ and any $\gamma<\zeta$, one has

$$
\begin{equation*}
\mathbb{E}\left\|M_{\varepsilon} Z^{\varepsilon} ; \widehat{Z}\right\|_{\gamma ; \mathfrak{K}} \lesssim \varepsilon^{\theta} \tag{6.48}
\end{equation*}
$$

uniformly over $\varepsilon \in(0,1]$. Here $\zeta$ is such that (3.26) holds for all polynomials $P$ of parabolic degree less or equal $\zeta$.

Proof: As already stated, the proof follows along the lines of [19, Thm. 10.22], with some changes due to the presence of the additional integration operator $\mathcal{E}$. Thus we will only comment on those parts of the proof that need to be adapted.

The result will follow if we are able to choose the constants $C_{1}(\tau)$ and $C_{2}(\tau)$ in (6.10) for each $\varepsilon \in(0,1]$ in such a way that the bounds (6.15) hold for all $\tau \in \mathcal{F}_{-}$; equivalently, the bounds (6.19) should hold for all these $\tau$. The limiting model $\widehat{Z}$ will be constructed step by step during this process.

The first step is to set

$$
\left(\widehat{\Pi}_{0} \Xi\right)(z)=\xi(z),
$$

and the fact that the required bounds are satisfied is a consequence of Proposition 9.5 in [19] (see also Lemma 10.2).

In the case of $\tau=\boldsymbol{}$, the limiting model is defined by (6.22), and the conclusion follows as shown in [19, Thm. 10.22] by applying the bounds on convolutions in [19, Lemmas 10.14 and 10.17]. We thus turn to the case $\tau=\boldsymbol{\beta}$. By (6.26), we have

$$
\begin{aligned}
\left(\left(\widehat{\mathcal{W}}^{(\varepsilon ; 1)} \dot{\xi}\right)(z)\right)\left(z_{1}\right) & =K_{\varepsilon}^{Q}\left(z-z_{1}\right), \\
\left(\left(\widehat{\mathcal{W}}^{(1)} ;\right)(z)\right)\left(z_{1}\right) & =K^{Q}\left(z-z_{1}\right) .
\end{aligned}
$$

The first bound in (6.19) then follows from the fact that

$$
\left\langle\left(\widehat{\mathcal{W}}^{(1)} \dot{\varphi}\right)(z),\left(\widehat{\mathcal{W}}^{(1)} \dot{)}\right)(\bar{z})\right\rangle=\int K^{Q}\left(z-z_{1}\right) K^{Q}\left(\bar{z}-z_{1}\right) \mathrm{d} z_{1}=\mathcal{Q}_{2}^{0}(z-\bar{z})
$$

(via the change of variables $z_{1} \mapsto z-z_{1}$ ). The right-hand side being bounded as shown in Lemma 6.2, the bound (6.19) is satisfied for $\zeta=0$ provided $\bar{\kappa}$ is small enough. In order to prove the second bound in (6.19), we use the fact that $K^{Q}$ is singular of order -1 in the sense of [19, Def. 10.12]. It thus follows from Lemma 10.17 in [19] that

$$
\left|K^{Q}(z)-K_{\varepsilon}^{Q}(z)\right| \lesssim \frac{\varepsilon^{\theta}}{\|z\|_{\mathfrak{s}}^{1+\theta}}
$$

for any $\theta>0$. As a consequence, we have the bound

$$
\left\langle\left(\delta \widehat{\mathcal{W}}^{(\varepsilon ; 1)} \mathfrak{\xi}\right)(z),\left(\delta \widehat{\mathcal{W}}^{(\varepsilon ; 1)} \dot{\varphi}\right)(\bar{z})\right\rangle \lesssim \varepsilon^{2 \theta} \int \frac{1}{\left\|z_{1}\right\|_{\mathfrak{s}}^{1+\theta}\left\|z-\bar{z}-z_{1}\right\|_{\mathfrak{s}}^{1+\theta}} \mathrm{d} z_{1} .
$$

The integral is bounded by a constant times $\|z-\bar{z}\|_{\mathfrak{s}}^{-\bar{\zeta}}$ for any $\bar{\zeta}>0$ provided $\theta<\frac{3}{2}$. This can be seen by noting that $z \mapsto\|z\|_{\mathfrak{s}}^{-(1+\theta)}$ is singular of order $-1-\theta$, and thus a fortiori singular of order $-\frac{5}{2}$, so that the claimed bound follows from [19, Lemma 10.14]. The second bound in (6.19) is thus fulfilled with $\zeta=0$ provided $\bar{\kappa}$ and $\theta$ are small enough.

