FRONT MOTION IN THE ONE-DIMENSIONAL STOCHASTIC CAHN–HILLIARD EQUATION∗

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Abstract. In this paper, we consider the one-dimensional Cahn–Hilliard equation perturbed by additive noise and study the dynamics of interfaces for the stochastic model. The noise is smooth in space and defined as a Fourier series with independent Brownian motions in time. Motivated by the work of Bates and Xun on slow manifolds for the integrated Cahn–Hilliard equation, our analysis reveals the significant difficulties and differences in comparison to the deterministic problem. New higher order terms that we estimate appear due to Itô calculus and stochastic integration and dominate the exponentially slow deterministic dynamics. Using a local coordinate system and defining the admissible interface positions as a multidimensional diffusion process, we derive a first order linear system of stochastic ordinary differential equations approximating the equations of front motion. Furthermore, we prove stochastic stability of the approximate slow manifold of solutions over a very long time scale and evaluate the noise effect.

Key words. one-dimensional stochastic Cahn–Hilliard equation, slow manifold, interface motion, additive noise, dynamics, stability

AMS subject classifications. 35K55, 35K40, 60H30, 60H15

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1. Introduction.

1.1. The problem. The standard Cahn–Hilliard equation is a simple model for the phase separation of a binary alloy at a fixed temperature proposed in [18, 19]. This model was extended by Cook [25, 43] in order to incorporate thermal fluctuations in the form of an additive noise. In this paper, we consider the one-dimensional Cahn–Hilliard equation posed on (0, 1) with an additive stochastic term:

\[(SC-H) \quad u_t = (-\varepsilon^2 u_{xx} + f(u))_{xx} + \partial_x \dot{W}_x, \quad 0 < x < 1, \quad t > 0,\]

with no-flux boundary conditions of Neumann type:

\[(1.1) \quad u_x = u_{xxx} = 0 \text{ at } x = 0, 1.\]

The nonlinearity \(f = f(u)\) is the derivative of a smooth double equal-well potential \(F\) taking its global minimum value 0 at \(u = \pm 1\) [1] with nondegenerate minima. A typical example is \(F(u) := \frac{1}{4}(u^2 - 1)^2\) with \(f(u) := u^3 - u\). The parameter \(\varepsilon > 0\) is a small atomistic interaction length modeling the width of layers that develop during the initial phase separation of spinodal decomposition (cf. [13, 14]). In the later

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stages of the separation process \( \varepsilon \) measures the width of transitions between the pure
phases \( u = \pm 1 \). Here, \( W_\varepsilon \) is a space-time noise smooth in space and is defined as the
formal derivative of an \( \varepsilon \)-dependent Wiener process \( W_\varepsilon \). As is common in stochastic
phase-field models, the noise scales with \( \varepsilon \). See, for example, the work of Funaki [35]
or Shardlow [47] on the stochastic Allen–Cahn equation. Here the noise strength is controlled by \( \varepsilon \); more specifically, it is bounded by \( O(\varepsilon^\delta) \) for some \( \delta > 9/2 \). For details see Assumption 2.3.

A characteristic feature of the Cahn–Hilliard equation model is the conservation of total mass
\( \int_0^1 u(t, x)dx \), which we now fix to be \( M \in (-1, 1) \). Substituting \( \tilde{u}(t, x) := \int_0^x u(t, y)dy \) we obtain the equivalent integrated stochastic Cahn–Hilliard equation

\[
\text{(ISC-H)} \quad \tilde{u}_t = -\varepsilon^2 \tilde{u}_{xxxx} + (f(\tilde{u}_x))_x + W_\varepsilon, \quad 0 < x < 1, \quad t > 0,
\]

associated with the boundary conditions

\[
\tilde{u}(t, 0) = 0, \quad \tilde{u}(t, 1) = M, \quad \tilde{u}_{xx}(t, 0) = \tilde{u}_{xx}(t, 1) = 0.
\]

Carr and Pego in [22, 23] presented a detailed analysis of the slow evolution of patterns of the singularly perturbed Ginzburg–Landau equation. They proved existence and persistence of metastable patterns and analyzed the equations governing their motion. These metastable states have been characterized in terms of the global unstable manifolds of equilibria. In [8, 9], Bates and Xun extended their argument and studied the dynamics of the one-dimensional Cahn–Hilliard equation in a neighborhood of an equilibrium having \( N + 1 \) transition layers, using several estimates presented in [22, 23]. They determined the exponentially slow speed of the layer motion and described precisely the layer motion directions. In addition, they established existence of an \( N \)-dimensional unstable invariant manifold attracting solutions exponentially fast uniformly in \( \varepsilon \). Related works in this direction are [10, 36, 45].

Motivated by the work of Bates and Xun for the deterministic problem, we study dynamics for the stochastic model. Due to stochastic integration, new higher order terms appear that we estimate using techniques and ideas of [8, 9, 22, 23]. In what follows, we shall refer frequently to some important definitions and results presented in the aforementioned articles; therefore, we give some details concerning our notation. Following [22, 23], we use the letter \( f \) for the nonlinearity in (SC-H) and denote by \( F \) the double equal well potential. In [8, 9] the symbol \( W' \) is used in place of \( f \); we avoided such a notation since we denote by the standard symbol \( \dot{W} \) the additive noise.

1.2. The effect of noise. The stochastic Cahn–Hilliard equation being one of the important examples of the nonlinear Langevin equations, it is based on a field-theoretic approach to the nonequilibrium dynamics of metastable states (see, for example, [25, 40, 43]). The multidimensional generalized stochastic Cahn–Hilliard equation associated with Neumann boundary conditions posed on bounded domains contains a time dependent noise in the chemical potential and an additive noise defined as the formal derivative of a Wiener process. The chemical potential noise describes external fields [38, 40, 42], while the free-energy independent noise may describe thermal fluctuations or external mass supply [25, 38, 40, 43].

Existence and uniqueness of solution for the stochastic problem was first studied in [26], where the nonlinearity \( f \) is a polynomial of odd degree and the problem is posed on multidimensional rectangular domains. Further, in [26], the author proved existence of solution and of its density for the stochastic Cahn–Hilliard equation with
additive noise (in the sense of Walsh; cf. [48]) posed on cubic domains. When the trace of the Wiener process is finite, existence was analyzed in [30]. In [5], existence for the generalized stochastic Cahn–Hilliard equation was derived for general convex or Lipschitz domains; the main novelty was the derivation of space-time Hölder estimates for the Green’s kernel of the stochastic problem, by using the domain’s geometry, which can be very useful in many other circumstances. The polynomial nonlinearity which forces the solution to stay between the pure phases ±1 has been analyzed in [13, 14, 20, 21, 26, 30], while in [29, 28, 37] a stochastic Cahn–Hilliard equation with reflection was considered.

In [13, 14] (see [15] for a review), the effect of noise on evolving interfaces during the initial stage of phase separation is analyzed. The evolution of these interfaces is stochastic and not yet fully understood. In [13], the authors show that for a solution starting at the homogeneous state, the probability of staying near a certain finite-dimensional space of pattern is high as long the solution stays within the distance of the homogeneous state. Further, in [14], the dynamics of a nonlinear partial differential equation perturbed by additive noise are considered. Under the assumption that the underlying deterministic equation has an unstable equilibrium, the authors show that the nonlinear stochastic partial differential equation exhibits essentially linear dynamics even far from equilibrium.

On the other hand, interface motion has been studied for many related models like Allen–Cahn or Ginzburg–Landau and phase-field models; see, for example, [4, 12, 16] for a rigorous analysis or the results of [32] for formal arguments, which describe the interfaces as interacting Brownian motions. Numerical results for interface motion are presented in [39, 47]. The problem of singular perturbation for a reaction-diffusion stochastic partial differential equation of Ginzburg–Landau type is investigated in [34]. The motion of interfaces for Cahn–Hilliard equation was studied only in an unpublished note by Brassesco in 2003, where she studied a solution with a single interface on \( \mathbb{R} \). When properly rescaled, the interface is driven by non-Markovian dynamics (cf. [12] for a similar result). In [46], the authors present a numerical study of the late stages of spinodal decomposition with noise.

The deterministic Cahn–Hilliard equation was proposed by Cahn and Hilliard [18, 17] as a model for the phase separation of a binary alloy at a fixed temperature, with \( u(t,x) \) defining the mass concentration of one of the phases at a point \( x \) at time \( t \). For a more physical background, derivation, and discussion of the deterministic Cahn–Hilliard equation and related equations, we refer to [7, 17, 18, 31, 33] and the references therein. Results for the noisy Cahn–Hilliard equation are of great interest for the study of Ostwald ripening [2, 3, 41] and nucleation [11]. For a survey, including numerical results and conjectures concerning the nucleation problem, see [15].

1.3. The approximate slow manifold. The space-time noise that we introduce is smooth in space, allowing for the application of the Itô formula. For our study of the dynamics of transition layers for the stochastic model, we closely follow the approach of Bates and Xun and of Carr and Pego based on the analysis of an approximate invariant manifold \( \mathcal{M} \). Although constructed in a different way, it can be thought of as piecing together a rescaled one kink (or front) of steady state solutions on the whole real line. The elements of the manifold are parametrized by the position of the fronts given by \( h \in \mathbb{R}^{N+1} \). Nevertheless, in our case the dependency on time is stochastic. This fact leads to the very interesting and difficult problem of further investigating the properties of \( \mathcal{M} \) by means of deriving higher order estimates related to the stationary problem.
We first present the details necessary for the steady state solutions $\phi$, the parameters $h$, and the manifold $\mathcal{M}$. Given $\varepsilon > 0$, we consider $a$ such that $f'(u) > 0$ for all $u$ satisfying $|u + 1| < a$. Then (cf. [22]), there exists $\rho > 0$ such that if $\ell$ satisfies $\varepsilon < \rho$, a unique solution $\phi = \phi(x, \ell, \pm 1)$ exists for the stationary Dirichlet problem
\[
\varepsilon^2 \phi_{xx} - f(\phi) = 0, \quad -\ell/2 < x < \ell/2,
\]
\[
\phi = 0, \quad x = \pm \ell/2,
\]
that satisfies
\begin{align*}
\text{(a) } & \phi(x, \ell, +1) > 0 \text{ for } |x| < \ell/2 \text{ and } |\phi(0) - 1| < a, \\
\text{(b) } & \phi(x, \ell, -1) < 0 \text{ for } |x| < \ell/2 \text{ and } |\phi(0) + 1| < a.
\end{align*}
For sufficiently small $\varepsilon > 0$, it is known that $\phi \approx \pm 1$ with transition layers of order $O(\varepsilon)$ near $x = \pm \ell/2$.

Following [9], we consider the slowly evolving solutions with $N + 1$ layers well separated and bounded away from the boundary $x = 0, 1$ and define the set of admissible positions $h$ of the interfaces
\[
\Omega_\rho := \{ h \in R^{N+1} : 0 < h_1 < \cdots < h_{N+1} < 1 \text{ and } \frac{\varepsilon}{\rho} < h_j - h_{j-1}, \ j = 1, \ldots, N + 2 \}
\]
with $h_0 := -h_1$, $h_{N+2} := 2 - h_{N+1}$. These interfaces evolve in time, and we expect them to have a width of order $\varepsilon$. Thus, the distance is bounded below by $\varepsilon/\rho$ for some small $\rho$. Later we fix $\rho = \varepsilon^\kappa$ for any small $\kappa > 0$.

Let $h \in \Omega_\rho$ be given as above, and denote the midpoints between interfaces by $m_j := \frac{h_{j-1} + h_j}{2}$ for $j = 1, \ldots, N + 2$ with $m_0 = 0$ and $m_{N+1} = 1$. Moreover, we define the function $u^h : I_j := [m_j, m_{j+1}] \to \mathbb{R}$ for the interfaces $h$ by
\[
\begin{aligned}
u^h(x) &= \left[ 1 - \chi \left( \frac{x - h_j}{\varepsilon} \right) \right] \cdot \phi \left( x - m_j, h_j - h_{j-1}, (-1)^j \right) \\
&\quad + \chi \left( \frac{x - h_j}{\varepsilon} \right) \cdot \phi \left( x - m_{j+1}, h_{j+1} - h_j, (-1)^{j+1} \right),
\end{aligned}
\]
where $\chi : \mathbb{R} \to [0, 1]$ is a $C^\infty$ cutoff function such that $\chi = 1$ on $[1, \infty)$ and $\chi = 0$ on $(-\infty, -1)$. (See Figure 1.1.)

\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{fig1.png}
\caption{Gluing together positive and negative solutions of (1.3) to obtain $u^h \in \mathcal{M}$. Note that $m_1 = 0$, $m_{N+2} = 1$, and $I_j = [m_j, m_{j+1}]$.}
\end{figure}

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Definition 1.1 (approximate slow manifold). The first approximate manifold of the stochastic Cahn–Hilliard equation solution is defined by
\[ M_1 := \{ u^h : h \in \Omega_\rho \} \]
Fixing a mass \( M \in (-1, 1) \), we define as the second approximate manifold the submanifold \( \mathcal{M} \) of \( M_1 \) where mass conservation holds, i.e.,
\[ \mathcal{M} := \{ u^h \in M_1 : \int_0^1 u^h \, dx = M \} \]
For the integrated equation, we consider the manifold \( \tilde{\mathcal{M}} := \{ \tilde{u}^h : u^h \in \mathcal{M}, \tilde{u}^h(x) = \int_0^x u^h \, dx \} \).

Remark 1.2. In view of the initial stochastic equation (SC-H), conservation of mass holds if and only if formally
\[ \int_0^1 \partial_x W_{\varepsilon} \, dy = W_{\varepsilon}(1) - W_{\varepsilon}(0) = 0. \]
This is later ensured by our assumptions on \( W_{\varepsilon} \), which impose Dirichlet boundary conditions for \( W_{\varepsilon} \) (cf. Definition 2.2 and Assumption 2.3). A very simple rigorous example is the following: consider \( \tilde{u}_{\varepsilon} := \delta_{\varepsilon} g(x) \tilde{\beta}(t) \), where \( \tilde{\beta}(t) \) is a white noise in time and \( g \) a smooth function satisfying \( g(1) = g(0) \). Then by integrating in space the equation (SC-H) and using the fact that
\[ \int_0^1 \partial_x \tilde{W}_{\varepsilon} \, dy = \delta_{\varepsilon} \tilde{\beta}(t) \int_0^1 g_{\varepsilon}(y) \, dy = 0, \]
we obtain mass conservation even with the noise. This example extends to infinite series of such terms.

Throughout the entire paper we assume that the additive noise in (SC-H) satisfies (1.6), and therefore the proposed stochastic model exhibits mass conservation.

1.4. The new coordinate system. Along \( \tilde{\mathcal{M}} \) the natural coordinate system would be to use the parameters \( h \in \Omega_\rho \) for the position in \( \mathcal{M} \) (where \( N \) of them are sufficient due to mass conservation), together with the orthogonal projection onto \( \mathcal{M} \). In order to relate the coordinate system to the deterministic flow of (ISC-H), one approximates the tangent space of \( \mathcal{M} \) by the span of some functions \( E_{\xi}^i, i = 1, \ldots, N \), related to eigenfunctions of the linearization to be defined later in the paper. Here, we follow [8].

We denote the \( L^2(0,1) \) inner product by \( \langle u, v \rangle := \int_0^1 uv \, dx \) and the induced \( L^2 \)-norm by \( \| \cdot \| \) and we introduce the symbol \( \tilde{g}(t, x) := \int_0^x g(t, y) \, dy \) for any \( g \), which is spatially integrable.

Due to mass conservation, we reduce the parameter space \( \Omega_\rho \) by one dimension. Define
\[ \xi := (\xi_1, \ldots, \xi_N) = (h_1, \ldots, h_N) \]
and consider \( h_{N+1} \) as a function of \( \xi \). Thus, for \( \tilde{u}_{\xi}^h := \frac{\partial \tilde{u}^h}{\partial h_{N+1}} \) and \( \tilde{u}_{\xi}^\xi := \frac{\partial \tilde{u}^\xi}{\partial \xi} \), we obtain
\[ \tilde{u}_{\xi}^j = \frac{\partial \tilde{u}^h}{\partial h_{N+1}} \cdot \frac{\partial h_{N+1}}{\partial h_j} + \frac{\partial \tilde{u}^h}{\partial h_j}. \]
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Fig. 1.2. The local coordinate system \( \tilde{u} = \tilde{\xi} + \tilde{v} \) around \( \tilde{M} \) for \( N = 1 \) (two interfaces). Note that \( E_1^\xi \approx \tilde{\alpha}_1^\xi \), which is the tangential vector along the manifold.

We use \( \tilde{u} \rightarrow (\xi, \tilde{v}) \) as coordinate system around \( \tilde{M} \). As sketched in Figure 1.2, let us split a solution \( \tilde{u} \) of (ISC-H) into a sum of stochastic processes

\[
\tilde{u}(t) := \tilde{\xi}(t) + \tilde{v}(t).
\]

Here the position on \( \tilde{M} \) is given by \( \tilde{\xi} \in \tilde{M} \), while the distance from \( \tilde{M} \) is given by \( \tilde{v} \), which is defined as the following projection such that

\[
\langle \tilde{v}, E_1^\xi \rangle = 0 \quad \text{for} \quad j = 1, \ldots, N.
\]

It turns out that the functions \( E_j^\xi \) are good approximations to the first eigenfunctions of the linearized integrated Cahn–Hilliard operator, which in turn are good approximations to the tangent space of \( \tilde{M} \). They are defined as follows:

\[
E_j^\xi(x) := \tilde{w}_j(x) - Q_j(x), \quad \tilde{w}_j(x) := \tilde{u}_j(x) + \tilde{u}_j^{(1)}(x),
\]

\[
Q_j(x) := \left( -\frac{1}{6} x^3 + \frac{1}{2} x^2 - \frac{1}{3} x \right) \tilde{w}_j^{xx}(0) + \frac{1}{6} (x^3 - x) \tilde{w}_j^{xx}(1) + x \tilde{w}_j(1), \quad j = 1, \ldots, N.
\]

The \( Q_j \) are exponentially small terms (cf. [8, pp. 437–439]), taking care of the boundary values of \( E_j^\xi \). More precisely, \( \tilde{w}_j \) are good approximations of these eigenfunctions, while \( \tilde{w}_j(0) = 0 \), and \( \tilde{w}_j(1), \tilde{w}_j^{xx}(0), \tilde{w}_j^{xx}(1) \) are exponentially small quantities. Introducing the polynomial correction terms \( Q_j \) in the definition of \( E_j^\xi(x) \) modifies the \( \tilde{w}_j \) so that \( E_j^\xi \) are good approximations and satisfy exactly the boundary conditions of the linearized integrated Cahn–Hilliard operator, i.e.,

\[
E_j^\xi = E_j^{\xi xx} = 0 \quad \text{for} \quad x = 0, 1.
\]

For short-hand notation, we also define higher derivatives using indices:

\[
E_{i\ell}^\xi := \frac{\partial E_i^\xi}{\partial \xi_\ell}, \quad E_{i\ell k}^\xi := \frac{\partial^2 E_i^\xi}{\partial \xi_\ell \partial \xi_k}, \quad \text{and} \quad \tilde{u}_{ik}^\xi := \frac{\partial^2 \tilde{u}_i^\xi}{\partial \xi_k \partial \xi_i}.
\]

1.5. Assumptions on the noise and layers. Throughout this paper the following three fundamental assumptions are considered for the noise and transition layers:

(a) The noise is sufficiently regular in space and of small strength. As derived, the manifold used in this paper is stable and attractive for a long time-scale with high probability and thus is a good approximation of the stochastic Cahn–Hilliard
equation solution, if the noise $\dot{W}_\varepsilon$ is sufficiently regular in space (cf. Assumption 2.3) and its strength is bounded by $O(\varepsilon^\delta)$ for some $\delta > 9/2$. The noise is presented in detail in Definition 2.2 as the formal derivative of a Wiener process $W_\varepsilon$ given by a Fourier series of independent Brownian motions in the sense of Da Prato and Zabczyk [27].