Next we consider the term $\mathcal{V}^{\text {. Then a decomposition analogous to the one given }}$ in (6.29) shows that we have

$$
\begin{aligned}
\left(\left(\widehat{\mathcal{W}}^{(\varepsilon ; 2)} \mathfrak{\bullet}\right)(z)\right)\left(z_{1}, z_{2}\right) & =K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}^{Q}\left(z-z_{2}\right) \\
\left(\widehat{\mathcal{W}}^{(\varepsilon ; 0)} \boldsymbol{\beta}\right)(z) & =\int K_{\varepsilon}\left(-z_{1}\right) K_{\varepsilon}^{Q}\left(-z_{1}\right) \mathrm{d} z_{1}-C_{1}^{\prime}(\varepsilon) .
\end{aligned}
$$

Note that the second line is independent of $z$, as a consequence of translation invariance, so that we can drop the argument $z$ from the notation. The integral in the second line
is equal to $\mathcal{Q}_{1}^{\varepsilon}(0)$, which is bounded uniformly in $\varepsilon$ by Lemma 6.2 . We may thus simply choose $C_{1}^{\prime}(\varepsilon)=0$, which amounts to setting

$$
\begin{aligned}
\left(\widehat{\Pi}_{0}^{(\varepsilon)} \mathfrak{\vartheta}^{\bullet}\right)(z) & =\iint K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}^{Q}\left(z-z_{2}\right) \xi\left(z_{1}\right) \diamond \xi\left(z_{2}\right) \mathrm{d} z_{1} \mathrm{~d} z_{2}+\widehat{\mathcal{W}}^{(\varepsilon ; 0)} \\
& =\chi_{\varepsilon}(z) \chi_{\varepsilon}^{Q}(z)
\end{aligned}
$$

Of course, another possibility would be to set $C_{1}^{\prime}(\varepsilon)=\widehat{\mathcal{W}}^{(\varepsilon ; 0)} \mathcal{\beta}^{\circ}$, which would result in $\widehat{\Pi}_{0} \mathcal{B}^{\circ}$ belonging to the second homogeneous Wiener chaos. Since $\mathcal{Q}_{1}^{\varepsilon}(0)$ is bounded, however, the choice of $C_{1}^{\prime}(\varepsilon)$ does not really matter, since it only results in a shift of the random variable defining $\widehat{Z}$. Whatever the choice, the bounds (6.19) can be checked by similar arguments as in the previous case.
 the same way. It thus remains to consider the symbols in the families of $\mathbb{V}^{\circ},{ }^{\circ}{ }^{\bullet \bullet}$ and $\bullet^{\bullet \bullet}$. As shown in the proof of [19, Thm. 10.22], the symbol ${ }^{\circ} \bullet \cdot$ already fulfils the bounds (6.19) as a consequence of the choice of $C_{1}(\varepsilon)$. Since the kernel $K^{Q}$ is less singular than $K$, the same conclusion holds for all symbols of the same family.

The symbol $\cup \cdot \sim$ requires the introduction of a further renormalisation constant $C_{2}(\cdot \cup)$ which has order $|\log \varepsilon|$. We thus have to check what happens to the other symbols in the same family. Consider for instance the case of $\mathfrak{\nu}^{\circ}$. A computation similar to the one made in (6.35) and (6.36) yields the expression

$$
\left(\widehat{\Pi}_{0}^{\varepsilon} \bullet \stackrel{\sim}{\bullet}\right)(z)=\left[\left(K *\left(\chi_{\varepsilon} \diamond \chi_{\varepsilon}^{Q}\right)\right)(z)-\left(K *\left(\chi_{\varepsilon} \diamond \chi_{\varepsilon}^{Q}\right)\right)(0)\right] \chi_{\varepsilon}^{\diamond 2}(z)-C_{2}(\bullet \bullet \cdot)
$$

Writing out the integrals and applying (6.40), we can decompose this expression into a sum of terms in the Wiener chaoses of order 0,2 and 4 . In particular, the term in the 0th Wiener chaos is given by

$$
\begin{aligned}
\left(\widehat{\mathcal{W}}^{(\varepsilon ; 0)} \mathfrak{V}^{\circ}\right)(z)= & 2 \iiint K(z-\bar{z}) K_{\varepsilon}\left(\bar{z}-z_{1}\right) K_{\varepsilon}^{Q}\left(\bar{z}-z_{2}\right) K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \mathrm{d} \bar{z} \mathrm{~d} z_{1} \mathrm{~d} z_{2} \\
& +2 \iiint K(-\bar{z}) K_{\varepsilon}\left(\bar{z}-z_{1}\right) K_{\varepsilon}^{Q}\left(\bar{z}-z_{2}\right) K_{\varepsilon}\left(z-z_{1}\right) K_{\varepsilon}\left(z-z_{2}\right) \mathrm{d} \bar{z} \mathrm{~d} z_{1} \mathrm{~d} z_{2} \\
& -C_{2}\left(\mathfrak{N}^{\circ}\right)
\end{aligned}
$$

The first integral on the right-hand side does not depend on $z$, as can be seen by the change of variables $\left(\bar{z}, z_{1}, z_{2}\right) \mapsto\left(\bar{z}+z, z_{1}+z, z_{2}+z\right)$. As a consequence, it is equal to $2 I_{01}(\varepsilon)$, and is thus uniformly bounded in $\varepsilon$ as shown in Lemma 6.2. The second integral is also bounded, so that we may choose $C_{2}\left(\stackrel{\nu}{*}^{\circ}\right)=0$ (or any other finite value). The terms in the other two Wiener chaoses are already bounded in the case of ${ }^{\circ} \cdot$, so they remain bounded in the present case. It is now quite obvious that the other symbols of the family do not need to be renormalised either, as the relevant integrals are obtained form the above ones by substituting the appropriate kernels $K$ by $K^{Q}$.

Finally, we have to deal with symbols in the family of ${ }^{\circ \circ} \cdot$. However, it is shown in the proof of [19, Thm. 10.22] that the constant $C_{2}\left(\sum_{0}\right)$ suffices to renormalise ${ }^{\circ} \cdot{ }^{\circ} \cdot($ as well, so the situation is completely analogous for the other symbols of the family.

### 6.3 Computation of the renormalised equations

In order to complete the proof of the main results, it remains to derive the classical equation with mollified noise that corresponds to the renormalised models. In other words, for a
given $M=M_{\varepsilon} \in \mathfrak{R}$, we have to characterise the solution map $\overline{\mathcal{S}}_{M}^{L}$ satisfying

$$
\begin{equation*}
\mathcal{R}^{M} \mathcal{S}_{M}^{L}\left(u_{0}, v_{0}, M Z^{\varepsilon}\right)=\overline{\mathcal{S}}_{M}^{L}\left(u_{0}, v_{0}, \xi^{\varepsilon}\right), \tag{6.49}
\end{equation*}
$$

where $\mathcal{R}^{M}$ is the reconstruction operator of the model $\Pi^{M}\left(=\widehat{\Pi}^{(\varepsilon)}\right)$ and $\mathcal{S}_{M}^{L}$ is the map giving the fixed point of the renormalised equation in terms of initial data and forcing term. Relation (6.49) can also be represented as the following commutative diagram:


In order to determine the classical equation, we will need the following algebraic result.
Lemma 6.4. For all $\tau \in \mathcal{F}_{0}$, one has

$$
\begin{equation*}
\left(\Pi_{z}^{M} \tau\right)(z)=\left(\Pi_{z} M \tau\right)(z) \tag{6.50}
\end{equation*}
$$

Proof: The claim can be checked directly by comparing the expressions given in (6.11) and (6.12). More generally, the first relation in (6.4) shows that (6.50) holds for all $\tau$ such that $\Delta^{M}(\tau)=M \tau \otimes \mathbf{1}$, which happens to be the case for all symbols excepts those of the families $\mathbb{M}, \mathscr{P}$. and $\mathbb{*}$. In the exceptional cases, the relation has to be checked by an explicit computation; in general it holds only when the model is evaluated at the base point, i.e. we may have $\left(\Pi_{z}^{M} \tau\right)(\bar{z}) \neq\left(\Pi_{z} M \tau\right)(\bar{z})$ if $z \neq \bar{z}$.

As a consequence of (6.50), when $f$ is a function, as is the case for $\varepsilon>0$, we have

$$
\begin{equation*}
\left(\mathcal{R}^{M} f\right)(z)=\left(\Pi_{z}^{M} f\right)(z)=\left(\Pi_{z} M f\right)(z)=(\mathcal{R} M f)(z) \tag{6.51}
\end{equation*}
$$

This directly implies the following result, which is the equivalent of Proposition 5.1 for the renormalised equation.

Proposition 6.5 (Renormalised equation). Assume that $\varepsilon>0$ and that (5.3) admits a fixed point $(U, V)$ where $U, V \in \mathcal{D}^{\gamma, \eta}$ for some $\gamma, \eta \in \mathbb{R}$. Then $(\hat{u}, \hat{v})=\left(\mathcal{R}^{M} U, \mathcal{R}^{M} V\right)$ satisfies the classical SPDE

$$
\begin{align*}
\partial_{t} \hat{u} & =\Delta \hat{u}+\widehat{F}(\hat{u}, \hat{v})+\xi^{\varepsilon}, \\
\partial_{t} \hat{v} & =A_{1} \hat{u}+A_{2} \hat{v}, \tag{6.52}
\end{align*}
$$

provided $\widehat{F}$ is such that

$$
\begin{equation*}
\widehat{F}(M U, M V)=M F(U, V)+\varrho(U, V), \tag{6.53}
\end{equation*}
$$

with @ having only components of strictly positive homogeneity.

Proof: Applying $M$ to the expansions of $U$ and $V$ given in (5.7), we see that $M U$ differs from $U$ and $M V$ differs from $V$ only by components of homogeneity $\frac{3}{2}-\kappa$ at least. Applying $\mathcal{R}$ to $M U$ and $M V$ and using (6.51), we obtain that as in the proof of Proposition 5.1

$$
\begin{aligned}
& \hat{u}(z)=\chi_{\varepsilon}(z)+\varphi(z) \\
& \hat{v}(z)=\chi_{\varepsilon}^{Q}(z)+\psi(z)
\end{aligned}
$$

On the other hand, applying $M$ to the right-hand side of (5.5) and using the fact that $M$ commutes with $\mathcal{I}$ and $\mathcal{E}$ and leaves polynomials invariant, we get

$$
M U=\mathcal{I} \boldsymbol{R}^{+}[\Xi+M F(U, V)]+\widehat{R}_{u}(M U, M V)+G u_{0}
$$

and a similar equation for $M V$. Applying the reconstruction operator, we thus obtain

$$
\hat{u}(z)=G * \boldsymbol{R}^{+}\left[\xi^{\varepsilon}(z)+(\mathcal{R} M F(U, V))(z)\right]+G u_{0}(z)
$$

and a similar expression for $\hat{v}(z)$. The claim then follows from the fact that (6.53) implies

$$
(\mathcal{R} M F(U, V))(z)=(\mathcal{R} \widehat{F}(M U, M V))(z)=\widehat{F}(\hat{u}(z), \hat{v}(z))
$$

where the last equality follows from the same argument as the one used in the proof of Proposition 5.1.