(b) We analyze local solutions of the (ISC-H): The coordinates $\xi$ of the projection onto the manifold perform a diffusion process. A main difference in stochastic dynamics of interfaces in comparison to the deterministic problem is that due to the noise, the movements of the layers are correlated, and thus the resulting stochastic ODE system given by (2.1) may be nonlinear for a general noise definition.

In order to make the analysis tractable, when we derive the equations of motion for the interface we assume that the coordinates $\xi$ of the projection onto the manifold perform a multidimensional diffusion process. By this natural assumption, we consider that the interfaces solve a very general stochastic ordinary equation driven by a Wiener process.

To be more precise, let $\tilde{u}$ be the solution of (ISC-H), where $W_\varepsilon$ is an $\varepsilon$-dependent Wiener process defined in Definition 2.2. We assume that the projection coordinates $\xi(t)$ (positions of the interface) is a stochastic diffusion process in $\mathbb{R}^N$. Since the specific $W_\varepsilon$ is introduced in (ISC-H), then the only underlying probability space is the Wiener space corresponding to $W_\varepsilon$. Therefore, diffusion is driven by $W_\varepsilon$ and is defined for any $k = 1, \ldots, N$ by

$$d\xi_k = b_k(\xi)dt + \langle \sigma_k(\xi), dW_\varepsilon \rangle$$

for some unknown vector field $b: \mathbb{R}^N \to \mathbb{R}^N$ and some variance $\sigma$ on $\mathbb{R}^N$. The unknown functions $b$, $\sigma$ might depend not only on $\xi$ but also on time $t$ and the distance from the manifold $\tilde{v}$.

As a result, we apply Itô calculus to the general system (2.1) in order to calculate explicitly the correlations of layers movements and derive finally closed forms of $b$ and $\sigma$. The assumption of $\xi$ being a diffusion process is justified later in Theorem 3.2 after the derivation of the SDE for the motion of the interfaces. More specifically, the diffusion process $\xi$ exists locally as a solution of the SDE defined up to a stopping time since the nonlinearities are only locally Lipschitz. It is possible to continue solutions until they leave the domain of definition of the equations close to $\tilde{M}$. In addition, as long as $\xi$ is well defined and $\|\tilde{v}\|$ sufficiently small, then $\tilde{u}$ given by $\tilde{u} := \tilde{u}^\xi + \tilde{v}$ is well defined and solves the initial (ISC-H) equation (cf. Theorem 3.2). Further, by attractivity and stochastic stability, we derive that the time of existence is with high probability larger than the exit time from some slow channel (neighborhood of the approximate manifold), in which we study the stability result.

(c) The number of transition layers is fixed. This is a natural assumption, which is also present in the work of Carr and Pego [22] and Bates and Xun in [8, 9]. Suggested by Fusco and Hale in [36] and further analyzed in [22, 8, 9], a geometric method was adopted and developed for the construction of a slow manifold of functions approximating a metastable state. This construction is valid for a fixed number of transition layers.
In [8], the study of dynamics of the one-dimensional Cahn–Hilliard equation considers the slow evolution of patterns in a neighborhood of an equilibrium having \( N + 1 \) transition layers. Further, in the aforementioned paper, the authors constructed an \( N \)-dimensional approximate invariant manifold consisting of states with a fixed number of \( N + 1 \) transition layers and a narrow tubular neighborhood or channel around this manifold. Solutions starting nearby approach this channel exponentially fast. In addition, [9] verifies the existence of an \( N \)-dimensional invariant manifold and all solutions inside the slow channel are attracted exponentially fast to this invariant manifold. The change of numbers of layers is only possible either by a rare stochastic event or when the solution leaves the slow channel after moving slowly along the manifold.

In our analysis, we study the dynamics for the stochastic problem locally in time, i.e., as long as the number of transition layers is fixed and thus indeed the layer locations are well separated and bounded away from the boundary points 0, 1 (cf. [22, 8, 9] for the deterministic problem). This is also justified by the fact that, as we prove, for a sufficiently bounded noise strength stability and attractivity of the manifold hold in the stochastic case also, at least for a very long time scale and with high probability. Of course the solution can leave the manifold at the boundary by a layer breaking down.

Moreover, due to rare stochastic events an extra “bump” (layer) could be formed. In our case this interesting event is rather unlikely, since the strength of the additive noise is sufficiently low so that the manifold \( \tilde{M} \) is stable and attractive with high probability. Apart from large deviation results the rigorous mathematical study on extra layers generation is highly not trivial. See, for example, the work of Chen [24] on generation, propagation, and annihilation of metastable patterns for the deterministic Allen–Cahn equation. Therefore, this is not analyzed in the present paper.

2. Main results. The SDE system for the motion of fronts is given by the projection onto the manifold \( \tilde{M} \), using the coordinate system of section 1.4. We then prove that \( \tilde{M} \) is locally exponentially attracting and show that solutions stay with high probability in a small slow tube around \( \tilde{M} \), until large times or until one of the layers becomes small. The flow along \( \tilde{M} \) is well described by the SDE for the interfaces \( \xi \). Depending on the strength of the noise, we investigate what the equation of motion of the fronts looks like and evaluate the noise effect. In addition, we study extensively the case \( N = 1 \), where the motion of the second interface is determined by the first. Here the motion is given by the Wiener process \( \dot{W}_c \) projected onto \( \tilde{M} \). Finally, the case of space-time white noise is discussed. In the last section, we present the proofs of the estimates used in our analysis concerning all the higher order terms that appear in the stochastic setting. These are technical results that are independent of the other sections.

We first explain briefly how the equations of motions along \( \tilde{M} \) are derived in section 3. For details we refer to subsection 3.2. If \( \tilde{u} \) is the solution of (ISC-H), then applying the Itô formula in differentiating with respect to \( t \), we get for \( i = 1, \ldots, N \) the following system in \( d\xi_1, \ldots, d\xi_N \) for the stochastic Cahn–Hilliard equation:

\[
\sum_j \left[ \langle \tilde{u}_j^\xi, E^\xi_i \rangle - \langle \tilde{v}, E^\xi_{ij} \rangle \right] d\xi_j = \int -\varepsilon^2 (\tilde{u}_j^{xx} + \tilde{v}_j^{xxxx}) + (f(\tilde{u}_j^x + \tilde{v}_j)) \right]_x, E^\xi_idt
\]
\[
+ \sum_{i,k} \left[ \frac{1}{2} \langle \tilde{v}, E^\xi_{ik} \rangle - \frac{1}{2} \langle \tilde{u}_k^x, E^\xi_i \rangle - \langle \tilde{u}_k^x, E^\xi_{ik} \rangle \right] d\xi_id\xi_k
\]
\[
+ \sum_j \langle dW_c, E^\xi_{ij} \rangle d\xi_j + \langle E^\xi_i, dW_c \rangle.
\]

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We note that the last three additive terms above at the right-hand side are not present in [8, 9], where the deterministic Cahn–Hilliard equation was studied.

Remark 2.1. In view of (2.1), we observe that the analysis of the stochastic dynamics is a much more complicated and difficult problem compared to the deterministic one.

1. Deterministic case: The system is linear in $d\xi_j$, therefore by estimating the inverse matrix on the left-hand-side (possibly close to $\mathcal{M}$) and the right-hand-side terms, the motion of interfaces is obtained; see [9].

2. Stochastic case: Obviously, for a general noise definition the system is non-linear due to the appearance of $d\xi_\ell d\xi_k$, which as we shall prove will dominate the exponentially small deterministic dynamics. In what follows, in order to get rid of the correlations $d\xi_\ell d\xi_k$, we make the ansatz that $\xi$ performs a diffusion process, which is justified later. Further, we need estimates for the additional higher order terms $E^\xi_{ij}$, $E^\xi_{ik}$, and $\hat{u}^\xi_{ik}$. Therefore, we need to improve the estimates of [8].

The sufficiently regular noise $W_\varepsilon$ is the formal derivative of a Wiener process $W_t$ defined as follows.

Definition 2.2 (the Wiener process $W_t$, [27]). Let $W_t$ be a $Q_\varepsilon$-Wiener process in the underlying Hilbert space $H = L^2(0,1)$, $Q_\varepsilon$ a symmetric operator, and $(e_k)_{k \in \mathbb{N}}$ an orthonormal basis with corresponding eigenvalues $\alpha_{\varepsilon,k}^2$ such that

$$Q_\varepsilon e_k = \alpha_{\varepsilon,k}^2 e_k$$

for a sequence of independent real-valued standard Brownian motions $\{\beta_k(t)\}_{t \geq 0}$.

We always rely on the following assumption, which implies mass conservation and regularity.

Assumption 2.3. Suppose that the $e_k$ are the eigenfunctions of the Dirichlet–Laplacian. Moreover, assume for some $\delta_\varepsilon > 0$

$$\begin{align*}
(1) \quad &\|Q_\varepsilon\| < C\delta_\varepsilon^2, \\
(2) \quad &\sum_{k=1}^{\infty} \alpha_{\varepsilon,k}^2 B_\varepsilon(e_k) < C\delta_\varepsilon^2, \\
(3) \quad &\|\partial_\varepsilon Q_\varepsilon\| < C\delta_\varepsilon^2,
\end{align*}$$

where we assume additionally that $\delta_\varepsilon < \varepsilon^{\delta/2 - \kappa}$ for some small $\kappa > 0$. $B_\varepsilon$ is defined as

$$B_\varepsilon(e) = \varepsilon^2 \|e_{xx}\|^2 + \|e_x\|^2,$$

while for $g = \sum_{k=1}^\infty \gamma_k e_k \in L^2(0,1)$ the linear operator $\partial_\varepsilon Q_\varepsilon$ is defined as

$$(\partial_\varepsilon Q_\varepsilon) g := \sum_{k=1}^{\infty} \gamma_k \partial_\varepsilon (Q_\varepsilon e_k) = \sum_{k=1}^{\infty} \gamma_k \alpha_{\varepsilon,k}^2 \partial_\varepsilon e_k.$$

The first assumption on the norm of $Q_\varepsilon$ as an operator in $H$ implies that the strength of the noise is bounded by $O(\delta_\varepsilon)$, while the second and third are additional assumptions on the noise regularity. Note that $B_\varepsilon(\cdot)^{1/2}$ is equivalent to the standard $H^2$-norm (see (3.15)).

For the calculation of the motion of the interfaces, we will assume that $\xi(t)$ a diffusion process in $\mathbb{R}^N$ (see section 1.5) is defined for any $k = 1, \ldots, N$ by

$$d\xi_k = b_k(\xi) dt + \langle \sigma_k(\xi), dW_\varepsilon \rangle,$$
Fig. 2.1. The stability of the slow manifold $\tilde{M}$ for two interfaces ($N = 1$). A small tubular neighborhood $\Gamma$, the slow channel, is attracting over long time scales. Solutions tend to exit at the end of $\Gamma$ by a breakdown of an interface.

for some vector field $b : \mathbb{R}^N \to \mathbb{R}^N$ and some variance $\sigma : \mathbb{R}^N \to H^N$. Later in Theorem 3.2 we justify this ansatz.

Following [9] we define the matrix

$$A_{ij}(\xi) = \langle \tilde{u}_j^\xi, E^\xi_i \rangle - \langle \tilde{v}, E^\xi_i \rangle,$$

which is invertible close to the slow manifold. The assumptions on the noise combined with (2.1) gives the following SDE system for the motion of interfaces:

$$\sum_j A_{ij}(\xi) d\xi_j = \langle -\epsilon^2 (\tilde{u}_{jxxxx}^\xi + \tilde{v}_{xxxx}) + (f(\tilde{u}_x^\xi + \tilde{v}_x))_x, E^\xi_i \rangle dt$$

$$+ \sum_{l,k} \left[ \frac{1}{2} \langle \tilde{v}, E^\xi_{ikl} \rangle - \frac{1}{2} \langle \tilde{u}_k^{\xi l}, E^\xi_i \rangle - \langle \tilde{u}_k^{\xi l}, E^\xi_{il} \rangle \right] \langle \varphi_{\epsilon} \sigma_k(\xi), \sigma_l(\xi) \rangle dt$$

$$+ \sum_j \langle \varphi_{\epsilon} E^\xi_{ij}, \sigma_j(\xi) \rangle dt + \langle E^\xi_i, dW_\epsilon \rangle$$

(cf. also the equivalent presentation (3.11)). From this we can easily read off $b$ and $\sigma$. Moreover, the flow along $\tilde{M}$ is described by the interface positions. It is now easy to check, by construction, that the difference $\tilde{v} = \tilde{u} - \tilde{u}^\xi$ is actually the $\tilde{v}$ of the coordinate system (see section 1.4). In addition, a solution of (2.2) together with a corresponding equation for $\tilde{v}$ (see (3.14)) describes a solution $\tilde{u}$ of (ISC-H); see Theorem 3.2.

Further, in section 3, the variance $\sigma$ of the multidimensional diffusion process $\xi$ of the interfaces is computed first explicitly and then is estimated in terms of $\epsilon$. A main result is the stochastic analysis of the stability of the approximate manifold, which is presented in Theorem 3.6 (see also Figure 2.1). Over a long time scale of order $O(\epsilon^{-q})$ for any large $q > 0$, we show that, with high probability, the solution of the stochastic Cahn–Hilliard equation stays in a small neighborhood $\Gamma$ of the integrated manifold $\tilde{M}$, unless an interface breaks down.

In section 4, we first present Theorem 4.1, in which we approximate the terms in (2.2) and derive the equations of motion of interfaces. Further, we consider several examples where Theorem 4.1 is simplified. If the noise is exponentially small, then we recover the slow motion results of [8, 9]. There is a slow channel given by a neighborhood of $\tilde{M}$, in which with high probability the motion of the interfaces is described by the deterministic regime. There is also an interesting intermediate regime of still exponentially small noise, which for simplicity of presentation we do not consider in this article. Here, due to the presence of noise, additional deterministic and stochastic
In view of (1.5), we also define fronts and show that the approximate manifold is locally attracting. The solution of (SC-H), then ˜ξ is approximately the projection of the Wiener process Wε onto M.

Finally in this section, we also discuss the case of nonsmooth in space space-time white noise (Qε = δt I d), which is unfortunately not covered by our assumptions. Here ξ1 would be close to a Brownian motion with variance ε2/(4ℓ2).

Section 5 provides estimates for the second order derivatives 2ξh,j,1,ε, for the higher order derivatives of Eξ and ˜wξ, and a bound for the quantity (Lξ ˜w, 2ξh,j,1) (needed in the proof of the stability theorem). Here the operator Lξ acting on a general smooth in space function ϕ is given by

\[ Lξ(ϕ) := -ε2ϕxxxx + (f′(u)ϕ)x. \]

The results of this section are quite technical since their proof involves extensive calculations related to properties of solutions of the stationary problem (1.3). The new terms to estimate appear only in the stochastic setting due to the frequent application of the Itô formula and were therefore not treated in the work of Bates and Xun [8, 9] or Carr and Pego [22, 23].

3. Front motion. In this section, we derive the equations of motions of the fronts and show that the approximate manifold is locally attracting.

3.1. Preliminaries and definitions. Let us first recall some notation. If u is the solution of (SC-H), then ˜u(t, x) := \int_0^t u(t, y)dy is the solution of the integrated one, i.e., of (ISC-H). Let a, ε, ρ, N be given; for some ℓ such that ε/ℓ < ρ, we consider the unique solution \( \phi \) of (1.3) which satisfies the properties (a) and (b). Let also \( (h_1, \ldots, h_{N+1}) \in \Omega_ρ \) be the admissible interface positions and take \( h_0 := -h_1, h_{N+2} := 2 - h_{N+1} \).

Let \( ℓ_j := h_j - h_{j-1} \) be the distance between interfaces and \( ℓ := \min\{ℓ_1, \ldots, ℓ_N\} \) the lower bound on them. Note that by the construction of \( \Omega_ρ \) the functions \( \phi \) are always well defined. Let

\[ r := ε/ℓ, \quad β±(r) := 1 ± ϕ(0, ℓ, ±) \quad \text{and} \quad α±(r) := F(ϕ(0, ℓ, ±)) . \]

In view of (1.5), we also define

\[ ϕ^j(x) := ϕ(x - m_j, ℓ_j, (-1)^j) , \]

and \( u^h_j := ∂u^h_j/∂h_j \) for \( j = 1, \ldots, N + 1 \). Considering \( r_j := ε/ℓ_j \), let

\[ β^j(r) := \begin{cases} β+(r_j) & \text{for } j \text{ even}, \\ β-(r_j) & \text{for } j \text{ odd}, \end{cases} \quad \text{and} \quad β(r) := \max_j β^j(r) . \]
We recall that in [9], as an application of the implicit function theorem,

\begin{equation}
\frac{\partial h_{N+1}}{\partial h_j} = (-1)^{N-j} + O(\varepsilon^{-1} \beta(r)).
\end{equation}

In addition, let

\[ \alpha^j(r) := \begin{cases} 
\alpha_+(r_j) & \text{for } j \text{ even}, \\
\alpha_-(r_j) & \text{for } j \text{ odd}, 
\end{cases} \quad \text{and} \quad \alpha(r) := \max_j \alpha^j(r). \]

We see later that both \( \alpha \) and \( \beta \) are exponentially small in \( \varepsilon \) if we consider \( r_j \leq \rho \leq \varepsilon^\kappa \) for some small positive \( \kappa \).

### 3.2. The SDE for the front motion.

Let \( \tilde{u} \) be a solution of (ISC-H). We assume that the \( N \) front positions, i.e., the coordinates of \( \xi(t) = (\xi_1(t), \ldots, \xi_N(t)) \), define a multidimensional diffusion process which is given by

\begin{equation}
\begin{aligned}
\frac{d\xi_k}{dt} &= b_k(\xi)dt + \langle \sigma_k(\xi), dW_\xi \rangle, \\
&\quad k = 1, \ldots, N,
\end{aligned}
\end{equation}

for some vector field \( b : \mathbb{R}^N \to \mathbb{R}^N \) and some variance \( \sigma : \mathbb{R}^N \to H^N \). The main aim of this paragraph is to identify \( b \) and \( \sigma \), which might also depend on \( \tilde{v} \), i.e., on the distance from the manifold.

We use the Itô formula, in order to differentiate \( \tilde{u}^{\xi} \) with respect to \( t \), and get

\begin{equation}
\begin{aligned}
\frac{d\tilde{u}^{\xi}}{dt} &= \sum_{j=1}^N \tilde{u}_j^{\xi} d\xi_j + \frac{1}{2} \sum_{1 \leq k, l \leq N} \tilde{u}_{kl}^{\xi} d\xi_k d\xi_l + d\tilde{v} \\
&\quad \text{with } \tilde{u}_{kl}^{\xi} = \frac{\partial^2 \tilde{u}^{\xi}}{\partial \xi_k \partial \xi_l}.
\end{aligned}
\end{equation}

We take as in [9, p. 175] the inner product in the space of equation (ISC-H) with \( E_i^{\xi} \) to get for any \( i = 1, \ldots, N \)

\begin{equation}
\langle E_i^{\xi}, d\tilde{u} \rangle = \langle \mathcal{L}^{\xi}(\tilde{u}), E_i^{\xi} \rangle dt + \langle E_i^{\xi}, dW_\xi \rangle,
\end{equation}

where we defined the nonlinear ICH operator as

\[ \mathcal{L}^{\xi}(u) := -\varepsilon^2 u_{xxxx} + (f(u_x))_x \]

for short-hand notation.

On the other hand, if we take the inner product of (3.3) with \( E_i^{\xi} \), we derive

\begin{equation}
\langle E_i^{\xi}, d\tilde{u} \rangle = \sum_{j=1}^N \langle \tilde{u}_j^{\xi}, E_i^{\xi} \rangle d\xi_j + \frac{1}{2} \sum_{1 \leq k, l \leq N} \langle \tilde{u}_{kl}^{\xi}, E_i^{\xi} \rangle d\xi_k d\xi_l + \langle E_i^{\xi}, d\tilde{v} \rangle.
\end{equation}

Throughout the rest of this paper, any summation is on 1, 2, \ldots, \( N \) for any index. In order to eliminate \( d\tilde{v} \), we apply the Itô formula to the orthogonality condition \( \langle \tilde{v}, E_i^{\xi} \rangle = 0 \) and arrive at

\[ \langle E_i^{\xi}, d\tilde{v} \rangle = -\langle \tilde{v}, E_i^{\xi} \rangle - \langle d\tilde{v}, dE_i^{\xi} \rangle \]

\[ = -\sum_j \langle \tilde{v}, E_i^{\xi} \rangle d\xi_j - \frac{1}{2} \sum_{j,k} \langle \tilde{v}, E_i^{\xi} \rangle d\xi_j d\xi_k - \sum_j \langle E_i^{\xi}, d\tilde{v} \rangle d\xi_j. \]
Now, we use that \( d\tilde{v} = d\tilde{u} - d\tilde{u}^\xi \) and the fact that \( dt dt = 0 \) and \( dW_z dt = 0 \). In detail,

\[
\begin{align*}
\sum_j \langle E^\xi_{ij}, d\tilde{v} \rangle d\xi_j &= - \sum_j \langle E^\xi_{ij}, d\tilde{u} \rangle d\xi_j + \sum_j \langle E^\xi_{ij}, d\tilde{u}^\xi \rangle d\xi_j \\
&= - \sum_j \langle E^\xi_{ij}, \mathcal{L}^\xi(\tilde{u}) \rangle dt d\xi_j - \sum_j \langle E^\xi_{ij}, dW_z \rangle d\xi_j + \sum_{j,k} \langle E^\xi_{ij}, \tilde{u}^\xi_{jk} \rangle d\xi_k d\xi_j \\
&= - \sum_j \langle E^\xi_{ij}, dW_z \rangle d\xi_j + \sum_{j,k} \langle E^\xi_{ij}, \tilde{u}^\xi_{jk} \rangle d\xi_k d\xi_j,
\end{align*}
\]

where we took the inner product in space of equation (ISC-H) with \( E^\xi_{ij} \), and used that

\[
d\xi_j dt = b_j(\xi) dt + \langle \sigma_j(\xi), dW_z \rangle dt = 0.
\]

Therefore, by (3.6) it follows that

\[
\langle E^\xi_{ij}, d\tilde{v} \rangle = - \sum_j \langle \tilde{v}, E^\xi_{ij} \rangle d\xi_j - \frac{1}{2} \sum_{j,k} \langle \tilde{v}, E^\xi_{ij} \rangle d\xi_j d\xi_k
\]

\[
\sum_j \sum_{l,k} \left[ \langle \tilde{u}^\xi_{ij}, E^\xi_{lk} \rangle - \frac{1}{2} \langle \tilde{u}^\xi_{ij}, E^\xi_{lj} \rangle - \langle \tilde{u}^\xi_{ij}, E^\xi_{lk} \rangle \right] d\xi_l d\xi_k
\]

\[
+ \sum_j \langle dW_z, E^\xi_{ij} \rangle d\xi_j + \langle E^\xi_{ij}, dW_z \rangle.
\]

**Lemma 3.1.** For all \( 1 \leq k, l \leq N \) it holds that

\[
\langle \sigma_k(\xi), dW_z \rangle \langle \sigma_l(\xi), dW_z \rangle = \langle Q_x \sigma_k(\xi), \sigma_l(\xi) \rangle dt.
\]

**Proof.** Since \( d\beta_j d\beta_i = \delta_{ij} dt \) and \( W_z(t) = \sum_{k=1}^\infty \alpha_{e_k, \beta_k}(t) e_k \) we obtain, using Parceval’s identity,

\[
\langle \sigma_k(\xi), dW_z \rangle \langle \sigma_l(\xi), dW_z \rangle = \sum_{i,j} \alpha_{e_i, e_j} \langle \sigma_k(\xi), e_i \rangle \langle \sigma_l(\xi), e_j \rangle d\beta_j d\beta_i
\]

\[
= \sum_j \alpha_{e_j}^2 \langle \sigma_k(\xi), e_j \rangle \langle \sigma_l(\xi), e_j \rangle dt
\]

\[
= \sum_j \langle Q_x \sigma_k(\xi), e_j \rangle \langle \sigma_l(\xi), e_j \rangle dt = \langle Q_x \sigma_k(\xi), \sigma_l(\xi) \rangle dt.
\]

Analogously to this lemma we easily obtain (using \( dt dW_z = 0 \))

\[
\langle E^\xi_{ij}, dW_z \rangle d\xi_j = \langle E^\xi_{ij}, dW_z \rangle \langle \sigma_j(\xi), dW_z \rangle = \langle Q_x E^\xi_{ij}, \sigma_j(\xi) \rangle dt.
\]
Moreover, for short-hand notation, as in [8], we define the matrix \( A(\xi) = (A_{ij}(\xi)) \in \mathbb{R}^{N \times N} \) by

\[
A_{ij}(\xi) = \langle \hat{u}^\xi_j, E_i^\xi \rangle - \langle \hat{v}, E_{ij}^\xi \rangle ,
\]

which is invertible, provided that we are near the slow manifold (cf. Lemma 3.4). Let us denote the inverse matrix of \( A \) by \( A^{-1}(\xi) = (A_{ij}^{-1}(\xi)) \in \mathbb{R}^{N \times N} \).

Therefore, for all \( i \in \{1, \ldots, N\} \) we arrive at

\[
\sum_j A_{ij}(\xi)d\xi_j = \langle \mathcal{L}^c(\check{u}^\xi + \check{v}), E_i^\xi \rangle dt + \sum_{i,l,k} \left[ \frac{1}{2} \langle \check{v}, E_{il}^\xi \rangle - \frac{1}{2} \langle \check{u}_l^\xi, E_i^\xi \rangle - \langle \hat{u}_l^\xi, E_{il}^\xi \rangle \right] \langle Q_{\varepsilon} \sigma_k(\xi), \sigma_l(\xi) \rangle dt + \sum_j \langle Q_{\varepsilon} E_{ij}^\xi, \sigma_j(\xi) \rangle dt + \langle E_i^\xi, dW_\varepsilon \rangle .
\]

To obtain the equation for \( d\xi \) we use that \( d\xi = A(\xi)^{-1} A(\xi) d\xi \). Thus, the final equation for \( \xi \) (as long as \( \check{u} \) is near the manifold) is given for any \( r = 1, \ldots, N \) by

\[
d\xi_r = \sum_i A_{ri}^{-1}(\xi) \langle \mathcal{L}^c(\check{u}^\xi + \check{v}), E_i^\xi \rangle dt + \sum_{i,l,k} \left[ \frac{1}{2} \langle \check{v}, E_{il}^\xi \rangle - \frac{1}{2} \langle \check{u}_l^\xi, E_i^\xi \rangle - \langle \hat{u}_l^\xi, E_{il}^\xi \rangle \right] \langle Q_{\varepsilon} \sigma_k(\xi), \sigma_l(\xi) \rangle dt + \sum_i A_{ri}^{-1}(\xi) \sum_j \langle Q_{\varepsilon} E_{ij}^\xi, \sigma_j(\xi) \rangle dt + \sum_i A_{ri}^{-1}(\xi) \langle E_i^\xi, dW_\varepsilon \rangle .
\]

We can now recover \( \sigma \) and \( b \) from (3.11). The only term that does involve noise is the last one. Thus, in view of (3.2) we derive

\[
\sigma_r(\xi) = \sum_i A_{ri}^{-1}(\xi) E_i^\xi .
\]

After we obtained \( \sigma \) we can proceed in order to determine \( b(\xi) \) from the remaining terms (cf. (3.2)). So, we get for \( r = 1, \ldots, N \) that

\[
b_r(\xi) = \sum_i A_{ri}^{-1}(\xi) \langle \mathcal{L}^c(\check{u}^\xi + \check{v}), E_i^\xi \rangle + \sum_{i,l,k} \left[ \frac{1}{2} \langle \check{v}, E_{il}^\xi \rangle - \frac{1}{2} \langle \check{u}_l^\xi, E_i^\xi \rangle - \langle \hat{u}_l^\xi, E_{il}^\xi \rangle \right] \langle Q_{\varepsilon} \sigma_k(\xi), \sigma_l(\xi) \rangle + \sum_i A_{ri}^{-1}(\xi) \sum_j \langle Q_{\varepsilon} E_{ij}^\xi, \sigma_j(\xi) \rangle .
\]

### 3.3. Justification of the ansatz.

Let us first give an equation for \( \check{v} = \check{u} - \check{u}^\xi \) describing the flow “orthogonal” to \( \mathcal{M} \). Following [8, p. 449], we consider equation

\[
d\check{v} = d\check{u} - \sum_{j=1}^N \check{u}^\xi_j d\xi_j - \frac{1}{2} \sum_{kl} \check{u}^\xi_{kl} d\xi_k d\xi_l ,
\]
and thus the key equation for the distance from the manifold $\mathcal{M}$ is described by

\begin{equation}
\dot{v} = L^c(\bar{u}) dt - \sum_j \tilde{a}_j^2 b_j(\xi) dt - \sum_j \tilde{a}_j^2 (\sigma_j(\xi), dW_z) - \frac{1}{2} \sum_{kl} \tilde{a}_{kl}^2 (\sigma_k(\xi), \sigma_l(\xi)) dt + dW_z.
\end{equation}

The following theorem is straightforwardly verified.

**Theorem 3.2.** Consider the pair of functions $(\xi, \bar{v})$ as local solutions of the system given by (3.14) and the ansatz (3.2), where $\sigma$ and $b$ are given by (3.12) and (3.13).

As long as $\|\bar{v}\| = O(\varepsilon^{3/2})$ and $\xi(t) \in \Omega_\rho$, the function $\bar{u} = \bar{u}_\xi + \bar{v}$ is well defined and solves (ISC-H) with $\langle \bar{v}, E_j^\xi \rangle = 0$.

The orthogonality condition follows directly from (3.14) as the differential $d\langle v, E_j^\xi \rangle = 0$. The fact that $\bar{u}$ is a solution follows from a lengthy calculation. Basically, one reverses the calculation of the previous section leading to (3.2).

**3.4. Stability and attractivity of the manifold.** In this section, we prove the stability and discuss the attractivity of $\mathcal{M}$. Considering the stability, we show that with high probability (over a long time scale) the solution stays close to $\mathcal{M}$, unless an interface breaks down. See Figure 2.1.

In [8, Theorem B], Bates and Xun show that in the deterministic setting the slow manifold is exponentially attracting in an $O(\varepsilon^{3/2})$-neighborhood in $\mathbb{R}^2$, until the solution reaches an exponentially small neighborhood, where the motion of the solution along the manifold is exponentially slow. Using large deviation estimates, it is straightforward to verify for small noise that the stochastic solution follows the deterministic one up to error terms of the order of the noise strength. Hence, the exponential attraction of $\mathcal{M}$ still holds for (ISC-H), until the solution reaches a neighborhood of the manifold that is determined by the strength of the noise.

Here, for simplicity of presentation we will focus only on the stability of $\mathcal{M}$. The proof can be modified, in order to show attraction, too. Once we are in the slow channel around $\mathcal{M}$, we cannot exit with high probability for a long time scale $T_\varepsilon$, unless one of the interfaces breaks down.

We define $A_\varepsilon$ and $B_\varepsilon$ as

\begin{equation}
A_\varepsilon(\bar{v}) = \int_0^1 [\varepsilon^2 \bar{v}_{xx}^2 + f'(u_\xi) \bar{v}_{xx}^2] dx \quad \text{and} \quad B_\varepsilon(\bar{v}) = \int_0^1 [\varepsilon^2 \bar{v}_{xx}^2 + \bar{v}_{xx}^2] dx.
\end{equation}

Obviously, it holds that

$$\|\partial_x \bar{v}\|_{L^2}^2 = \int_0^1 \bar{v}_{xx}^2 dx \leq \int_0^1 [\varepsilon^2 \bar{v}_{xx}^2 + \bar{v}_{xx}^2] dx = B_\varepsilon(\bar{v}).$$

Observe also that even if the function $f'(u_\xi)$ appearing in the definition of $A_\varepsilon$ changes sign, then provided $\rho$ is small for $\bar{v} \in C^2$ satisfying $\bar{v} = 0$ at $x = 0, 1$ and $\langle \bar{v}, E_j^\xi \rangle = 0$ for any $j = 1, \ldots, N$, there exists $C$ independent of $\varepsilon$, $\bar{v}$ such that

$$CA_\varepsilon(\bar{v}) \geq \varepsilon^2 B_\varepsilon(\bar{v}).$$

This estimate, which depends heavily on the properties of $\bar{v}$, is established in [8] after an extensive analysis of the spectrum of the linearized integrated Cahn–Hilliard operator (see [8, pp. 434–446, Lemmas 4.2, 3.2, 3.4]). More specifically, Bates and Xun proved that for $\rho$ small the spectrum consists of exactly $N$ exponentially small
eigenvalues, while all the other eigenvalues are negative and bounded away from 0 uniformly in \( \varepsilon \). Under weaker assumptions, such as \( \hat{\nu} \in H^2 \), the same estimate follows; cf. Theorem A.1 of [9, pp. 209–211]. Further, since \( f'(u^\varepsilon) \) is bounded we get 
\[ A_\varepsilon(\hat{\nu}) \leq c B_\varepsilon(\hat{\nu}), \]
while by definition and for \( \varepsilon < 1 \) it follows that 
\[ B_\varepsilon(\hat{\nu}) \leq \| \hat{\nu} \|^2_{H^2}. \]
Hence, the next relation holds true:
\[ \| \partial_x \hat{\nu} \|^2_L \leq B_\varepsilon(\hat{\nu}) \leq C \varepsilon^{-2} A_\varepsilon(\hat{\nu}) \leq c \varepsilon^{-2} B_\varepsilon(\hat{\nu}) \leq c \varepsilon^{-2} \| \hat{\nu} \|^2_{H^2}. \]
In addition, by Lemma 4.1 of [8] at p. 445, we have
\[ \| \hat{\nu} \|^2_\infty \leq B_\varepsilon(\hat{\nu}), \quad \| \hat{\nu}_x \|^2_\infty \leq \frac{1 + \varepsilon}{\varepsilon} B_\varepsilon(\hat{\nu}). \]

**Definition 3.3 (cf. [8, p. 452]).** Define a neighborhood \( \Gamma' \) of \( \mathcal{M} \) by
\[ \Gamma' = \{ \hat{\nu}^\varepsilon + \hat{\nu} : \xi \in \Omega_\rho, \ B_\varepsilon(\hat{\nu}) < \varepsilon^3 \}, \]
and we define the slow tube \( \Gamma \) by
\[ \Gamma := \{ \hat{\nu}^\varepsilon + \hat{\nu} : \xi \in \Omega_\rho, \ A_\varepsilon(\hat{\nu}) < \delta_\varepsilon^{2-\kappa} \}, \]
where \( 0 < \kappa \ll 1 \) presented in the definition of the noise (cf. Assumption 2.3) and \( \delta_\varepsilon \) estimates the noise strength.