Finally, we compute the expression of $\widehat{F}$. Even though we have seen that only two parameters are needed to renormalise the equation, we give here the result for an arbitrary value of the constants $C_{1}(\tau)$ and $C_{2}(\tau)$ in (6.10). Indeed, the computation is not really more difficult in the general case, and the result can be used to determine the constant $C_{3}$ in (2.6) and (2.7).

Proposition 6.6 (Expression for the renormalised $F$ ). Let $F$ be the cubic polynomial given in (5.6) and let $M$ be the renormalisation transformation defined in (6.10). Then

$$
\begin{equation*}
\widehat{F}(u, v)=F(u, v)-c_{0}(\varepsilon)-c_{1}(\varepsilon) u-c_{2}(\varepsilon) v, \tag{6.54}
\end{equation*}
$$

where the coefficients are given by

$$
\left(\begin{array}{c}
c_{0}(\varepsilon)  \tag{6.55}\\
c_{1}(\varepsilon) \\
c_{2}(\varepsilon)
\end{array}\right)=\left(\begin{array}{ccc}
\beta_{1} & \beta_{2} & \beta_{3} \\
3 \gamma_{1} & 2 \gamma_{2} & \gamma_{3} \\
\gamma_{2} & 2 \gamma_{3} & 3 \gamma_{4}
\end{array}\right)\left(\begin{array}{c}
C_{1}(\varepsilon)+C_{2}(\varepsilon) \\
C_{1}^{\prime}(\varepsilon)+C_{2}^{\prime}(\varepsilon) \\
C_{1}^{\prime \prime}(\varepsilon)+C_{2}^{\prime \prime}(\varepsilon)
\end{array}\right) .
$$

with $C_{1}(\varepsilon)=C_{1}(\mathcal{V}), C_{1}^{\prime}(\varepsilon)=C_{1}(\mathfrak{\vartheta}), C_{1}^{\prime \prime}(\varepsilon)=C_{1}(\stackrel{\otimes}{*})$ and

The constant $C_{2}^{\prime}(\varepsilon)$ has a similar definition with $\mathfrak{\sim}$. replaced by $\mathfrak{P}$. and so on, while $C_{2}^{\prime \prime}(\varepsilon)$ has a similar definition with 2 replaced by $\overbrace{}^{2}$. and so on.

Proof: The fact that $(U, V)$ differs from $(M U, M V)$ only by terms of order $\frac{3}{2}-\kappa$ implies

$$
\widehat{F}(U, V)=M F(U, V)+\varrho^{\prime}(U, V)
$$

where all components of $\varrho^{\prime}(U, V)$ have strictly positive homogeneity. Thus we only have to compute the difference $M F(U, V)-F(U, V)$, which reduces to determining the terms of
$F(U, V)$ which are modified by the renormalisation map, or, more precisely, which terms of homogeneity less or equal 0 appear when applying $M$.

The computation is straightforward, though somewhat lengthy. In practice, it suffices to expand the quadratic terms up to order $\frac{1}{2}-\kappa$, and cubic terms up to order $0-2 \kappa$. The quadratic terms yield the additional contribution

$$
-\left[\beta_{1} C_{1}(\varepsilon)+\beta_{2} C_{1}^{\prime}(\varepsilon)+\beta_{3} C_{1}^{\prime \prime}(\varepsilon)\right] \mathbf{1}
$$

It remains to treat terms of homogeneity $-\frac{1}{2}-\kappa$ and $0-2 \kappa$ in the cubic parts. Consider first those of homogeneity $0-2 \kappa$. Only 18 such terms need to be considered, which all belong to the family of $\mathbb{\sim}$. Their renormalisation yields a term

$$
-\left[b_{1} C_{2}(\varepsilon)+b_{2} C_{2}^{\prime}(\varepsilon)+b_{3} C_{2}^{\prime \prime}(\varepsilon)\right] \mathbf{1}
$$

There are 24 terms of homogeneity $-\frac{1}{2}-\kappa$ that are affected by the renormalisation map, which are all in the family of $\boldsymbol{*}$. Their renormalisation yields

$$
-\left[3 \gamma_{1} C_{2}(\varepsilon)+2 \gamma_{2} C_{2}^{\prime}(\varepsilon)+\gamma_{3} C_{2}^{\prime \prime}(\varepsilon)\right] \bullet-\left[\gamma_{2} C_{2}(\varepsilon)+2 \gamma_{3} C_{2}^{\prime}(\varepsilon)+3 \gamma_{3} C_{2}^{\prime \prime}(\varepsilon)\right] \risingdotseq .
$$

Replacing the $b_{i}$ by their expressions (5.8), and adding the last two expressions, we see that we can always factor out the expressions $\varphi \mathbf{1}+\boldsymbol{i}$ and $\psi \mathbf{1}+\boldsymbol{i}$. The result follows.

## 7 Proof of the main results

Since Theorem 2.1 is a particular case of Theorem 2.2, we proceed directly to proving the latter result.

Proof of Theorem 2.2. Consider first the situation on a fixed, sufficiently small time interval $[0, T]$. Fix $\varepsilon>0$ and let

$$
\begin{aligned}
\left(U^{\varepsilon}, V^{\varepsilon}\right) & =\mathcal{S}_{T}\left(u_{0}, v_{0}, M_{\varepsilon} Z^{\varepsilon}\right) \\
(U, V) & =\mathcal{S}_{T}\left(u_{0}, v_{0}, \widehat{Z}\right)
\end{aligned}
$$

where $\mathcal{S}_{T}$ is the solution map given in Proposition 5.4 (we have suppressed the dependence on the $W_{i}$, which are the same for both fixed points). For any $\delta>0$, the bound (5.17), Markov's inequality and Proposition 6.3 imply

$$
\mathbb{P}\left\{\left\|U^{\varepsilon} ; U\right\|_{\gamma, \eta ; T}+\left\|V^{\varepsilon} ; V\right\|_{\gamma, \eta ; T}>C_{0} \delta\right\} \leqslant \frac{1}{\delta} \mathbb{E}\left\|M_{\varepsilon} Z^{\varepsilon} ; \widehat{Z}\right\|_{2 \gamma+\alpha ; O} \lesssim \frac{\varepsilon^{\theta}}{\delta}
$$

It follows that $\left(U^{\varepsilon}, V^{\varepsilon}\right)$ converges to $(U, V)$ in probability as $\varepsilon \rightarrow 0$.
Next, let $(u, v)=(\mathcal{R} U, \mathcal{R} V)$ and $\left(u^{\varepsilon}, v^{\varepsilon}\right)=\left(\mathcal{R} U^{\varepsilon}, \mathcal{R} V^{\varepsilon}\right)$ (where we do not indicate the model-dependence of the reconstruction operators $\mathcal{R}$ ). The fact that ( $u^{\varepsilon}, v^{\varepsilon}$ ) converges to $(u, v)$ in probability as $\varepsilon \rightarrow 0$ is a consequence of the bound [19, (3.4)] in the reconstruction theorem, combined with the definitions (3.20) of $\|\Pi\|_{\gamma ; \mathfrak{\Re}}$ and (4.28) of $\|\cdot\| \|_{\gamma, \eta ; \mathfrak{\xi}}$.

The convergence result can be extended from a fixed time interval $[0, T]$ to the random time interval up to the first exit from a ball of radius $L$ in exactly the same way as in [19, Section 7.3]. In particular, the required continuous dependence on the model is proved in [19, Cor. 7.12].

Proposition 6.5 shows that $\left(u^{\varepsilon}, v^{\varepsilon}\right)$ satisfies the system (2.11), where the explicit expression of $\widehat{F}$ is a consequence of Proposition 6.6. The expressions (2.5) and (2.12) for the
renormalisation constants in dimension $d=3$ follow from (6.55), where the only nonzero terms are $C_{1}(\varepsilon)$, defined in (6.32), and $C_{2}(\varepsilon)$, defined in (6.43).

In the case $d=2$, Table 1 shows that only the terms $\boldsymbol{\bullet}, \boldsymbol{v}$ and $\boldsymbol{\vartheta}$ need to be renormalised, to that the expression (2.7) for the renormalisation constant follows by taking $C_{2}(\varepsilon)=0$ in (6.55).

Proof of Theorem 2.3. The only difference with the previous case is that we have to take into account the presence of the $n$ different operators $\mathcal{E}_{i}, i=1, \ldots, n$, introduced in Remark 4.1. This means that the fixed-point equation for $V$ in (5.3) is replaced by $n$ different equations of the form

$$
V_{i}=\mathcal{I E} \mathcal{E}_{i} \boldsymbol{R}^{+}[\Xi+F(U, V)]+\widehat{R}_{v, i}(U, V)+G^{Q_{i}} u_{0}+\widehat{Q}_{i} v_{0}, \quad i=1, \ldots, n
$$

Since $n$ is finite, it is straightforward to check that the proofs of Proposition 5.1 and Proposition 5.4 carry over to this situation.

It remains to check the expressions (2.17) for the renormalisation constants. To compute the diverging part of $\widehat{F}(U, V)-F(U, V)$, it is in fact sufficient to determine the coefficients of $\mathcal{V}, \dot{\sim}$ and $\dot{\sim}$ in $F(U, V)$, since these are the only ones that need to be renormalised. These can only result from the terms $\beta_{1} U^{2}+\gamma_{1} U^{3}$ in $F(U, V)$, since all other terms will contain at least one of the $\mathcal{E}_{i}$. It is now sufficient to check that $U$ has the same expansion as in (5.7), except that the first relation in (5.8) becomes

$$
b_{1}=\beta_{1}+3 \varphi \gamma_{1}+\sum_{i=1}^{n} \psi_{i} \gamma_{2, i},
$$

where $\psi_{i}$ is the coefficient of $\mathbf{1}$ in $V_{i}$. The result then follows as in the proof of Proposition 6.6 , taking into account the modified value of $b_{1}$.