The small tube \( \Gamma' \) is a neighborhood of the slow manifold, where the coordinate system (cf. (1.7)) is well defined, while the slow tube \( \Gamma \) is a neighborhood in which with high-probability solutions do not exit for long times unless one of the interfaces breaks down. Recall that \( \Gamma \subset \Gamma' \) by definition of \( \delta_\varepsilon \). We even have \( B_\varepsilon(\hat{\nu}) < C \delta_\varepsilon^{2-\kappa} \varepsilon^{-2} \leq C \varepsilon^7 \), which we need in the proof of stability.

As indicated in the introduction, the first term at the right-hand side of the flow given by (3.11) is identical to the right-hand side of the deterministic flow and has been estimated in [8]. In our stochastic case, in order to approximate the flow, we need to bound also the additional higher order terms and estimate the contribution of the noise. In section 4, we will identify the dominant terms in (3.11).

Using (4.27) of [9] and the fact that \( \| E_\varepsilon^\nu \| = O(\varepsilon^{-1/2}) \) [9, p. 187], we obtain in \( \Gamma' \) considering the matrix \( A \) the following invertibility result.

**Lemma 3.4.** Suppose that \( h \in \Omega_\rho \) and \( \| \hat{\nu} \| = O(\varepsilon^{3/2}) \); then
\[ A_{ij}(\xi) = O(\varepsilon) + \begin{cases} (-1)^{i+j+4} \ell_{j+1} & \text{if } i \geq j, \\ 0 & \text{if } i < j, \end{cases} \]
and the matrix is invertible with
\[ A_{ij}^{-1}(\xi) = O(\varepsilon) + \begin{cases} \frac{1}{\tau_{j+1}} & \text{if } i = j, j - 1, \\ 0 & \text{otherwise}, \end{cases} \]
where \( 1 > \ell_i > \varepsilon/\rho \) denotes the length of the \( i \)-th interface.

As the equation is deterministically stable, we can show that \( \hat{\nu} \) stays small for a long time (depending on the noise strength). To be more precise, we show a bound on \( A_\varepsilon(\hat{\nu}) \) for solutions near \( \mathcal{M} \). Compare also (86) of [8].

Fix some large time \( T_\varepsilon \) and define \( \tau^* > 0 \) as the first exit time (below the threshold \( T_\varepsilon \)) of \( \hat{\nu} \) from \( \Gamma' \). This is the stopping time
\[ \tau^* = T_\varepsilon \wedge \inf \{ t > 0 : \xi(t) \not\in \Omega_\rho \text{ or } A_\varepsilon(\hat{\nu}(t)) \geq \delta_\varepsilon^{2-\kappa} \}. \]
Note that for \( t \leq \tau^* \) also \( B_\varepsilon(\hat{\nu}(t)) \leq C \varepsilon^7 \), as discussed above.

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**Definition 3.5.** We say that a term is $O(\varepsilon^r)$ if it is asymptotically smaller than any polynomial in $\varepsilon$ uniformly for times $t \leq \tau^*$. 

Note that $\alpha, \beta$ are $O(\varepsilon^r)$ if $\rho = \varepsilon^r$.

**Theorem 3.6.** Suppose $\rho = \varepsilon^\kappa$ for some small $\kappa > 0$, $\delta_{e} \geq C\varepsilon^q$ for any $q > 0$, and suppose that for all $p > 0$ there exists a constant $c_p > 0$ such that $E A_{e}(\hat{v}(0))^p \leq c_p \delta_{e}^{2p}$.

Then for all $p > 0$ there exists a constant $C_p > 0$ such that

$$E A_{e}(\hat{v}(\tau^*))^p \leq C_p(T_{e} + 1)\delta_{e}^{2p}.$$  

Therefore, we can show that the probability that the solution exits from the slow tube before $T_{e}$ (i.e., $\tau^* = T_{e}$) or an interface is breaking down (i.e., $\xi(\tau^*) \notin \Omega_p$) is bounded above by

$$P \left( A_{e}(\hat{v}(\tau^*)) \geq \delta_{e}^{2-\kappa} \right) \leq E A_{e}(\hat{v}(\tau^*))^p \delta_{e}^{-(2-\kappa)} \leq C_p(T_{e} + 1)\delta_{e}^{\kappa p}$$

for any $p > 0$. Thus the probability that the solution exits from the slow tube before $T_{e}$ is of the order of $O(\varepsilon^r)$ provided $T_{e} \ll \delta_{e}^{-q}$ for some large $q > 0$. The typical case for applications would be to consider a noise strength polynomial in $\varepsilon$, where we can take $T_{e} = \varepsilon^{-q}$ for any $q > 0$.

**Remark 3.7** (exponentially small noise strength $\delta_e$). If we want to have exponentially long times $T_{e}$, then we need to take exponentially small noise strength $\delta_{e}$ and look closer at the various error terms in the proof of Theorem 3.6. This is straightforward, but for simplicity of presentation, we refrain from stating details here.

On the other hand, assuming that $\delta_{e}$ is exponentially small, the probability of the solution leaving the slow tube $\Gamma$ before time $T_{e}$, without an interface breaking down, is exponentially small, even for some exponentially large time $T_{e}$.

**3.5. Bounds on the SDE.** The following lemmas replace the bound on $\hat{\xi}$, used in the deterministic setting (cf. Lemma 4.3. in [8]).

**Lemma 3.8.** Let $\hat{\xi} + \hat{v} \in \Gamma'$ and $r = 1, \ldots, N$; then (with $E_{N+1}^{\xi} = 0$ for shorthand notation)

$$\sigma_{r}(\xi) = \frac{1}{4\ell_{r+1}}(E_{r}^{\xi} + E_{r+1}^{\xi}) + O(\varepsilon) \quad \text{and} \quad \|\sigma_{r}(\xi)\| \leq C/\ell < C\rho/\varepsilon.$$  

**Proof.** Note that $\|\hat{v}\| \leq B_{e}(\hat{v})^{1/2}$. Thus from the definition of $\sigma$ (cf. (3.12)), Lemma 3.4, and the bound on $E_{r}^{\xi}$ one obtains

$$\|\sigma_{r}(\xi)\| \leq \sum_{i} |A_{r}^{-1}(\xi)||E_{i}^{\xi}| \leq C/\ell.$$  

Moreover,

$$\sigma_{r}(\xi) = A_{r}^{-1}E_{r}^{\xi} + A_{r+1}^{-1}E_{r+1}^{\xi} + O(\varepsilon),$$

and the claim follows from Lemma 3.4. \(\square\)

The next lemma estimates the vector field $b$ of the diffusion process $\xi$.

**Lemma 3.9.** Let $\hat{\xi} + \hat{v} \in \Gamma'$ and assume that $\rho = \varepsilon^\kappa$ for some small $\kappa > 0$; then there is a constant $c > 0$ such that

$$|b_{r}(\xi)| \leq c\|Q_{r}\|\left(\varepsilon^{3\kappa-7/2} + \varepsilon^{2\kappa-5/2}\right) + O(\varepsilon)$$

for any $r = 1, \ldots, N$. 

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Proof. We recall first \( b_{\tau} \),

\[
(3.18) \quad b_{\tau}(\xi) = \sum_i A_{\tau i}^{-1}(\xi) \langle L^\varepsilon(\tilde{u}^\varepsilon + \tilde{v}), E_i^\varepsilon \rangle 
+ \sum_{i,l,k} A_{\tau i}^{-1}(\xi) \left[ \frac{1}{2} \langle \tilde{v}, E_{\tilde{u}kl}^\varepsilon \rangle - \frac{1}{2} \langle \tilde{u}^\varepsilon, E_{\tilde{u}il}^\varepsilon \rangle - \langle \tilde{u}^\varepsilon, E_{\tilde{l}i}^\varepsilon \rangle \right] \langle Q_\varepsilon \sigma_k(\xi), \sigma_l(\xi) \rangle 
+ \sum_i A_{\tau i}^{-1}(\xi) \sum_j \langle Q_\varepsilon E_j^\varepsilon, \sigma_j(\xi) \rangle.
\]

Then we use Lemma 3.4 and the bound on \( \sigma \). Moreover, in section 5, after tedious computations the next estimates are derived (cf. (5.44), (5.45), (5.46), (5.41), and (5.42), respectively):

\[
\begin{aligned}
|\langle \tilde{u}^\varepsilon_{kl}, E_i^\varepsilon \rangle| &\leq O(\varepsilon^{-1/2}) \left[ 4\ell_{i+1} + O(\varepsilon^{-3}\beta) \right], \\
|\langle \tilde{u}^\varepsilon, E_{\tilde{u}kl}^\varepsilon \rangle| &\leq O(\varepsilon^{-1/2} + \varepsilon^{-4}\beta), \\
|\langle \tilde{v}, E_{\tilde{u}il}^\varepsilon \rangle| &\leq O(\varepsilon^{-3/2} + \varepsilon^{-5}\beta) \| \tilde{v} \| \leq c + O(\varepsilon^{-7/2}\beta) \| \tilde{v} \|, \\
|\langle Q_\varepsilon \sigma_k(\xi), \sigma_l(\xi) \rangle| &\leq O(\varepsilon^{3/2} + \varepsilon^{5}\beta), \\
|\langle Q_\varepsilon E_j^\varepsilon, \sigma_j(\xi) \rangle| &\leq O(\varepsilon^{-1/2} + \varepsilon^{-3}\beta). 
\end{aligned}
\]

In addition, we observe that (cf. [9])

\[
\left| \sum_i A_{\tau i}^{-1}(\xi) \langle L^\varepsilon(\tilde{u}^\varepsilon + \tilde{v}), E_i^\varepsilon \rangle \right| = O(\alpha/\ell) + O(\varepsilon) = O(e_\varepsilon). 
\]

In this way, since \( \sigma = O(\rho e^{-1}) \) and \( A_{ij}^{-1} = O(\rho e^{-1}) \), we obtain

\[
|b_{\tau}(\xi)| \leq c|Q_\varepsilon|\rho^3\varepsilon^{-1/2} + c|Q_\varepsilon|\rho^5\varepsilon^{-5/2} + O(e_\varepsilon) 
\leq c|Q_\varepsilon| \left\{ \varepsilon^{-4/3} + \varepsilon^{-5/3} + O(e_\varepsilon) \right\}.
\]

3.6. Proof of stability. Now let us turn to the proof of Theorem 3.6. Considering the linearized Cahn–Hilliard operator and using the Itô formula we arrive at

\[
(3.19) \quad dA_\varepsilon(\tilde{v}) = d\langle -L^\varepsilon\tilde{v}, \tilde{v} \rangle + 2\langle -L^\varepsilon\tilde{v}, \tilde{d}\tilde{v} \rangle + dR
\]

with

\[
dR = \int_0^1 \tilde{v}_x^2 f''(u^\varepsilon) du^\varepsilon dx + \int_0^1 \tilde{v}_x f'''(u^\varepsilon) \tilde{v}_x dx + \int_0^1 \tilde{v}_x^2 f'''(u^\varepsilon) (du^\varepsilon)^2 dx.
\]

All terms in \( R \) appear, because \( A_\varepsilon \) itself depends on \( \xi \) through \( f'(u^\varepsilon) \). Using the Itô formula and (3.2) and (3.14) for \( \xi \) and \( \tilde{v} \), we expand all terms

\[
dR = \sum_j \int_0^1 \tilde{v}_x^2 f'''(u^\varepsilon) u_j^\varepsilon dx b_j dt + \sum_{i,j} \int_0^1 \tilde{v}_x^2 f'''(u^\varepsilon) u_{ij}^\varepsilon dx \langle Q_\varepsilon \sigma_j, \sigma_i \rangle dt
+ \sum_j \int_0^1 \tilde{v}_x f'''(u^\varepsilon) u_j^\varepsilon dx(\sigma_j, dW_\varepsilon) + \frac{1}{2} \sum_{i,j} \int_0^1 \tilde{v}_x^2 f'''(u^\varepsilon) u_{ij}^\varepsilon u_j^\varepsilon dx \langle Q_\varepsilon \sigma_j, \sigma_i \rangle dt
+ \sum_{i,j} \int_0^1 \tilde{v}_x f''(u^\varepsilon) u_j^\varepsilon u_i^\varepsilon dx \langle Q_\varepsilon \sigma_j, \sigma_i \rangle dt + \sum_j \int_0^1 \tilde{v}_x f''(u^\varepsilon) u_j^\varepsilon \partial_x \langle Q_\varepsilon \sigma_j \rangle dx dt.
\]
Now we use Theorem 5.8 in $\xi$ variables to obtain that $\|u_{ij}^\xi\|_\infty = o(\varepsilon^{-1})$, $\|u_{ij}^\xi\|_\infty = o(\varepsilon^{-2})$ and $\|u_{ij}^\xi\| = o(\varepsilon^{-1/2})$. Moreover, by definition it holds that $\|\tilde{e}_x\|^2 \leq B_\varepsilon(\tilde{v})$, so using Lemmas 3.9 and 3.8 we have $b_j = o(\varepsilon^{2-3/2})$ and $\sigma_j = o(\varepsilon^{-1})$. Finally, as $u^\xi$ is uniformly bounded, we can bound the nonlinearity $f$ by a constant and get

$$dR = o(\varepsilon^{9/2})dt + o(\varepsilon^{7/2})dt + (I_R, dW\varepsilon)$$

with

$$I_R = \sum_j \int_0^1 \tilde{e}_x^2 f''(u^\xi) u_{ij}^\xi \, dx \sigma_j = o(\varepsilon^{-2}).$$

As we are in the slow channel, we obtain

$$dR = o(\varepsilon^2)dt + (I_R, dW\varepsilon).$$

(3.20)

This is the crucial and only point where we need $B_\varepsilon(\tilde{v}) = o(\varepsilon^2)$ in order to estimate the fifth term of $R$.

Now we turn to the other terms in (3.19). Lemma 3.1 gives

$$dA_\varepsilon(\tilde{v}) - dR = 2\langle -L^\varepsilon \tilde{v}, L^\varepsilon(\tilde{u}) \rangle dt$$

(3.21)

$$- \sum_j 2\langle -L^\varepsilon \tilde{v}, \tilde{u}_{ij}^\xi \rangle b_j(\xi) dt$$

(3.22)

$$- \sum_j 2\langle -L^\varepsilon \tilde{v}, \tilde{u}_{ij}^\xi \rangle (\sigma_j(\xi), dW\varepsilon)$$

(3.23)

$$- \sum_{k_l} \langle -L^\varepsilon \tilde{u}_{kl}^\xi, \tilde{u}_{ij}^\xi \rangle (Q_{\xi} \sigma_k(\xi), \sigma_l(\xi)) dt$$

(3.24)

$$+ \sum_{ij} \langle -L^\varepsilon \tilde{u}_{ij}^\xi, \tilde{u}_{ij}^\xi \rangle (Q_{\xi} \sigma_i(\xi), \sigma_j(\xi)) dt$$

(3.25)

$$+ \sum_i \langle -L^\varepsilon \tilde{u}_{ij}^\xi, Q_{\xi} \sigma_i(\xi)\rangle dt$$

$$- 2\langle L^\varepsilon \tilde{v}, dW\varepsilon \rangle$$

(3.26)

$$+ \text{trace}(Q_{\xi}^{1/2} L^\varepsilon Q_{\xi}^{1/2}) dt.$$
where we used that for some constant \(a > 0\) independent of \(\varepsilon\) and \(r\) (cf. [8, Lemma 3.2, p. 434 and Lemma 4.2, p. 446])

\[
(3.28) \quad B_\varepsilon(\check{v}) < C\varepsilon^{-2}A_\varepsilon(\check{v}) < \frac{C}{\varepsilon^2}\varepsilon^{-2}\|L^c\check{v}\|^2.
\]

Using \(B_\varepsilon(\check{v}) = \mathcal{O}(\varepsilon^{6+\kappa})\) in the slow channel, we obtain

\[2\langle -L^c\check{v}, L^c(\check{u}) \rangle \leq -\frac{1}{2}\|L^c\check{v}\|^2 - aA_\varepsilon(\check{v}) + C\varepsilon^{-2}\alpha(\check{v})^2.
\]

Now consider the remaining four deterministic integrals. For the term in (3.22), notice that

\[
\langle L^c\check{v}, \check{u}_j^\xi \rangle = \langle \check{v}, L^c\check{u}_j^\xi \rangle = \langle \check{v}, \partial_x \partial_j L^b(u^\xi) \rangle.
\]

Thus integration by parts and Lemma 5.2 of [8] yields

\[
(3.29) \quad |\langle L^c\check{v}, \check{u}_j^\xi \rangle| \leq C\|\partial_x \check{v}\|\varepsilon^{-2}\beta(r) = \mathcal{O}(\varepsilon).
\]

We use now (3.29) to arrive at

\[
(3.30) \quad \left| \sum_j (L^c\check{v}, \check{u}_j^\xi)b_j(\xi) \right| \leq C\varepsilon^{-5/2}\beta(r)B_\varepsilon(\check{v})^{1/2} \sup_j \|b_j(\xi)\| = \mathcal{O}(\varepsilon),
\]

which is exponentially small in \(\varepsilon\) by Lemma 3.9. By Definition 3.5, a term is \(\mathcal{O}(\varepsilon)\) if it is asymptotically smaller than any polynomial in \(\varepsilon\) uniformly for times \(t \leq \tau^*\).

Now let us turn to (3.24). Similarly, we get

\[
|\langle -L^c\check{u}_i^\xi, \check{u}_j^\xi \rangle| = |\langle \check{u}_i^\xi, \partial_x \partial_j L^b(u^\xi) \rangle| \leq \|\check{u}_i^\xi\|_{L^1} \|\partial_x \partial_j L^b(u^\xi)\|_{L^\infty} \leq C\varepsilon^{-4}\beta(r),
\]

where we used Lemma 5.1 of [8] and the bound \(\|\check{u}_i^\xi\|_{L^1} = \mathcal{O}(1)\) (cf. (5.38) for \(\beta\) bounded). Thus we obtain for the term in (3.24)

\[
(3.31) \quad \left| \sum_{ij} \langle -L^c\check{u}_i^\xi, \check{u}_j^\xi \rangle \langle Q_\varepsilon \sigma_i(\xi), \sigma_j(\xi) \rangle \right| \leq C\varepsilon^{-4}\beta(r)\|Q_\varepsilon\|\varepsilon^{-2} = \mathcal{O}(\varepsilon).
\]

For the term in (3.23) we use the bounds on \(\langle -L^c\check{v}, \check{u}_k^\xi \rangle\) provided by Theorem 5.47. Thus, we get

\[
|\langle L^c\check{v}, \check{u}_k^\xi \rangle \langle Q_\varepsilon \sigma_k(\xi), \sigma_1(\xi) \rangle| \leq C\|Q_\varepsilon\|\varepsilon^{-2}C\varepsilon^{-2}\beta(\check{v})\|\check{v}\| = \mathcal{O}(\varepsilon).
\]

Using similar estimates and Lemma 3.8 the term in (3.25) is also \(\mathcal{O}(\varepsilon)\).