## A Proof of Lemma 6.2

It follows from the definition of $K$ that it satisfies $|K(z)| \lesssim\|z\|_{\mathfrak{s}}^{-3}$. Computing the derivatives of $K$, one sees that it is singular of order -3 in the sense of [19, Def. 10.12]. Thus the first bound in (6.44) follows from [19, Lemma 10.17].

To show the second bound in (6.44), we use the facts that $Q(t, x)$ is bounded and supported on $\{t>0\}, K_{\varepsilon}(t, x)$ is supported on $\left\{t>-\varepsilon^{2}\right\}$, and thus

$$
K_{\varepsilon}^{Q}(t, x) \asymp \int_{-\varepsilon^{2}}^{t} K_{\varepsilon}(s, x) \mathrm{d} s \lesssim \int_{-\varepsilon^{2}}^{t} \frac{1}{\left(|s|^{1 / 2}+|x|+\varepsilon\right)^{3}} \mathrm{~d} s \lesssim \frac{1}{|x|+\varepsilon}
$$

(recall the notation $f \asymp g$ indicating that we have both $f \lesssim g$ and $g \lesssim f$ ). The first bound in (6.45) follows from

$$
\int K_{\varepsilon}(z)^{2} \mathrm{~d} z \lesssim \int \frac{1}{\left(\left\|z_{1}\right\|_{\mathfrak{s}}+\varepsilon\right)^{6}} \mathrm{~d} z_{1} \asymp \int_{0}^{1} \int_{0}^{1} \frac{r^{2} \mathrm{~d} r}{\left(t^{1 / 2}+r+\varepsilon\right)^{6}} \mathrm{~d} t \asymp \frac{1}{\varepsilon} .
$$

To obtain the second bound in (6.45), we note that

$$
\int K_{\varepsilon}(t, x) K_{\varepsilon}^{Q}(t, x) \mathrm{d} x \lesssim \int_{0}^{1} \frac{r^{2} \mathrm{~d} r}{\left(r+|t|^{1 / 2}+\varepsilon\right)^{3}(r+\varepsilon)} \asymp \frac{1}{|t|^{1 / 2}+\varepsilon}
$$

so that

$$
\iint K_{\varepsilon}(t, x) K_{\varepsilon}^{Q}(t, x) \mathrm{d} x \mathrm{~d} t \lesssim \int_{-\varepsilon^{2}}^{1} \frac{1}{|t|^{1 / 2}+\varepsilon} \mathrm{d} t \lesssim 1
$$

The third bound in (6.46) follows from the fact that

$$
\int K_{\varepsilon}^{Q}(t, x)^{2} \mathrm{~d} x \lesssim \int_{0}^{1} \frac{r^{2} \mathrm{~d} r}{(r+\varepsilon)^{2}} \asymp 1
$$

To obtain the bound on $\mathcal{Q}_{0}$, we distinguish the regimes $\|z\|_{\mathfrak{s}} \leqslant \varepsilon$ and $\|z\|>\varepsilon$. If $\|z\|_{\mathfrak{s}} \leqslant \varepsilon$, the same computation as the one made for the integral of $K_{\varepsilon}(z)^{2}$ shows that $\left|\mathcal{Q}_{0}(z)\right| \lesssim 1 / \varepsilon$, while for $\|z\|_{\mathfrak{s}}>\varepsilon$, we can use [19, Lemma 10.14] to obtain

$$
\left|\mathcal{Q}_{0}(z)\right| \lesssim \int \frac{1}{\left\|z_{1}\right\|_{\mathfrak{s}}^{3}} \frac{1}{\left\|z-z_{1}\right\|_{\mathfrak{s}}^{3}} \mathrm{~d} z_{1} \lesssim \frac{1}{\|z\|_{\mathfrak{s}}}
$$

Using spherical coordinates, one can show that for all $a, b \geqslant 0$ and all $n \geqslant 1$, one has

$$
I_{n}:=\int \frac{1}{\left(\left|x_{1}\right|+a\right)^{n}} \frac{1}{\left|x-x_{1}\right|+b} \mathrm{~d} x_{1} \asymp \int_{0}^{1} \frac{r^{2} \mathrm{~d} r}{(r+a)^{n}(r+|x|+b)} .
$$

From this we easily get that whenever $a \geqslant b \geqslant 0$,

$$
I_{1} \lesssim 1+|x|+a^{2}, \quad I_{3} \lesssim \frac{1+|\log a|}{|x|+a}
$$

and thus

$$
\begin{aligned}
& \left|\mathcal{Q}_{1}^{\varepsilon}(t, x)\right| \lesssim \int_{0}^{1} \frac{1+\left|\log \left(t^{1 / 2}+\varepsilon\right)\right|}{|x|+t^{1 / 2}+\varepsilon} \mathrm{d} t \lesssim 1 \\
& \left|\mathcal{Q}_{2}^{\varepsilon}(t, x)\right| \lesssim 1+|x|+\varepsilon^{2} \lesssim 1
\end{aligned}
$$

It remains to bound the terms $I_{i j}(\varepsilon)$ and $I_{i j}^{Q}(\varepsilon)$. It is in fact sufficient to check that

$$
\begin{aligned}
& I_{00}(\varepsilon) \lesssim \int_{0}^{1} \int_{0}^{1} \frac{r^{2} \mathrm{~d} r}{\left(t^{1 / 2}+r\right)^{3}\left(t^{1 / 2}+r+\varepsilon\right)^{2}} \mathrm{~d} t \lesssim|\log \varepsilon| \\
& I_{01}(\varepsilon) \\
& \lesssim \int_{0}^{1} \int_{0}^{1} \frac{r^{2} \mathrm{~d} r}{\left(t^{1 / 2}+r\right)^{3}\left(t^{1 / 2}+r+\varepsilon\right)} \mathrm{d} t \lesssim 1 \\
& I_{00}^{Q}(\varepsilon)
\end{aligned} \int_{0}^{1} \int_{0}^{1} \frac{r \mathrm{~d} r}{\left(t^{1 / 2}+r+\varepsilon\right)^{2}} \mathrm{~d} t \lesssim 1, ~ \$
$$

since the other terms can be bounded above by either $I_{01}(\varepsilon)$ or $I_{00}^{Q}(\varepsilon)$. Elementary computations show that these bounds indeed hold true.

## References

[1] H. Alzubaidi and T. Shardlow. Numerical simulations of SDEs and SPDEs. In C. Laing and T. Shardlow, editors, Stochastic Methods in Neuroscience, pages 344-366. OUP, 2009.
[2] D. Barkley. A model for fast computer simulation of waves in excitable media. Physica D, 49:61-70, 1991.
[3] F. Barret. Sharp asymptotics of metastable transition times for one dimensional SPDEs. Ann. Inst. Henri Poincaré Probab. Stat., 51(1):129-166, 2015.
[4] N. Berglund and B. Gentz. Sharp estimates for metastable lifetimes in parabolic SPDEs: Kramers' law and beyond. Electron. J. Probab., 18:no. 24, 58, 2013.
[5] N. Berglund, B. Gentz, and C. Kuehn. From random Poincaré maps to stochastic mixed-mode-oscillation patterns. J. Dyn. Diff. Equat., 27(1):83-136, 2015.
[6] P. C. Bressloff. Spatiotemporal dynamics of continuum neural fields. J. Phys. A, 45(3):033001, 109, 2012.
[7] P. C. Bressloff. Waves in neural media. From single neurons to neural fields. Lecture Notes on Mathematical Modelling in the Life Sciences. Springer, New York, 2014.
[8] R. Cantrell and C. Cosner. Spatial Ecology via Reaction-Diffusion Equations. John Wiley \& Sons, 2004.
[9] S. Cerrai. Elliptic and parabolic equations in $\mathbf{R}^{n}$ with coefficients having polynomial growth. Comm. Partial Differential Equations, 21(1-2):281-317, 1996.
[10] G. Da Prato and A. Debussche. Strong solutions to the stochastic quantization equations. Ann. Probab., 31(4):1900-1916, 2003.
[11] G. Da Prato and J. Zabczyk. Stochastic equations in infinite dimensions, volume 44 of Encyclopedia of Mathematics and its Applications. Cambridge University Press, Cambridge, 1992.
[12] W. G. Faris and G. Jona-Lasinio. Large fluctuations for a nonlinear heat equation with noise. J. Phys. A, 15(10):3025-3055, 1982.
[13] J. Feldman. The $\lambda \varphi_{3}^{4}$ field theory in a finite volume. Comm. Math. Phys., 37:93-120, 1974.
[14] J. S. Feldman and K. Osterwalder. The Wightman axioms and the mass gap for weakly coupled $\left(\Phi^{4}\right)_{3}$ quantum field theories. Ann. Physics, 97(1):80-135, 1976.
[15] R. FitzHugh. Impulses and physiological states in models of nerve membrane. Biophys. J., 1:445-466, 1961.
[16] A. Gierer and H. Meinhardt. A theory of biological pattern formation. Kybernetic, 12:30-39, 1972.
[17] M. Hairer. Solving the KPZ equation. Ann. Math., 178(2):559-664, 2013.
[18] M. Hairer. Introduction to regularity structures. Braz. Jour. Prob. Stat., 2014. accepted, to appear.
[19] M. Hairer. A theory of regularity structures. Inventiones Mathematicae, 198(2):269-504, 2014.
[20] M. Hairer and C. Labbé. A simple construction of the continuum parabolic Anderson model on $\mathbb{R}^{2}$. Preprint arXiv/1501.00692, 2015.
[21] M. Hairer and E. Pardoux. A Wong-Zakai theorem for stochastic PDEs. Preprint arXiv/1409.3138, 2014.
[22] M. Hairer and H. Shen. The dynamical sine-Gordon model. Preprint arXiv/1409.5724, 2014.
[23] M. Hairer and H. Weber. Large deviations for white-noise driven, nonlinear stochastic PDEs in two and three dimensions. Preprint arXiv/1404.5863, 2014.
[24] A. L. Hodgkin and A. F. Huxley. A quantitative description of ion currents and its applications to conduction and excitation in nerve membranes. J. Physiol. (Lond.), 117:500-544, 1952.
[25] E. Keller and S. Segel. Initiation of slime mold aggregation viewed as an instability. J. Theor. Biol., 26:399-415, 1970.
[26] M. Koper. Bifurcations of mixed-mode oscillations in a three-variable autonomous Van der Pol-Duffing model with a cross-shaped phase diagram. Physica D, 80:72-94, 1995.
[27] J. Krüger and W. Stannat. Front propagation in stochastic neural fields: a rigorous mathematical framework. SIAM J. Appl. Dyn. Syst., 13(3):1293-1310, 2014.
[28] C. Kuehn. Multiple Time Scale Dynamics. Springer, 2015.
[29] C. Kuehn and M. Riedler. Large deviations for nonlocal stochastic neural fields. J. Math. Neurosci., 4(1):1-33, 2014.
[30] A. Kupiainen. Renormalization group and stochastic PDE's. arXiv:1410.3094, 2014.
[31] G. Lord, C. Powell, and T. Shardlow. An Introduction to Computational Stochastic PDEs. CUP, 2014.
[32] C. Morris and H. Lecar. Voltage oscillations in the barnacle giant muscle fiber. Biophys. J., pages 193-213, 1981.
[33] J. S. Nagumo, S. Arimoto, and S. Yoshizawa. An active pulse transmission line simulating nerve axon. Proc. IRE, 50:2061-2070, 1962.
[34] M. Sauer and W. Stannat. Lattice approximation for stochastic reaction diffusion equations with one-sided Lipschitz condition. Math. Comp., 84:743-766, 2015.
[35] M. E. Sweedler. Hopf algebras. Mathematics Lecture Note Series. W. A. Benjamin, Inc., New York, 1969.
[36] H. Tuckwell. Analytical and simulation results for the stochastic spatial Fitzhugh-Nagumo model neuron. Neural Computation, 20(12):3003-3033, 2008.
[37] H. Tuckwell. Stochastic partial differential equations in neurobiology: linear and nonlinear models for spiking neurons. In Stochastic Biomathematical Models, pages 149-173. Springer, 2013.
[38] H. C. Tuckwell. Stochastic Processes in the Neurosciences. SIAM, Philadelphia, PA, 1989.
[39] H. C. Tuckwell and J. Jost. The effect of various spatial distributions of weak noise on rhythmic spiking. J. Comput. Neurosci., 30(2):361-371, 2011.
[40] J. Tyson and J. Keener. Singular perturbation theory of traveling waves in excitable media (a review). Physica D, 32(3):327-361, 1988.