For the term in (3.26), we use the eigenfunctions \(e_k\) of \(Q_\varepsilon\) and the uniform bound on \(f'(u^\xi)\), in order to obtain

\[
\text{trace}(Q_\varepsilon^{1/2}L^cQ_\varepsilon^{1/2}) = \sum_{k=1}^\infty \alpha_{k,k}^2 \langle L^c e_k, e_k \rangle \leq C \sum_{k=1}^\infty \alpha_{k,k}^2 B_\varepsilon(e_k) \leq C\delta^2.
\]

This is the largest deterministic term, as the other ones are all \(\mathcal{O}(\varepsilon)\). This term comes directly from the Itô correction of the additive noise.
Consider now (3.21)–(3.26) with all deterministic integrals already estimated and include the bound on $R$ from (3.20). For $t \leq \tau^*$

$$
(3.32) \quad dA_\varepsilon(\tilde{v}(t)) \leq C\delta^2 \dot{\varepsilon}^2 dt - (\frac{1}{2}\|L^c\tilde{v}\|^2 + aA_\varepsilon(\tilde{v}))dt + \langle I, dW_\varepsilon \rangle,
$$

where

$$
I = \sum_j 2\langle -L^c\tilde{v}, \tilde{u}_j^c \rangle \sigma_j(\xi) - 2L^c\tilde{v} + I_R
$$

with $I_R = O(B_\varepsilon(\tilde{v})\varepsilon^{-3/2})$.

In order to bound $I$, we use (3.29) and the asymptotic formula for $\sigma_j(\xi)$ of Lemma 3.8 combined with (54)–(55) of [8] to obtain that $(L^c\tilde{v}, \tilde{u}_j^c)\sigma_j(\xi) = O(\varepsilon)$ and thus

$$
|\langle I, Q_\varepsilon I \rangle| \leq O(\varepsilon) + C\|Q_\varepsilon\| \|L^c\tilde{v}\|^2 + B_\varepsilon(\tilde{v})^2\varepsilon^{-4}.
$$

Now from (3.28) as in the slow channel at least $B_\varepsilon(\tilde{v}) = O(\varepsilon^6)$ we obtain $B_\varepsilon(\tilde{v})^2\varepsilon^{-3} \leq C\|L^c\tilde{v}\|^2B_\varepsilon(\tilde{v})\varepsilon^{-6} \leq C\|L^c\tilde{v}\|^2$ and thus

$$
|\langle I, Q_\varepsilon I \rangle| \leq O(\varepsilon) + C\|Q_\varepsilon\| \|L^c\tilde{v}\|^2.
$$

Now we can bound powers of $A_\varepsilon$ for $t \leq \tau^*$

$$
(3.33) \quad \frac{1}{p}dA_\varepsilon(\tilde{v})^p = A_\varepsilon(\tilde{v})^{p-1}dA_\varepsilon(\tilde{v}) + \frac{p-1}{2}A_\varepsilon(\tilde{v})^{p-2}(dA_\varepsilon(\tilde{v}))^2
$$

$$
\leq C\varepsilon^{2\delta}A_\varepsilon(\tilde{v})^{p-1}dt - (\frac{1}{2}\|L^c\tilde{v}\|^2 + aA_\varepsilon(\tilde{v}))A_\varepsilon(\tilde{v})^{p-1}dt + A_\varepsilon(\tilde{v})^{p-1}\langle I, dW_\varepsilon \rangle + \frac{p-1}{2}A_\varepsilon(\tilde{v})^{p-2}\langle I, Q_\varepsilon I \rangle dt.
$$

Taking integrals up to $\tau^*$ and expectation, we easily obtain from (3.32) and (3.33) (using that the expectation of a stochastic integral is 0)

$$
EA_\varepsilon(\tilde{v}(\tau^*)) + \frac{1}{2}E\int_0^{\tau^*}\|L^c\tilde{v}\|^2 dt + aE\int_0^{\tau^*}A_\varepsilon(\tilde{v})dt \leq A_\varepsilon(\tilde{v}(0)) + CT\varepsilon^2,
$$

and for $p \geq 2$

$$
\frac{1}{p}EA_\varepsilon(\tilde{v}(\tau^*))^p + \frac{1}{2}E\int_0^{\tau^*}\|L^c\tilde{v}\|^2A_\varepsilon(\tilde{v})^{p-1}dt + aE\int_0^{\tau^*}A_\varepsilon(\tilde{v})^p dt
$$

$$
\leq \frac{1}{p}EA_\varepsilon(\tilde{v}(0))^p + C\delta^2E\int_0^{\tau^*}A_\varepsilon(\tilde{v})^{p-1}dt + O(\varepsilon)E\int_0^{\tau^*}A_\varepsilon(\tilde{v})^{p-2}\|L^c\tilde{v}\|^2 dt + C\|Q_\varepsilon\| \cdot E\int_0^{\tau^*}A_\varepsilon(\tilde{v})^{p-2}\|L^c\tilde{v}\|^2 dt.
$$

Now (using $\delta \geq C\varepsilon^9$) it is easy to verify by induction on $p$ that

$$
\frac{1}{p}EA_\varepsilon(\tilde{v}(\tau^*))^p + \frac{1}{2}E\int_0^{\tau^*}\|L^c\tilde{v}\|^2A_\varepsilon(\tilde{v})^{p-1}dt + aE\int_0^{\tau^*}A_\varepsilon(\tilde{v})^p dt \leq C(T\varepsilon + 1)\delta^2\varepsilon^p.
$$

This implies the claim.
4. Motion of the interfaces. In this section, we investigate some important special cases in detail to see what the SDE (2.2) for \( \xi \) actually implies for the motion of the interfaces. We assume first that the noise is exponentially small. Then considering the two interfaces problem (i.e., when \( N = 1 \)) we discuss the case of noise strength being polynomial in \( \varepsilon \). Although not covered by our theorems, we present some comments on what the equation would look like for nonsmooth space-time white noise, which means that \( Q_\varepsilon \) is proportional to the identity. Finally, we present the approximate SDE system for the front motion considering the general precise system (2.1).

Let us first state the result we achieved so far. The motion of the interfaces for the stochastic model is given by the following theorem.

**Theorem 4.1.** Let \( \tilde{u}_\varepsilon + \tilde{v} \in \Gamma \) and assume that \( \rho \) is small. Then the equations dominating the flow of the stochastic Cahn–Hilliard equation within the slow channel are given by

\[
\begin{align*}
    d\xi_1 &= \frac{1}{4\ell_2}(\alpha^3 - \alpha^1)dt + O(\varepsilon\alpha)dt + dA_s^{(1)}, \\
    d\xi_2 &= \frac{1}{4\ell_2}(\alpha^3 - \alpha^1)dt + \frac{1}{4\ell_3}(\alpha^4 - \alpha^2)dt + O(\varepsilon\alpha)dt + dA_s^{(2)}, \\
    d\xi_3 &= \frac{1}{4\ell_3}(\alpha^4 - \alpha^2)dt + \frac{1}{4\ell_4}(\alpha^5 - \alpha^3)dt + O(\varepsilon\alpha)dt + dA_s^{(3)}, \\
    &\vdots \\
    d\xi_N &= \frac{1}{4\ell_N}(\alpha^{N+1} - \alpha^{N-1})dt + \frac{1}{4\ell_{N+1}}(\alpha^{N+2} - \alpha^N)dt + O(\varepsilon\alpha)dt + dA_s^{(N)},
\end{align*}
\]

where

\[
\alpha^j = \frac{1}{2}K_+^2\Delta_2^2 \exp(-A_\pm\ell_j/\varepsilon) \left[ 1 + O\left( \frac{\ell_j}{\varepsilon} \exp\left( \frac{-A_\pm\ell_j}{2\varepsilon} \right) \right) \right], \quad j = 1, 2, \ldots, N + 2,
\]

for

\[
A_\pm := f^i(\pm 1) \quad \text{and} \quad K_\pm := 2 \exp\left( \int_0^1 \left[ \frac{A_\pm}{2F(\pm t)^{1/2}} - \frac{1}{1-T} \right] dt \right).
\]

Here, the stochastic processes \( A^{(r)} \), \( r = 1, \ldots, N \), are related to the noise; they depend on the symmetric operator \( Q_\varepsilon \) and the variance \( \sigma \) and are given by the formula

\[
dA_s^{(r)} := \sum_{i,j,k} A_{r_i}^{-1}(\xi) \left[ \frac{1}{2}(\tilde{v}, E_{ik}^\xi) - \frac{1}{2}(\tilde{u}_k^\xi, E_{ij}^\xi) - (\tilde{u}_k^\xi, E_{ij}^\xi) \right] (Q_\varepsilon \sigma_k(\xi), \sigma_i(\xi)) dt
\]

\[
+ \sum_{j} A_{r_i}^{-1}(\xi) \langle Q_\varepsilon E_{ij}^\xi, \sigma_j(\xi) \rangle dt + \sum_{i} A_{r_i}^{-1}(\xi) \langle E_{ij}^\xi, dW_{\varepsilon} \rangle.
\]

The quantities \( K_\pm \) are constants introduced by Carr and Pego in [22].

**Proof.** Recall that as long as \( \tilde{u} \) is near the manifold, then by (3.11) we obtained for any \( r = 1, \ldots, N \)

\[
d\xi_r = \sum_{i} A_{r_i}^{-1}(\xi) \langle \mathcal{L}^r(\tilde{u}_\varepsilon + \tilde{v}), E_{ij}^\xi \rangle dt + dA_s^{(r)}.
\]

Lemma 3.4 gives that the matrix \( A^{-1} \) and therefore the terms \( \sum_i A_{r_i}^{-1}(\xi) \langle \mathcal{L}^r(\tilde{u}_\varepsilon + \tilde{v}), E_{ij}^\xi \rangle \) are identical to those presented in [8, 9] for the deterministic case (i.e., when \( dA_s^{(r)} = 0 \) for any \( r \)). Hence, using (4.32) of [9] we obtain the result. \( \square \)
Remark 4.2. Note that using the relation $t_j = h_j - h_{j-1}$ and the asymptotic formula for $\frac{\partial N_{j+1}}{\partial n_j}$ we can derive an analogous system in $h_j$ or in $t_j$ (cf. [9]).

Remark 4.3. In view of the assumptions of Theorem 4.1, and as mentioned throughout our analysis, it is sufficient that $\rho = \varepsilon^\kappa$ for some small $\kappa > 0$. In this case $\alpha^j$ are exponentially small. So, by stability, for a sufficiently bounded noise strength, the distance $\|\tilde{v}\|$ will remain small and thus the matrix $A^{-1}$ will remain well defined.

We observe
\begin{equation}
(4.5) \quad dA^{(r)}_s := A^{(r)}_s dt + \sum_i A^{-1}_{ri}(\xi_0)\langle E^\xi_i, dW_t \rangle
\end{equation}
for
\begin{equation}
(4.6) \quad A^{(r)}_s := \sum_{i,l,k} A^{-1}_{li}(\xi_0) \left[ \frac{1}{2} \langle \tilde{v}, E^\xi_{il} \rangle - \frac{1}{2} \langle \tilde{u}^\xi_{lk}, E^\xi_{li} \rangle - \langle \tilde{u}^\xi_{ik}, E^\xi_{li} \rangle \right] \langle Q, \sigma_k(\xi), \sigma_l(\xi) \rangle + \sum_i A^{-1}_{ri}(\xi_0) \sum_j \langle Q, E^\xi_{ij}, \sigma_j(\xi) \rangle.
\end{equation}

Following Lemma 3.9 we obtain in the slow channel
\begin{equation}
(4.7) \quad |A^{(r)}_s| \leq c\|Q\|\rho^2(\rho \varepsilon^{-3/2} + \varepsilon^{-5/2}) \quad \text{for all } r = 1, \ldots, N.
\end{equation}

Thus, in the case of $\|Q\| = O(\varepsilon^{4+1/2} \alpha)$, since $\rho$ is at least bounded, we can show that $A^{(r)}_s = O(\varepsilon \alpha)$. It is not hard to show that we can also neglect the stochastic term from (4.1) in order to recover the result of Bates and Xun on metastable slow motion, at least with high probability.

An interesting case arises when the additional terms in $A^{(r)}_s$ are of the order of $O(\alpha)$. Then we obtain additional terms in (4.1). Nevertheless, for simplicity of presentation, we refrain from stating details here. Obviously, for a polynomial noise strength the extra drift $A^{(r)}_s dt$ coming from stochastic dynamics would dominate the exponentially small terms involving $\alpha^j$ and $\alpha$.

4.1. Polynomial noise strength. For the remainder of this section we fix $N = 1$, which is the case of two interfaces, and a noise strength $\delta_\varepsilon = \varepsilon^\delta$ for some $\delta > 9/2$. To be more precise suppose $Q_\varepsilon = Q_0\varepsilon^\delta$ with $Q_0 = O(1)$.

Using (4.1), we notice that the equation of motion for the first interface is given by
\begin{equation}
d\xi_1 = O(\alpha) dt + dA^{(1)}_s,
\end{equation}
and the motion of the second interface is fixed due to mass conservation.

Recall that $\ell_2$ is the distance between the two interfaces, and fix $\rho = \varepsilon^\kappa$, which means that the lower bound on $\ell_2$ is $\varepsilon^{1-\kappa}$. Let us now first look at (3.12):
\begin{equation}
\sigma_1(\xi) = A_{11}^{-1} E^\xi_1.
\end{equation}
Since $\tilde{u}^\xi_1 = \tilde{u}^\xi_2 \frac{\partial h_2}{\partial n_1} + \tilde{u}^\xi_1$ while $\frac{\partial h_2}{\partial n_1} = 1 + O(e_\varepsilon)$ and $E^\xi_1 = \tilde{u}^\xi_1 + \tilde{u}^\xi_2 + O(e_\varepsilon)$, it follows that
\begin{equation}
E^\xi_1 = \tilde{u}^\xi_1 + O(e_\varepsilon),
\end{equation}
and again the error term remains of the same order under differentiation with respect to $\xi_1$. Second, from (4.24) in [9] there is a constant $c_*$ such that $||\tilde{u}^\xi_1||^2 = 4\ell_2 + c_*\varepsilon + \varepsilon$. 

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\( \mathcal{O}(e^\varepsilon) \), and the error term remains \( \mathcal{O}(e^\varepsilon) \) under differentiation. (In our case \( N = 1 \) we have that \( \tilde{u}_1 \) used in [9] is up to errors of order \( \mathcal{O}(e^\varepsilon) \) equal to \( \tilde{u}_1^\varepsilon \)). Moreover, by definition

\[
(4.8) \quad A_{11} = \langle \tilde{u}_1^\varepsilon, E_1^\varepsilon \rangle - \langle \tilde{v}, E_1^\varepsilon \rangle = ||\tilde{u}_1^\varepsilon||^2 + ||\tilde{v}||_\infty \mathcal{O}(\varepsilon^{-1/2}) + \mathcal{O}(e^\varepsilon),
\]

where we used (5.42) (cf. also [9], where the same estimate is used, though never presented analytically) for \( E_{11}^\varepsilon = \mathcal{O}(\varepsilon^{-1/2}) \). Recall that in the slow channel \( \Gamma \)

\[
(4.9) \quad ||v||_\infty \leq (B_\varepsilon(v))^{1/2} \leq C\varepsilon\langle A_\varepsilon(v) \rangle^{1/2} \leq C\varepsilon^{-1}(\delta^2 - \kappa)^{1/2} \leq C\varepsilon^{-1+\delta(1-\kappa/2)}.
\]

Thus we proved

\[
A_{11} = 4\ell_2 + c_\varepsilon + \mathcal{O}(\varepsilon^{(1-\kappa/2) - \frac{3}{2}}) \quad \text{and} \quad \sigma_1(\xi) = \frac{1}{4\ell_2 + c_\varepsilon + \mathcal{O}(\varepsilon^{(1-\kappa/2) - \frac{3}{2}})} E_1^\varepsilon + \mathcal{O}(e^\varepsilon).
\]

Now we can consider the deterministic drift

\[
A^{(1)}_Q = A^{-1}_{11}(\xi) \left[ \frac{1}{2}(\tilde{v}, E_{11}^\varepsilon) - \frac{1}{2}(\tilde{u}_1^\varepsilon, E_1^\varepsilon) - \langle \tilde{u}_1^\varepsilon, E_1^\varepsilon \rangle \right] \langle Q_\varepsilon \sigma_1(\xi), \sigma_1(\xi) \rangle
+ A^{-1}_{11}(\xi) \langle Q_\varepsilon E_1^\varepsilon, \sigma_1(\xi) \rangle
= A^{-3}_{11} \mathcal{O}(\varepsilon^{-3/2}) ||\tilde{v}|| - \frac{3}{4} \frac{\partial}{\partial \xi} \| E_1^\varepsilon \|^2 \| Q_\varepsilon^{1/2} E_1^\varepsilon \|^2 + A^{-1}_{11} \frac{1}{2} \frac{\partial}{\partial \xi} \| Q_\varepsilon^{1/2} E_1^\varepsilon \|^2 + \mathcal{O}(e^\varepsilon).
\]

Thus, in the slow channel \( \Gamma \) (cf. (4.9)) the equation of motion for the interface is reduced to

\[
d\xi = A^{-3}_{11} \mathcal{O}(\varepsilon^{(1-\kappa/2) - 5/2}) \| Q_\varepsilon^{1/2} E_1^\varepsilon \|^2 dt - \frac{3}{4} A^{-3}_{11} \left( \frac{\partial}{\partial \xi} \| E_1^\varepsilon \|^2 \right) \| Q_\varepsilon^{1/2} E_1^\varepsilon \|^2 dt
+ A^{-1}_{11} \frac{1}{2} \frac{\partial}{\partial \xi} \| Q_\varepsilon^{1/2} E_1^\varepsilon \|^2 dt + A^{-1}_{11}(E_1^\varepsilon, dW_\varepsilon) + \mathcal{O}(e^\varepsilon) dt.
\]

By (45) of [8] we know that

\[
\tilde{u}_1^\varepsilon = 1 - u^\varepsilon + \mathcal{O}(e^\varepsilon) \quad \text{and} \quad u_1^\varepsilon = -u^\varepsilon + \mathcal{O}(e^\varepsilon)
\]

(as \([0, 1] = I_1 \cup I_2 \) and \( u^\varepsilon(m_1) = u^\varepsilon(0) = -1 + \mathcal{O}(e^\varepsilon) \)). Furthermore, the error terms remain \( \mathcal{O}(e^\varepsilon) \), under differentiation with respect to \( \xi \). Thus, we obtain

\[
||\tilde{u}_1^\varepsilon||^2 = ||1 - u^\varepsilon||^2 + \mathcal{O}(e^\varepsilon) = 1 - 2M + ||u^\varepsilon||^2 + \mathcal{O}(e^\varepsilon).
\]

Differentiation yields

\[
\frac{\partial}{\partial \xi} ||\tilde{u}_1^\varepsilon||^2 = 2\langle u_1^\varepsilon, u_1^\varepsilon \rangle + \mathcal{O}(e^\varepsilon) = -2\langle u_1^\varepsilon, u_1^\varepsilon \rangle + \mathcal{O}(e^\varepsilon) = u_1^\varepsilon(0)^2 - u_1^\varepsilon(1)^2 + \mathcal{O}(e^\varepsilon) = \mathcal{O}(e^\varepsilon).
\]

Thus we verified that

\[
\frac{\partial}{\partial \xi} ||E_1^\varepsilon||^2 = \mathcal{O}(e^\varepsilon).
\]

Therefore, the equation of motion for \( \xi \) simplifies to

\[
(4.11) \quad d\xi = \mathcal{O}(\varepsilon^{(3-\kappa/2) - 11/2}) dt + A^{-1}_{11} \frac{1}{2} \frac{\partial}{\partial \xi} \| Q_\varepsilon^{1/2} E_1^\varepsilon \|^2 dt + A^{-1}_{11}(E_1^\varepsilon, dW_\varepsilon).
\]
Remark 4.4. Let us comment in more detail what this formula implies for the motion of the interface. First, $A_{11}$ is approximately the constant $4\ell^2$ with very small derivatives. Moreover, from (4.8) we see that $A_{11}^{-1}E_1^\xi$ is a normalized tangent vector at $\mathcal{M}$. So the deterministic drift in (4.11) is an Itô–Stratonovic correction and the motion of $\xi$ is approximately the Wiener process $\dot{W}_{\varepsilon}$ projected onto $\tilde{\mathcal{M}}$.

Although this is not covered by our assumptions, as a final example we consider a space-time white noise with $Q_{\varepsilon} = \varepsilon \delta \text{Id}$. In this case

$$d\xi = \mathcal{O}(\varepsilon^{3\delta-7/2})dt + \varepsilon^\delta A_{11}^{-1}(E_1^\xi, d\dot{W}_{\varepsilon}),$$

which is a rescaled equation valid on the time scale $\mathcal{O}(\varepsilon^{-\delta})$. Up to the small deterministic error terms, $\xi$ is a stochastic process with mean zero and linear quadratic variation. More specifically,

$$\int_0^t \varepsilon^{2\delta} A_{11}^{-2}(E_1^\xi, E_1^\xi)dt = \varepsilon^{2\delta} \int_0^t A_{11}^{-2}\|\dot{u}_1^\xi\|^2 dt + \mathcal{O}(\varepsilon t)$$

$$= \varepsilon^{2\delta} \int_0^t A_{11}^{-1} dt + \mathcal{O}(\varepsilon^{\delta-3/2+\kappa}) \frac{t}{\ell^2}$$

$$= \frac{\varepsilon^{2\delta}}{4\ell^2} t + \mathcal{O}(\varepsilon^{2\delta+1}) t + \mathcal{O}(\varepsilon^{3\delta-7/2+\kappa}) t.$$

Recalling Levy’s characterization of Brownian motion, in first approximation for times not too large the interface behaves similar to a Brownian motion with variance $\varepsilon^{2\delta}/(4\ell^2)$.

4.2. Conclusions. Let us summarize the results of our analysis:

1. There exists a slow tube $\Gamma$ (around the slow manifold $\Gamma'$) where the coordinate system (cf. (1.7)) is well defined and from which solutions with high probability do not exit for long times $T_{\varepsilon}$ unless one of the interfaces breaks down (stochastic stability).

More specifically, according to Theorem 3.6 this probability is bounded below by

$$1 - C_p(T_{\varepsilon} + 1)\delta_{\varepsilon}^p$$

for any $p > 0$, where $\delta_{\varepsilon}$ measures the noise strength (less than $\varepsilon^{3/2}$). So if the noise is exponentially small, then this probability is large for an exponentially long time, while when the noise being polynomially small the probability is large for any polynomially long time.

2. In $\Gamma$ the approximate SDE of front motion for the stochastic Cahn–Hilliard is given by (4.1). Further, the extra stochastic terms from correlations of the interface motions are important since the deterministic dynamics are exponentially small.

5. Higher order estimates.

5.1. Preliminaries. This section deals with the estimation of all the following higher order terms that appear due to stochastic integration when deriving the equations of motion in the slow channel:

$$\langle \dot{u}_1^\xi, E_{il}^{\xi} \rangle, \quad \langle \dot{u}_k^\xi, E_{il}^\xi \rangle, \quad \text{and} \quad \langle \dot{u}_k^\xi, E_{il}^\xi \rangle.$$

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In addition, we bound the quantity \( \langle L^c \tilde{v}, \tilde{u}_{kl} \rangle \), where for a general smooth in space function \( \phi \) the operator \( L^c \) is defined by

\[
L^c \phi := -\varepsilon^2 \phi_{xxxx} + (f'(u^h)\phi_x)_x.
\]

In order to achieve rigorous estimates for all these terms, we investigate the properties of the stationary problem (1.3). Our analysis admits extensive calculations and is based on the ideas and techniques presented in [8, 9, 22, 23] for the deterministic case where analogous terms of lower order have been estimated already.

Note first that for the construction of the approximate manifold of solutions for the stochastic Cahn–Hilliard equation we use a local coordinate system when presenting the admissible interface positions. The \( h_{N+1} \) variable depends on \( h_i = \xi_i, i = 1, \ldots, N \); therefore, when differentiating two times in \( \xi \) variables and applying the chain rule the second order term \( \frac{\partial^2 h_{N+1}}{\partial \xi_i \partial \xi_j} \) appears. More specifically, for a general function \( f \) smooth in space and any \( i, j = 1, \ldots, N \), we obtain

\[
\begin{align*}
\frac{\partial f}{\partial \xi_i} &= \frac{\partial f}{\partial h_i} + \frac{\partial f}{\partial h_{N+1}} \frac{\partial h_{N+1}}{\partial h_i}, \quad \text{and} \\
\frac{\partial^2 f}{\partial \xi_i \partial \xi_j} &= \frac{\partial^2 f}{\partial h_i \partial h_j} + \left( \frac{\partial^2 f}{\partial h_{N+1}^2 \partial h_j} + \frac{\partial^2 f}{\partial h_{N+1} \partial h_j} \right) \frac{\partial h_{N+1}}{\partial h_i} \\
&\quad + \frac{\partial f}{\partial h_{N+1}} \left( \frac{\partial^2 h_{N+1}}{\partial h_i \partial h_j} + \frac{\partial^2 h_{N+1}}{\partial h_i \partial h_{N+1}} \frac{\partial h_{N+1}}{\partial h_j} \right).
\end{align*}
\]

By the next lemma considering \( \rho = \varepsilon^\kappa \) for some small \( \kappa > 0 \) and thus \( \alpha, \beta \) are exponentially small, we estimate \( |\frac{\partial^2 h_{N+1}}{\partial \xi_i \partial \xi_j}| \). As in [8], where the analogous first order estimate has been derived, we use an implicit function theorem argument combined with the mass conservation constraint. If \( u^h \) is in the second approximate manifold \( M \), then, by definition, mass conservation holds, i.e.,

\[
M = M(h) = \int_0^1 u^h(x)dx.
\]

Differentiating twice with respect to \( h \) variables, we get

\[
\frac{d^2}{dh_i \partial h_j} M(h) = \int_0^1 u_{ij}^h dx,
\]

where \( u_{ij}^h := \frac{\partial^2 u^h}{\partial h_i \partial h_j} = \frac{\partial u^h}{\partial h_j}. \)

**Lemma 5.1.** For any \( i, j = 1, \ldots, N \) the next inequality follows:

\[
|\frac{\partial^2 h_{N+1}}{\partial h_i \partial h_j}| \leq O(\varepsilon).
\]

**Proof.** Let \( \ell \) be a generic positive variable. According to the analysis presented in [22], when comparing the \( x \) and \( \ell \) derivatives of the solution \( \phi \) of the stationary problem (1.3), we obtain a residual function \( w \) given by the following relation:

\[
2\phi_{x}(x, \ell, \pm 1) = -(\text{sgn } x)\phi_x(x, \ell, \pm 1) + 2w(x, \ell, \pm 1).
\]
Let us define \( I_j := [m_j, m_{j+1}] \), \( \chi^j(x) := \chi(\frac{x-h_j}{h}) \). If \( w^j(x) := w(x-m_j, h_j-h_{j-1}, (-1)^j) \), then the interval \([h_{j-1} - \varepsilon, h_{j+1} + \varepsilon]\) contains the support of \( u^h_j \) and

\[
\frac{\partial u^h_j}{\partial x} = \begin{cases} 
\chi^{j-1}w^j & \text{for } x \in I_{j-1}, \\
(1 - \chi^j)(-\phi^j + w^j) + \chi^j(-\phi^{j+1} - w^{j+1}) + \chi^j(\phi^j - \phi^{j+1}) & \text{for } x \in I_j, \\
(1 - \chi^{j+1})w^{j+1} & \text{for } x \in I_{j+1},
\end{cases}
\]

where \( \chi^j = \partial_x(\chi(\frac{x-h_j}{h})) \) and \( \phi^j = \phi(x-m_j, l_j-l_{j-1}, (-1)^j) \) (cf. \([8, p. 430]\)). We note that in \( I_j \) (cf. \([8, p. 430]\))

\[
u^h_j = -u^h_x + (1 - \chi^j)w^j - \chi^j w^{j+1}
\]

and thus

\[
u^h_{ji} = -\frac{\partial u^h_j}{\partial h_i} + (\delta_{j,i} \chi^i_x w^j + (1 - \chi^j)(A_{j,i} u^j_x + B_{j,i} w^j_i)) - \delta_{j,i} \chi^i w^{j+1} - \chi^i (A_{j,i+1} u^{j+1}_x + B_{j,i+1} w^{j+1}_i) \]

in \( I_j \),

where \( w^j_i = w_i(x-m_j, l_j-l_{j-1}, (-1)^j) \) and \( w^j_{i'} = w_{i'}(x-m_j, l_j-l_{j-1}, (-1)^j) \) with \( \delta_{j,i} \) being the Kronecker delta. Moreover,

\[
A_{j,i} := \frac{\partial(x-m_j)}{\partial h_i} = \begin{cases} 0 & \text{for } i \neq j, j-1, \\
-1/2 & \text{for } i = j, j-1,
\end{cases}
\]

and

\[
B_{j,i} := \frac{\partial(h_j - h_{j-1})}{\partial h_i} = \begin{cases} 0 & \text{for } i \neq j, j-1, \\
1 & \text{for } i = j, \\
-1 & \text{for } i = j-1.
\end{cases}
\]

In a similar way we obtain

\[
u^h_{ji} = \delta_{j-1,i} \chi^{j-1}_x w^j + \chi^{j-1}(A_{j,i} u^j_x + B_{j,i} w^j_i) \]

in \( I_{j-1} \),

and

\[
u^h_{ji} = \delta_{j+1,i} \chi^{j+1}_x w^{j+1} - (1 - \chi^{j+1})(A_{j,i} u^{j+1}_x + B_{j,i} w^{j+1}_i) \]

in \( I_{j+1} \).

Using now the bounds on \( w^j, w^j_x, \) and \( w^j_{i'} \) (cf. \([22, 8, p. 172]\)), we obtain for \( r > 0 \)
sufficiently small

\[
\left| \int_{I_{j-1} \cup I_{j+1}} \nu^h_{ji}(x) dx \right| \leq C \varepsilon^{-2}(r^{-1} + 1) \beta(r) \mathcal{K}_{j,i} + \mathcal{O}(\varepsilon r) \delta_{j-1,i} + \mathcal{O}(\varepsilon r) \delta_{j+1,i}
\]

with \( \mathcal{K}_{j,i} = |A_{j,i}| + |A_{j,i+1}| + |B_{j,i}| + |B_{j+1,i}| \) and

\[
\left| \int_{I_j} \left( \delta_{j,i} \chi^i_x w^j + (1 - \chi^j)(A_{j,i} u^j_x + B_{j,i} w^j_i) - \delta_{j,i} \chi^i w^{j+1}
- \chi^j (A_{j+1,i} u^{j+1}_x + B_{j+1,i} w^{j+1}_i) \right) dx \right| \leq C \varepsilon^{-2}(r^{-1} + 1) \beta(r) \mathcal{K}_{j,i} + \mathcal{O}(\varepsilon r) \delta_{j,i}.
\]
Therefore, using the estimates for \( u^i \) it follows that
\[
d^2 \frac{d}{dh_j dh_i} M(h) = \int_0^1 u_j^h dx \\
= \int_j - \partial^2 u^h_{x \partial h_i} dx + O(\varepsilon^{-2}(r^{-1} + 1)\beta(r))K_{j,i} \\
+ O(e^{-l})(\delta_{j-1,i} + \delta_{j,i} + \delta_{j+1,i}) \\
= \int_j \left(- \frac{\partial u^h_{i,j}}{\partial x} \right) dx + O(\varepsilon^{-2}(r^{-1} + 1)\beta(r))K_{j,i} \\
+ O(e^{-l})(\delta_{j-1,i} + \delta_{j,i} + \delta_{j+1,i}) \\
= -(u^h_i(m_{j+1}) - u^h_i(m_j)) + O(\varepsilon^{-2}(r^{-1} + 1)\beta(r))K_{j,i} \\
+ O(e^{-l})(\delta_{j-1,i} + \delta_{j,i} + \delta_{j+1,i}).
\]

Since the support of \( u^h_i \) is \( I_i-1 \cup I_i \cup I_{i+1} \supseteq m_{i-1}, m_i, m_{i+1}, m_{i+2} \) we get that \( d^2 \frac{d^2}{dh_j dh_i} M = 0 \) if \( j \neq i-1, i, i+1, i+2 \), while, for example, \( u^h_i(m_i) = \chi^{i-1}w|m_i = \chi^{i+1}|w(0, l_i, \pm 1) \) and \( u^h_i(m_{i+1}) = -(1-\chi^{i+1})w^{i+1}|m_{i+1} = -(1-\chi^{i+1})w(0, l_{i+1}, \pm 1) \). But \( w(0) = O(\varepsilon^{-1})\alpha_+^*(r) \) \( [22, p. 558] \), since \( \phi(x)(0)^{-1} = \varepsilon^2/f(\phi(0)) \) and \( \varepsilon/l \) is uniformly bounded, while \( \chi \) is \( C^\infty \).

Let us now for simplicity consider \( N = 1 \). Then \( M(h_1, y) = \text{constant} \), when \( y = h_2 \), where \( h_2 \) is a function of \( h_1 \), so
\[
\frac{\partial M}{\partial h_1} + \frac{\partial M}{\partial y} \frac{\partial y}{\partial h_1} = 0
\]
and thus
\[
\frac{\partial^2 M}{\partial h_1 \partial h_1} + \frac{\partial^2 M}{\partial y \partial h_1} + \frac{\partial^2 M}{\partial h_1 \partial h_1} \frac{\partial y}{\partial h_1} = 0.
\]

We set \( y = h_2 \) to get, using the estimate \( \frac{\partial h_2}{\partial h_1} = O(1) \),
\[
O(e^{-l}) + O(e^{-l})O(1) + O(1) \frac{\partial^2 h_2}{\partial h_1^2} = 0
\]
and thus
\[
\frac{\partial^2 h_2}{\partial h_1^2} = O(e^{-l}).
\]
The case \( N > 1 \) follows in a similar way. \( \square \)

5.2. The estimates. We define \( I_s := [-\ell/2 - \varepsilon, \ell/2 + \varepsilon] \). Then for any \( x \in I_s \) it holds that \( [9, 22, 23] \)
\[
|w| \leq c\varepsilon^{-1}\beta_\pm(r), \\
|w_x| \leq c\varepsilon^{-2}r^{-1}\beta_\pm(r), \\
|w_t| \leq c\varepsilon^{-2}\beta_\pm(r), \\
|w_{xt}| \leq c\varepsilon^{-3}r^{-1}\beta_\pm(r), \\
|w_{xx}| \leq c\varepsilon^{-3}\beta_\pm(r).
\]
(5.7)
For the purposes of our proof we will need estimates for the terms

\[ |w_{\ell}|, |w_{xxx}|, |w_{xx\ell}|, |w_{x\ell\ell}|, |w_{xxxx}|, |w_{x\ell\ell\ell}|, \text{ and } |w_{xx\ell\ell}|. \]

It is sufficient to estimate these terms in \( I := [0, \ell/2 + \varepsilon] \) or in \( (0, \ell/2 + \varepsilon] \). We write \( I = [0, \ell/2 - \varepsilon H] \cup [\ell/2 - \varepsilon H, \ell/2 + \varepsilon] \) for a positive \( H \) to be defined in what follows. We set

\[ I_H := [0, \ell/2 - \varepsilon H], \quad \text{and} \quad J := [\ell/2 - \varepsilon H, \ell/2 + \varepsilon] \]

and prove the following lemma bounding the second derivative of \( w \) in \( \ell \) on \( I_s \).