## Contents

1 Introduction ..... 1
2 Results ..... 4
2.1 The standard FitzHugh-Nagumo equation ..... 4
2.2 General cubic nonlinearities ..... 5
2.3 Vectorial gating variables ..... 6
3 Regularity structure for the Allen-Cahn equation ..... 7
3.1 Regularity structures ..... 8
3.2 Models and reconstruction theorem ..... 10
3.3 Lifting the convolution maps ..... 14
4 Extension of the regularity structure ..... 15
4.1 Extension theorem ..... 16
4.2 Multilevel Schauder estimates ..... 20
4.3 Extension to the space $\mathcal{D}^{\gamma, \eta}$ ..... 21
5 The fixed-point equation ..... 22
5.1 The set-up ..... 22
5.2 Local existence and uniqueness of the fixed point ..... 25
6 Renormalisation of the equation ..... 28
6.1 The renormalisation group ..... 28
6.2 Convergence of the renormalised models ..... 30
6.3 Computation of the renormalised equations ..... 36
7 Proof of the main results ..... 39
A Proof of Lemma 6.2 ..... 40

Nils Berglund<br>Université d'Orléans, Laboratoire MAPMO<br>CNRS, UMR 7349<br>Fédération Denis Poisson, FR 2964<br>Bâtiment de Mathématiques, B.P. 6759<br>45067 Orléans Cedex 2, France<br>E-mail address: nils.berglund@univ-orleans.fr<br>Christian Kuehn<br>Vienna University of Technology<br>Institute for Analysis and Scientific Computing<br>1040 Vienna, Austria<br>E-mail address: ck274@cornell.edu


[^0]:    ${ }^{1}$ Below, the renormalisation constants will rather be defined in terms of a truncated version of the heat kernel. This does however not affect their singular parts (terms of order $\varepsilon^{-1}$ and $\log \left(\varepsilon^{-1}\right)$ ).

[^1]:    ${ }^{2}$ The extended models are different, but $\mathcal{E}$ does not depend on the model. The situation is thus much simpler as in [19, Thm 5.12], where the $\mathcal{K}_{\gamma}$ associated with the two models are different.

[^2]:    ${ }^{3}$ See also the fact that $\Pi_{x} \tau(y) \leqslant C\|x-y\|_{\mathfrak{s}}^{|\tau|}$, which vanishes in $x=y$ if $|\tau|>0$, cf. [19, Paragraph after (8.28)].