**Lemma 5.2.** For any \( x \in I_s \) it holds that

\[ |w_{\ell\ell}(x)| \leq c\varepsilon^{-3} \beta(r). \tag{5.8} \]

**Proof.** Motivated by the proof of [22] for the estimate of \( |w_\ell| \), we use that

\[ \varepsilon^2 w_{xx} = f'(\phi(x))w \quad \text{in} \quad (0, \ell/2 + \varepsilon) \supset I_H^0 \]

and differentiate twice with respect to \( \ell \) to obtain

\[ \varepsilon^2(w_{\ell\ell})_{xx} - f'(\phi)w_{\ell\ell} = \mathcal{F} \]

for \( \mathcal{F} := f'''(\phi)\phi_\ell^2 w + f''(\phi)\phi_\ell w + 2f''(\phi)\phi_{\ell\ell} w_{\ell} \). By the maximum principle it follows that

\[ |w_{\ell\ell}(x)| \leq \max \left\{ |w_{\ell\ell}(0)|, |w_{\ell\ell}(\ell/2 - \varepsilon H)|, \sup_{x \in I_H^0} \left| \mathcal{F}/f'(\phi) \right| \right\} \text{ for any } x \in I_H. \tag{5.9} \]

Following Carr and Pego (cf. [22, p. 560]), we choose \( H \) such that \( f'(\phi(x)) \geq c_0 > 0 \) for \( 0 < x < \ell/2 - \varepsilon H \). Since \( \varepsilon^2 \phi_\ell^2 = 2(F(\phi) - \alpha) \), there exists a constant \( C > 0 \) such that \( \frac{1}{|\phi_\ell|} \leq \frac{C}{\varepsilon^2} \) for any \( x \in J = [\ell/2 - \varepsilon H, \ell/2 + \varepsilon] \) (cf. [22, pp. 560, 557]).

First we estimate \( |w_{\ell\ell}(x, \ell, -1)| \) in \( J \). It holds that (cf. [22, p. 558])

\[ w(x, \ell, -1) = \varepsilon^{-1} \ell^{-2} \alpha_-'(r) \phi_\ell(|x|, \ell, -1) \int_{\ell/2}^{x} \frac{ds}{\phi_\ell(s, \ell, -1)^2}. \tag{5.10} \]

Let us define \( \mathcal{A} := \int_{\ell/2}^{x} \frac{ds}{\phi_\ell(s, \ell, -1)^2} \). With a slight abuse of notation and for simplicity of notation, we neglect the index "-" in \( \alpha_- \) by using \( \alpha \). Differentiation of (5.10) yields

\[ w_{\ell\ell} = \varepsilon^{-1} \left\{ (\ell^{-2} \alpha'(r))_{\ell\ell}\phi_\ell \mathcal{A} + 2(\ell^{-2} \alpha'(r))_{\ell}\phi_{\ell\ell} \mathcal{A} + 2(\ell^{-2} \alpha'(r))_{\ell\ell}\phi_{\ell\ell} \mathcal{A} \right\} \tag{5.11} \]

According to [22, 23] it follows that

\[ |\alpha'(r)| \leq cr^{-2} \alpha \quad \text{and} \quad |\alpha''(r)| \leq cr^{-4} \alpha. \]

Analogously we obtain

\[ |\alpha'''(r)| \leq cr^{-6} \alpha. \]
Observing that $r = \varepsilon/\ell$ is bounded, i.e., $\ell^{-1} \leq \varepsilon^{-1}$, we derive
\begin{equation}
|\ell^{-2}\alpha'(r)| \leq \varepsilon^{-2}\alpha(r), \quad |(\ell^{-2}\alpha'(r))_r| \leq \varepsilon^{-3}\alpha(r), \quad \text{and} \quad |(\ell^{-2}\alpha'(r))_{rr}| \leq \varepsilon^{-4}\alpha(r).
\end{equation}

Obviously since $x \in J$ one has $|A| \leq \varepsilon^{2+1}$. It holds that (cf. [22, p. 552])
\begin{equation}
\varepsilon^2 \phi_x^2 = 2(F(\phi) - \alpha),
\end{equation}
while
\begin{equation}
\varepsilon^2 \phi_{xx} = f(\phi).
\end{equation}

Since $\int_{\ell/2}^{\ell/2} |\phi_x|dx \leq 2$ (cf. [22, p. 558]), and $\phi$ satisfies a Dirichlet problem, we get by trace inequality that $\phi$ is uniformly bounded. Therefore, we obtain
\begin{align*}
|\phi_x| &\leq \varepsilon^{-1}, \\
|\phi_{xx}| &\leq \varepsilon^{-2}, \quad \text{and} \quad |\phi_{xxx}| \leq \varepsilon^{-3}.
\end{align*}

Using now the definition (5.2) of $w$ and the fact that $|w| + |\phi_x| \leq \varepsilon^{-1}$, we arrive at
\begin{equation}
|\phi_\ell| \leq \varepsilon^{-1},
\end{equation}
while $|\phi_{\ell\ell}| \leq c|\phi_{xx}| + c|w_x|$. So, using $|w_x| \leq \varepsilon^{-2}$ (cf. [9]), we get
\begin{equation}
|\phi_{\ell\ell}| \leq \varepsilon^{-2}.
\end{equation}

By (5.14) it follows that
\begin{equation}
|\phi_{\ell\ell\ell}| \leq \varepsilon^{-3}.
\end{equation}

Finally, we also need an estimate for the term $\phi_{\ell\ell\ell}$. We differentiate (5.13) twice with respect to $\ell$ in order to obtain
\begin{equation}
|\varepsilon^2 \phi_{x}\phi_{\ell\ell}| \leq \varepsilon^{-2}.
\end{equation}

Hence using the bound $\frac{1}{|\phi_x|} \leq \varepsilon$ valid in $J$, it holds that
\begin{equation}
|\phi_{\ell\ell}| \leq \varepsilon^{-3} \quad \text{in} \quad J.
\end{equation}

In order to compute the derivatives of $A$ in (5.11), we apply the formulas
\begin{equation*}
\frac{d}{d\ell} \int_{s(\ell)}^{b} g(s, \ell)ds = \int_{s(\ell)}^{b} g_\ell(s, \ell)ds - s'(\ell)g(s(\ell), \ell)
\end{equation*}
and
\begin{align*}
\frac{d^2}{d\ell^2} \int_{s(\ell)}^{b} g(s, \ell)ds &= \int_{s(\ell)}^{b} g_{\ell\ell}(s, \ell)ds - s'(\ell)g_\ell(s(\ell), \ell) \\
& \quad - s''(\ell)g(s(\ell), \ell) - s'(\ell)^2g_x(s(\ell), \ell) - s'(\ell)^2g_x(s(\ell), \ell) - s'(\ell)g_\ell(s(\ell), \ell).
\end{align*}

After tedious calculations, using the estimates above and the fact that the length of the interval is of order $O(\varepsilon)$, we arrive at
\begin{align*}
|A_\ell| &\leq \varepsilon^2, \quad |A_{\ell\ell}| \leq \varepsilon\varepsilon.
\end{align*}
We note that $\varepsilon/\ell$ is bounded, i.e., $\ell^{-1} \leq c\varepsilon^{-1}$. Thus by (5.11) and (5.12) we obtain
\[(5.15)\quad |w_{\ell \ell}| \leq c\varepsilon^{-3}\alpha \quad \text{in} \quad J.
\]
So by (5.15), since $\ell/2 - \varepsilon H \in J$, it follows that
\[(5.16)\quad |w_{\ell \ell}(\ell/2 - \varepsilon H)| \leq c\varepsilon^{-3}\alpha.
\]
By the definition of $F$, the fact that $f' \geq c_0 > 0$ in $I_H$, and the first and third estimate of (5.7), we get (using $\beta := \beta_-$) that
\[
\sup_{I_H} F/f'(\phi) \leq c \left[ |\phi_{\ell}^2|w + |\phi_{\ell\ell}|w + |\phi_{\ell\ell}|w_{\ell} \right] \leq c\varepsilon^{-1}\beta \left[ |\phi_{\ell}^2| + |\phi_{\ell\ell}| + \varepsilon^{-1}|\phi_{\ell}| \right].
\]
In addition, since $|w_{\ell}| + |\phi_{\ell\ell}| \leq c\varepsilon^{-2}$ [22, 9], it follows that
\[
|\phi_{\ell}| \leq c\varepsilon^{-2}.
\]
Thus as we already proved $|\phi_{\ell}| \leq c\varepsilon^{-1}$, we derive
\[(5.17)\quad \sup_{I_H} F/f'(\phi) \leq c\varepsilon^{-3}\beta.
\]
Let us now turn to the missing estimate on $|w_{\ell \ell}(0)|$. In [23] by using $w(0) = -\partial\phi/\partial x (\varepsilon/\ell)$, it was demonstrated that $|w_{\ell}(0)| \leq c\varepsilon^{-2}\beta$. Analogously by differentiation in $\ell$, it follows that
\[(5.18)\quad |w_{\ell \ell}(0)| \leq c\varepsilon^{-3}\beta.
\]
Using now (5.9), (5.15), (5.16), (5.17), and (5.18) we obtain that $|w_{\ell \ell}(x)| \leq c\varepsilon^{-3}\beta$ for any $x$ in $I = I_H \cup J$. By symmetry we prove finally that $|w_{\ell \ell}| \leq c\varepsilon^{-3}\beta \pm r(\ell)$ in $I_s$. \(\square\)

The next three lemmas present bounds for the third or higher order terms.

**Lemma 5.3.** For any $x \in I_s^0 - \{0\}$ it holds that
\[(5.19)\quad |w_{xxx}(x)| \leq c\varepsilon^{-4}\beta \pm r(\ell)
\]
and
\[(5.20)\quad |w_{x \ell \ell}(x)| \leq c\varepsilon^{-4}\beta \pm r(\ell).
\]

**Proof.** We consider $x \in (0, \ell/2 + \varepsilon)$ and $\varepsilon^2 w_{xx} = f'(\phi)w$. By differentiation in $x$, using (5.7), and the bound on $|\phi_x|$, or by differentiating in $\ell$, using (5.7), and the bound on $|\phi_{\ell}|$, we get the following:
\[
|w_{xxx}| \leq c\varepsilon^{-2} \left[ |f'(\phi)||w_x| + |f''(\phi)||\phi_x||w| \right]
\leq c\varepsilon^{-2} \left[ c\varepsilon^{-2}\beta + c\varepsilon^{-1}|\phi_{\ell}| \right] \leq c\varepsilon^{-4}\beta
\]
and
\[
|w_{x \ell \ell}| \leq c\varepsilon^{-2} \left[ |f'(\phi)||w_{\ell}| + |f''(\phi)||\phi_{\ell}||w| \right]
\leq c\varepsilon^{-2} \left[ c\varepsilon^{-2}\beta + c\varepsilon^{-1}|\phi_{\ell}| \right] \leq c\varepsilon^{-4}\beta
\]
with $\beta = \beta_-$. Again by symmetry, we obtain the bounds for all $x$ in $I_s^0 - \{0\}$. \(\square\)
Lemma 5.4. For any $x \in I_s - \{0\}$ it holds that

\begin{equation}
|w_{x\ell\ell}(x)| \leq c \varepsilon^{-4} r^{-1} \beta_\pm(r).
\end{equation}

Proof. Consider $x \in (0, \ell/2 + \varepsilon]$ and write $w_{x\ell\ell}(\ell/2) - w_{x\ell\ell}(x) = \int_x^{\ell/2} w_{x\ell\ell}(s)\, ds$ in order to obtain

\begin{equation}
|w_{x\ell\ell}(x)| \leq |w_{x\ell\ell}(\ell/2)| + \int_x^{\ell/2} |w_{x\ell\ell}(s)|\, ds.
\end{equation}

We use the definition of $w$ given in (5.10), set $p = \varepsilon^{-1} \ell^2 \alpha'$, and recall that $A = \int_0^{\ell/2} \phi^{-2}_{\ell} \, ds$. Differentiating first with respect to $x$ and then twice with respect to $\ell$

\begin{equation}
w_{x\ell\ell} = p_{\ell} \phi_{xx} A + p_{\ell} \phi_{x\ell} A_t + p_{\ell} \phi_{xx\ell} A_t + p_{\ell} \phi_{x\ell\ell} A
\end{equation}

\begin{equation*}+
2p \phi_{x\ell} A_t + p \phi_{xx} A_t - p_{\ell} \phi_{x\ell} - p (\phi_{x\ell} \phi_{xx}^2 - 2 \phi_{x\ell} \phi_{xx} \phi_{x\ell} / \phi_x^2).
\end{equation*}

Observe that $A = 0$ at $x = \ell/2$, while

\begin{equation}
A_{t}(\ell/2) = -\frac{1}{2} \phi_x(\ell/2)^{-2} \quad \text{and} \quad A_{tt}(\ell/2) = \phi_x(\ell/2)^{-3} \phi_{xt}(\ell/2) + \phi_{x\ell}(\ell/2)^{-3} \phi_{x\ell}(\ell/2)^{-3}.
\end{equation}

We also note $\ell/2 \in J$; thus by the bounds of Lemma 5.2 we obtain $|\phi_{x\ell}(\ell/2)| \leq c \varepsilon^{-2}$ and $|\phi_x(\ell/2)^{-1}| \leq c \varepsilon$. Hence, as in Lemma 5.2 for a general $x \in J$, we get $\varepsilon^{-1}|A_{t}(\ell/2)| + |A_{tt}(\ell/2)| \leq c \varepsilon$.

In addition, using the third estimate from (5.12) yields $|p_{\ell\ell}(\ell/2)| \leq c \varepsilon^{-5} \alpha$. Furthermore, as $\ell/2 \in J$ by the proof of Lemma 5.2 we have $|\phi_{xx}(\ell/2)| \leq c \varepsilon^{-2}$ and $|\phi_{x\ell}(\ell/2)| \leq c \varepsilon^{-3}$. Therefore, we obtain finally (with $\alpha = \alpha_-$)

\begin{equation}
|w_{x\ell\ell}(\ell/2)| \leq c \varepsilon^{-4} \alpha.
\end{equation}

Recall $\varepsilon^2 \phi_{xx} = f'(\phi) w$ in $(0, \ell/2 + \varepsilon)$. By taking twice the $\ell$-derivative yields $|\beta = \beta_-$)

\begin{equation}
|w_{x\ell\ell\ell}| \leq c \varepsilon^{-2} \left[ |\phi_{\ell}|^2 |w| + |\phi_{\ell}||w_{\ell}| + |w_{\ell\ell}| \right] \leq c \varepsilon^{-5} \beta.
\end{equation}

Here, we used the bound $|\phi_{\ell}| \leq c \varepsilon^{-1}$ from the proof of Lemma 5.2, the first and third estimate of (5.7), the fact that $|w| \leq c \varepsilon^{-1} \beta$ and $|w_{\ell}| \leq c \varepsilon^{-2} \beta$, and the bound $|w_{\ell\ell}| \leq c \varepsilon^{-3} \beta$ of Lemma 5.2.

Using $r = \varepsilon/\ell$ and $x \in (0, \ell/2)$ we get $|x - \ell/2| \leq c (\ell/2 + \varepsilon) \leq c \varepsilon r^{-1}$. Therefore, (5.22), (5.23), and (5.24) imply

\begin{equation}
|w_{x\ell\ell}(x)| \leq c \varepsilon^{-4} r^{-1} \beta, \quad x \in (0, \ell/2 + \varepsilon).
\end{equation}

By symmetry the analogous result holds for any $x \in [-\ell/2 - \varepsilon, 0)$. \qed

Analogously the next lemma follows.

Lemma 5.5. For any $x \in I_s - \{0\}$ it holds that

\begin{equation}
|w_{x\ell\ell\ell}(x)| + |w_{x\ell\ell}(x)| + |w_{x\ell\ell\ell}(x)| \leq c \varepsilon^{-5} r^{-1} \beta_\pm(r).
\end{equation}

In order to estimate $E_{i1}^\xi$, $E_{ij}^\xi$, and $E_{ijk}^\xi$ we first need the following lemma for the correction terms $Q_j$. 

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Lemma 5.6. For any $i, j, k$ it follows that

\[
|Q_{ij}| \leq c \varepsilon^{-3} \beta, \\
|Q_{ij}| \leq c \varepsilon^{-4} r^{-1} \beta, \\
|Q_{ijk}| \leq c \varepsilon^{-5} r^{-1} \beta.
\]

(5.26)

Proof. We recall that

\[
u_h^j = \begin{cases} 
\chi^{j-1} w^j & \text{on } I_{j-1}, \\
(1 - \chi^j)(-\phi_x^j + w^j) + \chi^j(-\phi_x^j + w^j + 1) + \chi_x^j(\phi_x^j - \phi_x^{j+1}) & \text{on } I_j, \\
-(1 - \chi^{j+1})w^{j+1} & \text{on } I_{j+1}.
\end{cases}
\]

Consider the functions on $x = 0, 1$ in the first and last sets of their support. Using the bounds on $|w|, |w_{xx}|$, we arrive at

\[
|\tilde{u}^j_h| \leq c \varepsilon^{-1} \beta \text{ and thus } |\tilde{w}_j| \leq c \varepsilon^{-1} \beta, \\
|\tilde{u}^h_{jj}| \leq c \varepsilon^{-3} \beta \text{ and thus } |\tilde{w}_{jxx}| \leq c \varepsilon^{-3} \beta.
\]

The estimates of $|w_x|$ and $|w_t|$ and of $|w_{xxx}|$ and $|w_{txx}|$, respectively, give

\[
|\tilde{u}^h_{j1}| \leq c \varepsilon^{-2} r^{-1} \beta \text{ and thus } |\tilde{w}_{j1}| \leq c \varepsilon^{-2} r^{-1} \beta, \\
|\tilde{u}^h_{jxx}| \leq c \varepsilon^{-4} r^{-1} \beta \text{ and thus } |\tilde{w}_{jxx}| \leq c \varepsilon^{-4} r^{-1} \beta.
\]

Finally, using the estimates of $|w_{xx}|, |w_{tt}|$, and $|w_{ttx}|$ and of $|w_{xxxx}|, |w_{xxx}|$, and $|w_{txxx}|$, respectively, we obtain

\[
|\tilde{u}^h_{jik}| \leq c \varepsilon^{-3} r^{-1} \beta \text{ and thus } |\tilde{w}_{jik}| \leq c \varepsilon^{-3} r^{-1} \beta, \\
|\tilde{u}^h_{jikxx}| \leq c \varepsilon^{-5} r^{-1} \beta \text{ and thus } |\tilde{w}_{jikxx}| \leq c \varepsilon^{-5} r^{-1} \beta.
\]

Recall also

\[
\tilde{w}_j := \tilde{u}^h_j + \tilde{u}^h_{j+1}, \\
Q_j(x) := \left( -\frac{1}{6} x^3 + \frac{1}{2} x^2 - \frac{1}{3} x \right) \tilde{w}_{jxx}(0) + \frac{1}{6} (x^3 - x) \tilde{w}_{jxx}(1) + x \tilde{w}_j(1), \ j = 1, \ldots, N.
\]

This definition of $Q_j$ combined with the previously obtained bounds on $\tilde{w}_j$ imply the result. \( \square \)

Remark 5.7. By [22, pp. 557–565], the following estimates hold true. First,

\[
\int_{-\ell/2}^{\ell/2} \phi_x(x, \ell, -1)^2 + \int_{0}^{\ell/2} \phi_x(x, \ell, +1)^2 \leq \varepsilon^{-1} S_\infty + E(r),
\]

(5.27)

where $|E| \leq c \varepsilon^{-1} \beta$ and $S_\infty = \int_{-1}^{1} \sqrt{2F(u)}du$. Moreover,

\[
\int_{-\ell/2}^{\ell/2} |\phi_x| dx \leq 2
\]

(5.28)

and

\[
\int_{-\ell/2}^{\ell/2} |\phi_{xx}|^2 dx \leq c \varepsilon^{-3}.
\]

(5.29)
In addition, there exists a constant $c > 0$ such that for all $x \in [h_j - \varepsilon, h_j + \varepsilon], j = 0, \ldots, N + 1$, we have
\begin{align}
\phi^j(x) - \phi^{j+1}(x) &\leq c|\alpha^j - \alpha^{j+1}|, \\
\phi^{j+1}_x(x) - \phi^{j+1}_x(x) &\leq c\varepsilon^{-1}|\alpha^j - \alpha^{j+1}|,
\end{align}
and
\begin{align}
\phi^{j+1}_{xx}(x) - \phi^{j+1}_{xx}(x) &\leq c\varepsilon^{-2}|\alpha^j - \alpha^{j+1}|,
\end{align}
provided $\varepsilon/\ell_j, \varepsilon/\ell_{j+1} < r_0$ for some sufficiently small $r_0 > 0$ (cf. [8]).

Now, we are able to bound the terms $\tilde{u}^h$ and $u^h$.

**Theorem 5.8.** For any $i, j, k$ it holds that
\begin{align}
\|\tilde{u}^h_j\|_{\infty} &\leq \mathcal{O}(1) + \mathcal{O}(\|w\|_{\infty}), \\
\|\tilde{u}^h_{jk}\| &\leq c\varepsilon^{-1/2}(1 + S^{1/2}_\infty + \max(r_j\alpha^j, r_{j+1}\alpha^{j+1})^{1/2}) + c\|w_x\| + c\|w_\ell\|, \\
\|\tilde{u}^h_{jk1}\| &\leq c\varepsilon^{-3/2} + c\|w_x\| + c\|w_\ell\| + c\|w_x\| + c\|w_\ell\|, \\
\|u^h_j\|_{\infty} &\leq \mathcal{O}(\varepsilon^{-1}), \\
\|u^h_{j}\|_{\infty} &\leq \mathcal{O}(\varepsilon^{-2}), \\
\|u^h_{j}\| &\leq \mathcal{O}(\varepsilon^{-1/2}).
\end{align}

**Proof.** We use the definition of $u^h_j$ and get by (5.28) that
\begin{align}
\|\tilde{u}^h_j\| \leq c \int_0^x |\phi_x| dx + c\|w\|_{\infty} \leq c + c\|w\|_{\infty}.
\end{align}

Also, it follows that $\|u^h_j\| = \mathcal{O}(\|\phi_x\|) = \mathcal{O}(\varepsilon^{-1})$, and thus
\begin{align}
\|u^h_j\|_{\infty} \leq \mathcal{O}(\varepsilon^{-1}).
\end{align}

By [8, p. 38],
\begin{align}
u^h_j = -u^h_x + (1 - \chi^j)w^j - \chi^j w^{j+1} \quad \text{on } I_j.
\end{align}

Combining this with (5.3) we obtain
\begin{align}
u^h_{ji}(x) = \begin{cases} \mathcal{O}(w_x + w_\ell) & \text{for } x \in I_{j-1}, \\
-\mathcal{O}(w_x + w_\ell) & \text{for } x \in I_j, \\
\mathcal{O}(w_x + w_\ell) & \text{for } x \in I_{j+1}.
\end{cases}
\end{align}

Therefore, we arrive at
\begin{align}
\tilde{u}^h_{ji}(x) = \int_0^x u^h_{ji}(y)dy = \begin{cases} \mathcal{O}(w_x + w_\ell) & \text{for } x \in I_{j-1}, \\
\mathcal{O}(u^h_x + w_x + w_\ell) & \text{for } x \in I_j, \\
\mathcal{O}(w_x + w_\ell) & \text{for } x \in I_{j+1}.
\end{cases}
\end{align}

By [22, p. 563, (8.6)],
\begin{align}
\|u^h_j\| \leq \varepsilon^{-1/2}(S^{1/2}_\infty + \max(r_j\alpha^j, r_{j+1}\alpha^{j+1})^{1/2}).
\end{align}

Using this estimate we obtain
\begin{align}
\|\tilde{u}^h_{ji}\| \leq c\varepsilon^{-1/2}(1 + S^{1/2}_\infty + \max(r_j\alpha^j, r_{j+1}\alpha^{j+1})^{1/2}) + c\|w_x\| + c\|w_\ell\|.
\end{align}
For higher derivatives, observe now that

$$
\tilde{u}^{h}_{jik}(x) = \int_{0}^{x} u^{h}_{jik}(y) dy = \begin{cases} 
O(w_{x} + w_{x}w + w_{x}w_{t}) & \text{for } x \in I_{j-1}, \\
O(w_{x}^{2} + w_{x}w_{x} + w_{x}w_{t}) & \text{for } x \in I_{j}, \\
O(w_{x}w + w_{x}w_{t}) & \text{for } x \in I_{j+1}.
\end{cases}
$$

In addition, since $u^{h}_{j} = -u^{h}_{x} + (1 - \chi^{j})w_{x} - \chi^{j}w_{x}w_{x+1}$ on $I_{j}$, we obtain

$$
\|u^{h}_{x}\| \leq \|u^{h}_{x}\| + c\|w_{x}\|.
$$

The argument of Lemma 8.3 of [22, p. 562] applied to $u^{h}$ on $I_{j}$ using the support of $|\phi'_{x} - \phi^{j+1}_{x}|$ combined with (5.31), (5.28), and

$$
\phi'_{x} = O(|\phi_{x}|) + O(|\phi_{x}^{j} - \phi_{x}^{j+1}|)
$$

finally yields

(5.36)

$$
\|u^{h}_{x}\| \leq \|\phi_{x}\| + \sqrt{O(\epsilon^{-2}\epsilon)} \leq \epsilon^{-1/2}.
$$

Analogously, differentiating $u^{h}$ twice with respect to $x$, using the bounds of (5.32) and (5.29), and using the support of $|\phi'_{x} - \phi^{j+1}_{x}|$ yields

$$
u^{h}_{xx} = O(|\phi_{xx}|) + O(|\phi_{xx}^{j} - \phi_{xx}^{j+1}|).
$$

Thus

$$
\|u^{h}_{xx}\| \leq \|\phi_{xx}\| + \sqrt{O(\epsilon^{-3/2})} \leq \epsilon^{-3/2}.
$$

So, it follows that

(5.37)

$$
\|u^{h}_{x}\| \leq c\epsilon^{-3/2} + c\|w_{x}\|.
$$

Combining the previous estimates yields

$$
\|u^{h}_{jik}\| \leq c\epsilon^{-3/2} + c\|w_{x}\| + c\|w_{x}\| + c\|w_{x}\| + c\|w_{x}\| + c\|w_{x}\|.
$$

Using again (5.34) we obtain

$$
|u^{h}_{ij}| \leq O(u^{h}_{xj}) = O(u^{h}_{xx}) = O(\epsilon^{-2}).
$$

Therefore

$$
\|u^{h}_{ij}\| \leq O(\epsilon^{-2}).
$$

Further, by (5.34) and (5.36) it follows that

$$
\|u^{h}_{j}\| \leq O(\|u^{h}\|) = O(\epsilon^{-1/2}).
$$

Combining now the bound $|\tilde{u}^{h}_{j}| \leq O(1) + O(\|w\|)$ with the implicit function theorem, for the change of variables for $h$ to $\xi$, we obtain

(5.38)

$$
|\tilde{u}^{h}_{j}| \leq (O(1) + O(\|w\|))(O(1) + O(\epsilon^{-1/2}))
$$
Moreover, for the second derivative in $\xi$ variables
\[ \hat{u}_{jk}^\xi \leq |\hat{u}_{jk}^h| |O(1) + O(\varepsilon^{-1}\beta)|^2 + |\hat{u}_{jk}^h||O(1) + O(\varepsilon^{-1}\beta)| + |\hat{u}_{jk}^h|O(\varepsilon). \]
So we have verified the following lemma.

**Lemma 5.9.** For all $j, k$ it holds that
\[ \|\hat{u}_{jk}^\xi\| \leq (O(1) + O(\|w\|))[O(1) + O(\varepsilon^{-1}\beta)] \]
and
\[ \|\hat{u}_{jk}^\xi\| \leq [O(1) + O(\varepsilon^{-2}\beta^2) + O(\varepsilon^{-1}\beta)]|O(w_x + w_t) + \varepsilon^{-1/2} + \varepsilon^{-1/2}A] \]
with defined as $A = S_\infty^{1/2} + \max_j(r_j\alpha^j, r_j+1\alpha^{j+1})^{1/2}$.

The following theorem gives the bounds on $E_{ij}^\xi$ in the $L^2$-norm.

**Theorem 5.10.** For all $i, j, k$ the following inequalities hold:
\[ \|E_{ij}^\xi\| \leq 4\ell_{i+1} + O(\varepsilon^{-3}\beta), \]
\[ \|E_{ij}^\xi\| \leq O(\varepsilon^{-1/2}) + O(\varepsilon^{-r-1}\beta), \]
\[ \|E_{ij}^\xi\| \leq O(\varepsilon^{-3/2}) + O(\varepsilon^{-5r-1}\beta). \]

**Proof.** Using the bound $\|E_{ij}^\xi\| \leq \|\hat{w}_j\| + \|Q_j\|$, the estimate of $\|\hat{w}_j\|$ presented in (4.24) on p. 186 of [9], and our Lemma 5.6, we obtain (5.41). Also, observe
\[ E_{ji}^\xi = \hat{w}_{ji} + O(Q_{ji}) + O(Q_{ijx}) = O(w_x + w_t) + \int_0^x (-u_{xxi}^h)dy + O(Q_{ji}) + O(Q_{ijx}) \]
\[ \leq O(w_x + w_t) + O(u_{xxi}^h) + O(Q_{ji}). \]
Hence, by (5.35) and Lemma 5.6, we conclude
\[ \|E_{ji}^\xi\| \leq O(\varepsilon^{-1/2}) + O(\varepsilon^{-r-1}\beta). \]

Furthermore, using (5.34) we obtain
\[ E_{jik}^\xi = \hat{w}_{jik} + O(Q_{jik}) + O(Q_{jikx}) \]
\[ = O(w_{xx} + w_{xt} + w_{xt}) + \int_0^x (-u_{x}^h)dy + O(Q_{jik}) + O(Q_{jikx}) \]
\[ \leq O(w_{xx} + w_{xt} + w_{xt}) + O(u_{xx}^h) + O(Q_{jik}). \]
Thus, by (5.37) and Lemma 5.6 this implies
\[ \|E_{jik}^\xi\| \leq O(\varepsilon^{-3/2}) + O(\varepsilon^{-5r-1}\beta). \]

**Remark 5.11.** We note that the bound on $\|E_{ij}^\xi\|$ presented in Theorem 5.10 coincides in the main order term with the estimate that was used but not presented analytically in [9].

Using all the results of the previous analysis we are now ready to derive by Cauchy–Schwarz inequality all the desired estimates for the higher order derivatives. They are presented in the following main theorem of this section.
Theorem 5.12. These inequalities hold for all \( i, l, k \):
\[
\langle \tilde{u}^\xi, E^\xi_i \rangle \leq O(\varepsilon^{-1/2}) \left[ 4l_{i+1} + O(\varepsilon^{-3}\beta) \right],
\]
\[
\langle \tilde{u}^\xi, E^\xi_{il} \rangle \leq O(\varepsilon^{-1/2}) + O(\varepsilon^{-4}r^{-1}\beta),
\]
\[
\langle \tilde{v}, E^\xi_{dk} \rangle \leq O(\varepsilon^{-3/2}) + O(\varepsilon^{-5}r^{-1}\beta) \cdot \|\tilde{v}\|.
\]

It remains to analyze \( \langle L^\varepsilon \tilde{v}, \tilde{u}^\xi_{hl} \rangle \). Here we provide the following main result.

Theorem 5.13. For all \( k \) and \( l \) it holds that
\[
\langle L^\varepsilon \tilde{v}, \tilde{u}^\xi_{hl} \rangle \leq \varepsilon^{-5}\beta(r) \left( O(1) + O(\varepsilon^{-2}\beta(r)^2) \right) \|\tilde{v}\|.
\]

Proof. Note that by symmetry and definition
\[
\langle L^\varepsilon \tilde{v}, \tilde{u}^\xi_{hl} \rangle = -\langle \tilde{v}, \partial_{\xi} \partial_{\xi} \partial_{\xi} \partial_{\xi} \rangle \cdot L^b(u^\xi) \ .
\]

Recall that we defined \( L^b(\phi) = \varepsilon^2 \phi_{xx} - f(\phi) \). As in [8, pp. 452–453] for \( x \in [h_j - \varepsilon, h_j + \varepsilon], j = 1, 2, \ldots, N + 1 \), we write
\[
L^b(u^h) = f_1 + f_2 + G,
\]
where we defined
\[
f_1 := \varepsilon^2 \chi^1(x)(\phi^{l+1} - \phi^l), \quad f_2 := 2\varepsilon^2 \chi^1(x)(\phi^{l+1} - \phi^l),
\]
\[
G := (\phi^{l+1} - \phi^l)^2 \left\{ (1 - \chi^1) \int_0^{\chi^1} sf''(\theta)d\theta + \chi^1 \int_0^1 f''(\theta)d\theta \right\}
\]
with \( \theta = \theta(s) := (1-s)\phi^l(x) + s\phi^{l+1}(x) \). For all other \( x \), we have no contribution of \( L^b(u^h) \).

In Lemma 5.2 of [8, p. 454], after differentiating \( f_1 \), \( f_2 \), \( G \) with respect to \( h_j \) it is derived that
\[
\left| \frac{\partial}{\partial h_j} L^b u^h \right| \leq c\varepsilon^{-2}\beta(r) .
\]

Applying the analogous argument to the second differential with respect to \( h_j \) and \( h_i \) yields after some calculation
\[
\left| \frac{\partial^2}{\partial h_j \partial h_i} L^b u^h \right| \leq c\varepsilon^{-3}\beta(r) .
\]

Note that in this argument the worst term is \( |\phi_{xxx}^l(x) - \phi_{xxx}^{l+1}(x)| \). But as \( \varepsilon^2 \phi_{xxx} = f'(\phi)\phi_x \) with \( f(\phi) = \phi^3 - \phi \) and \( f'(\phi) = 3\phi^2 - 1 \), using the estimates of \( \phi, \phi_x \) and the results for the differences presented on p. 453 of [8], we get
\[
|\phi_{xxx}^l(x) - \phi_{xxx}^{l+1}(x)| = \varepsilon^{-2}|f'(\phi)\phi_x^l(x) - f'(\phi)\phi_x^{l+1}(x)|
\]
\[
= \varepsilon^{-2}|f'(\phi)\phi_x^l(x) - f'(\phi)\phi_x^{l+1}(x) - f'(\phi)\phi_x^{l+1}(x)|
\]
\[
\leq \varepsilon^{-2}|f'(\phi)||\phi_x^l(x) - \phi_x^{l+1}(x)| + \varepsilon^{-2}|\phi_x^{l+1}(x)||f'(\phi) - f'(\phi^{l+1})|
\]
\[
\leq c\varepsilon^{-2}|\phi_x(x) - \phi_x^{l+1}(x)| + c\varepsilon^{-2}|\phi_x^{l+1}(x)||f'(\phi) - f'(\phi^{l+1})|
\]
\[
\leq c\varepsilon^{-3}|\phi^{l+1} - \phi^{l+1}| + c\varepsilon^{-3}|3\phi^{l+1}(x) - 1 - 3\phi^{l+1}(x)^2 + 1|
\]
\[
\leq c\varepsilon^{-3}|\phi^{l+1} - \phi^{l+1}| + c\varepsilon^{-3}|\phi^{l+1}(x) + \phi^{l+1}||\phi_x(x) - \phi^{l+1}|
\]
\[
\leq c\varepsilon^{-3}|\phi^{l+1} - \phi^{l+1}| + c\varepsilon^{-3}|\phi^{l+1} - \phi^{l+1}| .
\]

Again as in [8, p. 456], by using that \( \varepsilon^2 w_{xx} = f'(\phi(x))w \) and differentiation in \( x, w \)
obtain

\[(5.50) \quad \left| \frac{\partial^2}{\partial h_j \partial h_i \partial x} \mathcal{L}^b u^h \right| \leq c\varepsilon^{-5}\beta(r).\]

Changing now to $\xi$ variables, using that the second derivative appears by applying the formula (5.1) to (5.50), and since (cf. [8, p. 454]) it holds that

\[(5.51) \quad \left| \frac{\partial}{\partial h_j} \frac{\partial}{\partial x} \mathcal{L}^b u^h \right| \leq c\varepsilon^{-4}\beta(r),\]

we finally obtain

\[(5.52) \quad \left| \frac{\partial^2}{\partial \xi_k \partial \xi_l \partial x} \mathcal{L}^b u^h \right| \leq \varepsilon^{-5}\beta(r) \left\{ (\mathcal{O}(1) + \varepsilon^{-1}\beta(r))^2 + (\mathcal{O}(1) + \varepsilon^{-1}\beta(r)) \right\} + \varepsilon^{-4}\beta(r)\mathcal{O}(\varepsilon),\]

we obtain

\[\varepsilon^{-5}\beta(r) \left\{ (\mathcal{O}(1) + \varepsilon^{-2}\beta(r))^2 \right\}.\]

So, the result follows.

\[\square\]

REFERENCES